# Gratings: Theory and Numeric Applications 

Tryfon Antonakakis, Fadi Issam Baida, Abderrahmane Belkhir, Kirill Cherednichenko, Shane Cooper, Richard Craster, Guillaume Demésy, John Desanto,, Gérard Granet, Boris Gralak, et al.

## To cite this version:

Tryfon Antonakakis, Fadi Issam Baida, Abderrahmane Belkhir, Kirill Cherednichenko, Shane Cooper, et al.. Gratings: Theory and Numeric Applications. E. Popov. AMU (PUP), pp.430, 2012, E. Popov, 978-2-8539-9860-4. hal-00785737v1

## HAL Id: hal-00785737 <br> https://hal.science/hal-00785737v1

Submitted on 6 Feb 2013 (v1), last revised 11 Feb 2013 (v2)

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## GRATINCES: THBORY AND NUNIBIBIC APPHCAYHONS

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ISBN: 2-85399-860-4

First Edition, 2012, PresseVuniversitaireVde Provence (PUP)

## World Wide Web:

www.fresnel.fr/numerical-grating-book

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Gratings: Theory and Numeric Applications, Evgeny Popov, editor (Institut Fresnel, CNRS, AMU, 2012)

## Editorial Preface

A typical question that almost all of us (the authors' team and other colleagues) has been asked not only once has in general the meaning (although usually being shorter): "What is the best method for modeling of light diffraction by periodic structures?" Unfortunately for the grating codes users, and quite fortunately for the theoreticians and code developers, the answer is quite short, there is no such a bird like the best method.

In the more than 30 years active studies on the subject, I have worked on the theory and numerical applications of several approximate methods, like Rayleigh expansion, coupled-wave theory, beam propagation method, first-order approximations, singular Green's function approximation, effective index medium theory, etc. My conviction is that they are quite useful (otherwise why to exist) for physical understanding, but my heart lies in what is considered as rigorous grating theories. Name 'rigorous' is used in the sense that in establishing the theories, exact vector macroscopic Maxwell equations and boundary conditions are applied without approximations. The approaches become approximate after computer implementation, due to the impossibility to work with infinite number of equations and unknowns, and due to the finite length of the computer word.

Of course, there are always initial approximations and assumptions, like the infinite dimensions of the grating plane, linearity of the optical response, etc. From physical point of view, the main feature of the methods, presented in this book are characterized by the use of optical parameters of different substances as something given by other physical optics theories and the experiment as an ultimate judge.

The necessity to use more than a single rigorous method comes from practice: different optogeometrical structures made of different materials and working in different spectral regions require a variety of methods, because each one is more effective in some cases, and less effective (or failing completely) in others. In addition, each approach is a subject of constant research and development. Grating modeling, grating manufacturing and grating use go hand in hand, and practice provides strong stimuli for the theory development. Vice versa, recent grating technologies and application cannot advance without proper theoretical and numerical support.

When I started my grating studies, the method of coordinate transformations that uses eigenvector technique to integrate the Maxwell equations (sometimes known as the Cmethod) has just been formulated. It worked perfectly for holographic grating whatever the polarization and the grating material, but failed completely for grooves with steep facets. It took more than 15 years to refine its formulation, so that now it can deal with echelles and pyramidal bumps (in the case of two-dimensional periodicity) with slopes up to 87 deg steepness. However, the method is not at all adapted to lamellar gratings. On the other hand, the Fourier modal method (also known as Rigorous coupled-wave approach, RCW) is perfect for such profiles, but its use in the case of arbitrary grating profiles (e.g., sinusoidal or triangular profiles) in case of metallic grating material causes problems when using a staircase approximation of the profile. The differential method does not use this approximation and could deal with arbitrary profiles, but it took more than 20 years to make it working with
metallic gratings in TM (p, or S) polarization. And quite ironically, the improvement came from advances in the competing RCW approach.

These methods are relatively easy for programing nowadays, after solving the numerical problems due to growing exponentials and factorization rules of the product of permittivity and electric field, however there are still some persisting problems for highly conduction metals. In addition, neither the differential, nor the Fourier modal methods can deal with infinitely conducting gratings.

Several methods are quite flexible concerning the geometry of the diffracting objects and the grating material. For example, the integral method can treat inverted profiles, rod gratings with arbitrary cross section, finitely or infinitely conducting materials in any polarization, but its programming require deep mathematical understanding of the singularities and integrability of the Green's functions. Other two flexible methods are quite famous and widely used, even in the form of commercially available codes. These are the finite-element method, and the finite-difference time domain method. The flexibility with respect to the geometrical structure, optical index inhomogeneity and anisotropy, etc. has to be paid by the necessity of sophisticated meshing algorithms and very large sparse matrix manipulations.

These few examples represent only the top of the iceberg, and are invoked to illustrate the basic idea that the best method has not been invented, yet. Probably never.

We have tried to gather a team of specialists in rigorous theories of gratings in order to cover as large variety of methods and applications as practically possible. The last such effort dates quite long ago, and it has resulted in the famous Electromagnetic Theory of Gratings (ed. R. Petit, Springer, 1980), a book that has long served the community of researchers and optical engineers, but that is now out of press and requires a lot of update and upgrade, something that we hope to achieve, at least partially with this new book.

Our choice of electronic publishing is determined by the desire to ensure larger free access that is not easily available through printed editions. I want to thank all the contributors to this Edition. Special thanks are due to my colleagues Fréderic Forestier and Boris Gralak for the technical efforts to make the electronic publishing possible.

Marseille, France
Evgeny Popov
December 2012

Chapter 1:
Introduction to Diffraction Gratings: Summary of Applications

Evgeny Popov

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## Chapter 1

# Introduction to Diffraction Gratings: Summary of Applications 

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Periodic systems play an important role in science and technology. Moreover, desire for order in Nature and human society has accompanied development of philosophy. Simple periodical oscillation is referenced as 'harmonic' in mechanics, optics, music, etc., the name deriving from the Greek $\dot{\alpha} \rho \mu o v i ́ \alpha$ (harmonía), meaning "joint, agreement, concord" [1.1]. It is not our purpose here to study harmony in general, nor harmony in physics. We are aiming to much more modest target: rigorous methods of modeling light propagation and diffraction by periodic media.

The methods presented in the book have already shown their validity and use from x-ray domain to MW region, for nonmagnetic and magnetic materials, metals and dielectrics, linear and nonlinear optical effects. The existing variety of these methods is due not only to historical reasons, but mainly to the absence of The Method, a universal approach that could solve all diffraction problems. Some of the approaches cover greater domain of problems, but more specialized ones are generally more efficient. The other reason of the great number of methods is the complexity and variety of their objects and applications.

### 1.1. Diffraction property of periodic media

The most important property of diffraction grating to create diffraction orders has been documented by Rittenhause for the first time in 1786 [1.2] due to the observation made by Francis Hopkinson through a silk handkerchief. The appearance of diffraction orders rather than the specularly reflected and transmitted beams was studied experimentally by Young in 1803 [1.3] with his discovery of the sine rule. A detailed presentation of the analytic properties of gratings can be found in Chapter 2.

Secondary-school pupils are supposed nowadays to know the Snell-Descartes law, which undergraduate students in universities are supposed to be able to demonstrate: singleray diffraction on a plane interface results in a single transmitted and single reflected rays. The critical advantage of periodic perturbation of the interface (variation of the refractive index or surface corrugation) changes the impulsion (wavevector surface component $\overrightarrow{\mathrm{k}}_{\mathrm{S}, \mathrm{m}}$ ) of the incident wave $\overrightarrow{\mathrm{k}}_{\mathrm{S}, \mathrm{i}}$ along the surface by adding or subtracting an integer number of grating impulses (grating vectors) $\overrightarrow{\mathrm{K}}$ :

$$
\begin{equation*}
\overrightarrow{\mathrm{k}}_{\mathrm{S}, \mathrm{~m}}=\overrightarrow{\mathrm{k}}_{\mathrm{S}, \mathrm{i}}+\mathrm{m} \overrightarrow{\mathrm{~K}} \tag{1.1}
\end{equation*}
$$

where $\overrightarrow{\mathrm{K}}=\frac{2 \pi}{\mathrm{~d}} \hat{\mathrm{~d}}$ and d is the grating period in a unit-vector direction $\hat{\mathrm{d}}$.
If the interface lies in the xy -plane and the periodicity is along the x -axis (Fig.1.1), and the incidence lies in a plane perpendicular to the grooves, the equation in reflection takes the form of the so-called grating equation:

$$
\begin{equation*}
\sin \theta_{\mathrm{m}}=\sin \theta_{\mathrm{i}}+\mathrm{m} \frac{\lambda}{\mathrm{~d}} \tag{1.2}
\end{equation*}
$$

where $\lambda$ is the wavelength of light.


Fig.1.1. Lamellar grating working in in-plane regime in TM (transverse magnetic) polarization, together with the coordinate system, incident wavevector.

The case of conical (off-plane) diffraction by a plane grating having one-dimensional periodicity can also be described by eq.(1.1), preserving the wavevector component parallel to the groove direction:

$$
\begin{align*}
& \mathrm{k}_{\mathrm{y}, \mathrm{~m}}=\mathrm{k}_{\mathrm{y}, \mathrm{i}} \\
& \mathrm{k}_{\mathrm{x}, \mathrm{~m}}=\mathrm{k}_{\mathrm{x}, \mathrm{i}}+\mathrm{mK}  \tag{1.3}\\
& \mathrm{k}_{\mathrm{z}, \mathrm{~m}}=\sqrt{\mathrm{k}^{2}-\mathrm{k}_{\mathrm{x}, \mathrm{~m}}^{2}-\mathrm{k}_{\mathrm{y}, \mathrm{i}}^{2}}
\end{align*}
$$

where k is the wavenumber in the cladding. As a result, the diffracted beams lie on a cone, thus the name conical diffraction.

Two-dimensional periodicity (having grating vectors $\overrightarrow{\mathrm{K}}_{1}$ and $\overrightarrow{\mathrm{K}}_{2}$ ) imposed on the optogeometrical properties of the plane interface creates two sets of diffraction orders together with their spatial combinations, subjected to the same rule formulated just before eq.(1.1):

$$
\begin{equation*}
\overrightarrow{\mathrm{k}}_{\mathrm{S}, \mathrm{mn}}=\overrightarrow{\mathrm{k}}_{\mathrm{S}, \mathrm{i}}+\mathrm{m} \overrightarrow{\mathrm{~K}}_{1}+\mathrm{n} \overrightarrow{\mathrm{~K}}_{2} \tag{1.4}
\end{equation*}
$$

Pure three-dimensional (3D) periodicity appears in crystallography and photonic crystals, if the substance is assumed to fill the entire space. The third-direction periodicity imposes additional condition to the wavenumber (said more precisely, to $\mathrm{k}_{\mathrm{z}, \mathrm{m}}$, which is already defined by the wave equation as given in the third equation of (1.3)), which leads to creation of discrete modes propagating in 3D periodic structures. This also leads to the appearance of propagating and forbidden zones structure.

### 1.2. Classical gratings in spectroscopy

The most common application of diffraction gratings is due to the fact that outside of the specular order $(\mathrm{m}=0)$, the diffraction order direction depends on the wavelength, as stated in eq.(1.2). The result is that the grating acts as a dispersive optical component, with several advantages when compared to the prisms:

1. The grating can be a plane device, while the prism is a bulk one that requires larger volumes of optically pure glass (to add the difficulties of weight and temperature expansion constrains).
2. Provided a suitable reflecting material, the grating can work in spectral regions, where there is no transparent 'glass' with sufficient dispersion.
3. Grating dispersion can be varied, as it depends on the groove period, while prism dispersion depends on the material choice and groove angle, which gives quite limited choices.
Since the first works of Young, followed by Fraunhofer's quite serious attention to diffraction gratings use and properties [1.4], there is rarely more important device in spectroscopy achievements that lay the basis of modern physics. An interested reader can find some important aspects of their history, properties, and application in [1.5]. Despite the above-listed advantages compared to prisms, the diffraction gratings never have sufficient performance for their spectroscopy customers:
4. There is no enough diffraction efficiency, defined as the ratio of the incident light diffracted in the order used by the application.
This is probably the problem that has been mostly treated by rigorous grating methods, because they are the only ones to provide feasible results on the energy distribution, while the other characteristics (spectral resolution, scatter, dispersion, order overlap, etc.) can be obtained by simpler approaches.


Fig.1.2. Diffraction efficiency of aluminum-made sinusoidal grating with period $d=0.5 \mu \mathrm{~m}$ as a function of the wavelength. Littrow mount $\left(-1^{s t}\right.$ order diffracted in the direction of the incident beam) for the two fundamental polarizations (transverse electric TE and transverse magnetic TM).

A typical spectral dependence of the diffraction efficiency (defined more precisely as the ratio between the energy flow in the corresponding order and the energy flow in the incident order in z-direction) of a surface-relief sinusoidal grating made of aluminum and working in the $-1^{\text {st }}$ Littrow mount ${ }^{1}$ is presented in Fig.1.2. As observed, the problems of spectral variation of efficiency adds to its polarization dependence.
2. There is no enough resolution, defined as the ratio between the working wavelength and the smallest distinguishable (as usual, defined using the Rayleigh criterion) spectral interval:

[^0]\[

$$
\begin{equation*}
\mathrm{R}=\frac{\lambda}{\Delta \lambda} \tag{1.5}
\end{equation*}
$$

\]

Classical diffraction theories show that the spectral resolution is proportional to the number of grooves illuminated by the incident beam, and inversely proportional to $\cos \left(\theta_{\mathrm{i}}\right)$, if we consider 1D periodicity in non-conical diffraction (speaking more precisely, the inverse proportionality implies to the sum of the cosines of the incident and the diffracted angles). The latter dependence gives advantages to grazing incidence for higher spectral resolution applications, namely astronomy.

## 3. The efficiency depends on the polarization.

Laser resonators usually use Brewster windows, so that the requirements are for high-efficient grating working in TM (transverse magnetic) polarization. However, spectroscopic applications do not like this at all. In astronomy, this property can be quite costly, because the loss of a half of the incident light intensity requires twice the exposer time. Low polarization dependent losses (PDL) are one of the most important criterions in optical communications, in general, and in grating applications for wavelength demultiplexing, in particular, necessary for multichannel optical connections. Fortunately, contrary to stellar spectroscopy, gratings used in optical communications work in only very small spectral interval, and are used in quite smaller sizes, so that there exist several solutions that provide high efficiency in unpolarized light over a limited spectral region. The idea is to shift the maxima in the spectral dependence of the two polarizations in Fig.1.2 in order to make them overlap at some required wavelength. One solution is to use a grating having two-dimensional (2D) periodicity, as shown further on in Fig.7.1. The period in the perpendicular direction is sufficiently small as not to introduce additional diffraction orders, and this additional corrugation can shift the position of the TE maximum to longer wavelengths [1.6].


Fig.1.3. SEM picture of the profile of a grating etched in Si wafer and used for wavelength demutiplexing (after [1.7], with the publisher's permission)

Another more conventional solution is to use a classical grating with a 1D periodicity, but having a deformed profile in order to perform the same shift of the TE maximum [1.7, 8]. This can be achieved by introducing a flat region at the bottom of the grooves and "sharpening" the groove triangle, which needs a sharper apex angle. While this is quite difficult to be made with grating ruling or holographic recording, etching in crystalline silicon naturally produces grooves with $70.5^{\circ}$ apex angle, as observed in Fig.1.3. A flat region on the top is made if the etching is not complete. When covered with gold, such grating can be used from the silicon side, which is transparent at wavelengths around $1.55 \mu \mathrm{~m}$. Unpolarized
efficiency greater than $80 \%$ can be kept over the communication interval of $50-60 \mathrm{~nm}$, as observed in Fig.1.4.


Fig.1.4. Spectral dependence of efficiency in $-1^{\text {st }}$ order for the grating shown in Fig.1.3 (after [1.7], with the publisher's permission).

## 4. There is an overlap of diffraction orders.

As stated by the grating equation, if the period is chosen to provide diffraction orders inside a large spectral interval, the second diffracted order has the same direction of propagation as the first one for wavelength of light twice shorter. This problem is avoided in many spectroscopic devices by three different methods: additional spectral filtering of undesired orders; interchange of grating having different spatial frequency (period); adding cross-dispersion grating for echelle applications.

In addition, there is a contradiction between some requirements as, for example, free spectral range (spectrum covered by the grating), dispersion and overlap of orders. In some cases, the only compromise is to use interchangeable gratings in order to cover larger spectral range without order overlap and maintaining higher dispersion. Some spectrographs are designed having multiple channels with splitting of the incident beam between them, however reducing the illumination power in each channel.

### 1.3. Echelle gratings in astronomy

In astronomy, measurements of very low signals coming from far cosmic objects require large periods of time, so that the loss of energy due to low efficiency or/and strong polarization dependence leads to a further growth of costs. When efficiency constraints are added to the independence of the polarization, the only known solution is the echelle grating (high grooveangle triangular groove profile used in grazing incidence, Fig.1.5) that has the advantages of almost equal and high efficiency in both fundamental polarizations [1.9].


Fig.1.5. Schematical presentation of echelle grating.

The high efficiency in unpolarized light is obtained when the diffraction is made as if light is reflected by the working facet in its normal direction. Of course, it is necessary that the reflections at consecutive facets are in phase. The problem is that the efficiency varies rapidly with the wavelength, and the maxima switch between consecutive orders (Fig.1.6). The separation of orders is usually made using another shorter-period grating with grooves perpendicular to the echelle (called cross-dispersion), so that the different orders are separated in direction perpendicular to the echelle dispersion direction.


Fig.1.6. Diffraction efficiency of an echelle made of aluminum with 31.6 grooves/mm and $64^{\circ}$ groove angle of the working facet. Incident angle is equal to $64^{\circ}$. The numbers of the diffraction orders are indicated in the figure.

Echelle gratings are also used in transmission, glued to the hypotenuse face of a prism that has a role to deviate back the beam diffracted by the grating, so that the principal diffracted order of the device propagates almost in the same direction as the incident beam for a chosen central wavelength. The device is known as a Carpenter prism or GRISM and gives the possibility to convert an imaging device (camera) into a long slit spectrograph [1.10], commonly used in airborne or space borne scientific missions.

UV excimer lasers used in photolithography at 193.3 nm wavelength can be equipped with an echelle grating for narrowing the spectral line. While in lasers working in the visible and IR, the resonators are equipped with diffraction gratings with symmetrical grooves working in $-1^{\text {st }}$ order than easily can be made holographically or lithographically, a grating working in the lowest order at 193.3 nm must have more than $9000 \mathrm{gr} / \mathrm{mm}$, very difficult for fabrication and impossible for replication. Echelles take longer to be made and are more expensive as they require mechanical ruling engines with temperature control and clean environment, but have large periods and can be replicated from the ruled masters and submasters to become available at acceptable prices.

### 1.4. Gratings as optical filters

Spectroscopic applications of gratings use one or more diffraction orders that differ from the specular reflected and transmitted ones, because of the required spectral dependence. There exist, however, several applications that use the zero order(s), even with corrugation of the
surface or modulation of the refractive index in the grating region. These are devices having refractive or reflection properties in the zeroth order that are modified by the grating surface.

### 1.4.1. Zero-order diffraction (ZOD) imaging

One of the applications uses the diffraction in higher orders to change the spectral dependency in the zero order, as light that is diffracted in the higher order(s) is absent in the zero order. The structure known as zero-order diffraction (ZOD) microimage [1.11] represents a transmission grating (usually with a rectangular or triangular groove form that can be replicated by relief printing in plastic sheet. Appropriately choosing the groove depth, one obtains broad-band color filters in transmission by using non-absorbing materials without colorants that can bleach.

Another type of gratings have periods, smaller than the wavelength in the substrate and the cladding, chosen to avoid the propagation of other that the zero orders. Such structures are known as subwavelength gratings, and they play important role in integrated optical devices and recently in plasmonics to transfer energy from one to another guided mode in dielectric or metallic waveguides, or to change the mode direction, or to focus guided light (see Section 1.5).

### 1.4.2. Surface/guided mode excitation

A subwavelength grating can serve to couple the incident light to surface or cavity resonances that can exist in the grating structure. The absence of higher orders means only that they are evanescent rather than propagating in the cladding or in the substrate. The horizontal component of their propagation constant ( $\mathrm{say}, \mathrm{k}_{\mathrm{x}}$ ) is larger than the wavenumber in the surrounding media, and thus it can excite a surface or waveguide mode that is subjected to the same requirement in order to stay confined to the surface. The grating action is the same as in the coupling between the incident wave and one of the diffraction orders, only that now $\mathrm{k}_{\mathrm{x},-1}$ is equal to the real part of the mode propagation constant $\mathrm{k}_{\mathrm{g}}$ :

$$
\begin{equation*}
\operatorname{Re}\left(\mathrm{k}_{\mathrm{g}}\right)=\mathrm{k}_{\mathrm{x}, \mathrm{i}}-\mathrm{K} \tag{1.6}
\end{equation*}
$$

To fail to see the quite small difference between $\mathrm{k}_{\mathrm{g}}$ and $\mathrm{k}_{0}$ in the case of surface plasmons on highly conducting metal surfaces is one of the rare failures of Lord Rayleigh [1.12] when trying to explain Wood's anomalies [1.13]. The case when $\mathrm{k}_{\mathrm{x}, \mathrm{m}}=\mathrm{k}_{0}$ represents a transition of the m-th diffracted order between a propagating and an evanescent type, as follows from eq.(1.3). This transition leads to a redistribution of energy between the propagating orders, observed in efficiency behavior as a phenomenon called cut-off anomaly.

Lord Rayleigh attributed Wood's anomalies to the cut-off of higher orders, whereas the true explanation is that Wood anomaly is due to surface plasmon excitation. It is easy to judge nowadays, but the difference in the spectral and angular positions of the cut-off and surface plasmon anomaly sometimes is smaller than the experimental error, or more important, the error in the knowledge of the grating period. The resonant nature of the surface-plasmon anomaly has much more pronounced features than the cut-off anomaly, a fact first noticed by Fano [1.14] and further developed by Hessel and Oliner [1.15]. The form of the so-called Fano-type anomaly can easily be derived from the interference with a non-resonant contribution (for example, non-resonant reflection by the grating layer) and a resonant effect (for example, excitation by the incident wave of a surface mode that is diffracted back into the direction of the non-resonant wave), as sketched in Fig.1.7.


Fig.1.7. Process of interference between the non-resonant and the resonant reflections (second term in eq.(1.7)) that is created due to the excitation of the surface or guided wave through the $-1^{s t}$ grating order, and then radiated into the cladding through the $+1^{\text {st }}$ diffraction order.

Here, by mode, we mean a surface or guided wave that represents an eigen (proper) solution of the homogeneous diffraction problem (assuming non-zero diffracted field without incident wave). The latter effect is commonly described in physics as having a Lorentzian character (e.g. electric circuit dipole oscillator). For example, the wavelength dependence of the total reflectivity r will be the sum of the non-resonant and the resonant Lorentzian contributions:

$$
\begin{equation*}
\mathrm{r}=\mathrm{r}_{\text {n.res. }}+\frac{\mathrm{c}_{\text {res. }}}{\mathrm{k}_{\mathrm{x}, \mathrm{i}}-\mathrm{k}_{\mathrm{x}}^{\mathrm{p}}} \tag{1.7}
\end{equation*}
$$

where the coefficients $r_{\text {n.res. }}$ and $c_{\text {res. }}$ are slowly varying functions of the incident wave parameters. The pole $\mathrm{k}_{\mathrm{x}}^{\mathrm{p}}$ of the resonant term is equal to the surface/guided wave propagation constant. The Fano-type equation is obtained by taking the common denominator in eq.(1.7):

$$
\begin{equation*}
\mathrm{r}=\mathrm{r}_{\text {n. } \mathrm{res}} \frac{\mathrm{k}_{\mathrm{x}, \mathrm{i}}-\mathrm{k}_{\mathrm{x}}^{\mathrm{z}}}{\mathrm{k}_{\mathrm{x}, \mathrm{i}}-\mathrm{k}_{\mathrm{x}}^{\mathrm{p}}} \tag{1.8}
\end{equation*}
$$

with $k_{x}^{z}=k_{x}^{p}-c_{\text {res. }} / r_{\text {n.res. }}$. representing a zero of the reflectivity. The formula implies that for each resonance, there exist an associated zero, that is expressed as a minimum of the anomaly. In reality, the pole has always non-zero imaginary part due to the interaction between the incident and the surface/guided wave. The zero takes, in general, complex values, but there are several important practically cases when the zero could become (and remains) real, i.e., the reflectivity can become zero. The imaginary part of the pole determines the quality factor (width) of the resonant maximum.

The same reasoning applies in transmission (and in any other existing propagating order), with the same pole (same resonance), but different zeros. Depending on the imaginary parts of the pole and the zero, sometimes the resonant anomaly can show itself as a Lorentzian maximum, sometimes as a pure minimum on otherwise highly-reflecting background, sometimes both maximum and minimum can manifest themselves. The important cases when the zeros are almost real are used in surface plasmon absorption detectors and in resonant dielectric grating filters (see the next two subsections).

### 1.4.3. Surface plasmon absorption detector

Wood anomaly can sometimes lead to a total absorption of light by shallow metallic gratings with groove depth that does not exceed $10 \%$ of the wavelength. This phenomenon was called Brewster effect in metallic gratings [1.16] and has found an important application in chemical
and biochemical surface-plasmon grating detectors [1.17]. The effect of total light absorption appears in a narrow spectral and angular interval (Fig.1.8).


Fig.1.8. Reflectivity of an optical surface plasmon detector as a function of the incident angle. Wavelength equals $850 \mathrm{~nm}, T M$ polarization. Cladding is glass, the substrate index for the two curves is indicated in the figure, and the grating layer is made of silver with thickness of 40 nm . The grating profile function contains two Fourier harmonics: $35 \sin (2 \pi \mathrm{x} / \mathrm{d})+59 \sin (4 \pi \mathrm{x} / \mathrm{d}+\pi / 2)$ with the depths given in nanometers (after [1.18], with the permission of the publisher).

The position of the anomaly depends on the propagation constant of the plasmon surface wave that is quite sensible to the variation of the refractive index of the surrounding dielectric, thus the determination of the position of the anomaly brings information about the composition of the surrounding medium, i.e., serves as an optical detector (Fig.1.8). Introduction of the second Fourier component of the profile function increases the direct interaction between the plasmon surface waves propagating in opposite directions (the grating vector of the second harmonic is twice longer that of the first one), interaction that increases the sensitivity of the device.

### 1.4.4. Resonant dielectric filters

There is a particular case when the maximum in the reflectivity due to the resonant waveguide mode excitation stays theoretically at $100 \%$ on a low-reflective non-resonant background value. This is the case of corrugated dielectric waveguides having symmetrical grooves. If, in addition, the substrate is identical to the cladding, the zero in reflection remains real whatever the other parameters, i.e. a $100 \%$ maximum is accompanied by a 0 value minimum in the reflectivity (Fig.1.9). This peculiarity was accidentally found in 1983 [1.19, 20] followed by a theoretical explication in 1984 [1.21]. The effect has been independently rediscovered in 1989 by Magnusson [1.22], who has done a lot for its analysis, and quite important practically, for its extension in transmission [1.23]. The advantage of the device is that the quality factor increases with the decrease of the groove depth, i.e. very narrow-line spectral filters can be obtained using shallow gratings, at least theoretically. However, as usual, the advantage is paid back somewhere, namely in the tight angular tolerances: when eq.(1.6) results in high spectral sensibility, it also is responsible for strong angular sensibility. Sentenac and Fehrembach [1.24] proposed to introduce a direct coupling between the waveguide modes propagating in the opposite direction (as was done in Fig.1.8), but this time aiming to a significant reduction of the angular sensibility of the effect.


Fig.1.9. Reflection by a corrugated dielectric waveguide. Symmetrical triangular profile with groove angle equal to $10^{\circ}$. Substrate and cladding have optical index equal tol, the layer has index 2.3 and its thickness is equal to 69 nm . Incident angle is equal to $26.7^{\circ}$. Excitation of TE waveguide mode leads to a narrow spectral anomaly.

The idea was based on the flattening of the angular dependence of the guided wave propagation constant on the boundaries of the forbidden gaps, created by the direct mode coupling due to the structure periodicity, one of the basic properties of photonic (and electronic) crystals.

### 1.4.5. Enhanced transmission through hole arrays in metallic screens

In 1998 Ebbesen et al. [1.25] reported an interesting effect on metallic screen perforated with circular holes (Fig.1.10a). They were surprised to observe peaks with relatively strong transmission in the spectral dependence (Fig.1.11).


Fig.1.10. (a) Schematical representation and notations of a two-dimensional hole array perforated in a metallic screen deposited on a glass substrate and illuminated from above with linearly polarized incident wave. (b) Surface-plasmon assisted energy flow inside the aperture close to the resonance.


Fig.1.11. Spectral dependence of the transmission of the structure presented in Fig.1.10a, with $d=$ $0.9 \mu \mathrm{~m}, t=0.2 \mu \mathrm{~m}$, and hole diameter of $0.2 \mu \mathrm{~m}$ (after [1.25], with the publisher's permission).

These peaks were identified as resulting from surface plasmon excitation on the upper, or the lower surface. Many theoretical and experimental works appeared as a result of this discovery, which revived the interest to surface plasmons and grating structure, leading to a new name of the plasmon studies called now plasmonics.

It seems nowadays that the explanation of the enhanced transmission observed in Fig.1.11 lies in the common action of two phenomena [1.26]. The first one is the surface plasmon excitation that enhances the electromagnetic field intensity near the entrance apertures, as represented schematically in Fig.1.10b. The second phenomenon is the tunneling of the mode in the vertical hollow waveguides inside the holes. Although for small holes even the lowest mode is evanescent (contrary to 1D lamellar gratings, where TEM mode exists without cut-off), it gives the possibility of energy transfer from the upper to the lower interface, much more efficiently than the tunneling through the non-perforated screen [1.27].

One unexpected consequence of the observation of Ebbesen et al. came in fluorescence single-molecule microscopy and in the biomembrane studies. Both require small measuring volumes in order to study only a very small number of molecules or a small portion of the membrane surface. However, small measuring volumes mean week signals. Here comes the role of the surface plasmon enhanced field and evanescent mode inside the aperture. When the mode is close to its cut-off, the real part of its propagation constant along the axis of the hole tends to zero, which leads to a compression of the electromagnetic field at the entrance aperture, increasing even more the field density and thus the measured signal. In addition, as the field is evanescent inside the hole, the effective measuring volume does not extend to the entire hole depth (see Fig.1.12). The small cross-section of the holes reduces further on the measuring volume to much smaller values than permitted by the diffraction limit, in several cases volumes of the order of several attolitres have been reported.


Fig.1.12. Schematic presentation of (a) single-molecule fluorescence backside microscopy field density, and (b) tiny-portion cell membrane microscopy inside a metallic aperture (after [1.28] with the publisher's permission).

### 1.4.6. Non-resonant filters

Resonant filtering has its advantages, when narrow spectral bands are aimed, but in some important applications it is necessary to have wider bands, accompanied by angular and polarization invariance with respect to the incident wave. Such are the requirements for color filters used to separate the RGB colors on each pixel of the CCD cameras. A promising example [1.29] contains a small subwavelength grating consisting of circular metallic bumps with different diameters and heighst (Fig.1.13a) that can filter light in different spectral regions, depending on their geometrical parameters, (Fig.1.13b).


Fig.1.13. (a) Nanostucture consisting of coaxial metallic cylinders, with a subwavelength period equal to 300 nm , external diameter 260 nm , and internal diameter 160 nm . (b) Transmission spectrum of the structure for different metals and cylinder heights, as shown in the inset (after [1.29] with the publisher's permission).

### 1.4.7. Flying natural gratings: butterflies, cicadas

The observation of Hopkinson (see Sec.1.1) of the diffraction on a handkerchief is far the less exotic grating that exists. As always, Nature had all the time to surpass humanity. A striking realization has been developed during the million-year long evolution of butterfly wings that
have so attractive coloring. The so called Morpho rhetenor butterflies have a non-pigment metallic blue color (Fig.1.14) that long has been attributed to multilayer reflection. However, Vukusic et al. [1.30] had the idea (and funds) to use an electron microscope to observe deeper in detail the structure of the scales, as reported in Fig.1.15a, a typical 3D structure that resembles a photonic crystal. The modeling with a 2D grating with a structure given in Fig.1.15b confirms the blue-spectrum reflection of the scales (Fig.1.16).


Fig.1.14. Entire view (to the left) and magnification of the scales (to the right) of a Morpho rhetenor butterfly (after [1.30] with the publisher's permission).


Fig.1.15. (a) Transmission electron microscope image showing the cross-section through a single Morpho rhetenor scale (after [1.30] with the publisher's permission). (b) Modeled structure; the two red lines define a grating layer, the optical index of white regions is 1 (after [1.31] with the publisher's permission).


Fig.1.16. Spectral dependence of the reflection of the butterfly scale (after [1.31], with the publisher's permission).

Another "application" is used by the cicadas (quite common in the southern Europe, see Fig.1.17a) to camouflage themselves on the tree branches. Their wings are covered with an anti-reflection nanostructure, first observed by Xie et al. [1.32] and shown in Fig.1.17b. Anyone that has entered the microwave measuring rooms can identify the cones on the walls that are about 1 million times larger, as scaled to the wavelength.


Fig.1.17. (a) A photo of a cicada, and (b) the nanostructure pattern on its wings (after [1.32] with the publisher's permission).

### 1.5. Gratings in Integrated optics and plasmonic devices

Gratings are used in integrated optical devices to deviate the direction of propagation of waveguide modes and surface waves, for their focusing inside or outside the guide, or for energy transfer between different modes. Let us consider the case with one-dimensional periodicity. The grating equation can be applied not only to free-space waves, but to the waveguides modes. If the period is suitably chosen, it is possible to couple one mode to another:

$$
\begin{equation*}
\operatorname{Re}\left(\mathrm{k}_{\mathrm{g}, 1}\right)=\operatorname{Re}\left(\mathrm{k}_{\mathrm{g}, 2}\right)-\mathrm{K} \tag{1.9}
\end{equation*}
$$

where $\mathrm{k}_{\mathrm{g}, 1}$ and $\mathrm{k}_{\mathrm{g}, 2}$ are the propagation constants of the modes.


Fig.1.18. Upper part: schematical presentation of Bragg relief type lamellar grating deposited on a waveguide (dielectric or plasmonic) with two propagating modes. Lower part: energy carried by each mode.

It is possible to couple the mode propagating in a given direction to the same but contrapropagative mode. This is the case of the so-called Bragg gratings that act as a distributed mirror forming a forbidden zone for the mode propagation, in which the mode field decreases exponentially without being radiated in the cladding and in the substrate. The result is that it is rejected back into a contra-propagative direction (Fig.1.18). Due to the limit size of the grating region, a small part of the incident mode 1 is transmitted to the right, thus the reflected mode 2 carries smaller amount of energy. The grating grooves can be made curvilinear in order to focus the mode. The same effect can be obtained by replacing the 1D structure by 2D periodicity having also a period in the transversal y-direction (see Section 1.8).

### 1.6. Beam-splitting applications

The fact that the periodicity creates diffraction orders can be used to create multiple beams from a single laser beam. The most-commonly used device is symmetrical groove transmission gratings used as beam splitters for optical disk readers, where the wavelength of the laser source is constant. As a rule, the zero order beam reads the track and the two first order beams read adjacent tracks to keep the head both centered and focused. By controlling the groove depth, the ratio of zero to first orders transmission can be varied over a factor of 10 , and a high degree of symmetry is inherent.

A special type of transmission grating can be used to generate an entire family of orders with a groove shape designed to make their intensities as equal as possible. This gives us multiple beam splitters, which may have 5 or even 20 orders on both sides of zero. A top view of such gratings under working conditions (with a laser input) gives rise to the term of "fanout gratings". Applications are found in scanning reference planes for construction use, optical computing, and others [1.5, 33]. Such gratings tend to have large groove spacings ( 10 to $100 \mu \mathrm{~m}$ ) and low depth modulations. The difficulty in making such gratings lies in achieving a groove shape that leads to a sufficient degree of efficiency uniformity among orders, especially if they are to function over a finite wavelength range. Two obvious candidates are cylindrical sections or an approximation of this shape in the form of a wide angle V with several segments of different angles. The choice may vary with the availability of the corresponding diamond tools. They have also been made by holographic methods, which are able to produce the parabolic groove form that gives the best energy uniformity
between the diffracted orders [1.34, 35]. Applications of fan-out gratings are found in scanning reference planes for construction use, in biophysics for simultaneous treatment of great number of samples, etc.

By combining two such grating at right angles to each other, an accurately defined 2D array of laser beams is generated to be used for calibrating the image field distortion of large precision lenses, in robotic vision systems, or in parallel optical computing. Instead of using two 1D periodical grating, it is possible to use a single large-period 2D crossed grating with specially optimized pattern inside each period [1.36].

### 1.7. Subwavelength gratings for photovoltaic applications

As explained in Sec.1.4, resonant excitation in metallic gratings can lead to a total absorption of incident light, effect necessary for the efficient work of photovoltaic devices. The problem with surface plasmon excitation and the accompanying light absorption is that its wave is not localized, thus it is characterized by a well-determined value of propagation constant along the surface, because non-local effects in the real space are localized in the inverse space. This is why the surface plasmon anomalies have very narrow angular and spectral width, an advantage in detector construction, but failing in photovoltaics. Evidently, effects that are less localized in the inverse space will be more localized in the real space. This leads us to cavity resonances, volume plasmonic excitation, and surface plasmons that propagate in the vertical direction but are localized in x-y direction. Cavity resonances in deep grooves or in closed cavities, like embedded dielectric spheres or cylinders inside a metallic sheet can absorb light within relatively large angular region (20-30 deg) [1.37-40].


Fig.1.19. Crossed metallic diffraction grating


Fig.1.20. Real and imaginary part of the effective refractive index of the structure of Fig.1.19 as a function of the filling ratio when the period is much shorter than the wavelength $\lambda=457 \mathrm{~nm}$.

However, the problem of cavity resonances is that they are strongly wavelengthsensitive. An alternative approach consists of using metamaterial behavior of small-feature structures. By mixing metallic and dielectric materials one can, in general, obtain strongly absorbing alloys with effective refractive index that does not exist for known materials. In addition, the equivalent metamaterial layer has a uniaxial anisotropy with axis perpendicular to the grating plane. Thus inside the xOy plane it has isotropic properties and its response is polarizationally independent, at least close to normal incidence.

For example, a crossed channel grating as presented in Fig.1.19 and made of silver bumps on a silver substrate. It can strongly absorb the incident light in much larger spectral
domain when compared with the surface plasmon excitation effects. As can be observed in Fig. 1.20, close to a filling ratio ( $\mathrm{f}=\mathrm{c}^{2} / \mathrm{d}^{2}$ ) equal to 0.7 , both the real and the imaginary part of the effective refractive index for small-period structure grow significantly to values that no existing material has in this spectral domain. This increases the effective optical thickness of the system, together with its absorption, so that a very thin grating ( $\mathrm{h}<10 \mathrm{~nm}$ ) can totally absorbed incident light [1.41]. Because of the 2D periodicity of the structure, it becomes polarization insensible. Moreover, this effect has no resonant nature and is extended angularly to almost the entire set of angles of incidence, as seen in Fig.1.21. In addition, the spectral domain of absorption stronger than $75 \%$ extends to a more than 100 nm interval [1.42].


Fig. 1.21. Reflectivity of the crossed graing in Fig.1.19 as a function of the wavelength and of the sine of the incident angle. Period $d=3$ nm, system thickness $h=6.2 \mathrm{~nm}$, filling factor $=0.7$, i.e. $c / d=$ 0.837.

However, the use of such structures for practical applications requires not only a larger spectrum of strong absorption, but also needs that the absorbed light is not immediately transformed into heat and that it is kept inside the photovoltaic layer that transforms light into electricity. The simple structure represented in Fig.1.19 is not sufficient to efficiently fulfill these purposes. It has to be complicated by a structuring in the depth of the device (Fig.1.22). Light is incident from the top of the device. Its entering part (in blue) is designed to act as an antireflector, and absorbs most of the blue part of the spectrum. The green part is absorbed in the middle of the rods, while the back-side part (in red) absorbs the longer wavelengths, and serves as a reflector to prevent the unabsorbed part to be lost in the substrate [1.43].

Fig.1.22. Grating rods as optimal photovoltaic absorber (after [1.43], with the publisher's permission)

### 1.8. Photonic crystals

Periodic structures in optics can serve for the photons in the same manner as semiconductor crystals for electrons. This common feature led in the ' 90 s to call such structures photonic crystals. As discovered by Yablonovich [1.44, 45], they present band-gaps that forbid propagation and thus guaranteeing $100 \%$ reflection inside the band. While this property is widely known and largely used in multilayer dielectric mirrors (Fig.1.23a), the band gap of 1D photonic crystals is limited in relatively smaller angular interval. Some other structures that have 2D periodicity combined with a nanostructuring in the third dimensions can be found in Figs.1.15, 17, and 22.


Fig.1.23. Schematic representation of (a) one-, (b) two- and (c, d) three-dimensional photonic crystals.


Fig.1.24. Forbidden bands for the photonic crystalmade of circular cylinders with $d=1.414 \mu \mathrm{~m}, r$ $=0.35 \mu \mathrm{~m}$ (after [1.46] with the publisher's permission).

A typical band-gap structure of a 2D photonic crystal presented in Fig.1.23b, but its surface cut at an angle of $45^{\circ}$ with respect to the vertical direction, is given in Fig. 1.24 for a system having period $\mathrm{d}=1.414 \mu \mathrm{~m}$ in both directions, and consists of circular cylindrical rods with radius equal to $0.35 \mu \mathrm{~m}$ and optical index of 2.9833 in air as a matrix. The figure presents the values of the smallest imaginary part of $\mathrm{k}_{\mathrm{z}}$. As can be seen, a forbidden gap in both TE and TM polarization exists for $\lambda \in[1.97,2.33 \mu \mathrm{~m}]$ [1.46, 47]. Inside the band gap electromagnetic field intensity diminishes exponentially and the entire incident light is reflected back. However, the choice of the ratio of the wavelength and the period allows for the propagation of the $-1^{\text {st }}$ order in reflection. This gives the possibility to guide the entire incident light into this order, thus perfect blazing in the $-1^{\text {st }}$ diffracted order can be obtained in both polarizations, a property that is strongly desirable in many applications. And indeed, Fig. 1.25 presents the diffraction efficiency of the system having cylinders with radii of 350 nm (a) and 150 nm (b). A well-defined spectral region with almost $100 \%$ efficiency in unpolarized light can be observed.


Fig.1.25. Diffraction efficiency in order -1 as a function of the wavelength lying inside the band gap. (a) $r=350 \mathrm{~nm}$, (b) $r=150 \mathrm{~nm}$ (after [1.46] with the publisher's permission).

The property of totally reflecting the incident light whatever the direction, can be of great importance for light guiding and manipulation. Light confinement and guiding in a single dimension is ensured by using planar waveguides. Channel waveguides and optical fibers confine light in two dimensions, but they suffer from two important limitations: dispersion and bending losses. High bending angles damage the guiding properties and lead to radiation losses. A waveguide constructed with photonic crystal walls can ensure bending, Fig.1.26, without losses even at $90^{\circ}$, as predicted numerically [1.48].

It is impossible even only to list here the potential applications of photonic crystal devices, as for example negative refraction, perfect lenses construction, photonic crystal fibers, nonlinear optical applications. The main problem that persists is purely technological: while it is relatively easy to fabricate 3D periodically structures working in the microwave and far-IR domain, scaling down to the visible and the near-IR presents a lot of challenges to optical industry. Fortunately, the process of fiber manufacturing enables literally such mechanical scaling of the dimension, transferring the initial large-diameter preform into a thin fiber by pulling it. If the preform is carefully drilled with macroscopic holes, they are preserved in the final fibers, but scaled to nanodimensions.

In the resulting photonic crystal fiber (Fig.1.27), light is preserved in the central hollow guide by repulsion from the surrounding structure that presents a forbidden gap for the working wavelength.


Fig.1.26. Light guiding in a cavity formed by photonic crystal walls consisting of cylindrical objects


Fig.1.27. Schematic representation of a portion of an optical fiber - photonic crystal hybrid structure. Small cylindrical holes run along the fiber length. A central hole (shown with a dashed line inside the fiber) serves as an energy propagator, which ensured low dispersion, low absorption losses and high damage threshold

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Chapter 2:

## Analytic Properties of Diffraction Gratings <br> Daniel Maystre

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## Chapter 2

# Analytic Properties of Diffraction Gratings 

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### 2.1 Introduction

Since the 80 's, specialists of gratings can rely on very powerful grating softwares [1-6]. These softwares are able to compute grating efficiencies for almost any kind of grating in any domain of wavelength, even though the progress of grating technologies needs endless extensions of grating theories to new kinds of structures. These softwares are based on elementary laws of Electromagnetics. Using mathematics, these laws lead to boundary value problems which can be solved on computers using adequate algorithms.

However, a grating user should not ignore some general properties of gratings which can derived directly from the boundary value problem without any use of computer. These analytic properties are valuable at least for two reasons. First, they strongly contribute to a better understanding of an instrument which puzzled and fascinated many specialists of Optics since the beginning of the 20th century. Secondly, they allow a theoretician to check the validity of a new theory or its numerical implementation, although one must be very cautious: a theory can fail while its results satisfy some analytic rules. Specially, this surprising remark apply to properties like energy balance or reciprocity theorem.

The first part of this chapter is devoted to the use of the elementary laws of Electromagnetics for stating the boundary value problems of gratings in various cases of materials and polarizations. Then, we deduce from the boundary value problems the most important analytic properties of gratings.

### 2.2 From the laws of Electromagnetics to the boundary-value problems

### 2.2.1 Presentation of the grating problem

Figure 2.1 represents a diffraction grating. Its periodic profile $\mathscr{P}$ of period $d$ along the $x$ axis separates air (region $\mathscr{R}_{0}$ ) from a grating material (region $\mathscr{R}_{1}$ ) which is generally a metal or a dielectric. The $y$ axis is the axis of invariance of the structure and the $z$ axis is perpendicular to the average profile plane. We denote by $z_{M}$ the ordinate of the top of $\mathscr{P}$, its bottom being located


Figure 2.1: Notations.
on the $x y$ plane by hypothesis. We suppose that the incident light can be described by a sum of monochromatic radiations of different frequencies. Each of these can in turn be described in a time-harmonic regime, which allows us to use the complex notation (with an $\exp (-i \omega t)$ time-dependence). In this chapter, we assume that the wave-vector of each monochromatic radiation lies in the cross-section of the grating ( $x z$ plane). In the following, we deal with a single monochromatic radiation.

The electromagnetic properties of the grating material (assumed to be non-magnetic) are represented by its complex refractive index $v$ which depends on the wavelength $\lambda=2 \pi c / \omega$ in vacuum ( $c=1 / \sqrt{\varepsilon_{0} \mu_{0}}$ being the speed of light, with $\varepsilon_{0}$ and $\mu_{0}$ the permittivity and the permeability of vacuum). This complex index respectively includes the conductivity (for metals) and/or the losses (for lossy dielectrics). It becomes a real number for lossless dielectrics.

In the air region, the grating is illuminated by an incident plane wave. The incident electric field $\overrightarrow{E^{i}}$ is given by :

$$
\begin{equation*}
\overrightarrow{E^{i}}=\vec{P} \exp \left(i k_{0} x \sin (\theta)-i k_{0} z \cos (\theta)\right) \tag{2.1}
\end{equation*}
$$

with $\theta$ being the angle of incidence, from the $z$ axis to the incident direction, measured in the counterclockwise sense, and $k_{0}$ being the wavenumber in the air ( $k_{0}=2 \pi / \lambda$, we take an index equal to unity for air). The wave-vector of the incident wave is given by:

$$
\overrightarrow{k_{0}^{i}}=\left[\begin{array}{c}
k_{0} \sin (\theta)  \tag{2.2}\\
0 \\
-k_{0} \cos (\theta)
\end{array}\right] .
$$

The physical problem is to find the total electric and magnetic fields $\vec{E}$ and $\vec{H}$ at any point of space.

### 2.2.2 Maxwell's equations

First, let us notice that the physical problem remains unchanged after translations of the grating or of the incident wave along the $y$ axis since they do not depend on $y$. Therefore, if $\vec{E}(x, y, z)$ and $\vec{H}(x, y, z)$ are the total fields for a given grating and a given incident wave, $\vec{E}\left(x, y+y_{0}, z\right)$ and $\vec{H}\left(x, y+y_{0}, z\right)$ will be solutions too, regardless of the value of $y_{0}$. Assuming, from the physical intuition, that the solution of the grating problem is unique, we deduce that $\vec{E}$ and $\vec{H}$ are independent of $y$.

In order to state the mathematical problem, we use the harmonic Maxwell equations in $\mathscr{R}_{0}$ :

$$
\begin{align*}
& \nabla \times \vec{E}=i \omega \mu_{0} \vec{H}  \tag{2.3}\\
& \nabla \times \vec{H}=-i \omega \varepsilon \vec{E} \tag{2.4}
\end{align*}
$$

with:

$$
\varepsilon= \begin{cases}\varepsilon_{0} & \text { in } \mathscr{R}_{0},  \tag{2.5}\\ \varepsilon_{1}=\varepsilon_{0} v^{2} & \text { in } \mathscr{R}_{1} .\end{cases}
$$

In the following, equations (2.3) and (2.4) will be called first and second Maxwell equations respectively. We note that Maxwell's equations $\nabla \cdot \vec{E}=0$ and $\nabla \cdot \vec{H}=0$ are the straightforward consequences of the first and second Maxwell equations (it suffices to take the divergence of both members).

We introduce the diffracted fields $\overrightarrow{E^{d}}$ and $\overrightarrow{H^{d}}$ defined by:

$$
\begin{align*}
& \overrightarrow{E^{d}}= \begin{cases}\vec{E}-\overrightarrow{E^{i}} & \text { in } \mathscr{R}_{0}, \\
\vec{E} & \text { in } \mathscr{R}_{1},\end{cases}  \tag{2.6}\\
& \overrightarrow{H^{d}}= \begin{cases}\vec{H}-\overrightarrow{H^{i}} & \text { in } \mathscr{R}_{0}, \\
\vec{H} & \text { in } \mathscr{R}_{1} .\end{cases} \tag{2.7}
\end{align*}
$$

The interest of the notion of diffracted field is that it satisfies the so-called radiation condition (or Sommerfeld condition, or outgoing wave condition), in contrast with the total field which does not satisfy this condition in $\mathscr{R}_{0}$ since it includes the incident field. This means that the diffracted fields must remain bounded and propagate upwards in $\mathscr{R}_{0}$ when $z \rightarrow+\infty$. The same property must be satisfied in $\mathscr{R}_{1}$, but that time the diffracted fields must remain bounded and propagate downwards in $\mathscr{R}_{1}$ when $z \rightarrow-\infty$. Since the incident fields satisfy Maxwell's equations in $\mathscr{R}_{0}$, the diffracted fields satisfy these equations as well. Introducing the components of the diffracted fields on the three axes, Maxwell's equations yield:

$$
\begin{align*}
& \partial E_{y}^{d} / \partial z=-i \omega \mu_{0} H_{x}^{d}  \tag{2.8a}\\
& \partial E_{y}^{d} / \partial x=i \omega \mu_{0} H_{z}^{d}  \tag{2.8b}\\
& \partial E_{z}^{d} / \partial x-\partial E_{x}^{d} / \partial z=-i \omega \mu_{0} H_{y}^{d} \tag{2.8c}
\end{align*}
$$

$$
\begin{align*}
& \partial H_{y}^{d} / \partial z=i \omega \varepsilon E_{x}^{d},  \tag{2.9a}\\
& \partial H_{y}^{d} / \partial x=-i \omega \varepsilon E_{z}^{d},  \tag{2.9b}\\
& \partial H_{z}^{d} / \partial x-\partial H_{x}^{d} / \partial z=i \omega \varepsilon E_{y}^{d} . \tag{2.9c}
\end{align*}
$$

### 2.2.3 Boundary conditions on the grating profile

On the grating profile, the tangential component of the electric and magnetic fields must be continuous ${ }^{11}$. Thus the boundary condition is given by:

$$
\begin{align*}
& \left(\overrightarrow{\left[E^{d}\right]_{0}}+\overrightarrow{\left[E^{i}\right]_{0}}\right) \times \vec{n}=\overrightarrow{\left[E^{d}\right]_{1}} \times \vec{n}  \tag{2.10}\\
& \left(\overrightarrow{\left[H^{d}\right]_{0}}+\overrightarrow{\left[H^{i}\right]_{0}}\right) \times \vec{n}=\overrightarrow{\left[H^{d}\right]_{1}} \times \vec{n} \tag{2.11}
\end{align*}
$$

with $\vec{n}$ being the unit normal to $\mathscr{P}$, oriented toward region $\mathscr{R}_{0}$ (figure 2.1 and the symbol $[\vec{F}]_{p}$ denoting the limit of $\vec{F}$ when a point of region $\mathscr{R}_{p}$ tends to the grating profile (with $p \in$ $(0,1)$ ). As for Maxwell's equations, we note that the other boundary conditions on the normal components of the fields are consequences of equations 2.10) and (2.11). It is worth noting that the linkage between these two boundary conditions is a typical example of an elementary property which is difficult to establish, at least for those who are not acquainted with the theory of distributions. Projecting equations (2.10) and (2.11) on the three axes yields:

$$
\begin{gather*}
{\left[E_{y}^{d}\right]_{0}-\left[E_{y}^{d}\right]_{1}=-\left[E_{y}^{i}\right]_{0}}  \tag{2.12a}\\
n_{x}\left[E_{z}^{d}\right]_{0}-n_{z}\left[E_{x}^{d}\right]_{0}-n_{x}\left[E_{z}^{d}\right]_{1}+n_{z}\left[E_{x}^{d}\right]_{1}=-n_{x}\left[E_{z}^{i}\right]_{0}+n_{z}\left[E_{x}^{i}\right]_{0}  \tag{2.12b}\\
{\left[H_{y}^{d}\right]_{0}-\left[H_{y}^{d}\right]_{1}=-\left[H_{y}^{i}\right]_{0}}  \tag{2.13a}\\
n_{x}\left[H_{z}^{d}\right]_{0}-n_{z}\left[H_{x}^{d}\right]_{0}-n_{x}\left[H_{z}^{d}\right]_{1}+n_{z}\left[H_{x}^{d}\right]_{1}=-n_{x}\left[H_{z}^{i}\right]_{0}+n_{z}\left[H_{x}^{i}\right]_{0} \tag{2.13b}
\end{gather*}
$$

### 2.2.4 Separating the general boundary-value problem into two separated scalar problems

The first conclusion to draw from equations (2.8), (2.9), (2.12) and (2.13) is that they can be separated into two independent sets. The first one, called TE case, includes equations (2.8a), 2.8b), 2.9c), 2.12a) and 2.13b). It only contains the transverse component (viz. the $y$-component) $E_{y}^{d}$ of the electric field and the $x z$ components (orthogonal to the $y$ axis) $H_{x}^{d}$ and $H_{z}^{d}$ of the magnetic field. It must be remembered that the incident field $\overrightarrow{E^{i}}$ is given by equation (2.1) and thus is not an unknown field. The same remark applies to the complementary set (TM case), but with the transverse component of the magnetic field and the $x z$ components of the electric field. As a consequence, the general problem of diffraction by a grating can be decomposed into two elementary mathematical problems.

### 2.2.4.1 The TE case problem

In the first one, the $x z$ components of the magnetic field can be expressed as functions of the transverse component of the electric field using equations (2.8a) and 2.8b). Inserting their expression in equation $(2.9 \mathrm{c})$ shows that $E_{y}^{d}$ satisfies a Helmholtz equation:

$$
\begin{equation*}
\nabla^{2} E_{y}^{d}+k^{2} E_{y}^{d}=0 \tag{2.14}
\end{equation*}
$$

[^1]with:
\[

k= $$
\begin{cases}k_{0} & \text { in } \mathscr{R}_{0}  \tag{2.15}\\ k_{1}=k_{0} v & \text { in } \mathscr{R}_{1}\end{cases}
$$
\]

The associated boundary condition on the diffracted electric field can be deduced from equations (2.12a) and (2.1):

$$
\begin{equation*}
\left[E_{y}^{d}\right]_{0}-\left[E_{y}^{d}\right]_{1}=-P_{y} \exp \left(i k_{0} x \sin (\theta)-i k_{0} z \cos (\theta)\right), \quad \text { with }(x, z) \in \mathscr{P} \tag{2.16}
\end{equation*}
$$

while the associated boundary condition on its normal derivative can be deduced from equations (2.13b), 2.8a) and 2.8b):

$$
\begin{align*}
& {\left[\frac{\mathrm{d} E_{y}^{d}}{\mathrm{~d} n}\right]_{0}-\left[\frac{\mathrm{d} E_{y}^{d}}{\mathrm{~d} n}\right]_{1}=-\left[\frac{\mathrm{d} E_{y}^{i}}{\mathrm{~d} n}\right]_{0},}  \tag{2.17}\\
& =-i P_{y} \vec{n} \cdot \overrightarrow{k_{0}^{i}} \exp \left(i k_{0} x \sin (\theta)-i k_{0} z \cos (\theta)\right), \quad \text { with }(x, z) \in \mathscr{P},
\end{align*}
$$

with $\frac{\mathrm{d} F}{\mathrm{~d} n}$ denoting the normal derivative $\overrightarrow{\mathrm{n}} . \nabla F$. It can be noticed that equation 2.17 entails the continuity of the normal derivative of the transverse component of the total electric field. Equations (2.14), (2.16) and 2.17) are not sufficient to define the boundary-value problem for TE case. A fourth condition must be added: the radiation condition:

$$
\begin{equation*}
E_{y}^{d} \text { must satisfy a radiation condition for } z \rightarrow \pm \infty \tag{2.18}
\end{equation*}
$$

The boundary value problem allows us to deduce a fundamental property of gratings. Let us suppose that the incident field is TE polarized, i.e. that the electric incident field is parallel to the $y$ axis ( $P_{x}=P_{z}=0$ ). In these conditions, the equations associated with the TM case are homogeneous: they do not contain the incident field since the right-hand member of equation (2.12b) vanishes. If we believe that the solution of the grating problem is unique, it must be concluded that the $x z$ component of the diffracted and total electric field vanish. On the other hand, the magnetic field is parallel to the $x z$ plane. In other words, in the TE case, the grating problem becomes scalar: we must determine the $y$-component of the diffracted electric field. The $x z$ components of the magnetic field deduce the $y$-component of the diffracted electric field using equations 2.8 a ) and 2.8 b .

### 2.2.4.2 The TM case problem

Now, let us deal with the TM case. As for the TE case, it can be shown that the $y$-component of the magnetic field satisfies a Helmholtz equation by using equations (2.8c), (2.9a) and (2.9b):

$$
\begin{equation*}
\nabla^{2} H_{y}^{d}+k^{2} H_{y}^{d}=0 \tag{2.19}
\end{equation*}
$$

The boundary conditions need the calculation of the incident magnetic field. From equation (2.1) and Maxwell equation (2.3), it turns out that:

$$
\begin{equation*}
\overrightarrow{H^{i}}=\vec{Q} \exp \left(i k_{0} x \sin (\theta)-i k_{0} z \cos (\theta)\right) \tag{2.20}
\end{equation*}
$$

with:

$$
\begin{equation*}
\vec{Q}=\frac{1}{\omega \mu_{0}} \overrightarrow{k_{0}^{\prime}} \cdot \vec{P} \exp \left(i k_{0} x \sin (\theta)-i k_{0} z \cos (\theta)\right) \tag{2.21}
\end{equation*}
$$

The associated boundary condition on the diffracted magnetic field can be deduced from equations (2.13a) and (2.20):

$$
\begin{equation*}
\left[H_{y}^{d}\right]_{0}-\left[H_{y}^{d}\right]_{1}=-Q_{y} \exp \left(i k_{0} x \sin (\theta)-i k_{0} z \cos (\theta)\right), \quad \text { with }(x, z) \in \mathscr{P}, \tag{2.22}
\end{equation*}
$$

while the boundary condition on its normal derivative is obtained by inserting the expressions of the $x z$ components of the electric field (equations (2.9a) and 2.9b) in equation 2.12b). Remarking that the incident field satisfies the same equations, we obtain finally:

$$
\begin{align*}
& \frac{1}{\varepsilon_{0}}\left[\frac{\mathrm{~d} H_{y}^{d}}{\mathrm{~d} n}\right]_{0}-\frac{1}{\varepsilon_{1}}\left[\frac{\mathrm{~d} H_{y}^{d}}{\mathrm{~d} n}\right]_{1}=-\frac{1}{\varepsilon_{0}}\left[\frac{\mathrm{~d} H_{y}^{i}}{\mathrm{~d} n}\right]_{0},  \tag{2.23}\\
& =-\frac{i Q_{y}}{\varepsilon_{0}} \vec{n} \cdot \overrightarrow{k_{0}^{i}} \exp \left(i k_{0} x \sin (\theta)-i k_{0} z \cos (\theta)\right), \quad \text { with }(x, z) \in \mathscr{P} .
\end{align*}
$$

It can be noticed that equation 2.23 has a simple interpretation: the product $\frac{1}{\varepsilon} \frac{\mathrm{~d} H_{y}}{\mathrm{~d} n}$ is continuous across the profile. Finally, the radiation condition yields:

$$
\begin{equation*}
H_{y}^{d} \text { must satisfy a radiation condition for } z \rightarrow \pm \infty . \tag{2.24}
\end{equation*}
$$

Equations (2.19), (2.22), (2.23) and radiation conditions for $z \rightarrow \pm \infty$ define the boundaryvalue problem for TM case. As for TE case, the uniqueness of the solution shows that that when the magnetic incident field is parallel to the $y$ axis ( $Q_{x}=Q_{z}=0$ ), the equations associated with the TE case are homogeneous: they do not contain the incident field. It can be concluded that the $x z$ components of the diffracted and total magnetic fields vanish. On the other hand, the electric field is parallel to the $x z$ plane. In other words, in the TM case, the grating problem becomes scalar: we must determine the y-component of the diffracted magnetic field. The $x z$ components of the electric field deduce from the $y$-component of the diffracted magnetic field using equations (2.9a) and 2.9b).

### 2.2.4.3 TE and TM cases: a unified presentation of the boundary-value problem

In order to deal with both cases simultaneously, we denote by $F^{d}$ the field defined by:

$$
F^{d}= \begin{cases}E_{y}^{d} & \text { for TE case }  \tag{2.25}\\ H_{y}^{d} & \text { for TM case }\end{cases}
$$

In the same way, by assuming that the incident field has a unit amplitude ( $P_{y}=1$ for TE case and $Q_{y}=1$ for TM case), the incident field in both cases is given by:

$$
\begin{equation*}
F^{i}=\exp \left(i k_{0} x \sin (\theta)-i k_{0} z \cos (\theta)\right), \tag{2.26}
\end{equation*}
$$

the total field $F$ being given by:

$$
F= \begin{cases}F^{d}+F^{i} & \text { in } \mathscr{R}_{0},  \tag{2.27}\\ F^{d} & \text { in } \mathscr{R}_{1} .\end{cases}
$$

Using equations (2.14), (2.16), (2.17), (2.18), (2.19), (2.22), (2.23) and (2.24), it is possible to gather both cases in a unique set of equations:

$$
\begin{align*}
& \nabla^{2} F^{d}+k^{2} F^{d}=0,  \tag{2.28}\\
& {\left[F^{d}\right]_{0}-\left[F^{d}\right]_{1}=-\exp \left(i k_{0} x \sin (\theta)-i k_{0} z \cos (\theta)\right) \quad \text { with }(x, z) \in \mathscr{P},}  \tag{2.29}\\
& \frac{1}{\tau_{0}}\left[\frac{\mathrm{~d} F^{d}}{\mathrm{~d} n}\right]_{0}-\frac{1}{\tau_{1}}\left[\frac{\mathrm{~d} F^{d}}{\mathrm{~d} n}\right]_{1},  \tag{2.30}\\
& =-\frac{i}{\tau_{0}} \vec{n} \cdot \overrightarrow{k_{0}^{t}} \exp \left(i k_{0} x \sin (\theta)-i k_{0} z \cos (\theta)\right), \quad \text { with }(x, z) \in \mathscr{P}, \\
& F^{d} \text { must satisfy a radiation condition for } y \rightarrow \pm \infty, \tag{2.31}
\end{align*}
$$

with:

$$
\tau_{i}=\left\{\begin{array}{lc}
1 & \text { for TE case, }  \tag{2.32}\\
\varepsilon_{i} & \text { for TM case }, i \in(0,1)
\end{array}\right.
$$

In the following, this boundary-value problem will be called normalized grating problem. It is worth noting that equations 2.29$)$ and (2.30) take a simpler form by introducing the total field $F$ :

$$
\begin{gather*}
{[F]_{0}=[F]_{1}}  \tag{2.33}\\
\frac{1}{\tau_{0}}\left[\frac{\mathrm{~d} F}{\mathrm{~d} n}\right]_{0}=\frac{1}{\tau_{1}}\left[\frac{\mathrm{~d} F}{\mathrm{~d} n}\right]_{1} \tag{2.34}
\end{gather*}
$$

### 2.2.5 The special case of the perfectly-conducting grating

The first grating theories were devoted to perfectly conducting gratings. This case is very important since it is realistic for metallic gratings in the microwave domain and far infrared regions. In the visible and infrared regions, it can provide qualitative results. However, in these regions, one must be very cautious. The existence of surface plasmons propagating at the vicinity of the grating surface generates strong resonance phenomena for TM case. Due to these phenomena, the properties of real metallic gratings and those of perfectly-conducting gratings may completely differ [2].Moreover, the perfect conductivity model allows one to simplify the grating theory, since the associated boundary-value problems are much simpler.

Basically, the equations associated to the perfect conductivity model are the same as for real metallic or dielectric gratings, except equations (2.4) and (2.11). Let us give a brief explanation to this property. In Maxwell equation (2.4), the right-hand member includes the volume current density $\vec{j}$ in the metal since this term is proportional to the electric field $(\vec{j}=\sigma \vec{E}$, $\sigma$ being the conductivity of the metal). When the conductivity tends to infinity, the volume current density and the total fields concentrate more and more on the grating surface since the skin depth tends to zero. As a consequence, at the limit when the conductivity tends to infinity, the fields are null in $\mathscr{R}_{1}$ while the volume current density $\vec{j}$ becomes a surface current density $\overrightarrow{j_{\mathscr{P}}}$. This surface current density cannot be included in the right-hand member of equation (2.4) since it is a singular distribution (for the specialist of Schwartz distributions [7], it writes $\stackrel{\overrightarrow{j \mathscr{P}}}{\overrightarrow{2}} \boldsymbol{\delta}_{\mathscr{P}}$ ). Finally, equation 2.4 becomes:

$$
\begin{equation*}
\nabla \times \vec{H}=-i \omega \tilde{\varepsilon} \vec{E}+\overrightarrow{j \mathscr{P}} \tag{2.35}
\end{equation*}
$$

with $\tilde{\varepsilon}$ being the permittivity of the material. Furthermore, taking into account that the total fields vanish inside $\mathscr{R}_{1}$, the boundary condition (equation (2.11)) becomes:

$$
\begin{equation*}
\vec{n} \times\left(\overrightarrow{\left[H^{d}\right]_{0}}+\overrightarrow{\left[H^{i}\right]_{0}}\right)=\overrightarrow{j_{\mathscr{P}}} \tag{2.36}
\end{equation*}
$$

This equation reduces to a relation between the surface current density on $\mathscr{P}$ and the limit of the magnetic field above $\mathscr{P}$. It does not constitute any more an element of the boundary-value problem.

In conclusion, for perfectly conducting gratings, the fields inside $\mathscr{R}_{1}$ vanish and, using equations (2.3), (2.4), (2.10), (2.6) and (2.7), the basic vector equations for the field in $\mathscr{R}_{0}$ can be written:

$$
\begin{gather*}
\nabla \times \overrightarrow{E^{d}}=i \omega \mu_{0} \overrightarrow{H^{d}},  \tag{2.37}\\
\nabla \times \overrightarrow{H^{d}}=-i \omega \varepsilon_{0} \overrightarrow{E^{d}}  \tag{2.38}\\
\left(\overrightarrow{\left[E^{d}\right]_{0}}+\overrightarrow{\left[E^{i}\right]_{0}}\right) \times \vec{n}=0 \tag{2.39}
\end{gather*}
$$

Following the same lines as in subsections 2.2.4.1 and 2.2.4.2, the boundary value problems for perfectly conducting gratings are given by:

For TE case:

$$
\begin{align*}
& \nabla^{2} E_{y}^{d}+k_{0}^{2} E_{y}^{d}=0,  \tag{2.40}\\
& {\left[E_{y}^{d}\right]_{0}=-P_{y} \exp \left(i k_{0} x \sin (\theta)-i k_{0} z \cos (\theta)\right), \quad \text { with }(x, z) \in \mathscr{P},} \\
& E_{y}^{d} \text { must satisfy a radiation condition for } z \rightarrow+\infty .
\end{align*}
$$

For TM case:

$$
\begin{align*}
& \nabla^{2} H_{y}^{d}+k_{0}^{2} H_{y}^{d}=0,  \tag{2.43}\\
& {\left[\frac{\mathrm{~d} H_{y}^{d}}{\mathrm{~d} n}\right]_{0}=-i Q_{y} \vec{n} \cdot \overrightarrow{k_{0}^{\prime}} \exp \left(i k_{0} x \sin (\theta)-i k_{0} z \cos (\theta)\right), \quad \text { with }(x, z) \in \mathscr{P},}  \tag{2.44}\\
& H_{y}^{d} \text { must satisfy a radiation condition for } z \rightarrow+\infty . \tag{2.45}
\end{align*}
$$

### 2.3 Pseudo-periodicity of the field and Rayleigh expansion

This section establishes the most famous property of diffraction gratings: the dispersion of light, which is a consequence of the well known grating formula. In general, this formula is demonstrated using heuristic considerations of physical optics. Here, we propose a rigorous demonstration based on the boundary-value problem stated in subsection 2.2.4.3. First, let us show that the field $F^{d}$ is pseudo-periodic, i.e. that:

$$
\begin{equation*}
F^{d}(x+d, z)=F^{d}(x, z) \exp \left(i k_{0} d \sin (\theta)\right) \tag{2.46}
\end{equation*}
$$

With this aim, we consider the function $G(x, z)$ defined by:

$$
\begin{equation*}
G(x, z)=F^{d}(x+d, z) \exp \left(-i k_{0} d \sin (\theta)\right) \tag{2.47}
\end{equation*}
$$

The pseudo-periodicity of $F^{d}$ is proved if we show that $F^{d}(x, z)=G(x, z)$. Owing to the uniqueness of the solution of the boundary-value problem defined by equations (2.28), (2.29), (2.30) and (2.31), this equation is satisfied if $G$ obeys the same equations. Obviously, $G$ satisfies these equations since $d$ is the grating period. Thus $F^{d}$ is pseudo-periodic, with coefficient of pseudoperiodicity $k_{0} \sin (\theta)$, as well as $F^{i}$ and $F$. Notice that in normal incidence $(\theta=0)$, pseudoperiodicity becomes ordinary periodicity, which in that case is a straightforward property since both grating and incident wave are periodic.

Using the pseudo-periodicity, let us show that the field above and below the grating is a sum of plane waves. With this aim, we notice from equation 2.28 that $F^{d}(x, z) \exp \left(-i k_{0} x \sin (\theta)\right)$ has a period $d$ and thus can be expanded in a Fourier series:

$$
\begin{equation*}
F^{d}(x, z) \exp \left(-i k_{0} x \sin (\theta)\right)=\sum_{n=-\infty}^{+\infty} F_{n}^{d}(z) \exp (2 i \pi n x / d) \tag{2.48}
\end{equation*}
$$

Multiplying both members of equation 2.48) by $\exp \left(i k_{0} x \sin (\theta)\right)$ yields :

$$
\begin{equation*}
F^{d}(x, z)=\sum_{n=-\infty}^{+\infty} F_{n}^{d}(z) \exp \left(i \alpha_{n} x\right) \tag{2.49}
\end{equation*}
$$

with:

$$
\begin{equation*}
\alpha_{n}=k_{0} \sin (\theta)+2 \pi n / d . \tag{2.50}
\end{equation*}
$$

Introducing this expression of $F^{d}(x, z)$ in Helmholtz equation 2.28, we find :

$$
\begin{equation*}
\sum_{n=-\infty}^{+\infty}\left(\mathrm{d}^{2} F_{n}^{d}(z) / \mathrm{d} z^{2}+\left(k^{2}-\alpha_{n}^{2}\right) F_{n}^{d}(z)\right) \exp \left(i \alpha_{n} x\right)=0 \tag{2.51}
\end{equation*}
$$

and multiplying both members by $\exp \left(-i k_{0} x \sin (\theta)\right)$,

$$
\begin{equation*}
\sum_{n=-\infty}^{+\infty}\left(\mathrm{d}^{2} F_{n}^{d}(z) / \mathrm{d} z^{2}+\left(k^{2}-\alpha_{n}^{2}\right) F_{n}^{d}(z)\right) \exp (2 i \pi n x / d)=0 \tag{2.52}
\end{equation*}
$$

It seems, at the first glance, that the left-hand member of equation (2.52) is a Fourier series, and thus that the coefficients of this Fourier series vanish. This is not correct. Indeed, we have to bear in mind that $k$, defined in equation (2.15) is not a constant. As a consequence, if $0<y<z_{M}$, a region called intermediate region in the following, $k^{2}$ depends on $x$ and the left-hand member of equation (2.52) is not a Fourier series. However, above and below this intermediate region, $k^{2}$ is constant and we can write that the Fourier coefficients vanish:

$$
\begin{array}{lll}
\forall n, & \mathrm{~d}^{2} F_{n}^{d}(z) / \mathrm{d} z^{2}+\gamma_{0, n}^{2} F_{n}^{d}(z)=0 & \text { if } y>z_{M}, \\
\forall n, & \mathrm{~d}^{2} F_{n}^{d}(z) / \mathrm{d} z^{2}+\gamma_{1, n}^{2} F_{n}^{d}(z)=0 & \text { if } y<0, \tag{2.53b}
\end{array}
$$

with:

$$
\begin{equation*}
\gamma_{i, n}=\sqrt{\left(k_{i}^{2}-\alpha_{n}^{2}\right)} \quad i \in(0,1) \tag{2.54}
\end{equation*}
$$

We deduce that:

$$
F_{n}^{d}(z)= \begin{cases}I_{0, n} \exp \left(-i \gamma_{0, n} z\right)+D_{0, n} \exp \left(+i \gamma_{0, n} z\right) & \text { if } y>z_{M}  \tag{2.55}\\ D_{1, n} \exp \left(-i \gamma_{1, n} z\right)+I_{1, n} \exp \left(+i \gamma_{1, n} z\right) & \text { if } y<0\end{cases}
$$

and therefore, using equation (2.49),

$$
F^{d}(x, z)=\left\{\begin{array}{cc}
\sum_{n=-\infty}^{+\infty}\left(I_{0, n} \exp \left(i \alpha_{n} x-i \gamma_{0, n} z\right)+\right. &  \tag{2.56}\\
\left.+D_{0, n} \exp \left(i \alpha_{n} x+i \gamma_{0, n} z\right)\right) & \text { if } z>z_{M} \\
\sum_{n=-\infty}^{+\infty}\left(D_{1, n} \exp \left(i \alpha_{n} x-i \gamma_{1, n} z\right)+\right. & \\
\left.+I_{1, n} \exp \left(i \alpha_{n} x+i \gamma_{1, n} z\right)\right) & \text { if } z<0
\end{array}\right.
$$

Let us remark that equation (2.54) does not assign to $\gamma_{i, n}$ a unique value. However, equation (2.56) shows that its determination can be chosen arbitrarily since a sign change does not modify the value of the field, provided that $I_{0, n}$ and $D_{0, n}$ are permuted. The determination of these constants will be given by:

$$
\begin{equation*}
\operatorname{Re}\left(\gamma_{i, n}\right)+\operatorname{Im}\left(\gamma_{i, n}\right)>0, \quad i \in(0,1) \tag{2.57}
\end{equation*}
$$

with $\operatorname{Re}(q)$ and $\operatorname{Im}(q)$ denoting the real and imaginary parts of $q$.
Equation (2.56) shows that the field above and below the intermediate region can be represented by plane wave expansions. The propagation constants of the plane waves along the $x$ and $z$ axes are respectively equal to $\alpha_{n}$ and $\pm \gamma_{i, n}$. In the physical problem, some of these plane waves must be rejected since they do not obey the radiation condition. This condition entails that $I_{0, n}=I_{1, n}=0$ since, according to equation (2.57), the associated plane waves propagate towards the grating profile. Finally, equations (2.56, 2.27) and the radiation condition allow us to express the total field by adding the incident field:

$$
F(x, z)=\left\{\begin{array}{c}
\exp \left(i \alpha_{0} x-i \gamma_{0,0} z\right)+  \tag{2.58}\\
\quad+\sum_{n=-\infty}^{+\infty} D_{0, n} \exp \left(i \alpha_{n} x+i \gamma_{0, n} z\right) \quad \text { if } z>z_{M} \\
\sum_{n=-\infty}^{+\infty} D_{1, n} \exp \left(i \alpha_{n} x-i \gamma_{1, n} z\right) \quad \text { if } z<0,
\end{array}\right.
$$

the sums being the expression of the scattered field in both regions. The unknown complex coefficients $D_{0, n}$ and $D_{1, n}$ are the amplitudes of the reflected and transmitted waves respectively.

The conclusion of this subsection is that above and below the intermediate region, the field reflected and transmitted by the grating takes the form of sums of plane waves (Rayleigh expansion [8]), each of them being characterized by its order $n$.

### 2.4 Grating formulae

According to equation (2.54), almost all the diffracted plane waves (an infinite number) are evanescent: they propagate along the $x$ axis at the vicinity of the grating profile since they decrease exponentially when $|z| \rightarrow+\infty$. For $z \rightarrow+\infty$, they correspond to the orders $n$ such that $\alpha_{n}^{2} \geq k_{0}^{2}$, thus rendering $\gamma_{0, n}=i \sqrt{\left(\alpha_{n}^{2}-k_{0}^{2}\right)}$ a purely imaginary number. Only a finite number of them, called $z$-propagative orders, propagate towards $z=+\infty$, with $\alpha_{n}^{2} \leq k_{0}^{2}$, thus $\gamma_{0, n}=\sqrt{\left(k_{0}^{2}-\alpha_{n}^{2}\right)}$ being real. Let us notice that among these orders, the $0^{\text {th }}$ order is always included, since $\gamma_{0, n}=k_{0} \cos (\theta)$. It propagates in the direction specularly reflected by the mean plane of the profile, whatever the wavelength may be. In contrast, the other $z$-propagative orders are dispersive. Indeed, their propagation constants along the $x$ and $z$ axes are equal to $\alpha_{n}$ and $\gamma_{0, n}$, in such a way that the diffraction angle $\theta_{0, n}$ of one of these waves, measured clockwise from the $z$ axis, can be deduced from $\alpha_{n}=k_{0} \sin \left(\theta_{0, n}\right)$. Using the expression of $\alpha_{n}$ given by equation (2.50), the angle of diffraction is given by :

$$
\begin{equation*}
\sin \left(\theta_{0, n}\right)=\sin (\theta)+n \frac{2 \pi}{k_{0} d}=\sin (\theta)+n \frac{\lambda}{d} . \tag{2.59}
\end{equation*}
$$

This is the famous grating formula, often deduced from heuristic arguments of physical optics.
For the field below the grooves, the wavenumber $k_{0}$ is replaced by $k_{1}=k_{0} v$. If the grating material is a lossless dielectric, the directions of propagation of the transmitted field obey a grating formula as well. This formula is similar to equation (2.59) but the angles of transmission $\theta_{1, n}$ can be deduced from $\alpha_{n}=k_{0} v \sin \left(\theta_{1, n}\right)$, which yields, using a counterclockwise convention:

$$
\begin{equation*}
v \sin \left(\theta_{0, n}\right)=\sin (\theta)+n \frac{2 \pi}{k_{0} d}=\sin (\theta)+n \frac{\lambda}{d} . \tag{2.60}
\end{equation*}
$$

The $0^{\text {th }}$ order is always included in the $z$-propagative orders ${ }^{2}$. It propagates in the direction of transmission by an air/dielectric plane interface, whatever the wavelength may be. In contrast, the other $z$-propagative orders are dispersive. When the grating material is metallic, the transmitted plane waves are absorbed by the metal and the $z$-propagating orders below the grooves no longer exist.

In conclusion of this section, the reflected and transmitted fields include, outside the grooves, a finite number of plane waves propagating to infinity with scattering angles given by the grating formulae. All the orders are dispersive, except the $0^{t h}$ orders. The reflected $0^{t h}$ order takes the specular direction while for a lossless material, the transmitted $0^{\text {th }}$ order takes the direction transmitted by an air/dielectric plane interface. Consequently, a polychromatic incident plane wave generates in a given order $n$ different from 0 a sum of plane waves scattered in different directions, i.e. a spectrum. The measurement of the intensity along this spectrum allows one to determine the spectral power of the incident wave. This dispersion phenomenon is the most important property of diffraction gratings. It explains why this optical component has been one of the most valuable tools in the history of Science.

### 2.5 Analytic properties of gratings

### 2.5.1 Balance relations

The mathematical balance relations established in this subsection will allow us to demonstrate very important general properties of gratings. These balance relations state mathematical links between characteristics of the field in two regions separated by large distances, without considering the fields in between. They can give a relation between the fields at $z=+\infty$ and the fields on the grating profile, or the fields at $z=-\infty$ and the fields on the grating profile, or the fields at $z=+\infty$ and $z=-\infty$.

### 2.5.1. 1 Lemma 1

We consider two pseudoperiodic functions $u$ and $v$ of the two variables $x$ and $z$, defined in $\mathscr{R}_{0}$, which belong to the class $G_{0}$ of functions having the following properties:

- They are pseudo-periodic, with the same coefficient of pseudo-periodicity $\alpha$, in other words, $u(x, z) \exp (-i \alpha x)$ and $v(x, z) \exp (-i \alpha x)$ are periodic,

[^2]- They are solutions of a Helmholtz equation:

$$
\begin{align*}
& \nabla^{2} u+k_{0}^{2} u=0  \tag{2.61a}\\
& \nabla^{2} v+k_{0}^{2} v=0 \tag{2.61b}
\end{align*}
$$

with $k_{0}$ being real.

- They are bounded for $z \rightarrow \infty$,
- They are square integrable in $x$ and locally square integrable in $z$,
- Their values on $\mathscr{P}$ are square integrable, as well as their normal derivatives.

We introduce the sesquilinear functional defined by:

$$
\begin{equation*}
\mathscr{F}_{0}=\int_{\mathscr{P}}\left(u \frac{\mathrm{~d} \bar{v}}{\mathrm{dn}}-\bar{v} \frac{\mathrm{~d} u}{\mathrm{dn}}\right) \mathrm{d} s . \tag{2.62}
\end{equation*}
$$

The symbol $\int_{\mathscr{P}}$ denotes a curvilinear integral on one period of the profile $\mathscr{P}$ of the grating, with $\mathrm{d} s$ being the differential of the curvilinear abscissa on $\mathscr{P}$. Obviously, the value in region $\mathscr{R}_{0}$ of the fields $F(x, z)$, solutions of the four boundary-value problems defined in subsection 2.2.4, belong to $G_{0}$, as well as the incident field $F^{i}$. It is to be noticed that we do not impose a boundary condition on $\mathscr{P}$ or a radiation condition at infinity, but we still impose that these functions must remain bounded at infinity.

Following the same lines as in section 2.3, it can be shown that above the top of the grooves, $u$ and $v$ can be represented by plane wave expansions, similar to that of equation (2.56): if $z>z_{M}$,

$$
\begin{align*}
& u(x, z)=\sum_{n=-\infty}^{+\infty}\left[I_{0, n} \exp \left(i \alpha_{n} x-i \gamma_{0, n} z\right)+D_{0, n} \exp \left(i \alpha_{n} x+i \gamma_{0, n} z\right)\right]  \tag{2.63a}\\
& v(x, z)=\sum_{n=-\infty}^{+\infty}\left[I_{0, n}^{\prime} \exp \left(i \alpha_{n} x-i \gamma_{0, n} z\right)+D_{0, n}^{\prime} \exp \left(i \alpha_{n} x+i \gamma_{0, n} z\right)\right] \tag{2.63b}
\end{align*}
$$

Let us notice that some terms must be eliminated in the Rayleigh expansions. Indeed, the field must remain bounded at infinity. It is not the case for the incident terms of coefficients $I_{0, n}$ and $I_{0, n}^{\prime}$ unless the corresponding plane waves are $z$-propagating waves. Thus we define the set $U_{0}$ of orders corresponding to $z$-propagating waves and equations (2.63) become:

$$
\begin{align*}
& u(x, z)= \sum_{n \in U_{0}} I_{0, n} \exp \left(i \alpha_{n} x-i \gamma_{0, n} z\right)+\sum_{n=-\infty}^{+\infty} D_{0, n} \exp \left(i \alpha_{n} x+i \gamma_{0, n} z\right)  \tag{2.64a}\\
& v(x, z)= \sum_{n \in U_{0}} I_{0, n}^{\prime} \exp \left(i \alpha_{n} x-i \gamma_{0, n} z\right)+\sum_{n=-\infty}^{+\infty} D_{0, n}^{\prime} \exp \left(i \alpha_{n} x+i \gamma_{0, n} z\right)  \tag{2.64b}\\
& \overline{v(x, z)}=\sum_{n \in U_{0}} \overline{I_{0, n}^{\prime}} \exp \left(-i \alpha_{n} x+i \gamma_{0, n} z\right)+ \\
&+\sum_{n=-\infty}^{+\infty} \overline{D_{0, n}^{\prime}} \exp \left(-i \alpha_{n} x-i \overline{\gamma_{0, n}} z\right) \tag{2.64c}
\end{align*}
$$



Figure 2.2: Balance relations.

Now, we show that $\mathscr{F}_{0}$ can be expressed as a function of the Rayleigh coefficients $I_{0, n}$, $D_{0, n}, I_{0, n}^{\prime}$ and $D_{0, n}^{\prime}$. With this aim, we multiply equation 2.61a by $\bar{v}$, the conjugate of equation (2.61b) by $u$ and we substract the first from the second, which yields:

$$
\begin{equation*}
u \nabla^{2} \bar{v}-\bar{v} \nabla^{2} u=0 \quad \text { in } \quad \mathscr{R}_{0} . \tag{2.65}
\end{equation*}
$$

Integrating equation (2.65) in the blue area of figure 2.2 and applying the second Green identity yields:

$$
\begin{equation*}
\int_{\Omega_{0}}\left(u \frac{\mathrm{~d} \bar{v}}{\mathrm{dn}}-\bar{v} \frac{\mathrm{~d} u}{\mathrm{dn}}\right) \mathrm{d} l=0, \tag{2.66}
\end{equation*}
$$

with $\Omega_{0}$ being the boundary of the blue area of figure 2.2 and $\mathrm{d} l$ denoting the differential of the curvilinear abscissa on $\Omega_{0}$. According to equations 2.64a and 2.64 c ,,$u \frac{\mathrm{~d} \bar{v}}{\mathrm{dx}}$ and $\bar{v} \frac{\mathrm{~d} u}{\mathrm{dx}}$ are periodic. Since the orientations of the normal on verticals $\mathrm{O} \mathrm{\Gamma}_{1}$ and $\mathrm{L} \mathrm{\Gamma}_{2}$ are opposite, the contributions of the integrals on these segments cancel each other. Furthermore, the normal to $\mathrm{O} \Gamma_{1}$ and $\mathrm{L} \Gamma_{2}$ is parallel to the $z$ axis and oriented downwords, then equation 2.66 becomes:

$$
\begin{equation*}
\int_{\mathscr{P}}\left(u \frac{\mathrm{~d} \bar{v}}{\mathrm{dn}}-\bar{v} \frac{\mathrm{~d} u}{\mathrm{dn}}\right) \mathrm{d} s=\int_{\Gamma_{1} \Gamma_{2}}\left(u \frac{\mathrm{~d} \bar{v}}{\mathrm{dz}}-\bar{v} \frac{\mathrm{~d} u}{\mathrm{dz}}\right) \mathrm{d} x . \tag{2.67}
\end{equation*}
$$

Introducing in the right-hand member of equation 2.66) the expressions of $u$ and $\bar{v}$ given by equations (2.64a) and (2.64c), separating the terms $n \in U_{0}$ from the other ones and taking into account that $\int_{x=0}^{d} \exp \operatorname{in} \frac{2 \pi}{d} x=\delta_{n, 0}$, with $\delta_{n, 0}$ being the Kronecker symbol, one can obtain, after some cumbersome but not difficult calculations that:

$$
\begin{equation*}
\int_{\mathscr{P}}\left(u \frac{\mathrm{~d} \bar{v}}{\mathrm{dn}}-\bar{v} \frac{\mathrm{~d} u}{\mathrm{dn}}\right) \mathrm{d} s=\sum_{n \in U_{0}} \gamma_{0, n}\left(I_{0, n} \overline{I_{0, n}^{\prime}}-D_{0, n} \overline{D_{0, n}^{\prime}}\right) \tag{2.68}
\end{equation*}
$$

### 2.5.1.2 Lemma 2

In this section, it is supposed that the grating material is lossless, in such a way that plane waves can propagate in $\mathscr{R}_{1}$. Lemma 2 is similar as lemma 1 , but for region $\mathscr{R}_{1}$. We denote by $U_{1}$ the set of orders corresponding to $z$-propagating waves in $\mathscr{R}_{1}$. The expressions of $u$ and $v$ below the x axis are given by:

$$
\begin{align*}
& u(x, z)= \sum_{n \in U_{1}} D_{1, n} \exp \left(i \alpha_{n} x-i \gamma_{1, n} z\right)+  \tag{2.69a}\\
& v(x, z)=\sum_{n=-\infty}^{+\infty} I_{1, n} \operatorname{Dexp}\left(i \alpha_{n} x+i \gamma_{1, n}^{\prime} z\right)  \tag{2.69b}\\
&\left.\overline{v(x, z)}=\sum_{n \in U_{1}} \overline{D_{1, n}^{\prime}} \exp \left(i \alpha_{n} x-i \gamma_{1, n} z\right)+\sum_{n=-\infty}^{+\infty} I_{1, n}^{\prime} \exp \left(i \alpha_{n} x+i \gamma_{1, n} z\right), i \gamma_{1, n} z\right)+ \\
&+\sum_{n=-\infty}^{+\infty} \overline{I_{1, n}^{\prime}} \exp \left(-i \alpha_{n} x-i \overline{\gamma_{1, n}} z\right) . \tag{2.69c}
\end{align*}
$$

Following the same lines as in section 2.5.1.1 but for the yellow area of figure 2.2 and noting that the normal is now oriented towards the exterior of the domain, it can be deduced that:

$$
\begin{equation*}
\int_{\mathscr{P}}\left(u \frac{\mathrm{~d} \bar{v}}{\mathrm{dn}}-\bar{v} \frac{\mathrm{~d} u}{\mathrm{dn}}\right) \mathrm{d} s=-\sum_{n \in U_{1}} \gamma_{1, n}\left(I_{1, n} \overline{I_{1, n}^{\prime}}-D_{1, n} \overline{D_{1, n}^{\prime}}\right) \tag{2.70}
\end{equation*}
$$

### 2.5.2 Compatibility between Rayleigh coefficients

In order to state a relation between the Rayleigh coefficients above and below the grating profile, we assume that the functions $u$ and $v$ satisfy the boundary conditions imposed on the total fields by equations (2.33) and (2.34). On the other hand, we do not impose radiation conditions at infinity, but the functions must remain bounded. In other words, $u$ and $v$ can be considered as solutions of the most general grating problem, in which the incident wave is not restricted to a single plane wave, but to the sum of all the plane waves generating diffracted waves in the same directions, with arbitrary amplitudes. It is straightforward to show from equations (2.33) and $(2.34)$ that the left-hand members of equations (2.68) and (2.70) are proportional, then to deduce a relation including the coefficients of the Rayleigh expansions of the field only:

$$
\begin{array}{r}
\frac{1}{\tau_{0}} \sum_{n \in U_{0}} \gamma_{0, n}\left(I_{0, n} \overline{I_{0, n}^{\prime}}-D_{0, n} \overline{D_{0, n}^{\prime}}\right)+  \tag{2.71}\\
\frac{1}{\tau_{1}} \sum_{n \in U_{1}} \gamma_{1, n}\left(I_{1, n} \overline{I_{1, n}^{\prime}}-D_{1, n} \overline{D_{1, n}^{\prime}}\right)=0
\end{array}
$$

This equation states the most general relation of compatibility between two solutions of the general diffraction grating problem associated to different sets of incident waves. When the grating material is perfectly conducting, it is easy to show that the compatibility equation holds, provided that the sum $n \in U_{1}$ is cancelled in equation (2.71).

Phenomenological theories of gratings make a wide use of the notion of scattering matrix (or $S$-matrix). The scattering matrix states the linear relation between the amplitudes of the diffracted and incident waves. We define the column matrix containing the amplitudes of the incident waves. More precisely, we define the normalized amplitudes of the incident and scattered waves by $\tilde{I}_{0, n}=\sqrt{\gamma_{0, n}} I_{0, n}, \tilde{D}_{0, n}=\sqrt{\gamma_{0, n}} D_{0, n}, \tilde{I}_{1, n}=\sqrt{\frac{\tau_{0}}{\tau_{1}} \gamma_{1, n}} I_{1, n}, \tilde{D}_{1, n}=\sqrt{\frac{\tau_{0}}{\tau_{1}} \gamma_{1, n}} D_{1, n}$, $n \in(0,1)$, and by definition, the scattering matrix is a square matrix defined by:

$$
\begin{equation*}
\mathbb{D}=\mathbb{S} \mathbb{I} \tag{2.72}
\end{equation*}
$$

with $\mathbb{I}$ being a column vector containing successively all the incident amplitudes $\tilde{I}_{0, n}$ for $n \in U_{0}$ and all the incident amplitudes $\tilde{I}_{1, n}$ for $n \in U_{1}, \mathbb{D}$ being a column vector containing successively all the diffracted amplitudes $\tilde{D}_{0, n}$, and all the incident amplitudes $\tilde{D}_{1, n}$ for $n \in U_{1}$. Thus, the order of column matrices $\mathbb{I}$ and $\mathbb{D}$ is the sum $\left|U_{0}\right|+\left|U_{1}\right|$ of the cardinals of $U_{0}$ and $U_{1}$ Using these notations, equation 2.71 can be expressed in the very simple form:

$$
\begin{equation*}
<\mathbb{D}\left|\mathbb{D}^{\prime}>=<\mathbb{I}\right| \mathbb{I}^{\prime}> \tag{2.73}
\end{equation*}
$$

the scalar product of two column matrices of order $N$ being defined by:

$$
\begin{equation*}
<P \mid Q>=\sum_{j=1}^{N} P_{j} \overline{Q_{j}} \tag{2.74}
\end{equation*}
$$

Using equation 2.72 to eliminate $\mathbb{D}$ in equation 2.77 yields:

$$
\begin{equation*}
<\mathbb{S} \mathbb{I}\left|\mathbb{S} \mathbb{I}^{\prime}>=<\left(\mathbb{S}^{*} \mathbb{S}\right) \mathbb{I}\right| \mathbb{I}^{\prime}>=<\mathbb{I} \mid \mathbb{I}^{\prime}> \tag{2.75}
\end{equation*}
$$

with $\mathbb{S}^{*}$ being the adjoint matrix of $\mathbb{S}$. Since equation 2.75 must be satisfied for any value of $\mathbb{I}$ and $\mathbb{I}^{\prime}$, we deduce that:

$$
\begin{equation*}
\mathbb{S}^{*} \mathbb{S}=\mathbb{1} \tag{2.76}
\end{equation*}
$$

with $\mathbb{1}$ being the identity matrix. Equation 2.76 shows that $\mathbb{S}$ is unitary.

### 2.5.3 Energy balance

The energy balance relation is obtain by taking $u=v$ in equation (2.77), which gives:

$$
\begin{equation*}
<\mathbb{D}|\mathbb{D}>=<\mathbb{I}| \mathbb{I}> \tag{2.77}
\end{equation*}
$$

or equivalently:

$$
\begin{equation*}
\|\mathbb{D}\|=\|\mathbb{I}\| \tag{2.78}
\end{equation*}
$$

Let us show why this equation is known as energy balance relation. To this end, it suffices to use the Poynting theorem and to calculate the flux of the Poynting vector $\vec{E} \times \overrightarrow{\vec{H}}$ through the rectangle $\Gamma_{1} \Gamma_{2} \Gamma_{4} \Gamma_{3}$ of figure 2.2 . Since the grating material is lossless, the flux of the Poynting vector through this rectangle (with now the normal oriented toward the exterior, in contrast with figure 2.2 ) must be null. The contributions of the vertical sides $\Gamma_{1} \Gamma_{3}$ and $\Gamma_{2} \Gamma_{4}$ cancel each other, thanks to the periodicity of the Poynting vector $(\vec{H}$ has a coefficient of pseudoperiodicity which is the opposite to that of $\vec{E}$ ). At the top of the rectangle, the calculation of the flux of the Poynting vector can be achieved by using the Rayleigh expansion given by
equations 2.64. Taking into account that $\int_{x=0}^{d} \exp \operatorname{in} \frac{2 \pi}{d} x=\delta_{n}$, elementary calculations show that the contributions to this flux of the different plane waves are decoupled and are proportional to $-\gamma_{0, n}\left|I_{0, n}\right|^{2}$ and $+\gamma_{0, n}\left|D_{0, n}\right|^{2}$. At the bottom of the rectangle, we use the Rayleigh expansion given by equations 2.69). The contributions of the plane waves are decoupled as well and are proportional to $-\frac{\tau_{0}}{\tau_{1}} \gamma_{1, n}\left|I_{1, n}\right|^{2}$ and $+\frac{\tau_{0}}{\tau_{1}} \gamma_{1, n}\left|D_{1, n}\right|^{2}$, with the same coefficient of proportionality as the contributions on the top of the rectangle. Therefore, the energy balance can be written:

$$
\begin{align*}
\sum_{n \in U_{0}} \gamma_{0, n}\left|D_{0, n}\right|^{2} & +\sum_{n \in U_{1}} \frac{\tau_{0}}{\tau_{1}} \gamma_{1, n}\left|D_{1, n}\right|^{2}=  \tag{2.79}\\
& =\sum_{n \in U_{0}} \gamma_{0, n}\left|I_{0, n}\right|^{2}+\sum_{n \in U_{1}} \frac{\tau_{0}}{\tau_{1}} \gamma_{1, n}\left|I_{1, n}\right|^{2}
\end{align*}
$$

The first and second terms in the left-hand member of equation (2.79) represent the energy diffracted upwards and downwords respectively and the corresponding terms in the right-hand member are the incident energy propagating downwords and upwards respectively.

Coming back to the physical problem where the incident wave is unique and has a unit amplitude (see equation (2.26), equation (2.79) becomes:

$$
\begin{equation*}
\sum_{n \in U_{0}} \gamma_{0, n}\left|D_{0, n}\right|^{2}+\sum_{n \in U_{1}} \frac{\tau_{0}}{\tau_{1}} \gamma_{1, n}\left|D_{1, n}\right|^{2}=\gamma_{0,0} \tag{2.80}
\end{equation*}
$$

the right-hand member representing the incident energy. In that case, the efficiency $\rho_{i, n}, i \in$ $(0,1)$ is defined as the ratio of the energy diffracted in a given order over the incident energy. Using equation 2.79 yields:

$$
\rho_{i, n}=\left\{\begin{array}{lc}
\frac{\gamma_{0, n}}{\gamma_{0,0}}\left|D_{0, n}\right|^{2} \quad \text { if } i=0  \tag{2.81}\\
\frac{\tau_{0}}{\tau_{1}} \frac{\gamma_{1, n}}{\gamma_{0,0}}\left|D_{1, n}\right|^{2} & \text { if } i=1
\end{array}\right.
$$

and the energy balance can be written:

$$
\begin{equation*}
\sum_{n \in U_{0}} \rho_{0, n}+\sum_{n \in U_{1}} \rho_{1, n}=1 \tag{2.82}
\end{equation*}
$$

The sum of efficiencies is equal to unity. When the grating is perfectly conducting, it is easy to show that the energy balance still holds, provided that the sum $n \in U_{1}$ is cancelled in equations (2.79), (2.80) and (2.82). When the grating material is lossy, the sum $n \in U_{1}$ must be cancelled as well and one can show that equation 2.82 becomes:

$$
\begin{equation*}
\sum_{n \in U_{0}} \rho_{0, n}<1 \tag{2.83}
\end{equation*}
$$

The sum of reflected efficiencies is smaller than one, a rather intuitive result if we bear in mind that a part of the incident energy is dissipated in the grating material.

### 2.5.4 Reciprocity

In order to demonstrate the well known reciprocity relation, we consider a function $u$, sum of the solution of the normalized grating problem (see equations(2.28), 2.29), (2.30) and (2.31))


Figure 2.3: The reciprocity theorem: The efficiency in the $p^{\text {th }}$ order is the same in the two cases symbolized by red and blue arrows.
and of the corresponding incident field (in other words, $u$ is the total field). In order to define $v$, we consider the $p^{t h}$ order of diffraction $\left(p \in U_{0}\right)$ in $\mathscr{R}_{0}$, with diffraction angle $\theta_{0, p}$.

Then, we consider a second problem, but with angle of incidence $\theta^{\prime \prime}=-\theta_{0, p}$, as shown ${ }^{3}$ in figure 2.3. The incident wave in this second case has a direction of propagation which is just the opposite of that of the $p^{t h}$ diffracted order in the first case and straightforward calculations show that the corresponding $p^{t h}$ order in $\mathscr{R}_{0}$ has a direction of propagation which is the opposite of that of the incident wave in the first case, which entails $\theta_{0, p}^{\prime \prime}=-\theta$. This geometrical property is known in optics as the reversion theorem. The constants of propagation of the $p^{t h}$ diffracted order in this second case are given by $\alpha_{p}^{\prime \prime}=-\alpha_{0}$ and $\gamma_{0, p}^{\prime \prime}=\gamma_{0,0}$ and more generally, the constants of propagation of an arbitrary $n^{\text {th }}$ diffracted order in this second case are given by $\alpha_{n}^{\prime \prime}=-\alpha_{p-n}$ and $\gamma_{0, n}^{\prime \prime}=\gamma_{0, p-n}$. Thus $v^{\prime \prime}$ can be written:

$$
\begin{equation*}
v^{\prime \prime}(x, z)=\exp \left(-i \alpha_{p} x-i \gamma_{0, p} z\right)+\sum_{n=-\infty}^{+\infty} D_{0, n}^{\prime \prime} \exp \left(-i \alpha_{p-n} x+i \gamma_{0, p-n} z\right) \tag{2.84}
\end{equation*}
$$

Functions $u$ and $v^{\prime \prime}$ do not satisfy the conditions of the equation of compatibility (equation (2.71) since they have not the same pseudo-periodicity. It is not so for $u$ and the function $v=\overline{v^{\prime \prime}}$ which is given by:

$$
\begin{equation*}
v(x, z)=\exp \left(i \alpha_{p} x+i \gamma_{0, p} z\right)+\sum_{n=-\infty}^{+\infty} \overline{D_{0, n}^{\prime \prime}} \exp \left(i \alpha_{p-n} x-i \overline{\gamma_{0, p-n}} z\right) . \tag{2.85}
\end{equation*}
$$

[^3]

Figure 2.4: Other reciprocity relations: The efficiency is the same in the two cases symbolized by red and blue arrows.

Identifying the incident and diffracted waves in equation (2.85) yields:

$$
\begin{align*}
& I_{0, n}^{\prime}=\overline{D_{0, p-n}^{\prime \prime}},  \tag{2.86a}\\
& D_{0, n}^{\prime}=\delta_{n-p}, \tag{2.86b}
\end{align*}
$$

and from equation (2.71), it turns out that:

$$
\begin{equation*}
\gamma_{0, p}^{\prime \prime} D_{0, p}^{\prime \prime}=\gamma_{0, p} D_{0, p} . \tag{2.87}
\end{equation*}
$$

This is the reciprocity theorem: the products of the amplitudes of the plane waves represented in figure 2.3 by their propagation constants along the $z$ axis is invariant. In order to state the reciprocity theorem in a form which is most widespread, we take the modulus square of both members of equation (2.87):

$$
\begin{equation*}
\gamma_{0, p}^{\prime \prime}{ }^{2}\left|D_{0, p}^{\prime \prime}\right|^{2}=\gamma_{0, p}^{2}\left|D_{0, p}\right|^{2} . \tag{2.88}
\end{equation*}
$$

Writing equation (2.88) in the form:

$$
\begin{equation*}
\frac{\gamma_{0, p}^{\prime \prime}}{\gamma_{0, p}}\left|D_{0, p}^{\prime \prime}\right|^{2}=\frac{\gamma_{0, p}}{\gamma_{0, p}^{\prime \prime}}\left|D_{0, p}\right|^{2}, \tag{2.89}
\end{equation*}
$$

and bearing in mind that $\gamma_{0, p}=\gamma_{0,0}^{\prime \prime}$ and $\gamma_{0, p}^{\prime \prime}=\gamma_{0,0}$, and using the definition of the efficiencies given in equation (2.81), equation (2.89) yields:

$$
\begin{equation*}
\rho_{0, p}^{\prime \prime}=\rho_{0, p} \tag{2.90}
\end{equation*}
$$

## The efficiency is invariant.

Figure 2.4 illustrates two other cases where the reciprocity theorem applies. These properties can be demonstrated by following the same lines as in the first part of this section. It is important to notice that the reciprocity theorem illustrated in figure 2.3 holds for lossy materials [9]. More surprisingly, the theorem can be generalized to evanescent waves [10].

### 2.5.5 Uniqueness of the solution of the grating problem

If two different solutions of the normalized grating problem exist, their difference $w(x, z)$ does not include any incident wave. We will show that such a field vanish. We assume here that the
grating material is lossless. First, using the compatibility equation (2.71) with $u=v=w$, it emerges that:

$$
\begin{equation*}
\frac{1}{\tau_{0}} \sum_{n \in U_{0}} \gamma_{0, n}\left|D_{0, n}\right|^{2}+\frac{1}{\tau_{1}} \sum_{n \in U_{1}} \gamma_{1, n}\left|D_{1, n}\right|^{2}=0 . \tag{2.91}
\end{equation*}
$$

Since $\tau_{0}, \tau_{1}, \gamma_{0, n}$ and $\gamma_{1, n}$ are positive, equation (2.91) implies that $D_{0, n}=D_{1, n}=0$. This is an important result since it means that if $w$ exists, it has no effect on the far field: the solution in the far field is unique. However, it could exist a function $w$ localized at the vicinity of the grating profile and tending to zero exponentially at infinity. The interested reader can find a complete and not straightforward demonstration of the uniqueness in [1], at least for the TE case.

### 2.5.6 Analytic properties of crossed gratings



Figure 2.5: A crossed grating with periods $d_{x}$ and $d_{z}$ on the $x$ and $z$ axes.
Now, we consider the diffraction problem schematized in figure 2.5. An incident wave of wavevector $\overrightarrow{k_{0}}$ is incident on a doubly-periodic structure separating air (region $\mathscr{R}_{0}$ ) from a grating material (region $\mathscr{R}_{1}$ ). We use all the notations defined in the preceding sections to characterize the materials. The incident field is schematized in figure 2.6. The direction of incidence is specified by the polar angles $\Phi$ and $\Psi$ (see figure 2.6). In order to define the polarization of the incident field, we construct the circle MNM' N ' in the plane perpendicular to $\overrightarrow{k_{0}}$, with the continuation of $\mathrm{NN}^{\prime}$ intersecting the $z$ axis and MM' being perpendicular to NN '. The polarization angle $\delta$ is the angle between M'M and the direction of the incident electric field $\vec{P}$. With these notations, the incident electric field is given by:

$$
\begin{equation*}
\overrightarrow{E^{i}}=\vec{P} \exp (i \alpha x+i \beta y-i \gamma z) \tag{2.92}
\end{equation*}
$$



Figure 2.6: Notations for the incident field.
with $\alpha=k_{0} \sin \Phi \cos \Psi, \beta=k_{0} \sin \Phi \sin \Psi$ and $\gamma=k_{0} \cos \Phi$. The projection of $\vec{P}$ on M'M is called transverse component of $\vec{P}$ and denoted by $P^{t}$. Its projection on N'N is called longitudinal (in plane) component and denoted $\mathrm{bz} P^{l}$, in such a way that $\vec{P}=P^{t} \frac{\overrightarrow{M M^{\prime}}}{M M^{\prime}}+P^{l} \frac{\overrightarrow{N N^{\prime}}}{N N^{\prime}}$.

As in the case of classical gratings, it is possible to show that above the top of the grating $\left(z>z_{M}\right)$, the field can be expanded in the form of a sum of plane waves:

$$
\vec{E}(x, z)=\left\{\begin{array}{c}
\sum_{n=-\infty}^{+\infty} \sum_{m=-\infty}^{+\infty}\left(\overrightarrow{I_{0, n, m}} \exp \left(i \alpha_{n} x+i \beta_{m} y-i \gamma_{0, n, m} z\right)+\right.  \tag{2.93}\\
\left.+\overrightarrow{D_{0, n}, m} \exp \left(i \alpha_{n} x+i \beta_{m} y+i \gamma_{0, n, m} z\right)\right), \quad \text { if } z>z_{M} \\
\sum_{n=-\infty}^{+\infty} \sum_{m=-\infty}^{+\infty}\left(\overrightarrow{D_{1, n, m}} \exp \left(i \alpha_{n} x+i \beta_{m} y-i \gamma_{1, n, m} z\right)+\right. \\
\left.+\overrightarrow{I_{1, n, m}} \exp \left(i \alpha_{n} x+i \beta_{m} y+i \gamma_{1, n, m} z\right)\right) \quad \text { if } z<0
\end{array}\right.
$$

The wavevectors of all these plane waves must be orthogonal to their vector amplitudes. As for the incident wave, we can define the transverse and longitudinal components of the vector amplitudes of the plane waves, the transverse component (for example $D_{0, n, m}^{t}$ ) being orthogonal to the $z$ axis in the plane perpendicular to the wavevector $\left(\alpha_{n}, \beta_{m}, \gamma_{0, n, m}\right)$ and the longitudinal (for example $D_{0, n, m}^{l}$ ) its component in the orthogonal direction of the same plane.

Using the Poynting theorem, it can be shown, as in section 2.5.3, that the efficiencies in the $z$-propagating orders are given by:

$$
\rho_{i, n, m}= \begin{cases}\frac{\gamma_{0, n, m}}{\gamma_{0,0}}\left|\vec{D}_{0, n, m}\right|^{2} & \text { if } i=0  \tag{2.94}\\ \frac{\gamma_{1, n, m}}{\gamma_{0,0}}\left(\frac{1}{v^{2}}\left|D_{1, n, m}^{l}\right|^{2}+\left|D_{1, n, m}^{t}\right|^{2}\right) & \text { if } i=1\end{cases}
$$

Of course, the line associated to $i=1$ in equation (2.94) must be cancelled if the grating material is lossy.

We define, as for classical gratings, the sets $U_{0}$ and $U_{1}$ of $z$-propagating orders in $\mathscr{R}_{0}$ and $\mathscr{R}_{1}$ respectively and, when the grating material is lossless, the energy balance can be written:

$$
\begin{equation*}
\sum_{(n, m) \in U_{0}} \rho_{0, n, m}+\sum_{(n, m) \in U_{1}} \rho_{1, n, m}=1 \tag{2.95}
\end{equation*}
$$

We will not demonstrate the reciprocity theorem, the interested reader can find the proof in [1]. This theorem, in the case of an order $(p, q)$ propagating in $\mathscr{R}_{0}$ can be expressed in the following form:

$$
\begin{equation*}
\gamma \vec{P} \cdot{\overrightarrow{D^{\prime}}}_{0, p, q}=\gamma^{\prime} \overrightarrow{P^{\prime}} \cdot \vec{D}_{0, p, q} \tag{2.96}
\end{equation*}
$$

In the first case, the incident electric field with vector amplitude $\vec{P}$ and propagation constant along the $z$ axis $-\gamma$ generates in $\mathscr{R}_{0}$ in the $(p, q)$ order, with $(p, q) \in U_{0}$, a plane wave of vector amplitude $\vec{D}_{0, p, q}$ and propagation constant along the $z$ axis $\gamma_{0, p, q}$. In the second case, we consider an incident wave which propagates in the direction which is just the opposite to that of the $(p, q)$ order in the first case. Thus its constant of propagation along the $z$ axis is $-\gamma^{\prime}=-\gamma_{0, p, q}$. The vector amplitude of this incident wave is equal to $\overrightarrow{P^{\prime}}$. It can be shown that in this second case, the $(p, q)$ order takes the direction which is the opposite of that of the incident wave in the first case and its vector amplitude is equal to ${\overrightarrow{D^{\prime}}}_{0, p, q}$. Thus, equation (2.96) can be expressed in the following form: the scalar product of the vector amplitudes of the incident and diffracted waves propagating in the opposite directions, multiplied by the propagation constant of the incident wave along the $z$ axis, is constant. It can be shown that this relation entails the reciprocity in natural light for the efficiencies:

$$
\begin{equation*}
\left.<\rho_{0, p, q}\right\rangle=<\rho_{0, p, q}^{\prime}> \tag{2.97}
\end{equation*}
$$

with $<\rho_{0, p, q}>$ being the average between the efficiencies in both cases of polarization ( $\delta=0$ and $\delta=\frac{\pi}{2}$ ).

### 2.6 Conclusion

We have established the mathematical bases of grating theories: the boundary-value problems. Most of the formalisms used for solving the grating problems numerically start from these boundary-value problems, for example the integral theory [1,2].Other theories use some conditions of these problems but deal directly with Maxwell equations, for example the RCWA method [5].

Without any doubt, the boundary-value problems are necessary to demonstrate the analytic properties of gratings. Very often, these properties are ignored or neglected. However, properties like energy balance or reciprocity are needed for a full understanding of the puzzling properties of this crucial component of optics and nanophotonics. These analytic properties are also widely used to check new grating softwares. However, they are not more than casting out nines. They can show that a software fails if they are nor satisfied on its numerical results. It must be emphasized that they can never prove its validity if they are satisfied.

Some important analytic properties of gratings have not been mentioned in this chapter. It is the case for example for the Marechal and Stroke theorem, the only grating property which allows one to know the field diffracted by a grating without any calculation. This theorem, which is restricted to perfecly conducting echelette gratings used for TM polarization in very special conditions will be given in the chapter devoted to the applications of grating properties.

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Chapter 3:

# Spectral Methods for Gratings <br> John A. DeSanto 

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## Chapter 3

# Spectral Methods for Gratings 

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#### Abstract

We present a unified formal treatment of spectral methods applied to scattering from penetrable gratings for both acoustic (scalar) and electromagnetic (vector) problems. These are derived from coordinate space representations for both acoustic problems for a one-dimensional grating, and full electromagnetic problems for a two-dimensional grating. The coordinate space representations are also derived here. By unified we mean that the electromagnetic results use a scalar analogy, in that the boundary unknowns are the electric field and its normal derivative.

In coordinate space, the kernels of either the integral representations or integral equations have two variables, the first relating to the field coordinate (which, for an integral equation, has been evaluated on the surface), and the second evaluated on the surface as part of the surface integration. We refer to this procedure as coordinate in both variables or simply a CC method. Partial spectral results involve a spectral replacement of the first coordinate variable, and we refer to these methods as $S C$. Full spectral methods involve an additional replacement of the second (always surface) coordinate variable by a spectral one and we refer to these as $S S$ methods, or in the case of a conjugate Rayleigh basis as $S S^{*}$.

For both scalar and electromagnetic cases, the partial spectral results are derived without the use of Green's functions. Instead we use plane wave states in Green's theorem. The partial spectral results are also used to generate surface inversion methods involving perturbation theory and the Kirchhoff approximation. For the full spectral scalar case three spectral expansions are considered, a physical optics modified Fourier expansion, a Floquet-Fourier expansion, and an expansion in conjugate Rayleigh basis functions. All are Floquet- or quasi-periodic. For the full spectral electromagnetic case only the conjugate Rayleigh basis expansion is presented.


### 3.1 Introduction

In this paper we present formal equations in spectral space to describe the scattering from periodic surfaces or gratings. We do this both for the acoustic case in one dimension for the direct scattering problem (Secs.4,5,7) and the inverse problem(Sec.6), and in two dimensions for the general electromagnetic problem (Secs. $9,10,11$ ). In order to do this and justify the validity of the spectral representations in various regions, we derive in each case the coordinate-space representations for the scattering (Secs.4,9). This includes the two- and three-dimensional periodic Green's function (Sec.3) necessary to describe the coordinate-space scattering, and the development of plane wave expansions in periodic media (Sec.2) which are used throughout the paper.

For the partial spectral-space equations it is not necessary to use the Green's function. Instead we describe a method using plane waves and Green's theorem in the periodic cell of the surface to derive the spectral equations directly and simply. One can describe the coordinatespace integral representations or integral equations to solve the boundary unknowns as equations in two coordinate spaces, the space of the source or integrated coordinate, and the space of the field or exterior coordinate. For clarity, we refer to this as a coordinate-coordinate (CC) representation. Spectral can then refer to one or both of these coordinate spaces transformed, a partial spectral space when one is transformed (Secs.5,10) where we use first the transform of the exterior coordinate (following from the plane waves and Green's theorem) so we are in a spectral-coordinate (SC) representation which we treat extensively, including its use in the inverse problem, the problem of surface reconstruction from (known) scattered field data. A second version (CS) is referred to ${ }^{033}$ but not extensively discussed. Using SC we can then represent the integrated coordinate in spectral space (SS) where we reference extensive computational results ${ }^{[663738 / 39}$. What has become of interest lately ${ }^{33}$ is the transform of the coordinate-space of integration into the conjugate spectral space $S^{\star}$, and we describe this $S S^{\star}$ method extensively (Secs.8,11) for acoustic and electromagnetic results respectively. Formally this is equivalent to a dual least squares method.

It is important to define what we mean by the word "spectral" ${ }^{28}$. In a general context, this could mean just Fourier space, and boundary function expansions in pure Fourier series. However, in our context, this is not correct. The reason is that the boundary fields are limits of Floquet- or quasi-periodic functions and must themselves be Floquet- or quasi-periodic. This reflects the nature of the incident field being, in general, incident off the normal. The periodic surface has right-left symmetry, but the boundary value problem does not (unless the field is incident normally). By "spectral" we thus mean here three kinds of expansions, a physicaloptics modified Fourier expansion (Sec.7) where a single plane wave modulates the Fourier series, an expansion which we term Floquet-Fourier which preserves the quasi-periodicity but without the modulating plane wave (Sec.7), and an expansion in conjugate plane wave states on the boundary (Secs.8,11) where plane waves modulate each term in the expansion. These latter are Rayleigh or Bloch wave type expansions and are useful because, at least in some degenerate cases, they lead to self-adjoint problems. All expansions are quasi-periodic. Relation of the results to Rayleigh and Waterman expansions is discussed.

Using the partial spectral methods developed in Sec. 5 for the direct scattering problem, we discuss in Sec. 6 how they can be used to find the periodic surface profile from the incident and scattered fields. Two methods are presented, one based on perturbation theory and the other on the Kirchhoff approximation, both for the scalar case. Both surface inversion methods were initially applied to truncated random surfaces ${ }^{[111 \mid 112}$ with good results, and the methodology is
here applied to gratings.
The electromagnetic equations we present are not done in the conventional formalism ${ }^{61}$ using boundary currents, but are based on earlier work of ours ${ }^{35]}$ which rely on a scalar analogue of the electromagnetic problem, so the boundary unknowns we use are the electric field and its normal derivative. The resulting equations become a direct scalar analogue in terms of different boundary unknowns for the electromagnetic problem. We also do not discuss the computational solutions of the equations but rely on references where available. Some related equations have been solved, and we point out the references where appropriate, but a full discussion of computational issues would require a separate paper.

A very large number of references are cited in the text. Many of these references, some not specifically attuned to spectral methods, nevertheless contain spectral components in the development or in reference to expansions, in particular to the Rayleigh expansion ${ }^{[11|12| 1448|50| 54|56| 57,}$
 the two expansions ${ }^{[62}$. The gratings we consider are infinite, although spectral methods have been applied to finite gratings also ${ }^{84}$, and our gratings are purely deterministic although random gratings have also been considered ${ }^{[83}$. Newer results on gratings apply surface integral methods to periodic nanostructures ${ }^{43}$, indicating the generality of the methods we describe. We are mainly interested in scattering methods which produce the scattered and transmitted fields and their use in inversion, although others prefer to consider the dispersion relation for surface plasmons and polaritons propagating along the grating ${ }^{4664763}$. Other rigorous theoretical and computational developments are also available ${ }^{2|677816 / 4445549167| 68177778 \mid 79871921961113}$, as well as approximations ${ }^{655110}$, applications ${ }^{913142 \mid 51711755}$, and other methodology ${ }^{[185]}$.

There are very many papers (our many references do not scratch the surface), reviews ${ }^{10140}$, 69888999198 and books ${ }^{5211031108}$ on scattering from gratings, not the least of them being the important book edited by Petit ${ }^{[90}$ in honor of which this paper is contributed.

### 3.2 Plane Waves in Periodic Media

In two dimensions $\vec{x}=(x, z)$, we write a plane wave as

$$
\begin{equation*}
\phi(\vec{x})=\exp \left[i k\left(\alpha_{0} x+\gamma_{0} z\right)\right], \tag{3.1}
\end{equation*}
$$

where $\alpha_{0}=\sin (\theta), \gamma_{0}=\cos (\theta)=\sqrt{1-\alpha_{0}^{2}}$, and $\theta$ is measured clockwise from the positive $z$ axis. With the time convention $\exp (-i \omega t)$, where $\omega$ is circular frequency, this is thus an up-going plane wave, and satisfies the two-dimensional Helmholtz equation

$$
\begin{equation*}
\left(\nabla_{2}^{2}+k^{2}\right) \phi(\vec{x})=0 \tag{3.2}
\end{equation*}
$$

where $k$ is the wavenumber (here considered to be strictly real). On a one-dimensional surface $z=h(x)$ we have

$$
\begin{equation*}
\phi\left(\vec{x}_{h}\right)=\exp \left[i k\left(\alpha_{0} x+\gamma_{0} h(x)\right)\right], \tag{3.3}
\end{equation*}
$$

where $\vec{x}_{h}=(x, h(x))$. This is referred to as a Rayleigh function. For one-dimensional periodic media (here the surface), $h(x+L)=h(x)$, where $L$ is the period, we have the relation

$$
\begin{equation*}
\phi(x+L, h(x+L))=\exp \left(i k \alpha_{0} L\right) \phi(x, h(x)) . \tag{3.4}
\end{equation*}
$$

The same result is true even off the surface, i.e.

$$
\begin{equation*}
\phi(x+L, z)=\exp \left(i k \alpha_{0} L\right) \phi(x, z) . \tag{3.5}
\end{equation*}
$$

These results are referred to as Floquet- or quasi-periodicity. The field scattered from this periodic surface, $\psi^{s c}$, satisfies the same Helmholtz equation, and is also quasi-periodic because the ratio $\frac{\psi^{s c}}{\phi}$ is periodic, i.e.

$$
\begin{equation*}
\frac{\psi^{s c}(x+L, z)}{\phi(x+L, z)}=\frac{\psi^{s c}(x, z)}{\phi(x, z)}, \tag{3.6}
\end{equation*}
$$

so that

$$
\begin{equation*}
\psi^{s c}(x+L, z)=\exp \left(i k \alpha_{0} L\right) \psi^{s c}(x, z), \tag{3.7}
\end{equation*}
$$

and the same is true on the surface $z=h(x)$. In general for any field function $\psi$ satisfying (3.2) and shifted an integer $n$ number of periods we have

$$
\begin{equation*}
\psi(x+n L, z)=\exp \left(i k \alpha_{0} n L\right) \psi(x, z) \tag{3.8}
\end{equation*}
$$

The same is of course true on the surface $z=h(x)$.
In three dimensions any field function which satisfies the three-dimensional Helmholtz equation

$$
\begin{equation*}
\left(\nabla_{3}^{2}+k^{2}\right) \psi(\vec{x})=0, \tag{3.9}
\end{equation*}
$$

where $\vec{x}=(x, y, z)$, and is quasi-periodic in $x$ with period $L_{1}$ and in $y$ with period $L_{2}$ satisfies

$$
\begin{equation*}
\psi\left(x+n_{1} L_{1}, y+n_{2} L_{2}, z\right)=\exp \left[i k\left(\alpha_{0} n_{1} L_{1}+\beta_{0} n_{2} L_{2}\right)\right] \psi(x, y, z), \tag{3.10}
\end{equation*}
$$

where $\alpha_{0}=\sin (\theta) \cos (\varphi)$ and $\beta_{0}=\sin (\theta) \sin (\varphi)$ with $\theta$ the polar angle, $\varphi$ the azimuthal angle, and $n_{1}$ and $n_{2}$ are integers. In three dimensions the up-going plane wave is now written as

$$
\begin{equation*}
\phi(\vec{x})=\exp \left[i k\left(\alpha_{0} x+\beta_{0} y+\gamma_{0} z\right)\right], \tag{3.11}
\end{equation*}
$$

where $\gamma_{0}=\sqrt{1-\alpha_{0}^{2}-\beta_{0}^{2}}$. Although we use some of the same notation in both two and three dimensions the interpretation will be clear from the context.

We consider the periodic surface $z=h(x)$ in one dimension and $z=h(x, y)=h\left(\vec{x}_{\perp}\right)$ in two dimensions which separates two media with wavenumbers $k_{1}$ for $z>h$ (the upper region 1) and $k_{2}$ for $z<h$ (the lower region 2). Notationally, subscripts are used to identify the region, e.g. $\phi$ becomes $\phi_{1}$ or $\phi_{2}, \gamma_{0}$ becomes $\gamma_{10}$ or $\gamma_{20}$, etc. The paper has a lot of notation, and we have tried to keep it as clear as possible.

### 3.3 Green's Functions in Periodic Media

In two dimensions the free-space Green's function is

$$
\begin{equation*}
G^{(2)}\left(\vec{x}^{\prime}, \vec{x}\right)=\frac{i}{4} H_{0}^{(1)}\left(k_{0}\left|\vec{x}^{\prime}-\vec{x}\right|\right), \tag{3.12}
\end{equation*}
$$

where $H_{0}^{(1)}$ is the Hankel function, and $k_{0}$ is a generic wave number. It satisfies the equation

$$
\begin{equation*}
\left(\nabla_{2}^{2}+k_{0}^{2}\right) G^{(2)}\left(\vec{x}^{\prime}, \vec{x}\right)=-\delta\left(\vec{x}^{\prime}-\vec{x}\right) . \tag{3.13}
\end{equation*}
$$

Its representation as a Fourier transform is

$$
\begin{equation*}
G^{(2)}\left(\vec{x}^{\prime}, \vec{x}\right)=\frac{1}{(2 \pi)^{2}} \iint \frac{\exp \left[i k_{x}\left(x^{\prime}-x\right)+i k_{z}\left(z^{\prime}-z\right)\right]}{k^{2}-k_{0+}^{2}} d k_{x} d k_{z}, \tag{3.14}
\end{equation*}
$$

where the integrals run from $-\infty$ to $\infty$. We have given $k_{0}$ a small positive imaginary part to define the integral, and $k^{2}=k_{x}^{2}+k_{z}^{2}$. If we choose a specific direction, here the fixed direction $z$, we can evaluate the $k_{z}$ integration using complex variables. The result is the Weyl representation ${ }^{34}$ for $G^{(2)}$

$$
\begin{equation*}
G^{(2)}\left(\vec{x}^{\prime}, \vec{x}\right)=\frac{i \pi}{(2 \pi)^{2}} \int_{-\infty}^{\infty} \frac{\exp \left[i k_{x}\left(x^{\prime}-x\right)+i K_{0}\left|z^{\prime}-z\right|\right]}{K_{0}} d k_{x}, \tag{3.15}
\end{equation*}
$$

where $K_{0}=\sqrt{k_{0}^{2}-k_{x}^{2}}$ for $k_{0}^{2}>k_{x}^{2}$, and $=i \sqrt{k_{x}^{2}-k_{0}^{2}}$ for $k_{x}^{2}>k_{0}^{2}$.
We have two regions. In region 1 , we let $k_{0}=k_{1}$ in (3.15), scale the integral using $k_{x}=$ $k_{1} \alpha$, and we get the Green's function for region 1 (subscript) in (2)-dimensions (superscript)

$$
\begin{equation*}
G_{1}^{(2)}\left(\vec{x}^{\prime}, \vec{x}\right)=\frac{i \pi}{(2 \pi)^{2}} \int_{-\infty}^{\infty} \frac{\exp \left[i k_{1}\left(\alpha\left(x^{\prime}-x\right)+\gamma_{1}(\alpha)\left|z^{\prime}-z\right|\right)\right]}{\gamma_{1}(\alpha)} d \alpha, \tag{3.16}
\end{equation*}
$$

where

$$
\begin{align*}
\gamma_{1}(\alpha) & =\sqrt{1-\alpha^{2}}, \quad\left(\alpha^{2}<1\right)  \tag{3.17}\\
& =+i \sqrt{\alpha^{2}-1}, \quad\left(\alpha^{2}>1\right) \tag{3.18}
\end{align*}
$$

In region 2 , let $k_{0}=k_{2}$ in (3.15) and scale using $k_{x}=k_{1} \alpha$ (the same scaling as in region 1) to get the Green's function in region 2 (subscript) in (2)-dimensions (superscript)

$$
\begin{equation*}
G_{2}^{(2)}\left(\vec{x}^{\prime}, \vec{x}\right)=\frac{i \pi}{(2 \pi)^{2}} \int_{-\infty}^{\infty} \frac{\exp \left[i k_{1}\left(\alpha\left(x^{\prime}-x\right)+\gamma_{2}(\alpha)\left|z^{\prime}-z\right|\right)\right]}{\gamma_{2}(\alpha)} d \alpha, \tag{3.19}
\end{equation*}
$$

where

$$
\begin{align*}
\gamma_{2}(\alpha) & =\sqrt{K^{2}-\alpha^{2}}, \quad\left(\alpha^{2}<K^{2}\right)  \tag{3.20}\\
& =+i \sqrt{\alpha^{2}-K^{2}}, \quad\left(\alpha^{2}>K^{2}\right) \tag{3.21}
\end{align*}
$$

and $K=k_{2} / k_{1}$, the ratio of wavenumbers. The same scaling in both regions can be thought of as simply a result of Snell's Law since the $x$-components of the phases of both functions must match at a flat interface.

The Green's functions above are for an infinite space or, in our case, an infinite surface. To find the periodic Green's function for a single cell of the surface we use the single or double layer potentials which occur in Sec.4. For example, define the single layer potential on an infinite surface $h(x)$ as

$$
\begin{equation*}
(S \psi)\left(\vec{x}_{h}^{\prime}\right)=\int_{-\infty}^{\infty} G^{(2)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right) \psi\left(\vec{x}_{h}\right) d x \tag{3.22}
\end{equation*}
$$

where $\psi$ is any field function (here it is the normal derivative). The result can be written as a sum over periodic cells

$$
\begin{equation*}
(S \psi)\left(\vec{x}_{h}^{\prime}\right)=\sum_{n=-\infty}^{\infty} I_{n}\left(x^{\prime}\right), \tag{3.23}
\end{equation*}
$$

where

$$
\begin{equation*}
I_{n}\left(x^{\prime}\right)=\int_{(2 n-1) L / 2}^{(2 n+1) L / 2} G^{(2)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right) \psi\left(\vec{x}_{h}\right) d x \tag{3.24}
\end{equation*}
$$

Use (3.16) or 3.19) in 3.24, shift the integration by defining $x^{\prime \prime}=x-n L$, and use the Floquet property of the field function to rewrite (3.22) as

$$
\begin{equation*}
(S \psi)\left(\vec{x}_{h}^{\prime}\right)=\int_{-L / 2}^{L / 2} G^{(2 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right) \psi\left(\vec{x}_{h}\right) d x \tag{3.25}
\end{equation*}
$$

where the two-dimensional periodic Green's function $((2 p)$-superscript $)$ is given by

$$
\begin{equation*}
G^{(2 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)=\frac{i \pi}{(2 \pi)^{2}} \int_{-\infty}^{\infty} \frac{\exp \left[i k_{1}\left(\alpha\left(x^{\prime}-x\right)+\gamma(\alpha)\left|h\left(x^{\prime}\right)-h(x)\right|\right)\right]}{\gamma(\alpha)} S(\alpha) d \alpha \tag{3.26}
\end{equation*}
$$

where the sum $S$ is given by

$$
\begin{equation*}
S(\alpha)=\sum_{n=-\infty}^{\infty} \exp \left[i n k_{1} L\left(\alpha_{0}-\alpha\right)\right] \tag{3.27}
\end{equation*}
$$

and can be evaluated using the Poisson sum ${ }^{955}$ to be

$$
\begin{equation*}
S(\alpha)=\frac{2 \pi}{i k_{1}} \sum_{j=-\infty}^{\infty} \delta\left(\alpha-\alpha_{j}\right) \tag{3.28}
\end{equation*}
$$

where $\delta$ represents the delta function and $\alpha_{j}=\alpha_{0}+j \lambda / L$ is the grating equation. The result substituted in (3.26) yields the periodic Green's function for region 1

$$
\begin{equation*}
G_{1}^{(2 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)=\frac{i}{2 k_{1} L} \sum_{j=-\infty}^{\infty} \frac{\exp \left[i k_{1}\left(\alpha_{j}\left(x^{\prime}-x\right)+\gamma_{1 j}\left|h\left(x^{\prime}\right)-h(x)\right|\right)\right]}{\gamma_{1 j}} \tag{3.29}
\end{equation*}
$$

where

$$
\begin{align*}
\gamma_{1 j} & =\sqrt{1-\alpha_{j}^{2}}, \quad\left(\alpha_{j}^{2}<1\right)  \tag{3.30}\\
& =+i \sqrt{\alpha_{j}^{2}-1}, \quad\left(\alpha_{j}^{2}>1\right) \tag{3.31}
\end{align*}
$$

and the periodic Green's function for region 2

$$
\begin{equation*}
G_{2}^{(2 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)=\frac{i}{2 k_{1} L} \sum_{j=-\infty}^{\infty} \frac{\exp \left[i k_{1}\left(\alpha_{j}\left(x^{\prime}-x\right)+\gamma_{2 j}\left|h\left(x^{\prime}\right)-h(x)\right|\right)\right]}{\gamma_{2 j}}, \tag{3.32}
\end{equation*}
$$

where

$$
\begin{align*}
\gamma_{2 j} & =\sqrt{K^{2}-\alpha_{j}^{2}}, \quad\left(\alpha_{j}^{2}<K^{2}\right)  \tag{3.33}\\
& =+i \sqrt{\alpha_{j}^{2}-K^{2}}, \quad\left(\alpha_{j}^{2}>K^{2}\right) \tag{3.34}
\end{align*}
$$

We have listed the Green's functions of both regions to stress the exterior scaling $k_{1}$ for both. The result is the residual of the $k_{1}$ in the phase of both terms, Snell's law, and the same Poisson sum.

The Green's functions satisfy the differential equations

$$
\begin{equation*}
\left(\nabla_{2}^{2}+k_{l}^{2}\right) G_{l}^{(2 p)}\left(\vec{x}^{\prime}, \vec{x}\right)=-\sum_{n=-\infty}^{\infty} \delta\left(x^{\prime}-x_{n}\right) \delta\left(z^{\prime}-z\right) \tag{3.35}
\end{equation*}
$$

where $x_{n}=x+n L$ and $l=1,2$. The periodic Green's function can also be written as a phased array of Hankel functions, e.g. for region 1

$$
\begin{equation*}
G_{1}^{(2 p)}\left(\vec{x}^{\prime}, \vec{x}\right)=\frac{i}{4} \sum_{n=-\infty}^{\infty} \exp \left(i k_{1} \alpha_{0} n L\right) H_{0}^{(1)}\left(k_{1} \sqrt{\left(x^{\prime}-x_{n}\right)^{2}+\left(z^{\prime}-z\right)^{2}}\right) \tag{3.36}
\end{equation*}
$$

The periodic Green's functions are also Floquet-periodic. For either Green's function, using (3.29) or (3.32) we have that

$$
\begin{equation*}
G^{(2 p)}\left(\vec{x}^{\prime}, \vec{x}_{n}\right)=\exp \left(-i k_{1} \alpha_{0} n L\right) G^{(2 p)}\left(\vec{x}^{\prime}, \vec{x}\right) \tag{3.37}
\end{equation*}
$$

where $\vec{x}_{n}=\vec{x}+\hat{i} n L$. Since the Floquet condition on any field function 3.8 has the conjugate phase of (3.37), the product of any Green's function times any field function $\psi$ is periodic,

$$
\begin{equation*}
G^{(2 p)}\left(\vec{x}^{\prime}, \vec{x}_{n}\right) \psi\left(\vec{x}_{n}\right)=G^{(2 p)}\left(\vec{x}^{\prime}, \vec{x}\right) \psi(\vec{x}) \tag{3.38}
\end{equation*}
$$

This result will be used later to cancel vertical integrals in Green's theorem for the coordinatespace representation.

The three-dimensional Green's function in free space is given by

$$
\begin{equation*}
G^{(3)}\left(\vec{x}^{\prime}, \vec{x}\right)=\frac{1}{4 \pi} \frac{\exp \left(i k_{0}\left|\vec{x}^{\prime}-\vec{x}\right|\right)}{\left|\vec{x}^{\prime}-\vec{x}\right|} \tag{3.39}
\end{equation*}
$$

where $k_{0}$ is a generic wave number. It satisfies the equation

$$
\begin{equation*}
\left(\nabla_{3}^{2}+k_{0}^{2}\right) G^{(3)}\left(\vec{x}^{\prime}, \vec{x}\right)=-\delta\left(\vec{x}^{\prime}-\vec{x}\right) \tag{3.40}
\end{equation*}
$$

where $\vec{x}=(x, y, z)$. Its Fourier representation is

$$
\begin{equation*}
G^{(3)}\left(\vec{x}^{\prime}, \vec{x}\right)=\frac{1}{(2 \pi)^{3}} \iiint \frac{\exp \left[i k_{x}\left(x^{\prime}-x\right)+i k_{y}\left(y^{\prime}-y\right)+i k_{z}\left(z^{\prime}-z\right)\right]}{k^{2}-k_{0+}^{2}} d k_{x} d k_{y} d k_{z} \tag{3.41}
\end{equation*}
$$

with $k^{2}=k_{x}^{2}+k_{y}^{2}+k_{z}^{2}$. Use complex integration on the preferred $z$-direction to yield

$$
\begin{equation*}
G^{(3)}\left(\vec{x}^{\prime}, \vec{x}\right)=\frac{i \pi}{(2 \pi)^{3}} \iint \frac{\exp \left[i k_{x}\left(x^{\prime}-x\right)+i k_{y}\left(y^{\prime}-y\right)+i K_{0}\left|z^{\prime}-z\right|\right]}{K_{0}} d k_{x} d k_{y} \tag{3.42}
\end{equation*}
$$

where

$$
\begin{align*}
K_{0} & =\sqrt{k_{0}^{2}-k_{x}^{2}-k_{y}^{2}}, \quad\left(k_{x}^{2}+k_{y}^{2}<k_{0}^{2}\right)  \tag{3.43}\\
& =+i \sqrt{k_{x}^{2}+k_{y}^{2}-k_{0}^{2}}, \quad\left(k_{x}^{2}+k_{y}^{2}>k_{0}^{2}\right) . \tag{3.44}
\end{align*}
$$

In the upper region, let $k_{0}=k_{1}$ in (3.42), and scale the wavenumbers as $k_{x}=k_{1} \alpha$ and $k_{y}=k_{1} \beta$. This yields the three-dimensional Green's function for region 1 in the Weyl representation

$$
\begin{equation*}
G_{1}^{(3)}\left(\vec{x}^{\prime}, \vec{x}\right)=\frac{i \pi k_{1}}{(2 \pi)^{3}} \iint \frac{\exp \left[i k_{1}\left(\alpha\left(x^{\prime}-x\right)+\beta\left(y^{\prime}-y\right)+\gamma_{1}(\alpha, \beta)\left|z^{\prime}-z\right|\right)\right]}{\gamma_{1}(\alpha, \beta)} d \alpha d \beta \tag{3.45}
\end{equation*}
$$

where

$$
\begin{align*}
\gamma_{1} & =\sqrt{1-\alpha^{2}-\beta^{2}}, \quad\left(\alpha^{2}+\beta^{2}<1\right)  \tag{3.46}\\
& =+i \sqrt{\alpha^{2}+\beta^{2}-1}, \quad\left(\alpha^{2}+\beta^{2}>1\right) \tag{3.47}
\end{align*}
$$

In the lower region, let $k_{0}=k_{2}$ in (3.42), scale the wavenumbers the same (two-dimensional Snell's law) to yield the three-dimensional Green's function for region 2 in the Weyl representation

$$
\begin{equation*}
G_{2}^{(3)}\left(\vec{x}^{\prime}, \vec{x}\right)=\frac{i \pi k_{1}}{(2 \pi)^{3}} \iint \frac{\exp \left[i k_{1}\left(\alpha\left(x^{\prime}-x\right)+\beta\left(y^{\prime}-y\right)+\gamma_{2}(\alpha, \beta)\left|z^{\prime}-z\right|\right)\right]}{\gamma_{2}(\alpha, \beta)} d \alpha d \beta \tag{3.48}
\end{equation*}
$$

where

$$
\begin{align*}
\gamma_{2} & =\sqrt{K^{2}-\alpha^{2}-\beta^{2}}, \quad\left(\alpha^{2}+\beta^{2}<K^{2}\right)  \tag{3.49}\\
& =+i \sqrt{\alpha^{2}+\beta^{2}-K^{2}}, \quad\left(\alpha^{2}+\beta^{2}>K^{2}\right) \tag{3.50}
\end{align*}
$$

The above results are for an infinite surface. To illustrate the reduction to a single cell of a two-dimensional periodic surface we choose a single layer potential (for either region) with density $\psi$ which is any field function (here the normal derivative of the velocity potential)

$$
\begin{equation*}
(S \psi)\left(\vec{x}_{h}^{\prime}\right)=\iint_{-\infty}^{\infty} G^{(3)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right) \psi\left(\vec{x}_{h}\right) d x d y . \tag{3.51}
\end{equation*}
$$

Here the surface is doubly periodic (period $L_{1}$ in $x$ and $L_{2}$ in $y$ )

$$
\begin{equation*}
h\left(\vec{x}_{\perp}+\vec{x}_{n_{1} n_{2}}\right)=h\left(\vec{x}_{\perp}\right), \tag{3.52}
\end{equation*}
$$

where $n_{1}$ and $n_{2}$ are integers and $\vec{x}_{n_{1} n_{2}}=\hat{i} n_{1} L_{1}+\hat{j} n_{2} L_{2}$. The field function is Floquet-periodic in two dimensions, see (3.10). We can thus write (3.51) as a double sum over periodic cells

$$
\begin{equation*}
(S \psi)\left(\vec{x}_{h}^{\prime}\right)=\sum_{n_{1}=-\infty}^{\infty} \sum_{n_{2}=-\infty}^{\infty} I_{n_{1} n_{2}}\left(\vec{x}_{h}^{\prime}\right), \tag{3.53}
\end{equation*}
$$

where

$$
\begin{equation*}
I_{n_{1} n_{2}}\left(\vec{x}_{h}^{\prime}\right)=\int_{\left(2 n_{2}-1\right) L_{2} / 2}^{\left(2 n_{2}+1\right) L_{2} / 2} \int_{\left(2 n_{1}-1\right) L_{1} / 2}^{\left(2 n_{1}+1\right) L_{1} / 2} G^{(3)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right) \psi\left(\vec{x}_{h}\right) d x d y \tag{3.54}
\end{equation*}
$$

Use the Weyl representation (3.45) or 3.48) for $G^{(3)}$, shift the integrations using $x^{\prime \prime}=x-n_{1} L_{1}$ and $y^{\prime \prime}=y-n_{2} L_{2}$ to yield

$$
\begin{equation*}
(S \psi)\left(\vec{x}_{h}^{\prime}\right)=\int_{-L_{2} / 2}^{L_{2} / 2} \int_{-L_{1} / 2}^{L_{1} / 2} G^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right) \psi\left(\vec{x}_{h}\right) d x d y \tag{3.55}
\end{equation*}
$$

where the three-dimensional periodic Green's function is given by

$$
\begin{equation*}
G^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)=\frac{i \pi k_{1}}{(2 \pi)^{3}} \iint \frac{\exp \left[i k_{1}\left(\alpha\left(x^{\prime}-x\right)+\beta\left(y^{\prime}-y\right)+\gamma\left|h\left(\vec{x}_{\perp}^{\prime}\right)-h\left(\vec{x}_{\perp}\right)\right|\right)\right]}{\gamma(\alpha, \beta)} P_{1} P_{2} d \alpha d \beta \tag{3.56}
\end{equation*}
$$

with the Poisson sums

$$
\begin{equation*}
P_{1}(\alpha)=\sum_{n_{1}=-\infty}^{\infty} \exp \left[i n_{1} k_{1} L_{1}\left(\alpha_{0}-\alpha\right)\right]=\frac{2 \pi}{k_{1} L_{1}} \sum_{j=-\infty}^{\infty} \delta\left(\alpha_{j}-\alpha\right) \tag{3.57}
\end{equation*}
$$

and

$$
\begin{equation*}
P_{2}(\beta)=\sum_{n_{2}=-\infty}^{\infty} \exp \left[i i_{2} k_{1} L_{2}\left(\beta_{0}-\beta\right)\right]=\frac{2 \pi}{k_{1} L_{2}} \sum_{j^{\prime}=-\infty}^{\infty} \delta\left(\beta_{j^{\prime}}-\beta\right) \tag{3.58}
\end{equation*}
$$

The grating equations are now $\alpha_{j}=\alpha_{0}+j \lambda / L_{1}$ and $\beta_{j^{\prime}}=\beta_{0}+j^{\prime} \lambda / L_{2}$. The result is the threedimensional periodic Green's function for region 1 with both coordinates on the surface

$$
\begin{equation*}
G_{1}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)=\frac{i}{2 k_{1} L_{1} L_{2}} \sum_{j=-\infty}^{\infty} \sum_{j^{\prime}=-\infty}^{\infty} \frac{\exp \left[i k_{1}\left(\alpha_{j}\left(x^{\prime}-x\right)+\beta_{j^{\prime}}\left(y^{\prime}-y\right)+\gamma_{1 j j^{\prime}}\left|h\left(\vec{x}_{\perp}^{\prime}\right)-h\left(\vec{x}_{\perp}\right)\right|\right)\right]}{\gamma_{1 j j^{\prime}}} \tag{3.59}
\end{equation*}
$$

where

$$
\begin{align*}
\gamma_{1 j j^{\prime}} & =\sqrt{1-\alpha_{j}^{2}-\beta_{j^{\prime}}^{2}}, \quad\left(\alpha_{j}^{2}+\beta_{j^{\prime}}^{2}<1\right)  \tag{3.60}\\
& =+i \sqrt{\alpha_{j}^{2}+\beta_{j^{\prime}}^{2}-1}, \quad\left(\alpha_{j}^{2}+\beta_{j^{\prime}}^{2}>1\right) \tag{3.61}
\end{align*}
$$

The three-dimensional periodic Green's function for region 2 is given by

$$
\begin{equation*}
G_{2}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)=\frac{i}{2 k_{1} L_{1} L_{2}} \sum_{j=-\infty}^{\infty} \sum_{j^{\prime}=-\infty}^{\infty} \frac{\exp \left[i k_{1}\left(\alpha_{j}\left(x^{\prime}-x\right)+\beta_{j^{\prime}}\left(y^{\prime}-y\right)+\gamma_{2 j j^{\prime}}\left|h\left(\vec{x}_{\perp}^{\prime}\right)-h\left(\vec{x}_{\perp}\right)\right|\right)\right]}{\gamma_{2 j j^{\prime}}}, \tag{3.62}
\end{equation*}
$$

where

$$
\begin{align*}
\gamma_{2 j j^{\prime}} & =\sqrt{K^{2}-\alpha_{j}^{2}-\beta_{j^{\prime}}^{2}}, \quad\left(\alpha_{j}^{2}+\beta_{j^{\prime}}^{2}<K^{2}\right)  \tag{3.63}\\
& =+i \sqrt{\alpha_{j}^{2}+\beta_{j^{\prime}}^{2}-K^{2}}, \quad\left(\alpha_{j}^{2}+\beta_{j^{\prime}}^{2}>K^{2}\right) \tag{3.64}
\end{align*}
$$

We use these Green's functions later for three-dimensional electromagnetic problems. Note again that the scaling is $k_{1}$ in front of both (3.59) and (3.62). Note also the obvious remark that for a two-dimensional surface $h\left(\vec{x}_{\perp}\right)$ we have two spectral parameters, $j$ and $j^{\prime}$.

Both Green's functions satisfy a two-dimensional Floquet condition

$$
\begin{equation*}
G^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}+\vec{x}_{n_{1} n_{2}}\right)=\exp \left[-i k_{1}\left(\alpha_{0} n_{1} L_{1}+\beta_{0} n_{2} L_{2}\right)\right] G^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right) . \tag{3.65}
\end{equation*}
$$

Combined with the two-dimensional Floquet condition on any field function (3.10), the product of any Green's function times a field function is periodic

$$
\begin{equation*}
G^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}+\vec{x}_{n_{1} n_{2}}\right) \psi\left(\vec{x}_{h}+\vec{x}_{n_{1} n_{2}}\right)=G^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right) \psi\left(\vec{x}_{h}\right) . \tag{3.66}
\end{equation*}
$$

We use this property later to cancel side integrals in Green's theorem. Techniques for computing these periodic Green's functions are available ${ }^{101102}$.

### 3.4 Integral Methods in Coordinate Space for Scalar Problems

We first present the coordinate-space representation of the scattering from a periodic surface as comparison and contrast to that of the spectral representations in later sections. In addition, these yield rigorous representations for the scattered field above the highest surface excursion and for the transmitted field below the lowest surface excursion, as well as projections on lines above and below the surface.

The total field in region $1, \psi_{1}$, equals the sum of incident plus scattered fields, $\psi_{1}=$ $\psi^{i n}+\psi^{s c}$, and it satisfies the scalar Helmholtz equation

$$
\begin{equation*}
\left(\nabla_{2}^{2}+k_{1}^{2}\right) \psi_{1}(\vec{x})=0, \tag{3.67}
\end{equation*}
$$

as do both incident and scattered fields. We do Green's theorem using $\psi_{1}$ and $G_{1}^{(2 p)}$. Cross multiply (3.67) and (3.35), multiply by the characteristic function of region 1

$$
\begin{equation*}
\Theta_{1}(\vec{x})=\theta(L / 2-x) \theta(x+L / 2) \theta(z-h(x)) \theta\left(H_{1}-z\right), \tag{3.68}
\end{equation*}
$$

where $\theta$ is the step function, $\theta(x)=1$ when $x>0$, and $\theta(x)=0$ when $x<0$, and integrate by parts. To express the results conveniently, introduce the bracket notation

$$
\begin{equation*}
\left[G_{1}^{(2 p)}, \psi_{1} ; \vec{x}^{\prime}, S\right]=\iint_{S}\left[G_{1}^{(2 p)}\left(\vec{x}^{\prime}, \vec{x}_{S}\right) \partial_{l} \psi_{1}\left(\vec{x}_{S}\right)-\partial_{l} G_{1}^{(2 p)}\left(\vec{x}^{\prime}, \vec{x}_{S}\right)\right] n_{l} d s \tag{3.69}
\end{equation*}
$$

where $\partial_{l}$ is the partial derivative ( $\partial_{x}$ for $l=1$ and $\partial_{z}$ for $l=2$ ), $n_{l}$ is the non-unit surface normal, and $d s$ the arc length along the surface. Repeated subscripts are summed. There are four surfaces $S$ : $x= \pm L / 2\left(h<z<H_{1}\right), z=h$, and $z=H_{1}$, both with $-L / 2<x<L / 2$. The result is
$\psi_{1}\left(\vec{x}^{\prime}\right) \Theta_{1}\left(\vec{x}^{\prime}\right)=\left[G_{1}^{(2 p)}, \psi_{1} ; \vec{x}^{\prime}, L / 2\right]-\left[G_{1}^{(2 p)}, \psi_{1} ; \vec{x}^{\prime},-L / 2\right]+\left[G_{1}^{(2 p)}, \psi_{1} ; \vec{x}^{\prime}, H_{1}\right]-\left[G_{1}^{(2 p)}, \psi_{1} ; \vec{x}^{\prime}, h\right]$.
Using (3.38), the first two brackets on the right hand side of (3.70) cancel by Floquet periodicity. For the moment assume the integral on $H_{1}$ represents the incident field (proof below), i.e.

$$
\begin{equation*}
\psi^{i n}\left(\vec{x}^{\prime}\right)=\left[G_{1}^{(2 p)}, \psi_{1} ; \vec{x}^{\prime}, H_{1}\right] . \tag{3.71}
\end{equation*}
$$

We thus have three results. Inside region $1, h\left(x^{\prime}\right)<z^{\prime}<H_{1}, \Theta_{1}=1$, we have

$$
\begin{equation*}
\psi_{1}\left(\vec{x}^{\prime}\right)=\psi^{i n}\left(\vec{x}^{\prime}\right)-\left[G_{1}^{(2 p)}, \psi_{1} ; \vec{x}^{\prime}, h\right] . \tag{3.72}
\end{equation*}
$$

Outside region 1, where $\Theta_{1}=0$, we have from (3.70)

$$
\begin{equation*}
\psi^{i n}\left(\vec{x}^{\prime}\right)=\left[G_{1}^{(2 p)}, \psi_{1} ; \vec{x}^{\prime}, h\right], \tag{3.73}
\end{equation*}
$$

which is an Extinction Theorem, and on the surface $z^{\prime}=h\left(x^{\prime}\right)$, taking into account the discontinuity of the double layer potential in (3.72), we have

$$
\begin{equation*}
\frac{1}{2} \psi_{1}\left(\vec{x}_{h}^{\prime}\right)=\psi^{i n}\left(\vec{x}_{h}^{\prime}\right)-\left[G_{1}^{(2 p)}, \psi_{1} ; \vec{x}_{h}^{\prime}, h\right] . \tag{3.74}
\end{equation*}
$$

In particular, we can write the scattered field above the highest surface excursion, $z^{\prime}>$ $\max (h)$, using (3.72). The absolute value in the Green's function in 3.72 is not present, i.e. we use the representation

$$
\begin{equation*}
G_{1}^{(2 p)}\left(\vec{x}^{\prime}, \vec{x}_{h}\right)=\frac{i}{2 k_{1} L} \sum_{j=-\infty}^{\infty} \frac{\exp \left[i k_{1}\left(\alpha_{j}\left(x^{\prime}-x\right)+\gamma_{1 j}\left(z^{\prime}-h(x)\right)\right)\right]}{\gamma_{1 j}} \tag{3.75}
\end{equation*}
$$

The result is that the scattered field above the highest surface excursion can be written exactly as a plane wave expansion of upgoing waves

$$
\begin{equation*}
\psi^{s c}\left(\vec{x}^{\prime}\right)=\sum_{j=-\infty}^{\infty} A_{j} \exp \left[i k_{1}\left(\alpha_{j} x^{\prime}+\gamma_{1 j} z^{\prime}\right)\right] \tag{3.76}
\end{equation*}
$$

where $A_{j}$ can be written as the integral

$$
\begin{equation*}
A_{j}=\frac{1}{L} \int_{-L / 2}^{L / 2} A(j, x) \exp \left[-i k_{1}\left(\alpha_{j} x+\gamma_{1 j} h(x)\right)\right] d x \tag{3.77}
\end{equation*}
$$

and the integrand $A(j, x)$ is in terms of the boundary unknowns

$$
\begin{equation*}
A(j, x)=\frac{-i}{2 k_{1} \gamma_{1 j}}\left\{\frac{\partial \psi_{1}}{\partial n}\left(\vec{x}_{h}\right)+i k_{1}\left(\gamma_{1 j}-\alpha_{j} h^{\prime}(x)\right) \psi_{1}\left(\vec{x}_{h}\right)\right\} . \tag{3.78}
\end{equation*}
$$

We have written $A(j, x)$ as a function of two variables, the first a discrete spectral (S) variable $j$ which has replaced the field coordinate variable, and the second a continuous coordinate (C) variable $x$, which is the surface integration variable. This is the basis for the spectral-coordinate (SC) approach used in Sec. 5.

There remains to prove (3.71). This time the representation for the Green's function evaluated on $z=H_{1}$ is with $\vec{x}_{1}=\left(x, H_{1}\right)$

$$
\begin{equation*}
G_{1}^{(2 p)}\left(\vec{x}^{\prime}, \vec{x}_{1}\right)=\frac{i}{2 k_{1} L} \sum_{j=-\infty}^{\infty} \frac{\exp \left[i k_{1}\left(\alpha_{j}\left(x^{\prime}-x\right)+\gamma_{1 j}\left(H_{1}-z^{\prime}\right)\right)\right]}{\gamma_{1 j}} . \tag{3.79}
\end{equation*}
$$

Using (3.79) in (3.71) and the representation (3.76) for the scattered field it is straighforward to show that

$$
\begin{equation*}
\left[G_{1}^{(2 p)}, \psi^{s c} ; \vec{x}^{\prime}, H_{1}\right]=0 \tag{3.80}
\end{equation*}
$$

Further, if we assume a general plane wave decomposition of the incident field in terms of downgoing waves

$$
\begin{equation*}
\psi^{i n}(\vec{x})=\sum_{n} I_{n} \exp \left[i k_{1}\left(\alpha_{n} x-\gamma_{1 n} z\right)\right] \tag{3.81}
\end{equation*}
$$

the relation

$$
\begin{equation*}
\left[G_{1}^{(2 p)}, \psi^{i n} ; \vec{x}^{\prime}, H_{1}\right]=\psi^{i n}\left(\vec{x}^{\prime}\right) \tag{3.82}
\end{equation*}
$$

follows immediately. Alternatively, one can view the integrand in (3.71)

$$
\begin{equation*}
\frac{\partial \psi_{1}}{\partial z}\left(x, H_{1}\right)-i k_{1} \gamma_{1 j} \psi_{1}\left(x, H_{1}\right) \tag{3.83}
\end{equation*}
$$

as projecting out only the downgoing waves and canceling the scattered waves. The combination of (3.80) and (3.82) is the proof of (3.71).

In the region below the surface, region 2, the total field $\psi_{2}$ satisfies the Helmholtz equation

$$
\begin{equation*}
\left(\nabla_{2}^{2}+k_{2}^{2}\right) \psi_{2}(\vec{x})=0, \tag{3.84}
\end{equation*}
$$

where the wavenumber $k_{2}=K k_{1}$ is written in terms of a scale factor $K$. The region is defined by the characteristic function

$$
\begin{equation*}
\Theta_{2}(\vec{x})=\theta(L / 2-x) \theta(x+L / 2) \theta(h(x)-z) \theta\left(z-H_{2}\right) . \tag{3.85}
\end{equation*}
$$

The Green's function $G_{2}^{(2 p)}$ is given by

$$
\begin{equation*}
G_{2}^{(2 p)}\left(\vec{x}^{\prime}, \vec{x}\right)=\frac{i}{2 k_{1} L} \sum_{j=-\infty}^{\infty} \frac{\exp \left[i k_{1}\left(\alpha_{j}\left(x^{\prime}-x\right)+\gamma_{2 j}\left|z^{\prime}-z\right|\right)\right]}{\gamma_{2 j}} \tag{3.86}
\end{equation*}
$$

where $\gamma_{2 j}$ is defined in (3.33). Green's theorem in this region, the cancellation of integrals along $x= \pm L / 2$ by Floquet periodicity, and the vanishing of the integral along $H_{2}$ since at this value of $z$ the total field consists of downward propagating waves, yields the result

$$
\begin{equation*}
\psi_{2}\left(\vec{x}^{\prime}\right) \Theta_{2}\left(\vec{x}^{\prime}\right)=\left[G_{2}^{(2 p)}, \psi_{2} ; \vec{x}^{\prime}, h\right] . \tag{3.87}
\end{equation*}
$$

On the boundary, the limit of 3.87) is

$$
\begin{equation*}
\frac{1}{2} \psi_{2}\left(\vec{x}_{h}^{\prime}\right)=\left[G_{2}^{(2 p)}, \psi_{2} ; \vec{x}_{h}^{\prime}, h\right] . \tag{3.88}
\end{equation*}
$$

For $\left.z^{\prime}<\min (h), 3.87\right)$ yields a representation of the total field in region 2 in terms of downward propagating plane waves

$$
\begin{equation*}
\psi_{2}\left(\vec{x}^{\prime}\right)=\sum_{j=-\infty}^{\infty} B_{j} \exp \left[i k_{1}\left(\alpha_{j} x^{\prime}-\gamma_{2 j} z^{\prime}\right)\right] \tag{3.89}
\end{equation*}
$$

where

$$
\begin{equation*}
B_{j}=\frac{1}{L} \int_{-L / 2}^{L / 2} B(j, x) \exp \left[-i k_{1}\left(\alpha_{j} x-\gamma_{2 j} h(x)\right)\right] d x, \tag{3.90}
\end{equation*}
$$

and $B(j, x)$ is in terms of the boundary values from region 2

$$
\begin{equation*}
B(j, x)=\frac{i}{2 k_{1} \gamma_{2 j}}\left\{\frac{\partial \psi_{2}\left(\vec{x}_{h}\right)}{\partial n}-i k_{1}\left(\gamma_{2 j}+\alpha_{j} h^{\prime}(x)\right) \psi_{2}\left(\vec{x}_{h}\right)\right\} \tag{3.91}
\end{equation*}
$$

We have assumed that $\psi$ is a velocity potential. Then the continuity conditions at the boundary are written as the continuity of velocity

$$
\begin{equation*}
\frac{\partial \psi_{2}}{\partial n}\left(\vec{x}_{h}\right)=\frac{\partial \psi_{1}}{\partial n}\left(\vec{x}_{h}\right) \doteq N\left(\vec{x}_{h}\right), \tag{3.92}
\end{equation*}
$$

and continuity of pressure

$$
\begin{equation*}
\rho_{2} \psi_{2}\left(\vec{x}_{h}\right)=\rho_{1} \psi_{1}\left(\vec{x}_{h}\right), \tag{3.93}
\end{equation*}
$$

where $\rho_{j}$ are the densities. In (3.92) we defined the normal derivative boundary unknown as $N$, and we define the surface field boundary unknown as $\psi\left(\vec{x}_{h}\right)=\psi_{1}\left(\vec{x}_{h}\right)$, so that $\psi_{2}\left(\vec{x}_{h}\right)=\frac{1}{\rho} \psi\left(\vec{x}_{h}\right)$ where $\rho=\rho_{2} / \rho_{1}$. Using these unknowns, (3.78) and (3.91) become

$$
\begin{equation*}
A(j, x)=\frac{-i}{2 k_{1} \gamma_{1 j}}\left[N\left(\vec{x}_{h}\right)+i k_{1}\left(\gamma_{1 j}-\alpha_{j} h^{\prime}(x)\right) \psi\left(\vec{x}_{h}\right)\right], \tag{3.94}
\end{equation*}
$$

and

$$
\begin{equation*}
B(j, x)=\frac{i}{2 k_{1} \gamma_{2 j}}\left[N\left(\vec{x}_{h}\right)-\frac{i k_{1}}{\rho}\left(\gamma_{2 j}+\alpha_{j} h^{\prime}(x)\right) \psi\left(\vec{x}_{h}\right)\right] . \tag{3.95}
\end{equation*}
$$

We can summarize these results using single $(S)$ and double $(D)$ layer potentials

$$
\begin{equation*}
\left(S_{j} u\right)\left(\vec{x}^{\prime}\right)=\int_{-L / 2}^{L / 2} G_{j}^{(2 p)}\left(\vec{x}^{\prime}, \vec{x}_{h}\right) u\left(\vec{x}_{h}\right) d x \tag{3.96}
\end{equation*}
$$

and

$$
\begin{equation*}
\left(D_{j} v\right)\left(\vec{x}^{\prime}\right)=\int_{-L / 2}^{L / 2} \frac{\partial G_{j}^{(2 p)}}{\partial n}\left(\vec{x}^{\prime}, \vec{x}_{h}\right) v\left(\vec{x}_{h}\right) d x \tag{3.97}
\end{equation*}
$$

and write the integral equations 3.74 and 3.88 in symbolic form as

$$
\begin{equation*}
\frac{1}{2} \psi=\psi^{i n}-\left(S_{1} N\right)+\left(D_{1} \psi\right) \tag{3.98}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{1}{2} \psi=\rho\left(S_{2} N\right)-\left(D_{2} \psi\right) \tag{3.99}
\end{equation*}
$$

Various combinations of these equations and integral equations formed by first taking the normal derivative of the field representations (3.72) and (3.87) and passing to the surface limit can be used to solve for the boundary unknowns $\psi$ and $N$. For the Dirichlet problem, $\psi=0$ and $\rho=0$ so (3.99) disappears and (3.98) is an integral equation of first kind for $N$. For the Neumann problem, first divide 3.99 by $\rho$, then let $\rho \rightarrow \infty$ and set $N=0$.

Direct integral equation methods have been used to computationally solve this problem ${ }^{21 / 22 \mid 677}$. Other integral equation solutions ${ }^{\sqrt{36137138139}}$ have been compared to the solutions of spectral methods presented later in this paper. Other methods have also been employed ${ }^{1811920 \mid 73174186}$. Point collocation questions arise ${ }^{\sqrt{1025553 / 60 \mid 72}}$ for any coordinate based method.

### 3.5 Partial Spectral Methods for Scalar Problems

In this section we use a direct method to generate integral equations in a partial spectral representation. The method uses Green's theorem again, but not the Green's function. Define the upand down-going plane wave states in region 1

$$
\begin{equation*}
\phi_{1 j}^{ \pm}(\vec{x})=\exp \left[i k_{1}\left(-\alpha_{j} x \pm \gamma_{1 j} z\right)\right], \tag{3.100}
\end{equation*}
$$

which satisfy the same Helmholtz equation as $\psi_{1}$, 3.67,

$$
\begin{equation*}
\left(\nabla_{2}^{2}+k_{1}^{2}\right) \phi_{1 j}^{ \pm}(\vec{x})=0 \tag{3.101}
\end{equation*}
$$

For convenience, in (3.100) we have chosen the conjugate in the $x$-coordinate. In the Green's function this occurs naturally. Cross multiply (3.67) and (3.101) and subtract the results, multiply by $\Theta_{1}$ from (3.68), integrate over all space, and then integrate by parts. Since all fields are Floquet periodic, the integrals along $x= \pm L / 2$ cancel. The results can be expressed with the collapsed bracket notation

$$
\begin{equation*}
[u, v ; S]=\int_{S}\left[u\left(\vec{x}_{S}\right) \partial_{l} v\left(\vec{x}_{S}\right)-v\left(\vec{x}_{S}\right) \partial_{l} u\left(\vec{x}_{S}\right)\right] n_{l} d s \tag{3.102}
\end{equation*}
$$

where, unlike the bracket notation in Sec.4, no exterior coordinate-space variable appears. There are two surfaces, $z=h$ with $-L / 2<x<L / 2$, and $z=H_{1}$ with $-L / 2<x<L / 2$. The result is

$$
\begin{equation*}
\left[\phi_{1 j}^{ \pm}, \psi_{1} ; h\right]=\left[\phi_{1 j}^{ \pm}, \psi_{1} ; H_{1}\right] . \tag{3.103}
\end{equation*}
$$

The result can be thought of as an analytic continuation from the periodic surface $z=h$ to a flat plane $z=H_{1}$ above the surface. The right hand side of (3.103) can be evaluated explicitly using (3.76) and (3.81) to give

$$
\left[\phi_{1 j}^{ \pm}, \psi_{1} ; H_{1}\right]=2 i k_{1} L \gamma_{1 j}\left\{\begin{array}{c}
-I_{A_{j}} \tag{3.104}
\end{array}\right\}
$$

Here the up-going plane waves $\phi_{1 j}^{+}$project out the down-going spectral components of the incident wave $I_{j}$, and the down-going plane waves $\phi_{1 j}^{-}$project out the up-going spectral components of the scattered waves $A_{j}$. Combining this with the left hand side of (3.103) we get the set of equations for region 1

$$
\frac{1}{L} \int_{-L / 2}^{L / 2} \phi_{1 j}^{ \pm}\left(\vec{x}_{h}\right) U_{j}^{ \pm}\left(\vec{x}_{h}\right) d x=\gamma_{1 j}\left\{\begin{array}{l}
-I_{j}  \tag{3.105}\\
A_{j}
\end{array}\right\},
$$

where

$$
\begin{equation*}
U_{j}^{ \pm}\left(\vec{x}_{h}\right)=\frac{1}{2 i k_{1}}\left[N\left(\vec{x}_{h}\right)-i k_{1}\left( \pm \gamma_{1 j}+\alpha_{j} h^{\prime}(x)\right) \psi\left(\vec{x}_{h}\right)\right] . \tag{3.106}
\end{equation*}
$$

We have incorporated the boundary unknowns defined following 3.92. Note that $U_{j}^{-}=$ $\gamma_{1 j} A(j, x)$ from (3.94).

In region 2, the up- and down-going plane waves are given by

$$
\begin{equation*}
\phi_{2 j}^{ \pm}(\vec{x})=\exp \left[i k_{1}\left(-\alpha_{j} x \pm \gamma_{2 j} z\right)\right] \tag{3.107}
\end{equation*}
$$

which satisfy the Helmholtz equation

$$
\begin{equation*}
\left(\nabla_{2}^{2}+k_{2}^{2}\right) \phi_{2 j}^{ \pm}(\vec{x})=0 \tag{3.108}
\end{equation*}
$$

Cross multiply (3.84) and (3.108), multiply the result by $\Theta_{2}$ from (3.85), integrate over all space, then integrate by parts. The Floquet periodicity cancels the integrals on $x= \pm L / 2$ and the result is the analytic continuation

$$
\begin{equation*}
\left[\phi_{2 j}^{ \pm}, \psi_{2} ; h\right]=\left[\phi_{2 j}^{ \pm}, \psi_{2} ; H_{2}\right] . \tag{3.109}
\end{equation*}
$$

The right hand side of (3.109) can be evaluated using (3.89) for $\psi_{2}$ to yield

$$
\left[\phi_{2 j}^{ \pm}, \psi_{2} ; H_{2}\right]=-2 i k_{1} L \gamma_{2 j}\left\{\begin{array}{c}
B_{j}  \tag{3.110}\\
0
\end{array}\right\} .
$$

Combined with 3.109, and incorporating the definitions of the boundary values following (3.92) yields the equations from the lower region

$$
\frac{1}{L} \int_{-L / 2}^{L / 2} \phi_{2 j}^{ \pm}\left(\vec{x}_{h}\right) L_{j}^{ \pm}\left(\vec{x}_{h}\right) d x=-\gamma_{2 j}\left\{\begin{array}{l}
B_{j}  \tag{3.111}\\
0
\end{array}\right\}
$$

where

$$
\begin{equation*}
L_{j}^{ \pm}\left(\vec{x}_{h}\right)=\frac{1}{2 i k_{1}}\left[N\left(\vec{x}_{h}\right)-\frac{i k_{1}}{\rho}\left( \pm \gamma_{2 j}+\alpha_{j} h^{\prime}(x)\right) \psi\left(\vec{x}_{h}\right)\right] . \tag{3.112}
\end{equation*}
$$

The lower equation in $(3.111)$ is a spectral version of the Extinction Theorem.
The procedure is to solve the combined $U^{+}$equation in (3.105) and the $L^{-}$equation in (3.111) for the boundary unknowns $N$ and $\psi$, and to evaluate the $U^{-}$and $L^{+}$equations for the scattered $\left(A_{j}\right)$ and transmitted $\left(B_{j}\right)$ amplitudes. The scattered and transmitted fields can be then found from (3.76) and (3.89) respectively. In order to find field values in the surface wells, we must use these boundary unknowns in (3.72) and (3.87) respectively.

The advantage of the method is that there are no Green's functions to compute. Instead, the results are projected onto plane wave based basis functions (really Rayleigh functions since they're on the surface). The Green's function does this in an alternate way.

It is useful with any theory to check simple special cases. It is also necessary that the general results reduce to simple solvable cases. Here we take the flat surface limit $(h=0)$, and derive from them the Fresnel reflection and transmission coefficients as a necessary check on the general results. For $h=0$, let $L \rightarrow \infty$, so that for any finite $j, \lim _{L \rightarrow \infty} \alpha_{j}=\alpha_{0}$, so that the only surviving waves are the $0^{t h}$ order reflection $\left(A_{0}\right)$ and transmission $\left(B_{0}\right)$ amplitudes. The surface fields $N$ and $\psi$ thus have two different flat-surface field representations which are

$$
\begin{equation*}
\psi_{1}(x, 0)=\left(I_{0}+A_{0}\right) \exp \left[i k_{1} \alpha_{0} x\right] \tag{3.113}
\end{equation*}
$$

$$
\begin{gather*}
N_{1}(x, 0)=-i k_{1} \gamma_{10}\left(I_{0}-A_{0}\right) \exp \left[i k_{1} \alpha_{0} x\right]  \tag{3.114}\\
\psi_{2}(x, 0)=B_{0} \exp \left[i k_{1} \alpha_{0} x\right] \tag{3.115}
\end{gather*}
$$

and

$$
\begin{equation*}
N_{2}(x, 0)=-i k_{1} \gamma_{20} B_{0} \exp \left[i k_{1} \alpha_{0} x\right] . \tag{3.116}
\end{equation*}
$$

From (3.105) and 3.111) we have

$$
\begin{equation*}
A_{0}=\frac{1}{2 i k_{1} \gamma_{10}} \lim _{L \rightarrow \infty} \frac{1}{L} \int_{-L / 2}^{L / 2}\left[N(x, 0)+i k_{1} \gamma_{10} \psi(x, 0)\right] \exp \left[-i k_{1} \alpha_{0} x\right] d x, \tag{3.117}
\end{equation*}
$$

and

$$
\begin{equation*}
B_{0}=\frac{-1}{2 i k_{1} \gamma_{20}} \lim _{L \rightarrow \infty} \frac{1}{L} \int_{-L / 2}^{L / 2}\left[N(x, 0)-\frac{i k_{1}}{\rho} \gamma_{20} \psi(x, 0)\right] \exp \left[-i k_{1} \alpha_{0} x\right] d x \tag{3.118}
\end{equation*}
$$

If we use the flat-surface field representations on the surface from region $1,\left(N_{1}\right.$ and $\left.\psi_{1}\right)$, in the $A_{0}$ equation, and the flat-surface field representations from region 2 , $\left(N_{2}\right.$ and $\left.\psi_{2}\right)$, in the $B_{0}$ equation, we just get identities. Instead, use the opposite procedure, i.e. write $A_{0}$ and $B_{0}$ as

$$
\begin{equation*}
A_{0}=\frac{1}{2 i k_{1} \gamma_{10}} \lim _{L \rightarrow \infty} \frac{1}{L} \int_{-L / 2}^{L / 2}\left[N_{2}(x, 0)+i k_{1} \gamma_{10} \psi_{2}(x, 0)\right] \exp \left[-i k_{1} \alpha_{0} x\right] d x \tag{3.119}
\end{equation*}
$$

and

$$
\begin{equation*}
B_{0}=\frac{-1}{2 i k_{1} \gamma_{20}} \lim _{L \rightarrow \infty} \frac{1}{L} \int_{-L / 2}^{L / 2}\left[N_{1}(x, 0)-\frac{i k_{1}}{\rho} \gamma_{20} \psi_{1}(x, 0)\right] \exp \left[-i k_{1} \alpha_{0} x\right] d x \tag{3.120}
\end{equation*}
$$

Using (3.113) through (3.116) in (3.119) and (3.120) we get two equations

$$
\begin{equation*}
A_{0}=\frac{\rho \gamma_{10}-\gamma_{20}}{2 \gamma_{10}} B_{0} \tag{3.121}
\end{equation*}
$$

and

$$
\begin{equation*}
B_{0}=\frac{\gamma_{10}\left[I_{0}-A_{0}\right]+\left(\gamma_{20} / \rho\right)\left[I_{0}+A_{0}\right]}{2 \gamma_{10}} \tag{3.122}
\end{equation*}
$$

These can be solved to yield the Fresnel reflection coefficient

$$
\begin{equation*}
\frac{A_{0}}{I_{0}}=\frac{\rho \gamma_{10}-\gamma_{20}}{\rho \gamma_{10}+\gamma_{20}} \tag{3.123}
\end{equation*}
$$

and the Fresnel transmission coefficient

$$
\begin{equation*}
\frac{B_{0}}{I_{0}}=\frac{2 \gamma_{10}}{\rho \gamma_{10}+\gamma_{20}} \tag{3.124}
\end{equation*}
$$

Finally we have that

$$
\begin{equation*}
1+\frac{A_{0}}{I_{0}}=\rho \frac{B_{0}}{I_{0}} \tag{3.125}
\end{equation*}
$$

as expected.

### 3.6 Surface Inversion Using the Partial Spectral Method

We can use the partial spectral results from Sec. 5 to develop simple algorithms to reconstruct the surface height $h(x)$ from the knowledge of the incident and scattered field amplitudes $I_{j}$ and $A_{j}$. For simplicity, we choose the Dirichlet problem, $\psi\left(\vec{x}_{h}\right)=0$. This is a perfectly reflecting case, and (3.111) vanishes identically (multiply it by $\rho$, and then set $\rho=0$ ). The resulting equations 3 3.105 become

$$
\frac{1}{L} \int_{-L / 2}^{L / 2} \phi_{1 j}^{ \pm}\left(\vec{x}_{h}\right) N\left(\vec{x}_{h}\right) d x=2 i k_{1} \gamma_{1 j}\left\{\begin{array}{l}
-I_{j}  \tag{3.126}\\
A_{j}
\end{array}\right\} .
$$

We describe two methods, the first is perturbation theory in the surface height, and the second is the use of the Kirchhoff approximation for the normal derivative $N$. The full details of both methods with numerical results were presented in ${ }^{[111112]}$. There the methods were applied to truncated rough surfaces. Some other methods can be found in ${ }^{59}$ for uniqueness questions and ${ }^{[17}$ for more detailed reconstruction algorithms.

For perturbation theory $\sqrt{3.100}$ ) is used on the surface, and becomes

$$
\begin{equation*}
\phi_{1 j}^{ \pm}\left(\vec{x}_{h}\right) \approx \exp \left[-i k_{1} \alpha_{j} x\right]\left(1 \pm i k_{1} \gamma_{1 j} h(x)\right) . \tag{3.127}
\end{equation*}
$$

Substituting (3.127) in (3.126), and adding and subtracting the resulting equations yields the two results

$$
\begin{equation*}
\frac{1}{L} \int_{-L / 2}^{L / 2} \exp \left[-i k_{1} \alpha_{j} x\right] N\left(\vec{x}_{h}\right) d x=-i k_{1} \gamma_{1 j}\left(I_{j}-A_{j}\right) \tag{3.128}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{1}{L} \int_{-L / 2}^{L / 2} \exp \left[-i k_{1} \alpha_{j} x\right] N\left(\vec{x}_{h}\right) h(x) d x=-\left(I_{j}+A_{j}\right) \tag{3.129}
\end{equation*}
$$

Fourier inverting both equations yields

$$
\begin{equation*}
N\left(\vec{x}_{h}\right)=-i k_{1} \sum_{j=-\infty}^{\infty} \exp \left[i k_{1} \alpha_{j} x\right]\left(I_{j}-A_{j}\right) \tag{3.130}
\end{equation*}
$$

and

$$
\begin{equation*}
N\left(\vec{x}_{h}\right) h(x)=-\sum_{j=-\infty}^{\infty} \exp \left[i k_{1} \alpha_{j} x\right]\left(I_{j}+A_{j}\right) \tag{3.131}
\end{equation*}
$$

Divide (3.131) by (3.130) (so that we factor out the boundary condition) and take the real part to get the approximation to the surface profile $h_{P T}$ produced by perturbation theory

$$
\begin{equation*}
h_{P T}(x)=\frac{1}{k_{1}} \operatorname{Im}\left\{\frac{\sum_{j=-\infty}^{\infty}\left(I_{j}+A_{j}\right] \exp \left[i k_{1} \alpha_{j} x\right]}{\sum_{j=-\infty}^{\infty}\left(I_{j}-A_{j}\right) \exp \left[i k_{1} \alpha_{j} x\right]}\right\} . \tag{3.132}
\end{equation*}
$$

where ( $\operatorname{Im}$ ) is the imaginary part. The equation simplifies for a single incident wave ( $I_{j}=\delta_{j 0} I_{0}$ ) to be

$$
\begin{equation*}
h_{P T}(x)=\frac{1}{k_{1}} \operatorname{Im}\left\{\frac{I_{0}+\sum_{j=-\infty}^{\infty} A_{j} \exp [i 2 \pi j x / L]}{I_{0}-\sum_{j=-\infty}^{\infty} A_{j} \exp [i 2 \pi j x / L]}\right\} . \tag{3.133}
\end{equation*}
$$

These equations (3.132) and (3.133) express the surface in terms of the amplitudes of the incident and scattered fields.

For the Kirchhoff approximation (KA), assume a single plane wave incidence

$$
\begin{equation*}
\psi^{i n}(\vec{x})=I_{0} \exp \left[i k_{1}\left(\alpha_{0} x-\gamma_{10} z\right)\right] \tag{3.134}
\end{equation*}
$$

and approximate the normal derivative on the surface in 3.126 by twice the normal derivative of the incident field

$$
\begin{equation*}
N\left(\vec{x}_{h}\right) \approx N^{K A}\left(\vec{x}_{h}\right)=2 n_{l} \partial_{l} \psi^{i n}\left(\vec{x}_{h}\right) \tag{3.135}
\end{equation*}
$$

For the lower equation in 3.126 this yields

$$
\begin{equation*}
\frac{1}{L} \int_{-L / 2}^{L / 2}\left[\gamma_{10}+\alpha_{0} h^{\prime}(x)\right] \exp \left[-i k_{1}\left(p_{j} x+q_{j} h(x)\right)\right] d x=-\gamma_{1 j} A_{j} / I_{0} \tag{3.136}
\end{equation*}
$$

where

$$
\begin{equation*}
p_{j}=\alpha_{j}-\alpha_{0} \tag{3.137}
\end{equation*}
$$

and

$$
\begin{equation*}
q_{j}=\gamma_{1 j}+\gamma_{0 j} \tag{3.138}
\end{equation*}
$$

The $h^{\prime}(x)$ term in 3.136) can be integrated by parts to yield

$$
\begin{equation*}
\frac{1}{L} \int_{-L / 2}^{L / 2} \exp \left[-i k_{1}\left(p_{j} x+q_{j} h(x)\right)\right] d x=-f_{j}^{-} A_{j} / I_{0} \tag{3.139}
\end{equation*}
$$

where

$$
\begin{equation*}
f_{j}^{-}=\frac{\gamma_{1 j}\left(\gamma_{1 j}+\gamma_{10}\right)}{\gamma_{1 j} \gamma_{10}+\left(1-\alpha_{j} \alpha_{0}\right)} \tag{3.140}
\end{equation*}
$$

We can re-express $p_{j}$ and $q_{j}$ using trig identities as

$$
\begin{equation*}
p_{j}=\sin \left(\theta_{j}^{s c}\right)-\sin \left(\theta^{i n}\right)=2 \cos \left\{\frac{\theta_{j}^{s c}+\theta^{i n}}{2}\right\} \sin \left\{\frac{\theta_{j}^{s c}-\theta^{i n}}{2}\right\} \tag{3.141}
\end{equation*}
$$

and

$$
\begin{equation*}
q_{j}=\cos \left(\theta_{j}^{s c}\right)+\cos \left(\theta^{i n}\right)=2 \cos \left\{\frac{\theta_{j}^{s c}+\theta^{i n}}{2}\right\} \cos \left\{\frac{\theta_{j}^{s c}-\theta^{i n}}{2}\right\} \tag{3.142}
\end{equation*}
$$

We thus have that $p_{j}$ and $q_{j}$ are confined to an Ewald circle

$$
\begin{equation*}
p_{j}^{2}+q_{j}^{2} \leqslant 4 \tag{3.143}
\end{equation*}
$$

and we further have $\left|p_{j}\right| \leqslant 2$ and $\left|q_{j}\right| \leqslant 2$. This restricts the acceptable $j$ values to a set $J$ and correspondingly restricts the acceptable scattering angles (modulo the incident angle), and thus the scattered amplitudes and fields used for the inversion. For fixed $q_{j}$, say $q_{j_{1}}, 3.139$ is a periodic Fourier transform restricted in $p_{j}$ and thus restricted in the set $J$. As $q_{j_{1}}$ increases, $p_{j}$ decreases, which is equivalent to a low-pass filter. As $q_{j_{1}}$ decreases, more data near grazing illumination and scattering is involved, where the Kirchhoff approximation gets worse. For fixed $q_{j_{1}}$, assume the integral in $\sqrt{3.139}$ can be approximately inverted to yield

$$
\begin{equation*}
\exp \left[-i k_{1} q_{j_{1}} h(x)\right]=\frac{-1}{I_{0}} \sum_{J} f_{j}^{-} A_{j} \exp \left[i k_{1} p_{j} x\right] \doteq \mathscr{R}(x) \tag{3.144}
\end{equation*}
$$

Taking the real $R e$ and imaginary Im parts of (3.144) (and neglecting periodic phase shifts) yields $h_{K A}$, the Kirchhoff approximation of the surface height

$$
\begin{equation*}
h_{K A}\left(x, q_{j_{1}}\right)=\frac{1}{k_{1} q_{j_{1}}} \arctan \left\{\frac{-\operatorname{Im}(\mathscr{R}(x))}{\operatorname{Re}(\mathscr{R}(x))}\right\}, \tag{3.145}
\end{equation*}
$$

which again produces the surface height function in terms of the scattered field amplitudes this time modulated by the Kirchhoff components. Each $q_{j_{1}}$ produces a different value of $h_{K A}$. For a non-periodic truncated random surface the method was used successfully to reconstruct ensemble surface height functions with approximately twice the rms height as for perturbation theory ${ }^{[112}$. The cited paper also contains a discussion of the various angle combinations for different reconstructions.

### 3.7 Full Spectral Methods for Scalar Problems: Physical Optics Modified Fourier Basis and Floquet-Fourier Expansions

In Sec.4, both the exterior and interior (integration) variables were in coordinate space. The equations generated were formally exact for the solution of the boundary values $\psi\left(\vec{x}_{h}\right)$ and $N\left(\vec{x}_{h}\right)$. Once the boundary values were found, the scattered and transmitted fields anywhere away from the surface wells could be evaluated via either direct transforms or summation methods in the resulting plane wave cum evanescent wave expansions. The periodic Green's function was used and had to be computed. Acceleration methods to do this are available ${ }^{[011102]}$.

In Sec.5, we used plane/evanescent waves to derive another set of equations, again formally exact, for the boundary unknowns which avoided the use of the periodic Green's functions. The equations to be solved were similar to the equations to be evaluated in the sense that both involved a close interplay between spectral and coordinate parameters in "parallel" as distinct from the "serial" presentation of methods in Sec. 4 .

Solution of the boundary unknowns in Secs. 4 and 5 using direct discretization methods involves matrix inversion where the rows and columns of the matrix are both sampled in coordinate space, and the sampling methods are flexible. In Sec. 5 the columns are sampled in coordinate space, but the rows are sampled in spectral space, and this is proscribed in terms of the Bragg waves. Convergence and the usefulness of the two sets of solutions have been discussed ${ }^{36137138139}$. The major point is that the limits of convergence, stability and errors are numerical and directly related to the solution of exact formal equations and not to any strictly "physical" approximations.

That changes when we attempt to approximate the surface fields in some spectral basis, and thus to write equations fully in spectral space. The first question is what do we mean by spectral space in this context? The second is what do we know about possible expansions? The main thing we know is that the surface fields are the limits of Floquet-periodic functions, so they must also be Floquet-periodic. In particular, they should not be expanded in a pure Fourier series (no matter the temptation) since the latter are only valid for normal incidence ( $\alpha_{0}=0$ ), where the Floquet periodicity reduces to ordinary periodicity. The validity of a pure Fourier expansion deteriorates for non-normal incidence.

In this section we briefly describe the use of a pure Floquet type expansion which defines "spectrum" in one particular way. It is also a physical optics (PO) expansion explained below. From this we are able to infer the results for what we refer to as a Floquet-Fourier (FF) expansion, and these latter results are presented at the end of the section. The PO expansions for the
boundary unknowns are

$$
\begin{equation*}
\psi\left(\vec{x}_{h}\right)=\exp \left[-i k_{1} \gamma_{10} h(x)\right] \sum_{j^{\prime}=-\infty}^{\infty} \psi_{j^{\prime}}^{(P O)} \exp \left[i k_{1} \alpha_{j^{\prime}} x\right] \tag{3.146}
\end{equation*}
$$

or, written in another form

$$
\begin{equation*}
\psi\left(\vec{x}_{h}\right)=\exp \left[i k_{1} \alpha_{0} x-i k_{1} \gamma_{10} h(x)\right] \sum_{j^{\prime}=-\infty}^{\infty} \psi_{j^{\prime}}^{(P O)} \exp \left[i 2 \pi j^{\prime} x / L\right] \tag{3.147}
\end{equation*}
$$

where the term outside the summation can be written using the complex conjugate of (3.100)

$$
\begin{equation*}
\exp \left[i k_{1} \alpha_{0} x-i k_{1} \gamma_{10} h(x)\right]=\bar{\phi}_{10}^{+}\left(\vec{x}_{h}\right) \tag{3.148}
\end{equation*}
$$

The term is the physical optics or Kirchhoff approximation of a down-going plane wave evaluated on the boundary. It serves to modulate the remaining Fourier series, and can be viewed as a precursor to more general Waterman-type expansions ${ }^{\sqrt{105}}$ in terms of down-going waves. (If the $h$ term is not present in (3.146), the expansion is still a Floquet-periodic expansion, and, since it is a generalization of the Fourier expansion, has the advantage of being invertible. Its result can be inferred from the results below, and are presented at the end of this section.) The normal derivative is similarly expanded

$$
\begin{equation*}
N\left(\vec{x}_{h}\right)=i k_{1} \exp \left[-i k_{1} \gamma_{10} h(x)\right] \sum_{j^{\prime}=-\infty}^{\infty} N_{j^{\prime}}^{(P O)} \exp \left[i k_{1} \alpha_{j^{\prime}} x\right] . \tag{3.149}
\end{equation*}
$$

Here we have scaled the normal derivative term by $i k_{1}$ for convenience. The expansion was initially introduced as a physical optics modified Fourier expansion ${ }^{3132] 33}$, and used by several others ${ }^{[15 / 26 / 2755 \mid 106 / 107]}$. The reference ${ }^{33]}$ can be viewed as the exact version of the approximate Rayleigh-Fano equations ${ }^{97}$ valid in perturbation theory for shallow surfaces. The expansions can be substituted into (3.105) and (3.111) to yield

$$
\sum_{j^{\prime}=-\infty}^{\infty} M_{1 j j^{\prime}}^{ \pm}(P O)\left[N_{j^{\prime}}^{(P O)} \mp \gamma_{1 j} \psi_{j^{\prime}}^{(P O)}\right]-\alpha_{j} \sum_{j^{\prime}=-\infty}^{\infty} \tilde{M}_{1 j j^{\prime}}^{ \pm}(P O) \psi_{j^{\prime}}^{(P O)}=2 \gamma_{1 j}\left\{\begin{array}{c}
-I_{j}  \tag{3.150}\\
A_{j}
\end{array}\right\}
$$

and

$$
\sum_{j^{\prime}=-\infty}^{\infty} M_{2 j j^{\prime}}^{ \pm}(P O)\left[N_{j^{\prime}}^{(P O)} \mp \frac{i k_{1}}{\rho} \gamma_{2 j} \psi_{j^{\prime}}^{(P O)}\right]-\frac{i k_{1}}{\rho} \alpha_{j} \sum_{j^{\prime}=-\infty}^{\infty} \tilde{M}_{2 j j^{\prime}}^{ \pm}(P O) \psi_{j^{\prime}}^{(P O)}=-2 \gamma_{2 j}\left\{\begin{array}{c}
B_{0}^{B_{j}}  \tag{3.151}\\
0
\end{array}\right\},
$$

where the physical optics $(P O)$ matrix elements are $(p=1,2)$

$$
\begin{equation*}
M_{p j j^{\prime}}^{ \pm}(P O)=\frac{1}{L} \int_{-L / 2}^{L / 2} \exp \left(i k_{1}\left[\left( \pm \gamma_{p j}-\gamma_{10}\right) h(x)+\left(\alpha_{j^{\prime}}-\alpha_{j}\right) x\right]\right) d x \tag{3.152}
\end{equation*}
$$

(note that $M_{10 j^{\prime}}^{+}(P O)=\delta_{j^{\prime} 0}$ ) and

$$
\begin{equation*}
\tilde{M}_{p j j^{\prime}}^{ \pm}(P O)=\frac{1}{L} \int_{-L / 2}^{L / 2} h^{\prime}(x) \exp \left(i k_{1}\left[\left( \pm \gamma_{p j}-\gamma_{10}\right) h(x)+\left(\alpha_{j^{\prime}}-\alpha_{j}\right) x\right]\right) d x \tag{3.153}
\end{equation*}
$$

The latter is written in such a way that integration by parts is obvious. Using integration by parts, the equations reduce to a simple form

$$
\begin{equation*}
\sum_{j^{\prime}=-\infty}^{\infty} M_{1 j j^{\prime}}^{ \pm}(P O)\left[N_{j^{\prime}}^{(P O)} \mp a_{1 j j^{\prime}}^{ \pm} \psi_{j^{\prime}}^{(P O)}\right]=2 \gamma_{1 j}\left\{\left\{_{A_{j}}^{-I_{j}}\right\},\right. \tag{3.154}
\end{equation*}
$$

where

$$
\begin{equation*}
a_{1 j j^{\prime}}^{ \pm}=\frac{ \pm\left(1-\alpha_{j} \alpha_{j^{\prime}}\right)-\gamma_{1 j} \gamma_{10}}{ \pm \gamma_{1 j}-\gamma_{10}} \tag{3.155}
\end{equation*}
$$

and

$$
\sum_{j^{\prime}=-\infty}^{\infty} M_{2 j j^{\prime}}^{ \pm}(P O)\left[N_{j^{\prime}}^{(P O)} \mp a_{2 j j^{\prime}}^{ \pm} \psi_{j^{\prime}}^{(P O)}\right]=-2 \gamma_{2 j}\left\{\begin{array}{l}
B_{j}  \tag{3.156}\\
0
\end{array}\right\}
$$

where

$$
\begin{equation*}
a_{2 j j^{\prime}}^{ \pm}=\frac{ \pm\left(K^{2}-\alpha_{j} \alpha_{j^{\prime}}\right)-\gamma_{2 j} \gamma_{10}}{\rho\left( \pm \gamma_{2 j}-\gamma_{10}\right)} . \tag{3.157}
\end{equation*}
$$

Finally, it is useful to rewrite the physical optics matrix elements as

$$
\begin{equation*}
M_{p j j^{\prime}}^{ \pm}(P O)=\frac{1}{L} \int_{-L / 2}^{L / 2} \exp \left[-i 2 \pi\left(j-j^{\prime}\right) x / L+i k_{1}\left( \pm \gamma_{p j}-\gamma_{10}\right) h(x)\right] d x, \tag{3.158}
\end{equation*}
$$

which displays the Fourier part explicitly. Note that for a flat surface, the only elements of 3.158 which survive are the diagonal elements $j=j^{\prime}$ (which equal 1 ). Further, $\tilde{M}_{p j j^{\prime}}^{ \pm}(P O)=0$, $a_{1 j j^{\prime}}^{ \pm}=\gamma_{1 j}, a_{2 i j^{\prime}}^{ \pm}=\gamma_{2 j} / \rho$, and, for a single plane wave incidence $\left(I_{j}=\delta_{j 0}\right)$, the usual flat surface limit of (3.154) and (3.156) follows directly.

A Floquet-Fourier (FF) expansion for the boundary unknowns can be written as

$$
\begin{equation*}
\psi\left(\vec{x}_{h}\right)=\sum_{j^{\prime}=-\infty}^{\infty} \psi_{j^{\prime}}^{(F F)} \exp \left(i k_{1} \alpha_{j^{\prime}} x\right) \tag{3.159}
\end{equation*}
$$

and

$$
\begin{equation*}
N\left(\vec{x}_{h}\right)=i k_{1} \sum_{j^{\prime}=-\infty}^{\infty} N_{j^{\prime}}^{(F F)} \exp \left(i k_{1} \alpha_{j^{\prime}} x\right) \tag{3.160}
\end{equation*}
$$

The equations corresponding to (3.154) and (3.156) are thus

$$
\sum_{j^{\prime}=-\infty}^{\infty} M_{1 j j^{\prime}}^{ \pm}(F F)\left[N_{j^{\prime}}^{(F F)} \mp a_{1 j j^{\prime}} \psi j^{\prime(F F)}\right]=2 \gamma_{1 j}\left\{\begin{array}{l}
-I_{j}  \tag{3.161}\\
A_{j}
\end{array}\right\}
$$

and

$$
\sum_{j^{\prime}=-\infty}^{\infty} M_{2 j j^{\prime}}^{ \pm}(F F)\left[N_{j^{\prime}}^{(F F)} \mp a_{2 j j^{\prime}} \psi_{j^{\prime}}^{(F F)}\right]=-2 \gamma_{2 j}\left\{\begin{array}{l}
B_{j}  \tag{3.162}\\
0
\end{array}\right\},
$$

where

$$
\begin{gather*}
a_{1 j j^{\prime}}=\frac{1-\alpha_{j} \alpha_{j^{\prime}}}{\gamma_{1 j}}  \tag{3.163}\\
a_{2 j j^{\prime}}=\frac{K^{2}-\alpha j \alpha_{j^{\prime}}}{\rho \gamma_{2 j}} \tag{3.164}
\end{gather*}
$$

and, for $p=1,2$, the matrix elements are

$$
\begin{equation*}
M_{p j j^{\prime}}^{ \pm}(F F)=\frac{1}{L} \int_{-L / 2}^{L / 2} \exp \left[-i 2 \pi\left(j-j^{\prime}\right) x / L \pm i k_{1} \gamma_{p j} h(x)\right] d x \tag{3.165}
\end{equation*}
$$

The FF equations follow from (3.154) through (3.158) by setting the $\gamma_{10}$ term to zero. These provide an alternative set of equations to solve for the alternative boundary function coefficients to produce the same coefficients for the scattered and transmitted fields ${ }^{588}$.

### 3.8 Full Spectral Methods for Scalar Problems: Conjugate Rayleigh Basis

A further spectral expansion consists in modifying the physical optics expansion by making the single physical optics plane wave dependent on the Bragg mode, so that the phase height term is dependent on the mode, and this leads to a conjugate Rayleigh (CR) expansion using the complex conjugate of the plane wave states 3.100 evaluated on the surface as

$$
\begin{equation*}
\psi\left(\vec{x}_{h}\right)=\sum_{j^{\prime}=-\infty}^{\infty} \psi_{j^{\prime}}^{(C R)} \exp \left[i k_{1} \alpha_{j^{\prime}} x-i k_{1} \bar{\gamma}_{1 j^{\prime}} h(x)\right]=\sum_{j^{\prime}=-\infty}^{\infty} \psi_{j^{\prime}}^{(C R)} \bar{\phi}_{1 j^{\prime}}^{+}\left(\vec{x}_{h}\right), \tag{3.166}
\end{equation*}
$$

and the scaled expansion for the normal derivative

$$
\begin{equation*}
N\left(\vec{x}_{h}\right)=i k_{1} \sum_{j^{\prime}=-\infty}^{\infty} N_{j^{\prime}}^{(C R)} \bar{\phi}_{1 j^{\prime}}^{+}\left(\vec{x}_{h}\right), \tag{3.167}
\end{equation*}
$$

where the overbar is complex conjugation. Substituting these expansions in 3.105) and 3.111, and carrying out the integration by parts necessary to simplify the slope terms as in Sec. 7 yields equations similar in form to (3.154) and (3.156). For the upper region equation we get

$$
\sum_{j^{\prime}=-\infty}^{\infty} M_{1 j j^{\prime}}^{ \pm}(C R)\left[N_{j^{\prime}}^{(C R)} \mp b_{1 j j^{\prime}}^{ \pm} \psi_{j^{\prime}}^{(C R)}\right]=2 \gamma_{1 j}\left\{\begin{array}{l}
-I_{j} \tag{3.168}
\end{array}\right\}
$$

where

$$
\begin{equation*}
b_{1 j j^{\prime}}^{ \pm}=\frac{1-\alpha_{j} \alpha_{j^{\prime}} \mp \gamma_{1 j} \bar{\gamma}_{1 j^{\prime}}}{\gamma_{1 j} \mp \bar{\gamma}_{1 j^{\prime}}} \tag{3.169}
\end{equation*}
$$

and the matrix elements are defined as

$$
\begin{equation*}
M_{1 j j^{\prime}}^{ \pm}(C R)=\frac{1}{L} \int_{-L / 2}^{L / 2} \exp \left[-i 2 \pi\left(j-j^{\prime}\right) x / L+i k_{1}\left( \pm \gamma_{1 j}-\bar{\gamma}_{1 j^{\prime}}\right) h(x)\right] d x=<\phi_{1 j}^{ \pm}, \phi_{1 j^{\prime}}^{+}> \tag{3.170}
\end{equation*}
$$

It is obvious that $M_{1 j j^{\prime}}^{+}(C R)$ is self-adjoint, positive definite and thus invertible, and this fact was used with success in solving the Dirichlet problem ${ }^{3}$. That is,

$$
\begin{equation*}
\left[M_{\left.1 j j^{\prime}\right]^{\star}}^{+}(C R)=M_{1 j j^{\prime}}^{+}(C R),\right. \tag{3.171}
\end{equation*}
$$

where the symbol $\star$ represents the adjoint.
The same expansion for the equations in region 2 yields the equations

$$
\sum_{j^{\prime}=-\infty}^{\infty} M_{2 j j^{\prime}}^{ \pm}(C R)\left[N_{j^{\prime}}^{(C R)} \mp b_{2 j j^{\prime}}^{ \pm} \psi_{j^{\prime}}^{(C R)}\right]=-2 \gamma_{2 j}\left\{\begin{array}{l}
B_{j}  \tag{3.172}\\
0
\end{array}\right\},
$$

where

$$
\begin{equation*}
b_{2 j j^{\prime}}^{ \pm}=\frac{1}{\rho} \frac{K_{2}-\alpha_{j} \alpha_{j^{\prime}} \mp \gamma_{2 j} \bar{\gamma}_{1 j^{\prime}}}{\gamma_{2 j} \mp \bar{\gamma}_{1 j^{\prime}}} \tag{3.173}
\end{equation*}
$$

and the matrix elements are

$$
\begin{equation*}
M_{2 j j^{\prime}}^{ \pm}(C R)=\frac{1}{L} \int_{-L / 2}^{L / 2} \exp \left[-i 2 \pi\left(j-j^{\prime}\right) x / L+i k_{1}\left( \pm \gamma_{2 j}-\bar{\gamma}_{1 j^{\prime}}\right) h(x)\right] d x=<\phi_{2 j}^{ \pm}, \phi_{1 j^{\prime}}^{+}> \tag{3.174}
\end{equation*}
$$

### 3.9 Integral Equation Methods in Coordinate Space for Electromagnetic Problems

Up to now we have considered one-dimensional surfaces and acoustic problems. These correspond directly to electromagnetic scattering problems where there is no change in polarization for the scattered and transmitted fields. The general electromagnetic problem for a periodic dielectric interface is for a two-dimensional surface $z=h\left(\vec{x}_{\perp}\right)=h(x, y)$ which separates media of different dielectric constants $\varepsilon_{j}$ for $j=1,2$ and permeability $\mu_{j}$. The wave numbers for the two regions are $k_{j}=k_{0} \sqrt{\varepsilon_{j} \mu_{j}}$ where $k_{0}=\omega / c, \omega$ is the circular frequency and $c$ the speed of light. There is now a change in polarization in the scattered and transmitted fields.

For the source-free electric field $\partial_{i} E_{i}=0$, and each component of the electric field $E_{i}$ with $i=1,2,3$ satisfies the same Helmholtz equation as the scalar field, viz. in region 1

$$
\begin{equation*}
\left(\nabla_{3}^{2}+k_{1}^{2}\right) E_{1 i}(\vec{x})=0, \tag{3.175}
\end{equation*}
$$

where $E_{1 i}$ is the $i^{t h}$ electric field component in region 1, and the Laplacian is three-dimensional. This is just the vector analogue of (3.67). We can use this to write the vector analogues of the scalar equations in Sec. 4 , using Green's theorem, the three-dimensional periodic Green's functions $G^{(3 p)}$ from (3.59) and (3.62), the two-dimensional Floquet periodicity of the field, and the characteristic function defining the region, e.g. for region 1

$$
\begin{equation*}
\Theta_{1}(\vec{x})=\theta\left(L_{1} / 2-x\right) \theta\left(x+L_{1} / 2\right) \theta\left(L_{2} / 2-y\right) \theta\left(y+L_{2} / 2\right) \theta\left(z-h\left(\vec{x}_{\perp}\right)\right) \theta\left(H_{1}-z\right) . \tag{3.176}
\end{equation*}
$$

In region 1 the result is

$$
\begin{equation*}
E_{1 i}\left(\vec{x}^{\prime}\right) \Theta_{1}\left(\vec{x}^{\prime}\right)=E_{i}^{i n}\left(\vec{x}^{\prime}\right)-\left[G_{1}^{(3 p)}, E_{1 i} ; \vec{x}^{\prime}, h\right], \tag{3.177}
\end{equation*}
$$

where $E_{i}^{i n}$ is the incident field, and the two-dimensional bracket is explicitly

$$
\begin{equation*}
\left[G_{1}^{(3 p)}, E_{1 i} ; \vec{x}^{\prime}, h\right]=\iint_{D}\left[G_{1}^{(3 p)}\left(\vec{x}^{\prime}, \vec{x}_{h}\right) N_{1 i}\left(\vec{x}_{h}\right)-N_{1}^{(3 p)}\left(\vec{x}^{\prime}, \vec{x}_{h}\right) E_{1 i}\left(\vec{x}_{h}\right)\right] d \vec{x}_{\perp} \tag{3.178}
\end{equation*}
$$

where $\vec{x}_{\perp}=(x, y)$, the domain of integration $D$ is $x \varepsilon\left[-L_{1} / 2, L_{1} / 2\right], y \varepsilon\left[-L_{2} / 2, L_{2} / 2\right]$, and the normal derivatives are

$$
\begin{equation*}
N_{1 i}\left(\vec{x}_{h}\right)=n_{l}\left(\vec{x}_{\perp}\right) \partial_{l} E_{1 i}\left(\vec{x}_{h}\right), \tag{3.179}
\end{equation*}
$$

and

$$
\begin{equation*}
N_{1}^{(3 p)}\left(\vec{x}^{\prime}, \vec{x}_{h}\right)=n_{l}\left(\vec{x}_{\perp}\right) \partial_{l} G_{1}^{(3 p)}\left(\vec{x}^{\prime}, \vec{x}_{h}\right), \tag{3.180}
\end{equation*}
$$

the normal derivatives of the boundary unknown and the periodic Green's function respectively. From (3.177), the field representation in $D$ is found by setting $\Theta_{1}=1$, the Extinction Theorem by setting $\Theta_{1}=0$, the scattered field is just the bracket term

$$
\begin{equation*}
E_{1 i}^{s c}\left(\vec{x}^{\prime}\right)=-\left[G_{1}^{(3 p)}, E_{1 i} ; \vec{x}^{\prime}, h\right], \tag{3.181}
\end{equation*}
$$

and the boundary integral equation is

$$
\begin{equation*}
\frac{1}{2} E_{1 i}\left(\vec{x}_{h}^{\prime}\right)=E_{i}^{i n}\left(\vec{x}_{h}^{\prime}\right)-\iint_{D}\left[G_{1}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right) N_{1 i}\left(\vec{x}_{h}\right)-N_{1}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right) E_{1 i}\left(\vec{x}_{h}\right)\right] d \vec{x}_{\perp} \tag{3.182}
\end{equation*}
$$

with the boundary unknowns $E_{1 i}$ and its normal derivative $N_{1 i}$.
Green's theorem in region 2 yields the representation for the total transmitted field

$$
\begin{equation*}
E_{2 i}\left(\vec{x}^{\prime}\right) \Theta_{2}\left(\vec{x}^{\prime}\right)=\left[G_{2}^{(3 p)}, E_{2 i} ; \vec{x}^{\prime}, h\right], \tag{3.183}
\end{equation*}
$$

where $\Theta_{2}=1-\Theta_{1}$, and the boundary integral equation becomes

$$
\begin{equation*}
\frac{1}{2} E_{2 i}\left(\vec{x}_{h}^{\prime}\right)=\iint_{D}\left[G_{2}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right) N_{2 i}\left(\vec{x}_{h}\right)-N_{2}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right) E_{2 i}\left(\vec{x}_{h}\right)\right] d \vec{x}_{\perp} \tag{3.184}
\end{equation*}
$$

with boundary unknowns $E_{2 i}$ and its normal derivative $N_{2 i}$. Equations (3.182) and (3.184) are similar to the scalar equations (3.74) and (3.88), but the fields and their normal derivatives are not the usual electromagnetic boundary values, the latter being typically written in terms of normal field components and currents ${ }^{61]}$. So to continue we must relate these usual boundary conditions to our boundary unknowns.

For the electric field on the boundary we have the continuity condition of the normal component of the displacement vector $\vec{D}=\varepsilon \vec{E}$ which becomes

$$
\begin{equation*}
\varepsilon \vec{n} \cdot \vec{E}_{2}=\vec{n} \cdot \vec{E}_{1} \tag{3.185}
\end{equation*}
$$

where $\varepsilon=\varepsilon_{2} / \varepsilon_{1}$, and the continuity of the magnetic current

$$
\begin{equation*}
\vec{n} \times \vec{E}_{2}=\vec{n} \times \vec{E}_{1} . \tag{3.186}
\end{equation*}
$$

These are four equations, three of which are independent. These three can be solved directly, or the four equations solved using a Moore-Penrose pseudo inverse to yield the boundary conditions on the electric field (in index notation) as

$$
\begin{equation*}
E_{2 i}\left(\vec{x}_{h}\right)=C_{i j}\left(\vec{x}_{h}\right) E_{1 j}\left(\vec{x}_{h}\right) \tag{3.187}
\end{equation*}
$$

with repeated subscripts summed from 1 to 3 and

$$
\begin{equation*}
C_{i j}\left(\vec{x}_{h}\right)=\delta_{i j}+\left(\varepsilon^{-1}-1\right) \hat{n}_{i} \hat{n}_{j}, \tag{3.188}
\end{equation*}
$$

with $\hat{n}$ representing the unit normal. These boundary conditions were introduced some time ago ${ }^{355}$ and used successfully for scattering from a body of revolution ${ }^{41}$.

The continuity conditions on the normal derivative components are more involved. The full details are in ${ }^{35}$. Briefly we introduce the bracket notation for when we set the field on the surface first and then differentiate

$$
\begin{equation*}
\left\{E_{m}\right\} \doteq E_{m}(x, y, h(x, y)) . \tag{3.189}
\end{equation*}
$$

Then the transverse (" t ") derivatives ( $x$ and $y$ ) are given by

$$
\begin{equation*}
\partial_{x}\left\{E_{m}\right\}=\left\{\partial_{x} E_{m}\right\}+h_{x}\left\{\partial_{z} E_{m}\right\} \tag{3.190}
\end{equation*}
$$

and

$$
\begin{equation*}
\partial_{y}\left\{E_{m}\right\}=\left\{\partial_{y} E_{m}\right\}+h_{y}\left\{\partial_{z} E_{m}\right\} . \tag{3.191}
\end{equation*}
$$

Using this notation and the continuity of the electric surface current $\vec{K}^{e}=-\vec{n} \times \vec{H}$, where $\vec{H}$ is the magnetic field, in index form

$$
\begin{equation*}
K_{2 i}^{e}\left(\vec{x}_{h}\right)=K_{1 i}^{e}\left(\vec{x}_{h}\right), \tag{3.192}
\end{equation*}
$$

we can write the continuity condition for the normal derivative as ${ }^{35}$

$$
\begin{equation*}
\left\{N_{2 i}\right\}=\mu\left\{N_{1 i}\right\}+\left(\varepsilon^{-1}-1\right) V_{i}\left(\vec{x}_{h}\right), \tag{3.193}
\end{equation*}
$$

where $V_{i}$ can be written in terms of transverse partial derivatives involving the normal components of the electric field as

$$
\begin{equation*}
V_{i}\left(\vec{x}_{h}\right)=n_{m} \partial_{i t}\left\{\hat{n}_{m} \hat{n}_{j} E_{1 j}\right\}-n_{i} \partial_{q t}\left\{\hat{n}_{q t} \hat{n}_{j} E_{1 j}\right\} . \tag{3.194}
\end{equation*}
$$

This $V_{i}$ term looks awkward, but it can be integrated by parts. First, choose the boundary unknowns as

$$
\begin{equation*}
E_{1 i}\left(\vec{x}_{h}\right)=\left\{E_{1 i}\right\} \doteq\left\{E_{i}\right\}, \tag{3.195}
\end{equation*}
$$

and

$$
\begin{equation*}
N_{1 i}\left(\vec{x}_{h}\right)=\left\{N_{1 i}\right\} \doteq\left\{N_{i}\right\} . \tag{3.196}
\end{equation*}
$$

Then we can write the equation for the upper region (3.182) as

$$
\begin{equation*}
\frac{1}{2}\left\{E_{i}^{\prime}\right\}+\iint_{D}\left[G_{1}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)\left\{N_{i}\right\}-N_{1}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)\left\{E_{i}\right\}\right] d \vec{x}_{\perp}=E_{i}^{i n}\left(\vec{x}_{h}^{\prime}\right) . \tag{3.197}
\end{equation*}
$$

Here $\left\{E_{i}^{\prime}\right\}$ means the exterior primed variable placed on the surface, i.e. $\vec{x}_{h}^{\prime}$. The equation (3.197) is diagonal in the index. The coupling is from the lower equation (3.184) written using (3.193) through (3.196) as

$$
\begin{equation*}
\frac{1}{2} C_{i j}\left(\vec{x}_{h}^{\prime}\right)\left\{E_{j}^{\prime}\right\}=\iint_{D}\left[G_{2}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)\left(\mu\left\{N_{i}\right\}+\left(\varepsilon^{-1}-1\right) V_{i}\left(\vec{x}_{h}\right)\right)-N_{2}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right) C_{i j}\left(\vec{x}_{h}\right)\left\{E_{j}\right\}\right] d \vec{x}_{\perp} \tag{3.198}
\end{equation*}
$$

The $V_{i}$ term can be integrated by parts to yield

$$
\begin{equation*}
\iint_{D} G_{2}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right) V_{i}\left(\vec{x}_{h}\right) d \vec{x}_{\perp}=\iint_{D} V_{i j}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)\left\{E_{i}\right\} d \vec{x}_{\perp} \tag{3.199}
\end{equation*}
$$

where

$$
\begin{equation*}
V_{i j}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)=\partial_{q t}\left\{n_{i} G_{2}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)\right\} \hat{n}_{q t} \hat{n}_{j}-\partial_{i t}\left\{n_{m} G_{2}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)\right\} \hat{n}_{m} \hat{n}_{j} . \tag{3.200}
\end{equation*}
$$

We can simplify 3.200 to yield

$$
\begin{equation*}
V_{i j}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)=N_{2}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right) \hat{n}_{i} \hat{n}_{j}-\left\{\partial_{i} G_{2}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)\right\} n_{j}, \tag{3.201}
\end{equation*}
$$

where now the derivative of the Green's function is taken first, and then the result set on the surface. Combining these results we can rewrite (3.198) as

$$
\begin{equation*}
\frac{1}{2} C_{i j}\left(\vec{x}_{h}^{\prime}\right)\left\{E_{j}^{\prime}\right\}=\iint_{D}\left[G_{2}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right) \mu\left\{N_{i}\right\}-W_{i j}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)\left\{E_{j}\right\}\right] d \vec{x}_{\perp}, \tag{3.202}
\end{equation*}
$$

where

$$
\begin{equation*}
W_{i j}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)=N_{2}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right) \delta_{i j}+\left(\varepsilon^{-1}-1\right)\left\{\partial_{i} G_{2}^{(3 p)}\left(\vec{x}_{h}^{\prime}, \vec{x}_{h}\right)\right\} n_{j} . \tag{3.203}
\end{equation*}
$$

Note that (3.197) and (3.202), if put in matrix form, are diagonal in three of the four matrix blocks multiplying the six-dimensional vector of boundary unknowns $\left[\left\{E_{i}\right\},\left\{N_{i}\right\}\right]^{T}$ where $T$ is transpose. The only coupling occurs in the single block of the electric fields from (3.202). We note this in contrast to pre-conditioning methods used to sparsify matrix inversion problems. Here the results are exact and highly sparse as formulated. They have been used computationally to treat the scattering from a body of revolution ${ }^{411}$.

We can use these representations to write plane wave representations for the scattered and transmitted fields, above and below the largest surface excursions. From (3.177) we can write the scattered field above the highest surface excursion $\left(z^{\prime}>\max (h)\right)$ as

$$
\begin{equation*}
E_{i}^{s c}\left(\vec{x}^{\prime}\right)=-\iint_{D}\left[G_{1}^{(3 p)}\left(\vec{x}^{\prime}, \vec{x}_{h}\right)\left\{N_{i}\right\}-N_{1}^{(3 p)}\left(\vec{x}^{\prime}, \vec{x}_{h}\right)\left\{E_{i}\right\}\right] d \vec{x}_{\perp} \tag{3.204}
\end{equation*}
$$

where now the Green's function is, following (3.59), with the field point above the surface

$$
\begin{equation*}
G_{1}^{(3 p)}\left(\vec{x}^{\prime}, \vec{x}_{h}\right)=\frac{i}{2 k_{1} L_{1} L_{2}} \sum_{j=-\infty}^{\infty} \sum_{j^{\prime}=-\infty}^{\infty} \frac{\exp \left[i k_{1}\left(\alpha_{j}\left(x^{\prime}-x\right)+\beta_{j^{\prime}}\left(y^{\prime}-y\right)+\gamma_{1 j^{\prime}}\left(z^{\prime}-h\left(\vec{x}_{\perp}\right)\right)\right]\right.}{\gamma_{1 j j^{\prime}}}, \tag{3.205}
\end{equation*}
$$

and

$$
\begin{equation*}
N_{1}^{(3 p)}\left(\vec{x}^{\prime}, \vec{x}_{h}\right)=n_{q} \partial_{q} G_{1}^{(3 p)}\left(\vec{x}^{\prime}, \vec{x}_{h}\right) . \tag{3.206}
\end{equation*}
$$

Combining these results we can write the scattered field exactly above the highest surface excursion as a plane wave expansion in terms of purely up-going waves as

$$
\begin{equation*}
E_{i}^{s c}\left(\vec{x}^{\prime}\right)=\sum_{j=-\infty}^{\infty} \sum_{j^{\prime}=-\infty}^{\infty} A_{i j j^{\prime}} \exp \left[i k_{1}\left(\alpha_{j} x^{\prime}+\beta_{j^{\prime}} y^{\prime}+\gamma_{1 j j^{\prime}} z^{\prime}\right)\right] \tag{3.207}
\end{equation*}
$$

where

$$
\begin{equation*}
A_{i j j^{\prime}}=\frac{1}{L_{1} L_{2}} \iint_{D} A_{i j j^{\prime}}\left(\vec{x}_{\perp}\right) \exp \left[-i k_{1}\left(\alpha_{j} x+\beta_{j^{\prime}} y+\gamma_{1 j j^{\prime}} h\left(\vec{x}_{\perp}\right)\right)\right] d \vec{x}_{\perp}, \tag{3.208}
\end{equation*}
$$

and

$$
\begin{equation*}
A_{i j j^{\prime}}\left(\vec{x}_{\perp}\right)=\frac{-i}{8 \pi^{2} k_{1} \gamma_{1 j j^{\prime}}}\left[\left\{N_{i}\right\}+i k_{1}\left(\gamma_{1 j j^{\prime}}-\alpha_{j} h_{x}-\beta_{j^{\prime}} h_{y}\right)\left\{E_{i}\right\}\right], \tag{3.209}
\end{equation*}
$$

in terms of the boundary unknowns.
Similarly, from (3.183), we have the transmitted field below the lowest surface excursion $\left(z^{\prime}<\min (h)\right)$

$$
\begin{equation*}
E_{2 i}\left(\vec{x}^{\prime}\right)=\iint_{D}\left[G_{2}^{(3 p)}\left(\vec{x}^{\prime}, \vec{x}_{h}\right) N_{2 i}\left(\vec{x}_{h}\right)-N_{2}^{(3 p)}\left(\vec{x}^{\prime}, \vec{x}_{h}\right) E_{2 i}\left(\vec{x}_{h}\right)\right] d \vec{x}_{\perp}, \tag{3.210}
\end{equation*}
$$

where now

$$
\begin{equation*}
G_{2}^{(3 p)}\left(\vec{x}^{\prime}, \vec{x}_{h}\right)=\frac{i}{8 \pi^{2} k_{1} L_{1} L_{2}} \sum_{j=-\infty}^{\infty} \sum_{j^{\prime}=-\infty}^{\infty} \frac{\exp \left[i k_{1}\left(\alpha_{j}\left(x^{\prime}-x\right)+\beta_{j^{\prime}}\left(y^{\prime}-y\right)-\gamma_{2 j j^{\prime}}\left(z^{\prime}-h\left(\vec{x}_{\perp}\right)\right)\right]\right.}{\gamma_{2 j j^{\prime}}}, \tag{3.211}
\end{equation*}
$$

and

$$
\begin{equation*}
N_{2}^{(3 p)}\left(\vec{x}^{\prime}, \vec{x}_{h}\right)=n_{q} \partial_{q} G_{2}^{(3 p)}\left(\vec{x}^{\prime}, \vec{x}_{h}\right) \tag{3.212}
\end{equation*}
$$

The result is the plane wave spectral representation for the transmitted field below the lowest surface excursion in terms of purely down-going waves as

$$
\begin{equation*}
E_{2 i}\left(\vec{x}^{\prime}\right)=\sum_{j=-\infty}^{\infty} \sum_{j^{\prime}=-\infty}^{\infty} B_{i j j^{\prime}} \exp \left[i k_{1}\left(\alpha_{j} x^{\prime}+\beta_{j^{\prime}} y^{\prime}-\gamma_{2 j j^{\prime}} z^{\prime}\right)\right] \tag{3.213}
\end{equation*}
$$

where

$$
\begin{equation*}
B_{i j j^{\prime}}=\frac{1}{L_{1} L_{2}} \iint_{D} B_{i j j^{\prime}}\left(\vec{x}_{\perp}\right) \exp \left[-i k_{1}\left(\alpha_{j} x+\beta_{j^{\prime}} y-\gamma_{2 j j^{\prime}} h\left(\vec{x}_{\perp}\right)\right)\right] d \vec{x}_{\perp}, \tag{3.214}
\end{equation*}
$$

and

$$
\begin{equation*}
B_{i j j^{\prime}}\left(\vec{x}_{\perp}\right)=\frac{i}{8 \pi^{2} k_{1} \gamma_{2 j j^{\prime}}}\left[N_{2 i}\left(\vec{x}_{h}\right)-i k_{1}\left(\gamma_{2 j j^{\prime}}+\alpha_{j} h_{x}+\beta_{j^{\prime}} h_{y}\right) E_{2 i}\left(\vec{x}_{h}\right)\right], \tag{3.215}
\end{equation*}
$$

written in terms of the boundary values from the lower region. Using the boundary conditions (3.187) and (3.193) and integration by parts we can rewrite (3.215) in terms of the boundary unknowns as

$$
\begin{equation*}
B_{i j j^{\prime}}\left(\vec{x}_{\perp}\right)=\frac{i}{8 \pi^{2} k_{1} \gamma_{2 j j^{\prime}}}\left[\mu\left\{N_{i}\right\}-W_{i j j^{\prime} l}\left(\vec{x}_{\perp}\right)\left\{E_{l}\right\}\right], \tag{3.216}
\end{equation*}
$$

where

$$
\begin{equation*}
\left.W_{i j j^{\prime} l}\left(\vec{x}_{\perp}\right)=i k_{1}\left[\alpha_{j} h_{x}+\beta_{j^{\prime}} h_{y}+\gamma_{2 j j^{\prime}}\right) \delta_{i l}-\left(\varepsilon^{-1}-1\right)\left(\delta_{i 1} \alpha_{j}+\delta_{i 2} \beta_{j^{\prime}}-\delta_{i 3} \gamma_{2 j j^{\prime}}\right) n_{l}\right] . \tag{3.217}
\end{equation*}
$$

Equations (3.207) and (3.213) are the exact plane wave representations in the appropriate regions. In the next section we write general partial spectral representations of the fields, and show the relations between them and the plane wave spectral representations here which are valid in limited domains.

### 3.10 Partial Spectral Methods for Electromagnetic Problems

We develop this section in analogy with the scalar results in Sec.5. This is the electromagnetic version of the Spectral-Coordinate approach. We define the three-dimensional plane wave states in the upper region 1 for up ${ }^{(+)}$- and down ${ }^{(-)}$-going waves as

$$
\begin{equation*}
\phi_{1 j j^{\prime}}^{ \pm}(\vec{x})=\exp \left[i k_{1}\left(-\alpha_{j} x-\beta_{j^{\prime}} y \pm \gamma_{1 j j^{\prime}} z\right),\right. \tag{3.218}
\end{equation*}
$$

where $\gamma_{1 j j^{\prime}}$ is defined following (3.59). The function satisfies the three-dimensional Helmholtz equation

$$
\begin{equation*}
\left(\nabla_{3}^{2}+k_{1}^{2}\right) \phi_{1 j j^{\prime}}^{ \pm}(\vec{x})=0 \tag{3.219}
\end{equation*}
$$

The incident electric field can be written as a general plane wave expansion of down-going waves

$$
\begin{equation*}
E_{i}^{i n}(\vec{x})=\sum_{j=-\infty}^{\infty} \sum_{j^{\prime}=-\infty}^{\infty} I_{i j j^{\prime}} \exp \left[i k_{1}\left(\alpha_{j} x+\beta_{j^{\prime}} y-\gamma_{1 j j^{\prime}} z\right)\right] \tag{3.220}
\end{equation*}
$$

Apply Green's theorem in the domain defined by $\Theta_{1}$ in 3.176 ) to $\phi_{1 j j^{\prime}}^{ \pm}$and $E_{1 i}$, use the twodimensional Floquet conditions to cancel the side integrals as in Sec. 9 and the result is

$$
\frac{1}{L_{1} L_{2}} \iint_{D} \phi_{1 j j^{\prime}}^{ \pm}\left(\vec{x}_{h}\right) U_{i j j^{\prime}}^{ \pm}\left(\vec{x}_{h}\right) d \vec{x}_{\perp}=\gamma_{1 j j^{\prime}}\left\{\begin{array}{l}
-I_{i j j^{\prime}} \tag{3.221}
\end{array}\right\}
$$

where $U$ is defined as

$$
\begin{equation*}
U_{i j j^{\prime}}^{ \pm}\left(\vec{x}_{h}\right)=\frac{1}{2 i k_{1}}\left[\left\{N_{i}\right\}-i k_{1}\left( \pm \gamma_{1 j^{\prime}}+\alpha_{j} h_{x}+\beta_{j^{\prime}} h_{y}\right)\left\{E_{i}\right\}\right] . \tag{3.222}
\end{equation*}
$$

The $A_{i j j^{\prime}}$ are the spectral coefficients of the scattered field from (3.207). Recall that the scattered field is evaluated on a flat surface $\left(z=H_{1}\right)$ above the highest surface excursion, so the representation (3.207) is rigorously valid and not a Rayleigh approximation. We also used the boundary values (3.195) and 3.196). In 3.221), the up-going plane wave states $\phi^{+}$project out the down-going incident field spectral components $I_{i j j^{\prime}}$, and the down-going plane wave states $\phi^{-}$project out the up-going scattered field spectral components $A_{i j j^{\prime}}$. Equations 3.221 and (3.222) are the vector generalizations of (3.105) and (3.106).

For the lower region 2, the three-dimensional up- and down-going plane wave states are defined as

$$
\begin{equation*}
\phi_{2 j j^{\prime}}^{ \pm}(\vec{x})=\exp \left[i k_{1}\left(-\alpha_{j} x-\beta_{j^{\prime}} y \pm \gamma_{2 j j^{\prime}} z\right)\right], \tag{3.223}
\end{equation*}
$$

where $\gamma_{2 j j^{\prime}}$ is defined following (3.62). The functions satisfy the three-dimensional Helmholtz equation

$$
\begin{equation*}
\left(\nabla_{3}^{2}+k_{2}^{2}\right) \phi_{2 j j^{\prime}}^{ \pm}(\vec{x})=0 \tag{3.224}
\end{equation*}
$$

Green's theorem on $\phi_{2 j j^{\prime}}^{ \pm}$and the total transmitted field $E_{2 i}$ in the domain defined by $\Theta_{2}(\vec{x})=$ $1-\Theta_{1}(\vec{x})$ yields the relations

$$
\frac{1}{L_{1} L_{2}} \iint_{D} \phi_{2 j j^{\prime}}^{ \pm}\left(\vec{x}_{h}\right) L_{i j j^{\prime}}^{ \pm}\left(\vec{x}_{h}\right) d \vec{x}_{\perp}=-\gamma_{2 j j^{\prime}}\left\{\begin{array}{l}
B_{i j j^{\prime}} \tag{3.225}
\end{array}\right\}
$$

where

$$
\begin{equation*}
L_{i j j^{\prime}}^{ \pm}\left(\vec{x}_{h}\right)=\frac{1}{2 i k_{1}}\left[N_{2 i}\left(\vec{x}_{h}\right)-i k_{1}\left( \pm \gamma_{2 j j^{\prime}}+\alpha_{j} h_{x}+\beta_{j^{\prime}} h_{y}\right) E_{2 i}\left(\vec{x}_{h}\right)\right], \tag{3.226}
\end{equation*}
$$

in terms of the boundary values from the lower region. Using the boundary values 3.195) and (3.196) and integration by parts, (3.226) can be rewritten as

$$
\begin{equation*}
L_{i j j^{\prime}}^{ \pm}\left(\vec{x}_{h}\right)=\frac{1}{2 i k_{1}}\left[\mu\left\{N_{i}\right\}-W_{i j j^{\prime} l}^{ \pm}\left(\vec{x}_{h}\right)\left\{E_{l}\right\}\right], \tag{3.227}
\end{equation*}
$$

with a sum over $l=(1,2,3)$ and where

$$
\begin{equation*}
W_{i j j^{\prime} l}^{ \pm}\left(\vec{x}_{h}\right)=i k_{1}\left[\left(\alpha_{j} h_{x}+\beta_{j^{\prime}} h_{y} \pm \gamma_{2 j j^{\prime}}\right) \delta_{i l}-\left(\varepsilon^{-1}-1\right)\left(\alpha_{j} \delta_{i 1}+\beta_{j^{\prime}} \delta_{i 2} \mp \gamma_{2 j j^{\prime}} \delta_{i 3}\right) n_{l}\right] . \tag{3.228}
\end{equation*}
$$

Note that $W_{i j j^{\prime} l}^{+}$is just $W_{i j j^{\prime} l}$ from $(3.217$ ). Equations $\sqrt{3.225}$ and 3.227 are the vector generalizations of (3.111) and (3.112). The procedure is to solve the upper equation (3.221) and the lower equation (3.225) for the boundary unknowns $\{N i\}$ and $\left\{E_{i}\right\}$ and evaluate the remaining equations for the scattered and transmitted amplitudes.

### 3.11 Full Spectral Methods for Electromagnetic Problems

In this section we develop the full spectral methods using the conjugate Rayleigh basis in analogy with Sec. 8 for the scalar case. In (3.221) and (3.225) we use the following expansions in the conjugate Rayleigh basis,

$$
\begin{equation*}
\left\{E_{i}\right\}=\sum_{l=-\infty}^{\infty} \sum_{l^{\prime}=-\infty}^{\infty} E_{i l l^{\prime}} \bar{\phi}_{1 l l^{\prime}}^{+}\left(\vec{x}_{h}\right), \tag{3.229}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\{N_{i}\right\}=i k_{1} \sum_{l=-\infty}^{\infty} \sum_{l^{\prime}=-\infty}^{\infty} N_{i l l^{\prime}} \bar{\phi}_{1 l l^{\prime}}^{+}\left(\vec{x}_{h}\right), \tag{3.230}
\end{equation*}
$$

with the normal derivative on the boundary scaled by $i k_{1}$ and where $\phi_{1 l l^{\prime}}^{+}$is from 3.218. The overbar is complex conjugation. Integrate the slope terms by parts as for example

$$
\begin{equation*}
<\phi_{1 j j^{\prime}}^{ \pm}, h_{x} \phi_{1 l l^{\prime}}^{+}>=\frac{\alpha_{j}-\alpha_{l}}{ \pm \gamma_{1 j j^{\prime}}-\bar{\gamma}_{1 l l l^{\prime}}}<\phi_{1 j j^{\prime}}^{ \pm}, \phi_{1 l l^{\prime}}^{+}> \tag{3.231}
\end{equation*}
$$

and the equations for the upper region can be written using (3.221) as

$$
\sum_{l=-\infty}^{\infty} \sum_{l^{\prime}=-\infty}^{\infty}<\phi_{1 j j^{\prime}}^{ \pm}, \phi_{1 l l^{\prime}}^{+}>\left[N_{i l l^{\prime}}-U^{ \pm}\left(j j^{\prime}, l l^{\prime}\right) E_{i l l l^{\prime}}\right]=2 \gamma_{1 j j^{\prime}}\left\{\begin{array}{l}
-I_{i j j^{\prime}} \tag{3.232}
\end{array}\right\},
$$

where

$$
\begin{equation*}
U^{ \pm}\left(j j^{\prime}, l l^{\prime}\right)=\frac{1-\alpha_{j} \alpha_{l}-\beta_{j^{\prime}} \beta_{l^{\prime}} \mp \gamma_{1 j j^{\prime}} \bar{\gamma}_{1 l l^{\prime}}}{ \pm \gamma_{1 j j^{\prime}}-\bar{\gamma}_{1 l l^{\prime}}} . \tag{3.233}
\end{equation*}
$$

We have written the double spectral values $j j^{\prime}$ and $l l^{\prime}$ as arguments of $U$ in illustration of the fact that they are each replacing coordinate sampling/integration along two-dimensional surfaces denoted by $\vec{x}_{h}^{\prime}$ and $\vec{x}_{h}$ respectively, as well as to indicate that the equations 3.232 ) are diagonal in the vector index " $i^{\prime \prime}$. That is, the $i^{t h}$ component of $A$ is related to the $i^{\text {th }}$ components of $N$ and $E$. There is no coupling in this index for the equations from region 1. It can be shown that the matrix $<\phi_{1 j j^{\prime}}^{+}, \phi_{1 l l^{\prime}}^{+}>$is self-adjoint, positive definite and hence invertible.

For the lower region 2 these same expansions and integration of the slope terms yields from (3.225)

$$
\begin{equation*}
\sum_{l=-\infty}^{\infty} \sum_{l^{\prime}=-\infty}^{\infty}<\phi_{2 j j^{\prime}}^{ \pm}, \phi_{1 l l^{\prime}}^{+}>\left[\mu N_{i l l^{\prime}}-L_{i p}^{ \pm}\left(j j^{\prime}, l l^{\prime}\right) E_{p l l^{\prime}}\right]=-2 \gamma_{2 j j^{\prime}}\left\{\left\{_{0}^{B_{i j j^{\prime}}}\right\},\right. \tag{3.234}
\end{equation*}
$$

where there is an implicit sum over the repeated subscript $p=(1,2,3)$. The full coupling of these equations resides in this summation. Here the $L$ term can be written as

$$
\begin{equation*}
L_{i p}^{ \pm}\left(j j^{\prime}, l l^{\prime}\right)=\frac{M_{i p}^{ \pm}\left(j j^{\prime}, l l^{\prime}\right)}{ \pm \gamma_{2 j j^{\prime}}-\bar{\gamma}_{1 l l^{\prime}}}, \tag{3.235}
\end{equation*}
$$

where $M$ can be written as a diagonal $(D)$ part and a full $(F)$ part, the latter of which contains the coupling,

$$
\begin{equation*}
M_{i p}^{ \pm}\left(j j^{\prime}, l l^{\prime}\right)=D^{ \pm}\left(j j^{\prime}, l l^{\prime}\right) \delta_{i p}+F_{i p}^{ \pm}\left(j j^{\prime}, l l^{\prime}\right), \tag{3.236}
\end{equation*}
$$

where

$$
\begin{equation*}
D^{ \pm}\left(j j^{\prime}, l l^{\prime}\right)=K^{2}-\alpha_{j} \alpha_{l}-\beta_{j^{\prime}} \beta_{l^{\prime}} \mp \gamma_{2 j j^{\prime}} \bar{\gamma}_{l l l^{\prime}}, \tag{3.237}
\end{equation*}
$$

and
$F_{i p}^{ \pm}\left(j j^{\prime}, l l^{\prime}\right)=\left(\varepsilon^{-1}-1\right)\left(\alpha_{j} \delta_{i 1}+\beta_{j^{\prime}} \delta_{i 2} \mp \gamma_{2 j j^{\prime}} \delta_{i 3}\right)\left[\left(\alpha_{j}-\alpha_{l}\right) \delta_{p 1}+\left(\beta_{j^{\prime}}-\beta_{l^{\prime}}\right) \delta_{p 2}-\left( \pm \gamma_{2 j j^{\prime}}-\bar{\gamma}_{1 l l^{\prime}}\right) \delta_{p 3}\right]$.
The procedure is to solve the upper equation (3.232) and the lower equation (3.234), which is a spectral extinction equation, for the unknown expansion coefficients $N_{i l l^{\prime}}$ and $E_{i l l^{\prime}}$, and evaluate the remaining equations for the scattered $A_{i j j^{\prime}}$ and transmitted $B_{i j j^{\prime}}$ spectral coefficients.

### 3.12 Summary

We have derived exact formal sets of equations, in both coordinate and various spectral domains, to describe the scattering from deterministic gratings. Both acoustic scalar one-dimensional problems and full electromagnetic two-dimensional problems were considered. Both involved a grating surface separating two homogeneous regions of space. Both involved coordinatespace representations from which proceeded rigorous plane wave spectral representations valid for the scattered field above the highest surface excursion and for the transmitted field below its lowest excursion. The electromagnetic development was treated in analogy with the scalar problem, with boundary conditions derived for the electric field and its normal derivative from the standard boundary conditions on currents and the normal components of the displacement vector.

From these coordinate representations we proceeded first to partial spectral representations where the word "partial" refers to the field variables. These could be derived in a straightforward way just using plane waves and Green's theorem, and without involving the Green's function explicitly. We stress again that the equations are exact. In addition, these led to surface inversion examples for the scalar case using perturbation theory (where the boundary values could be factored out), and the Kirchhoff approximation (where the boundary values were approximated).

The full spectral equations involved expanding the boundary unknowns in some set of functions, and it is here where the Rayleigh and Waterman assumptions come into play. For the scalar case we presented three expansions. The first was a physical optics modified Fourier expansion with a single plane wave modulating the surface fields. The second was what we referred to as a Floquet-Fourier basis which modulated the Fourier expansion by still preserving the Floquet-periodicity of the surface fields but without the full plane waves, and the third was an expansion in the conjugate Rayleigh basis where each term in the expansion could be thought of as modulated by a plane wave. For the electromagnetic case only the expansion in the conjugate Rayleigh basis was considered. Since we used a scalar analogy for the electromagnetic problem the resulting equations were formally analogous to the scalar equations with the additional complication being first a vector problem, and second the Bragg modal sampling in two twodimensional spaces, the spaces of boundary and field points.

We pointed out in the paper where any of these equations have been solved, but we repeat that the full computational results and the comparisons of different computational results for this problem require at least a separate paper if not a separate book.

## Appendix 3.A. A Note on Matrix Elements

For the fully spectral methods in Secs. 7 and 8 for the acoustic case and Sec. 11 for the electromagnetic case, the matrix elements have a characteristic form. In the one-dimensional case, after projection on the various basis sets considered in this paper, (see (3.158), (3.165), (3.170), and (3.174)), they have the general form

$$
\begin{equation*}
M(a, b)=\frac{1}{L} \int_{-L / 2}^{L / 2} \exp \left[i a x+i k_{1} b h(x)\right] d x \tag{3.239}
\end{equation*}
$$

For all the cases in question, $a=2 \pi\left(j^{\prime}-j\right) / L$ which is the Fourier part common to all, and $b=$ $\pm \gamma_{p j}-\gamma_{10}$ for the physical optics case with one overall plane wave, $b= \pm \gamma_{p j}$ for the FloquetFourier case with no plane waves modulating the field expansion, and $b= \pm \gamma_{p j}-\bar{\gamma}_{1 j^{\prime}}$ for the conjugate Rayleigh case with plane waves related to each Bragg mode in the sum. They have a general validity in surface scattering problems due to the presence of Green's functions or plane wave type expansions. For example, for a random surface, these functions $M$ were referred to as interaction functions in a Feynman diagram expansion ${ }^{[1142930}$, essentially a perturbation expansion in the functions.

In addition, for many surfaces, not necessarily analytic ones, the integral can be expressed in closed form in terms of special functions. For example, a cosine surface yields Bessel functions for $M$, a symmetric sawtooth function yields simple exponentials, a quadratic surface yields Fresnel integrals, a vortex-like surface involving a logarithm yields cosine integrals which can be evaluated in closed form (or, in a different form, confluent hypergeometric functions), a cycloid can be evaluated in terms of Bessel functions, a full-wave rectified surface in terms of a Bessel series, and a periodic array of semicircular cylinders (bosses) ${ }^{98}$ in terms of a Bessel series. These closed form solutions can be useful in computations or for approximations. The details can be found in ${ }^{32]}$. Two-dimensional integrals occurring in the electromagnetic problem can be developed in a similar way for egg-crate surfaces of the form $h(x, y)=h_{1}(x)+h_{2}(y)$.

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Chapter 4:

## Integral Method for Gratings

Daniel Maystre and Evgeny Popov

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# Chapter 4 

# Integral Method for Gratings 

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### 4.1. Introduction

Integral methods for scattering problems represent a class of mathematical methods based on integral equations. In this chapter, all the integral equations are deduced from boundary value problems of scattering and are classified as Fredholm integral equations. They can be written in the form:

$$
\begin{equation*}
\mathrm{cu}(\overrightarrow{\mathrm{r}})=\mathrm{v}(\overrightarrow{\mathrm{r}})+\int_{\mathbb{A}} \mathrm{W}\left(\overrightarrow{\mathrm{r}}, \overrightarrow{\mathrm{r}}^{\prime}\right) \mathrm{u}\left(\overrightarrow{\mathrm{r}}^{\prime}\right) \mathrm{d} \overrightarrow{\mathrm{r}}^{\prime}, \tag{4.1}
\end{equation*}
$$

in which $\mathbb{A}$ may represent for example the surface of a three-dimension diffracting object or the cross-section boundary of a two-dimension (cylindrical) object, v and u are continuous functions in $\mathbb{A}$. The kernel $W$ of the equation is also continuous in $\mathbb{A}$, and $c=0$ or 1 according to whether the equation is of the first or second kind. The mathematical problem is to determine $u$ if $v$ and W are known [1]. This kind of method is widely employed in many domains of physics [1,2], where it is usual to extend it to cases in which $u$, $v$, and $W$ are only peasewise continuous or can even be singular. A typical example of use in electromagnetism can be found by applying the second theorem of Green for expressing the field in a given volume in terms of its values and of the values of its normal derivative on the surrounding surface. In that case, the kernels of the integral equations include combinations of Green's functions and of their normal derivatives on the boundaries of the scattering objects.

The theory of integral equations can be described in a rigorous, elegant and concise way using distribution theory [3-5] that extends derivatives and other differential operators to discontinuous functions (a famous example is the so-called Dirac function that, for the mathematician, should not be called function). The interested reader can find a detailed presentation of rigorous use of the distribution theory in the electromagnetic theory of gratings in $[6,7]$.

The first application of the integral method in grating theory was proposed almost simultaneously for the case of perfectly conducting gratings by Petit and Cadilhac [8], Wirgin [9], and Uretski [10]. The first numerical implementation was reported by Petit [11,12] for TE polarization (called also P , or $\mathrm{E}_{/ /}$, or s polarization). The first extension of this integral equation to the other polarization (TM or $S$ or $p$ ) led to a non-integrable kernel. This problem was solved by Pavageau, who proposed for both polarizations new integral equations having continuous and bounded kernels [13].

Soon after, Wirgin [14], Neureuther and Zaki [15], and Van den Berg [16] gave formulations of the integral method applied to gratings made of metals with finite conductivity or dielectrics. The approach was based on the resolution of two coupled integral equations containing two unknown functions. Problems of limited memory storage and timecomputation on computers in the late ' 60 s restricted the numerical implementation of this theory to dielectric gratings only, for which very rare numerical results were published. This
restriction was not considered as dramatic by grating specialists at that time. Indeed, it was generally considered that in the visible and infrared regions in which the reflectivity of usual metals exceeds $90 \%$, the model of perfectly conducting grating is accurate. However, development of space optics and astronomy at that time required a precise treatement of the problem of metallic gratings used in the ultraviolet, a domain where the metal reflectivity starts to drop down as approaching electron plasma frequencies. This need required a new approach proposed by Maystre in 1972 [17] by using a single integral equation for a single unknown function. Solving the difficulties in the summation of the series in the kernels and in their integration, the integral method in this formulation was the first one to result in a computer code that was able to correctly model diffraction gratings behaviour over the entire spectrum for almost all commercial gratings profiles [18]. One of the most important conclusions for the practical applications and grating manufacturers was the definite demonstration of the inadequacy of the model of perfectly conducting grating in the nearinfrared, visible and ultraviolet regions [19].

However, the method was unable to treat some kinds of gratings, for example gratings having large periods and steep facets (echelle gratings, for example) or gratings with profile that cannot be represented by Fourier series (rod gratings, cavity gratings, etc.). A further development of this approach was proposed by Maystre [20]. It was numerically implemented in the early '90s in the code 'Grating 2000' by the author, which is used in many academic and industrial centers in the world. Other development was required for other exotic cases that started to find applications and thus required theoretical support for modeling. This development covered conical mountings and specially gratings with dielectric multilayer coatings [21,22], buried gratings and bimetallic gratings [23-26].

It must be emphasized that in this chapter, the authors use, without complete demonstrations, some analytic properties of gratings demonstrated in chapter 2. Thus, it is recommended to read this chapter before the present one.

First, we will deal with the most frequent problem: the bare metallic or dielectric grating. Then, extensions will be given to other kinds of gratings like perfectly conducting gratings, dielectric coated gratings or gratings in conical diffraction.

### 4.2. The integral method applied to a bare, metallic or dielectric grating.

### 4.2.1. The physical model

The grating surface S of period d separates a region $\mathrm{V}^{+}$with real relative electric permittivity and magnetic permeability $\varepsilon^{+}$and $\mu^{+}$respectively and a region $\mathrm{V}^{-}$with real or complex relative electric permittivity and magnetic permeability $\varepsilon^{-}$and $\mu^{-}$(figure 4.1). The indices $\mathrm{n}^{+}$and $\mathrm{n}^{-}$of these media are given by $\mathrm{n}^{+}=\sqrt{\varepsilon^{+} \mu^{+}}$and $\mathrm{n}^{-}=\sqrt{\varepsilon^{-} \mu^{-}}$. We consider the classical diffraction case with incident wavevector $\overrightarrow{\mathrm{k}}^{\mathrm{i}}$ lying in the xz plane, i.e. the plane perpendicular to the grooves. The incidence angle $\theta^{i}$ is measured in the counterclockwise sense from the z axis and $\lambda=\frac{2 \pi}{\mathrm{k}}$ denotes the wavelength of light in vacuum. The ordinate of the top of the profile is denoted by $\mathrm{z}_{0}$ and unit normal $\overrightarrow{\mathrm{N}}_{\mathrm{S}}$ is oriented towards $\mathrm{V}^{+}$. We denote by s the curvilinear abscissa on S , with origin being located at the origin of the Cartesian coordinates, and $s_{d}$ denoting the curvilinear abscissa of the point of $S$ of abscissa $x=d$.


Figure 4.1. Notations
In this chapter, we use the complex notation with a time-dependence in $\exp (-i \omega t)$. Let F denote the y-component of the electromagnetic field. In TE polarization, it stays for the ycomponent $\mathrm{E}_{\mathrm{y}}$ of the electric field, and in TM case for the y -component $\mathrm{H}_{\mathrm{y}}$ of the magnetic field.

### 4.2.2. The boundary value problem.

It is shown in chapter 2 that in that case of classical diffraction, the total field $\mathrm{F}^{\mathrm{T}}=\left\{\begin{array}{l}\mathrm{F}^{\mathrm{T}+} \text { in } \mathrm{V}^{+}, \\ \mathrm{F}^{\mathrm{T}-} \text { in } \mathrm{V}^{-}\end{array}\right.$is invariant along the y axis and that it is pseudo-periodic in x :

$$
\begin{equation*}
\mathrm{F}^{\mathrm{T}}(\mathrm{x}+\mathrm{d}, \mathrm{z})=\mathrm{F}^{\mathrm{T}}(\mathrm{x}, \mathrm{z}) \exp \left(\mathrm{i} \alpha_{0} \mathrm{~d}\right) \tag{4.2}
\end{equation*}
$$

with:

$$
\begin{equation*}
\alpha_{0}=\mathrm{kn}^{+} \sin \theta^{i}, \quad \mathrm{k}=\frac{2 \pi}{\lambda} . \tag{4.3}
\end{equation*}
$$

Moreover, the scattered field is defined by:

$$
\mathrm{F}=\left\{\begin{array}{c}
\mathrm{F}^{+}=\mathrm{F}^{\mathrm{T}+}-\mathrm{F}^{\mathrm{i}} \quad \text { in } \mathrm{V}^{+},  \tag{4.4}\\
\mathrm{F}^{-}=\mathrm{F}^{\mathrm{T}-} \quad \text { in } \mathrm{V}^{-},
\end{array}\right.
$$

with the incident field $\mathrm{F}^{\mathrm{i}}$ in $\mathrm{V}^{+}$being given by:

$$
\begin{equation*}
F^{\mathrm{i}}(\mathrm{x}, \mathrm{z})=\mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}-\mathrm{ikn} n^{+} \cos \theta^{\mathrm{i}} \mathrm{z}} \tag{4.5}
\end{equation*}
$$

The interest of the notion of scattered field is that it satisfies a radiation condition (also called Sommerfeld condition, or outgoing wave condition, see chapter 2) for $z \rightarrow \pm \infty$. The radiation
condition states that the scattered field at infinity must remain bounded and must propagate upward in $\mathrm{V}^{+}$and dawnward in $\mathrm{V}^{-}$. The scattered field also satisfies Helmholtz equations:

$$
\begin{equation*}
\nabla^{2} \mathrm{~F}^{ \pm}+\mathrm{k}^{2}\left(\mathrm{n}^{ \pm}\right)^{2} \mathrm{~F}^{ \pm}=0 \quad \text { in } \mathrm{V}^{ \pm} \tag{4.6}
\end{equation*}
$$

The invariance along the y axis allows one to reduce the scattering problem to a twodimension problem while the pseudo-periodicity restricts the study of the field to a single period of the grating. Consequently, it can be considered that $\mathrm{V}^{ \pm}$are no more volumes but surfaces extending on a single period of the grating.

In order to periodize the scattered field, we introduce the function $U=\left\{\begin{array}{l}\mathrm{U}^{+} \text {in } \mathrm{V}^{+}, \\ \mathrm{U}^{-} \text {in } \mathrm{V}^{-},\end{array}\right.$:

$$
\begin{equation*}
\mathrm{U}(\mathrm{x}, \mathrm{z})=\mathrm{e}^{-\mathrm{i} \alpha_{0} \mathrm{x}} \mathrm{~F}(\mathrm{x}, \mathrm{z}) \tag{4.7}
\end{equation*}
$$

We denote by $\psi^{ \pm}(\mathrm{s})$ the limit values of $\mathrm{U}^{ \pm}$on S and $\phi^{ \pm}(\mathrm{s})$ the values of $\mathrm{e}^{-\mathrm{i} \alpha_{0} \mathrm{x}} \frac{\mathrm{dF}}{\mathrm{F}^{ \pm}} \mathrm{dN}_{\mathrm{S}}$, with $\frac{\mathrm{dF}^{ \pm}}{\mathrm{dN}_{\mathrm{S}}}$ being the normal derivative of $\mathrm{F}^{ \pm}$on S .

The Helmholtz equations and the radiation condition are not sufficient to define the boundary value problem satisfied by the scattered field F. A third kind of condition must be added: the boundary conditions of the electromagnetic field components across S . The tangential components of the electric and magnetic fields are continuous across this interface, as far as the permittivities and permeabilities of the two media take finite values. For both polarizations, this property yields:

$$
\begin{equation*}
\psi^{+}(\mathrm{s})+\psi^{i}(\mathrm{~s})=\psi^{-}(\mathrm{s}), \tag{4.8}
\end{equation*}
$$

with $\psi_{\mathrm{i}}$ being the value of the periodized incident field, obtained from equation (4.5):

$$
\begin{equation*}
\psi^{\mathrm{i}}=\mathrm{F}^{\mathrm{i}}[\mathrm{x}(\mathrm{~s}), \mathrm{z}(\mathrm{~s})] \exp \left(-\mathrm{i} \alpha_{0} \mathrm{x}(\mathrm{~s})\right)=\exp \left[-\mathrm{ikn}{ }^{+} \cos ^{i} \mathrm{z}(\mathrm{~s})\right] . \tag{4.9}
\end{equation*}
$$

Using Maxwell equations, the continuity of the tangential component of the magnetic field (for TE polarization) and that of the electric field (for TM polarization) leads to the following relation:

$$
\begin{equation*}
\mathrm{q}^{+}\left[\phi^{+}(\mathrm{s})+\phi^{\mathrm{i}}(\mathrm{~s})\right]=\mathrm{q}^{-} \phi^{-}(\mathrm{s}), \tag{4.10}
\end{equation*}
$$

with

$$
\begin{align*}
& \phi^{\mathrm{i}}=\exp \left(-\mathrm{i} \alpha_{0} \mathrm{x}\right) \frac{\partial \mathrm{F}^{\mathrm{i}}}{\partial \mathrm{~N}_{\mathrm{S}}}= \\
& =-\mathrm{ikn}+\left(\frac{\mathrm{dz}(\mathrm{~s})}{\mathrm{ds}} \sin \theta^{i}+\frac{\mathrm{dx}(\mathrm{~s})}{\mathrm{ds}} \cos \theta^{i}\right) \exp \left(-\mathrm{ikn}{ }^{+} \cos \theta^{i} \mathrm{z}(\mathrm{~s})\right), \tag{4.11}
\end{align*}
$$

and

$$
\mathrm{q}^{ \pm}= \begin{cases}\frac{1}{\mu^{ \pm}}, & \text {for TE polarization, }  \tag{4.12}\\ \frac{1}{\varepsilon^{ \pm}}, & \text {for TM polarization. }\end{cases}
$$

### 4.2.3. Integral equation

The theoretical basis of the integral method lies on a general property of the electromagnetic field: the field inside a given surface of the xz plane can be expressed from the values of the field and of its normal derivatives on the curve surrounding the surface, according to the second Green's theorem. The value of $\mathrm{U}^{ \pm}$at a point of $\mathrm{V}^{ \pm}$of coordinates $x$ and $z$ be deduced from its values on $S$ using equation (4.139) of appendix 4.A:

$$
\begin{equation*}
\mathrm{U}^{ \pm}(\mathrm{x}, \mathrm{z})= \pm \int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{ \pm}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \phi^{ \pm}\left(\mathrm{s}^{\prime}\right)+\boldsymbol{N}^{ \pm}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \psi^{ \pm}\left(\mathrm{s}^{\prime}\right)\right] \mathrm{ds}^{\prime} \tag{4.13}
\end{equation*}
$$

with

$$
\begin{gather*}
\mathcal{G}^{ \pm}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{id}} \sum_{\mathrm{m}=-\infty}^{\infty} \frac{1}{\gamma_{\mathrm{m}}^{ \pm}} \exp \left\{\operatorname{imK}\left[\mathrm{x}-\mathrm{x}^{\prime}\left(\mathrm{s}^{\prime}\right)\right]+\mathrm{i} \gamma_{\mathrm{m}}^{ \pm} \mid \mathrm{z}-\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right\},  \tag{4.14}\\
\boldsymbol{N}^{ \pm}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{~d}} \sum_{\mathrm{m}=-\infty}^{\infty}\left\{\frac{\mathrm{dx}\left(\mathrm{~s}^{\prime}\right)}{\mathrm{ds} s^{\prime}} \operatorname{sgn}\left[\mathrm{z}-\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right]-\frac{\alpha_{\mathrm{m}}}{\gamma_{\mathrm{m}}^{ \pm}} \frac{\mathrm{dz} \mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)}{\mathrm{ds} s^{\prime}}\right\} \times  \tag{4.15}\\
\times \exp \left\{\mathrm{imK}\left[\mathrm{x}-\mathrm{x}^{\prime}\left(\mathrm{s}^{\prime}\right)\right]+\mathrm{i} \gamma_{\mathrm{m}}^{ \pm} \mid \mathrm{z}-\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right\}
\end{gather*}
$$

where:

$$
\begin{gather*}
\alpha_{\mathrm{m}}=\alpha_{0}+\mathrm{mK},  \tag{4.16}\\
\mathrm{~K}=\frac{2 \pi}{\mathrm{~d}},  \tag{4.17}\\
\gamma_{\mathrm{m}}^{ \pm}=\sqrt{\left(\mathrm{kn}^{ \pm}\right)^{2}-\alpha_{\mathrm{m}}^{2}}, \tag{4.18}
\end{gather*}
$$

with $\mathrm{s}^{\prime}$ being the curvilinear abscissa on a point of S with coordinates $\mathrm{x}^{\prime}\left(\mathrm{s}^{\prime}\right)$ and $\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)$.
According to section 4.A.4, the values of $\psi^{ \pm}\left(s^{\prime}\right)$ and $\phi^{ \pm}\left(s^{\prime}\right)$ are linked by a relation of compatibility. Using eqs. (4.146), (4.147) and (4.148) we obtain:

$$
\begin{align*}
& \int_{s^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \tilde{\phi}^{+}\left(\mathrm{s}^{\prime}\right)+\mathcal{N}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \tilde{\Psi}^{+}\left(\mathrm{s}^{\prime}\right)\right] \mathrm{ds} \mathrm{~s}^{\prime}-\frac{\tilde{\psi}^{+}\left(\mathrm{s}^{\prime}\right)}{2}=0  \tag{4.19}\\
& \int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{-}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \tilde{\phi}^{-}\left(\mathrm{s}^{\prime}\right)+\mathcal{N}^{-}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \tilde{\Psi}^{-}\left(\mathrm{s}^{\prime}\right)\right] \mathrm{ds} \mathrm{~s}^{\prime}+\frac{\tilde{\psi}^{-}\left(\mathrm{s}^{\prime}\right)}{2}=0 \tag{4.20}
\end{align*}
$$

with:

$$
\begin{align*}
& \boldsymbol{\mathcal { G }}^{ \pm}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{id}} \sum_{\mathrm{m}=-\infty}^{\infty} \frac{1}{\gamma_{\mathrm{m}}^{ \pm}} \exp \left[\operatorname{imK}\left(\mathrm{x}(\mathrm{~s})-\mathrm{x}^{\prime}\left(\mathrm{s}^{\prime}\right)\right)+\mathrm{i} \gamma_{\mathrm{m}}^{ \pm}\left|\mathrm{z}(\mathrm{~s})-\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right|\right]  \tag{4.21}\\
& \quad \boldsymbol{N}^{ \pm}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{~d}} \sum_{\mathrm{m}=-\infty}^{\infty}\left[\frac{\mathrm{dx}{ }^{\prime}}{\mathrm{ds} \mathrm{~s}^{\prime}} \operatorname{sgn}\left(\mathrm{z}(\mathrm{~s})-\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right)-\frac{\alpha_{\mathrm{m}}}{\gamma_{\mathrm{m}}^{ \pm}} \frac{\mathrm{dz}}{\mathrm{ds}}\right]  \tag{4.22}\\
& \quad \times \exp \left[\operatorname{imK}\left(\mathrm{x}(\mathrm{~s})-\mathrm{x}^{\prime}\left(\mathrm{s}^{\prime}\right)\right)+\mathrm{i} \gamma_{\mathrm{m}}^{ \pm}\left|\mathrm{z}(\mathrm{~s})-\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right|\right]
\end{align*}
$$

Introducing in eq. (4.20) the values of $\psi^{-}(\mathrm{s})=\psi^{+}(\mathrm{s})+\psi_{\mathrm{i}}(\mathrm{s})$ and $\phi^{-}(\mathrm{s})=\frac{\mathrm{q}^{+}}{\mathrm{q}^{-}}\left[\phi^{+}(\mathrm{s})+\phi_{\mathrm{i}}(\mathrm{s})\right]$ given by the continuity conditions on the grating profile (eqs (4.8), (4.9), (4.10), (4.11) and (4.12)) yields a second integral equations with two unknown functions $\psi^{+}$and $\phi^{+}$:

$$
\begin{gather*}
\int_{s^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left\{\frac{\mathrm{q}^{+}}{\mathrm{q}^{-}} \boldsymbol{\mathcal { G }}^{-}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)\left[\phi^{+}\left(\mathrm{s}^{\prime}\right)+\phi^{\mathrm{i}}\left(\mathrm{~s}^{\prime}\right)\right]+\mathcal{N}^{-}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)\left[\psi^{+}\left(\mathrm{s}^{\prime}\right)+\psi^{\mathrm{i}}\left(\mathrm{~s}^{\prime}\right)\right]\right\} \mathrm{ds} \mathrm{~s}^{\prime}  \tag{4.23}\\
+\frac{\psi^{+}\left(\mathrm{s}^{\prime}\right)+\psi^{\mathrm{i}}\left(\mathrm{~s}^{\prime}\right)}{2}=0
\end{gather*}
$$

Eqs (4.19) and (4.23) constitute a system of two integral equations with two unknown functions, which can be solved on a computer. The amplitudes $r_{m}$ and $t_{m}$ of the plane waves reflected and transmitted by the grating can be deduced from the solution of the integral equation using eqs. (4.184) and (4.185) of appendix 4.A:

$$
\begin{align*}
& r_{m}=\frac{1}{2 d} \int_{s=0}^{s_{d}} e^{-i m K x(s)-i \gamma_{m}^{+} z(s)}\left[\frac{-i \phi^{+}(s)}{\gamma_{m}^{+}}+\left(\frac{d x(s)}{d s}-\frac{\alpha_{m}}{\gamma_{m}^{+}} \frac{d z(s)}{d s}\right) \psi^{+}(s)\right] d s,  \tag{4.24}\\
& t_{m}=\frac{1}{2 d} \int_{s=0}^{s_{d}} e^{-i m K x(s)+i \gamma_{m}^{-} z(s)}\left[\frac{i \phi^{-}(s)}{\gamma_{m}^{-}}+\left(\frac{d x(s)}{d s}+\frac{\alpha_{m}}{\gamma_{m}^{+}} \frac{d z(s)}{d s}\right) \psi^{-}(s)\right] d s, \tag{4.25}
\end{align*}
$$

with $z_{0}$ being the ordinate of the top of the grating profile. For non-evenescent reflected orders, diffraction efficiencies $\rho_{\mathrm{m}}$ can be obtained using eq. (4.187):

$$
\begin{equation*}
\rho_{\mathrm{m}}=\frac{\gamma_{\mathrm{m}}^{+}}{\gamma_{0}^{+}}\left|\mathrm{r}_{\mathrm{m}}\right|^{2} \tag{4.26}
\end{equation*}
$$

For gratings made of a lossless dielectric material in $\mathrm{V}^{-}$, transmitted efficiencies can be defined as well:

$$
\begin{equation*}
\tau_{\mathrm{m}}=\frac{\mathrm{q}^{-}}{\mathrm{q}^{+}} \frac{\gamma_{\mathrm{m}}^{+}}{\gamma_{0}^{+}}\left|\mathrm{t}_{\mathrm{m}}\right|^{2} \tag{4.27}
\end{equation*}
$$

In that case the energy balance (see chapter 2) can be expressed by:

$$
\begin{equation*}
\sum_{\mathrm{m} \in \mathrm{P}^{+}} \rho_{\mathrm{m}}+\sum_{\mathrm{m} \in \mathrm{P}^{-}} \tau_{\mathrm{m}}=1 \tag{4.28}
\end{equation*}
$$

with $\mathrm{P}^{+}$and $\mathrm{P}^{-}$denoting respectively the set of non-evanescent reflected and transmitted orders. The numerical implementation of the integral equations will be described in section 6 . In contrast with the two coupled equations obtained in this section, the integral equation obtained by Maystre is unique. This feature requires the definition of a single and well adapted unknown function. The mathematical definition of this function needs the use of tools of applied mathematics described in appendix A. Appendix B contains a mathematical description of this mathematical function and of the integral equation. Here, we give a heuristic description of this function for TE polarization. First, we replace the material in $\mathrm{V}^{-}$ by the same material as in $\mathrm{V}^{+}$, the entire space being thus homogeneous. It can be shown that it exists one (and only one) distribution of surface current density $\Phi$ parallel to the y axis, placed on S (this surface separates now two identical media), which generates in $\mathrm{V}^{+}$a field equal to the actual diffracted field in the physical problem. Intuitively, it is easy, from this surface current density, to express in an integral form the actual scattered field in $\mathrm{V}^{+}$, this current distribution being nothing else than a set of current lines placed in a homogeneous medium. From the expression of the scattered field in $\mathrm{V}^{+}$, simple mathematical calculations allow one to deduce the scattered field and its normal derivative above S , thus $\Psi^{+}(\mathrm{s})$ and $\phi^{+}(\mathrm{s})$ from the unique unknown function $\Phi$.

Now, we abandon the field generated by the fictitious surface current density $\Phi$, except the integral expressions of $\psi^{+}(\mathrm{s})$ and $\phi^{+}(\mathrm{s})$ containing $\Phi$, and we come back to the actual physical grating problem. The continuity conditions for the tangential components of the field permit the calculation of $\psi^{-}(\mathrm{s})$ and $\phi^{-}$(s) thus, using the second Green theorem, of the actual physical field below $S$. At that point it has been shown that the four unknown functions contained in the classical theory previously described in this section can be derived from a single one and that, in some way, there is a redundancy in the use of multiple unknown functions. It is easy to understand that this single unknown function can be calculated from a single integral equation. This equation can be obtained for example by writing the continuity on S of the integral expressions of the field in $\mathrm{V}^{+}$and $\mathrm{V}^{-}$.

This method was the first one to show that, in contrast with the second Green theorem, it exists a formula that allows one to express the field inside a given domain from a single function defined on its boundary. This function is neither the field nor its normal derivative, but both can be deduced from it through simple integrals. These integrals automatically satisfy the compatibility condition.

### 4.3. The bare, perfectly conducting grating

Perfectly conducting gratings were historically the first gratings to be modeled using rigorous electromagnetic theories. They represent accurate models for metallic gratings working in far infrared and microwaves regions. The pioneering works appear in the ' 60 s [8] and were followed by many papers. Various formulations of the integral method have been published. They differ either in the form of the integral equation ot in the numerical implementation. A review of this matter may be found in [27, 6, 28].

### 4.3.1. Perfectly conducting gratings in TE polarization

Two different approaches will be presented in this section. The first one, published by R. Petit and M. Cadilhac in [8], leads to a Fredholm integral equation of the fist kind with a singular
kernel. A version leading to a Fredholm integral equation of the second kind with a nonsingular kernel was proposed by Pavageau et al. [13] using the ideas of Maue [29].

The total field in $\mathrm{V}^{+}$is pseudo-periodic, it satisfies the Helmholtz equation:

$$
\begin{equation*}
\nabla^{2} \mathrm{~F}^{\mathrm{T}+}+\mathrm{k}^{2}\left(\mathrm{n}^{+}\right)^{2} \mathrm{~F}^{\mathrm{T}+}=0 \quad \text { in } \mathrm{V}^{+} \tag{4.29}
\end{equation*}
$$

and a radiation condition at infinity. Thus we can apply the generalized compatibility condition of section 4.A.5:

$$
\begin{equation*}
\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \Phi^{+}\left(\mathrm{s}^{\prime}\right)+\mathcal{N}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \Psi^{+}\left(\mathrm{s}^{\prime}\right)\right] \mathrm{ds} \mathrm{~s}^{\prime}+\psi^{\mathrm{i}}(\mathrm{~s})=\frac{\Psi^{+}(\mathrm{s})}{2} \tag{4.30}
\end{equation*}
$$

with $\Psi^{+}\left(\mathrm{s}^{\prime}\right)$ and $\Phi^{+}\left(\mathrm{s}^{\prime}\right)$ denoting the limit of the total field on S and its normal derivative. The boundary condition on S is straightforward: the total electric field, which is parallel to the y axis thus tangential to the metal, vanishes on S. This property entails that $\Psi^{+}(\mathrm{s})=0$ and thus, eq. (4.30) becomes, in operator notation:

$$
\begin{equation*}
\mathcal{G}^{+} \Phi^{+}=-\psi^{\mathrm{i}} \tag{4.31}
\end{equation*}
$$

with $\boldsymbol{\mathcal { G }}^{+}(\mathrm{s}, \mathrm{s}$ ') given by eq. (4.147)

$$
\begin{equation*}
\boldsymbol{\mathcal { G }}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{id}} \sum_{\mathrm{m}=-\infty}^{\infty} \frac{1}{\gamma_{\mathrm{m}}^{+}} \exp \left[\operatorname{imK}\left(\mathrm{x}(\mathrm{~s})-\mathrm{x}^{\prime}\left(\mathrm{s}^{\prime}\right)\right)+\mathrm{i} \gamma_{\mathrm{m}}^{+}\left|\mathrm{z}(\mathrm{~s})-\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right|\right] \tag{4.32}
\end{equation*}
$$

This is a Fredholm integral equation of the first kind, with a singular kernel.
The amplitudes of the reflected waves are deduced from eq. (4.186):

$$
\begin{equation*}
\mathrm{r}_{\mathrm{m}}=\frac{1}{2 \mathrm{id} \gamma_{\mathrm{m}}^{+}} \int_{\mathrm{s}=0}^{\mathrm{s}_{\mathrm{d}}} \exp \left[-\mathrm{imKx}(\mathrm{~s})-\mathrm{i} \gamma_{\mathrm{m}}^{+} \mathrm{z}(\mathrm{~s})\right] \Phi^{+}(\mathrm{s}) \mathrm{ds} \tag{4.33}
\end{equation*}
$$

The efficiencies $\rho_{\mathrm{m}}=\frac{\gamma_{m}^{+}}{\gamma_{0}^{+}}\left|\mathrm{r}_{\mathrm{m}}\right|^{2}$ satisfy the energy balance relation:

$$
\begin{equation*}
\sum_{\mathrm{P}^{+}} \rho_{\mathrm{m}}=1, \tag{4.34}
\end{equation*}
$$

with $\mathrm{P}^{+}$denoting the set of non-evanescent orders.
An integral equation of the second kind with a regular continuous kernel can be found using the same function $\Phi^{+}$. It is shown in section 4.A5 that the normal derivative $\frac{\mathrm{dF}^{+}}{\mathrm{dN}_{\mathrm{S}}}$ can be calculated in that case (eq. (4.172)). This integral equation can be written either by writing that $\frac{\mathrm{dF}^{+}}{\mathrm{dN}_{\mathrm{S}}}=\phi^{+} \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}}$ is equal to $\left(\Phi^{+}-\phi^{\mathrm{i}}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}}$. The final equation is given by:

$$
\begin{equation*}
\frac{\Phi^{+}(\mathrm{s})}{2}=\phi^{\mathrm{i}}(\mathrm{~s})+\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \mathscr{K}\left(\mathrm{~s}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}} \Phi^{+}(\mathrm{s}) \mathrm{ds}^{\prime}, \tag{4.35}
\end{equation*}
$$

with:

$$
\begin{equation*}
\mathscr{K}\left(s, s^{\prime}\right)=\frac{1}{2 d} \sum_{m=-\infty,+\infty}\left[\operatorname{sgn}\left(z-z^{\prime}\right) \frac{d x}{d s}-\frac{\alpha_{m}}{\gamma_{m}^{+}} \frac{d z}{d s}\right] e^{\operatorname{imK}\left(x-x^{\prime}\right)+i \gamma_{m}^{+}\left|z-z^{\prime}\right|} \tag{4.36}
\end{equation*}
$$

### 4.3.2. Perfectly conducting gratings for TM polarization

Once again, the generalized compatibility condition is used (eq. (4.158):

$$
\begin{equation*}
\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \Phi^{+}\left(\mathrm{s}^{\prime}\right)+\mathcal{N}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \Psi^{+}\left(\mathrm{s}^{\prime}\right)\right] \mathrm{ds} \mathrm{~s}^{\prime}+\psi^{\mathrm{i}}(\mathrm{~s})=\frac{\Psi^{+}(\mathrm{s})}{2} \tag{4.37}
\end{equation*}
$$

In that case too, the tangential component of the electric field vanishes on the profile. It is shown in chapter 2 from Maxwell equations that this condition entails that the normal derivative of the total field vanishes on S , thus:

$$
\begin{equation*}
\Phi^{+}=0 \tag{4.38}
\end{equation*}
$$

so that

$$
\begin{equation*}
\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \mathcal{N}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \Psi^{+}\left(\mathrm{s}^{\prime}\right) \mathrm{ds} \mathrm{~s}^{\prime}+\psi^{\mathrm{i}}(\mathrm{~s})=\frac{\Psi^{+}(\mathrm{s})}{2} \tag{4.39}
\end{equation*}
$$

with $\boldsymbol{\mathcal { N }}^{+}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$ given by eq. (4.22) and (4.148):

$$
\begin{align*}
& \mathcal{N}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{~d}} \sum_{\mathrm{m}=-\infty}^{\infty}\left[\frac{\mathrm{dx} \mathrm{x}^{\prime}}{\mathrm{ds} s^{\prime}} \operatorname{sgn}\left(\mathrm{z}(\mathrm{~s})-\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right)-\frac{\alpha_{\mathrm{m}}}{\gamma_{\mathrm{m}}^{ \pm}} \frac{\mathrm{dz} \mathrm{z}^{\prime}}{\mathrm{ds}}\right] \times  \tag{4.40}\\
& \times \exp \left[\operatorname{imK}\left(\mathrm{x}(\mathrm{~s})-\mathrm{x}^{\prime}\left(\mathrm{s}^{\prime}\right)\right)+\mathrm{i} \gamma_{\mathrm{m}}^{+}\left|\mathrm{z}(\mathrm{~s})-\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right|\right] .
\end{align*}
$$

The Fredholm integral equation of the second kind with a regular continuous kernel is very close to that obtained for TE polarization (eq. (4.35)).

The amplitudes of the reflected waves are deduced from eq. (4.186):

$$
\begin{equation*}
\mathrm{r}_{\mathrm{m}}=\frac{1}{2 \mathrm{~d}} \int_{\mathrm{s}=0}^{\mathrm{s}_{\mathrm{d}}} \exp \left[-\mathrm{imKx}(\mathrm{~s})-\mathrm{i} \gamma_{\mathrm{m}}^{+} \mathrm{z}(\mathrm{~s})\right]\left(\frac{\mathrm{dx}(\mathrm{~s})}{\mathrm{ds}}-\frac{\alpha_{\mathrm{m}}}{\gamma_{\mathrm{m}}^{+}} \frac{\mathrm{dz}(\mathrm{~s})}{\mathrm{ds}}\right) \Psi(\mathrm{s}) \mathrm{ds} \tag{4.41}
\end{equation*}
$$

As for TE polarization, the efficiencies $\rho_{\mathrm{m}}=\frac{\gamma_{m}^{+}}{\gamma_{0}^{+}}\left|\mathrm{r}_{\mathrm{m}}\right|^{2}$ satisfy the energy balance relation:

$$
\begin{equation*}
\sum_{\mathrm{P}^{+}} \rho_{\mathrm{m}}=1 \tag{4.42}
\end{equation*}
$$

### 4.4. Multiprofile gratings

The use of dielectric coatings has many applications even for diffraction gratings use. For example, metallic gratings are covered by a thin layer of dielectric material in order to avoid oxidation of the metal. Dielectric gratings can require an antireflection coating consisting of a thin layer or a stack of layers. Conversely, a stack of layers is used to increase the metal
reflectivity or even to replace it, in order to reduce the absorption of light beams for high power laser applications.

Up to our knowledge, the first numerical results in the study of coated gratings was made by Van de Berg [16] for a single-layer perfectly conducting grating with the layer filling up the space between the grating surface and a plane surface. Although at that time the interest in such geometry remained mostly academic, further development of technology made it possible to fabricate such layer by dielectric coating and polishing. Another important application comes from the process of replication of dielectric gratings using an epoxy layer to transfer the replica to a plane surface of the substrate, or to have epoxy as grating layer itself.

After this initial work, two integral methods were proposed. The first one [20] is theoretically able to deal with an arbitrary multilayer grating without limitations concerning the shape of the profile or the conductivity of the layers. The second method [21] can deal with a multilayer grating without interpenetration of the profiles.Botten has solved the problem with a single-profile grating that has a stack of plane layers below and, eventually, above it [22, 23], by introducing a new form of Green's function, adapted to the multilayer system, which leeds to a single integral equation.

In what follows, we will at first describe the method for a single interface inside a stack, when the layer is relatively thin so that the upper and the lower interface interpenetrate. It is important to distinguish the two cases, with and without interpenetration, because in the latter case, it is possible to define a plane layer in between that does not cross the upper or the lower interface. This possibility enables one to use the plane-wave Rayleigh expansion of the electromagnetic field between the interfaces, whereas in the former case it is necessary to write and to solve a system of coupled integral equations that link the field components on the top and bottom of each layer.

### 4.4.1. Thin-layer gratings



Figure 4.2. Single layer inside a stack of a multilayer grating
Let us consider the case of a multilayer grating having profiles that interpenetrate. In other words, it is impossible to introduce inside the layer a plane surface that does not cross one of the profiles. Then it is impssible to use the plane-wave expansion between the interfaces and we are led to solve integral equations that are coupled on the two interfaces of each layer.

We introduce in figure 4.2 a grating made of M materials (numbered from 0 to M , separated by $\mathrm{M}-1$ profiles (numbered from 1 to M ). We introduce the following functions:

$$
\mathrm{U}_{\mathrm{j}}=\left\{\begin{array}{c}
\mathrm{F}_{\mathrm{j}} \mathrm{e}^{-\mathrm{i} \alpha_{0} \mathrm{x}}-\mathrm{F}^{\mathrm{i}} \mathrm{e}^{-\mathrm{i} \alpha_{0} \mathrm{x}} \delta_{\mathrm{j}, \mathrm{M}} \quad \text { in } \mathrm{V}_{\mathrm{j}} .  \tag{4.43}\\
0 \quad \text { elsewhere. }
\end{array}\right.
$$

The function $U_{j}$ inside each layer can be expressed from the fields and normal derivatives on the lower and upper interface, for $\mathrm{j}=1, \mathrm{M}-1$ :

$$
\begin{align*}
\mathrm{U}_{\mathrm{j}}(\mathrm{x}, \mathrm{y}) & =\int_{\mathrm{S}_{\mathrm{j}}} \mathcal{G}_{\mathrm{j}}^{+}\left(\mathrm{x}, \mathrm{y}, \mathrm{~s}_{\mathrm{j}}^{\prime}\right) \phi_{\mathrm{j}}^{+}\left(\mathrm{s}_{\mathrm{j}}^{\prime}\right) \mathrm{ds} \mathrm{~s}_{\mathrm{j}}{ }^{\prime}+\int_{\mathrm{S}_{\mathrm{j}}} \mathcal{N}_{\mathrm{j}}^{+}\left(\mathrm{x}, \mathrm{y}, \mathrm{~s}_{\mathrm{j}}^{\prime}\right) \psi_{\mathrm{j}}^{+}\left(\mathrm{s}_{\mathrm{j}}^{\prime}\right) \mathrm{ds} \mathrm{~s}_{\mathrm{j}}^{\prime} \\
& -\int_{\mathrm{S}_{\mathrm{j}+1}}^{\boldsymbol{\mathcal { G }}_{\mathrm{j}+1}^{-}\left(\mathrm{x}, \mathrm{y}, \mathrm{~s}_{\mathrm{j}+1}^{\prime}\right) \phi_{\mathrm{j}+1}^{-}\left(\mathrm{s}_{\mathrm{j}+1}^{\prime}\right) \mathrm{ds}_{\mathrm{j}+1}^{\prime}-\int_{\mathrm{S}_{\mathrm{j}+1}} \boldsymbol{\mathcal { N }}_{\mathrm{j}+1}^{-}\left(\mathrm{x}, \mathrm{y}, \mathrm{~s}_{\mathrm{j}+1}^{\prime}\right) \psi_{\mathrm{j}+1}^{-}\left(\mathrm{s}_{\mathrm{j}+1}^{\prime}\right) \mathrm{ds}_{\mathrm{j}+1}^{\prime} .} \tag{4.44}
\end{align*}
$$

The expression being limited to the second one if $\mathrm{j}=0$ and to the first one for $\mathrm{j}=\mathrm{M}$.
The functions derived from the Green functions in the various materials depend on the interface number:

$$
\begin{align*}
& \boldsymbol{\mathcal { G }}_{\mathrm{j}}^{ \pm}\left(\mathrm{x}, \mathrm{y}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{id}} \sum_{\mathrm{m}=-\infty}^{\infty} \frac{1}{\gamma_{\mathrm{j}, \mathrm{~m}}^{ \pm}} \exp \left\{\operatorname{imK}\left[\mathrm{x}-\mathrm{x}^{\prime}\left(\mathrm{s}_{\mathrm{j}}^{\prime}\right)\right]+\mathrm{i} \gamma_{\mathrm{j}, \mathrm{~m}}^{ \pm} \mid \mathrm{z}-\mathrm{z}^{\prime}\left(\mathrm{s}_{\mathrm{j}}^{\prime}\right)\right\},  \tag{4.45}\\
& \mathcal{N}_{\mathrm{j}}^{ \pm}\left(\mathrm{x}, \mathrm{y}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{~d}} \sum_{\mathrm{m}=-\infty}^{\infty}\left\{\frac{\mathrm{dx} \mathrm{x}^{\prime}\left(\mathrm{s}^{\prime}\right)}{\mathrm{ds} \mathrm{~s}^{\prime}} \operatorname{sgn}\left[\mathrm{z}-\mathrm{z}^{\prime}\left(\mathrm{s}_{\mathrm{j}}^{\prime}\right)\right]-\frac{\alpha_{\mathrm{m}}}{\gamma_{\mathrm{j}, \mathrm{~m}}^{ \pm}} \frac{\mathrm{dz} \mathrm{z}^{\prime}\left(\mathrm{s}_{\mathrm{j}}^{\prime}\right)}{\mathrm{ds}{ }^{\prime}}\right\} \times  \tag{4.46}\\
& \times \exp \left\{\operatorname{imK}\left[x-x^{\prime}\left(s_{j}^{\prime}\right)\right]+i \gamma_{\mathrm{j}, \mathrm{~m}}^{ \pm}\left|z-\mathrm{z}^{\prime}\left(\mathrm{s}_{\mathrm{j}}^{\prime}\right)\right|\right\},
\end{align*}
$$

with:

$$
\begin{equation*}
\gamma_{\mathrm{j}, \mathrm{~m}}^{+}=\gamma_{\mathrm{j}+1, \mathrm{~m}}^{-}=\sqrt{\left(\mathrm{kn}_{\mathrm{j}}\right)^{2}-\alpha_{\mathrm{m}}^{2}} \tag{4.47}
\end{equation*}
$$

It can be shown it the same manner as in eq. (4.146) that a compatibility condition on the $\mathrm{j}^{\text {th }}$ interfacewritten in a matrix form is given by:

$$
\begin{equation*}
\frac{\psi_{\mathrm{j}}^{+}}{2}=\boldsymbol{\mathcal { G }}_{\mathrm{j}}^{+} \phi_{\mathrm{j}}^{+}+\boldsymbol{\mathcal { N }}_{\mathrm{j}}^{+} \psi_{\mathrm{j}}^{+}-\boldsymbol{\mathcal { G }}_{\mathrm{j}, \mathrm{j}+1}^{-} \phi_{\mathrm{j}+1}^{-}-\boldsymbol{\mathcal { N }}_{\mathrm{j}, \mathrm{j}+1}^{-} \psi_{\mathrm{j}+1}^{-} \tag{4.48}
\end{equation*}
$$

Another compatibility equation is obtained on the $(\mathrm{j}+1)^{\text {th }}$ interface:

$$
\begin{equation*}
\frac{\psi_{\mathrm{j}+1}^{-}}{2}=\boldsymbol{\mathcal { G }}_{\mathrm{j}+1, \mathrm{j}}^{+} \phi_{\mathrm{j}}^{+}+\boldsymbol{\mathcal { N }}_{\mathrm{j}+1, \mathrm{j}}^{-} \psi_{\mathrm{j}}^{+}-\boldsymbol{\mathcal { G }}_{\mathrm{j}+1}^{-} \phi_{\mathrm{j}+1}^{-}-\boldsymbol{\mathcal { N }}_{\mathrm{j}+1}^{-} \psi_{\mathrm{j}+1}^{-} . \tag{4.49}
\end{equation*}
$$

In eqs. (4.48) and (4.49), we use the double-index Greens functions that are derived on two consecutive profiles: $\boldsymbol{\mathcal { G }}_{\mathrm{j}, \mathrm{j}+1}^{-} \equiv \boldsymbol{\mathcal { G }}_{\mathrm{j}}^{-}\left(\mathrm{s}_{\mathrm{j}}, \mathrm{s}_{\mathrm{j}+1}^{\prime}\right), \boldsymbol{\mathcal { G }}_{\mathrm{j}+1, \mathrm{j}}^{+} \equiv \boldsymbol{\mathcal { G }}_{\mathrm{j}}^{+}\left(\mathrm{s}_{\mathrm{j}+1}, \mathrm{~s}_{\mathrm{j}}^{\prime}\right)$, and similarly for $\mathcal{N}$.

The computability equation becomes, for the upper and lower media:

$$
\begin{align*}
& \frac{\psi_{\mathrm{M}}^{+}}{2}=\boldsymbol{\mathcal { G }}_{\mathrm{M}}^{+} \phi_{\mathrm{M}}^{+}+\boldsymbol{\mathcal { N }}_{\mathrm{M}}^{+} \psi_{\mathrm{M}}^{+},  \tag{4.50}\\
& \frac{\psi_{1}^{-}}{2}=-\boldsymbol{\mathcal { G }}_{1}^{-} \phi_{1}^{-}-\boldsymbol{\mathcal { N }}_{1}^{-} \psi_{1}^{-} . \tag{4.51}
\end{align*}
$$

By combining eqs (4.48) with eq. (4.49), we obtain the link between the unknows on the upper and on the lower interface of the $\mathrm{j}^{\text {th }}$ layer, for $\mathrm{j}=1, \mathrm{M}-1$ :

$$
\left(\begin{array}{cc}
\boldsymbol{\mathcal { N }}_{\mathrm{j}, \mathrm{j}+1}^{-} & \boldsymbol{\mathcal { G }}_{\mathrm{j}, \mathrm{j}+1}^{-}  \tag{4.52}\\
\boldsymbol{\mathcal { N }}_{\mathrm{j}+1}^{-}+\frac{\mathbb{I}}{2} & \boldsymbol{\mathcal { G }}_{\mathrm{j}+1}^{-}
\end{array}\right)\binom{\psi_{\mathrm{j}+1}^{-}}{\phi_{\mathrm{j}+1}^{-}}=\left(\begin{array}{cc}
\boldsymbol{\mathcal { N }}_{\mathrm{j}}^{+}-\frac{\mathbb{I}}{2} & \boldsymbol{\mathcal { G }}_{\mathrm{j}}^{+} \\
\boldsymbol{\mathcal { N }}_{\mathrm{j}+1, \mathrm{j}}^{+} & \boldsymbol{\mathcal { G }}_{\mathrm{j}+1, \mathrm{j}}^{+}
\end{array}\right)\binom{\psi_{\mathrm{j}}^{+}}{\phi_{\mathrm{j}}^{+}} .
$$

This equation gives the transmission operator of the unknown amplitudes across the $\mathrm{j}^{\text {th }}$ layer, i.e. from the $\mathrm{j}^{\text {th }}$ to the $(\mathrm{j}+1)^{\text {st }}$ interface, for $\mathrm{j}=1, \mathrm{M}-1$ :

$$
\mathrm{T}_{\mathrm{j}+1, \mathrm{j}}=\left(\begin{array}{cc}
\boldsymbol{\mathcal { N }}_{\mathrm{j}, \mathrm{j}+1}^{-} & \boldsymbol{\mathcal { G }}_{\mathrm{j}, \mathrm{j}+1}^{-}  \tag{4.53}\\
\boldsymbol{\mathcal { N }}_{\mathrm{j}+1}^{-}+\frac{\mathbb{I}}{2} & \boldsymbol{\mathcal { G }}_{\mathrm{j}+1}^{-}
\end{array}\right)^{-1}\left(\begin{array}{cc}
\boldsymbol{\mathcal { N }}_{\mathrm{j}}^{+}-\frac{\mathbb{I}}{2} & \boldsymbol{\mathcal { G }}_{\mathrm{j}}^{+} \\
\boldsymbol{\mathcal { N }}_{\mathrm{j}+1, \mathrm{j}}^{+} & \boldsymbol{\mathcal { G }}_{\mathrm{j}+1, \mathrm{j}}^{+}
\end{array}\right)
$$

The transmission operator includes an inverse operator. Numerically, this inversion leads to the inversion of a matrix, as we will see in section 4.6.

The transmission matrix across the $\mathrm{j}^{\text {th }}$ interface for $\mathrm{j}=1, \ldots, \mathrm{M}-1$ is obtained using the continuity of the tangential and normal field components, as given by eqs.(4.8) and (4.10):

$$
\binom{\psi_{\mathrm{j}}^{+}}{\phi_{\mathrm{j}}^{+}}=\mathrm{T}_{\mathrm{j}}^{+-}\binom{\psi_{\mathrm{j}}^{-}}{\phi_{\mathrm{j}}^{-}}, \quad \quad \mathrm{T}_{\mathrm{j}}^{+-}=\left(\begin{array}{cc}
\mathbb{I} & 0  \tag{4.54}\\
0 & \frac{\mathrm{q}_{\mathrm{j}}^{-}}{\mathrm{q}_{\mathrm{j}}^{+}}
\end{array}\right), \mathrm{j} \neq \mathrm{M}
$$

The advantage of this presentation is that there are no exponentially growing terms in the transmission matrices, since all the components of the two variable functions contain only scattered propagating or decreasing evanescent waves, in contrast with the other methods (differential, Fourier modal, Rayleigh, etc.). However, this formulation requires calculating the cross-layer functions between the interfaces, which leads to computation times equal to those of single-interface functions. As a consequence, the total computation time is almost multiplied by a factor 2 with respect to the case where cross-layer kernels can be avoided, as discussed in the next section.

Finally, it can be deduced from eqs. (4.53) and (4.54):

$$
\begin{equation*}
\binom{\psi_{\mathrm{M}}^{+}}{\phi_{\mathrm{M}}^{+}}+\binom{\psi^{\mathrm{i}}}{\phi^{\mathrm{i}}}=\mathrm{T}_{\mathrm{M}}^{+-} \mathrm{T}_{\mathrm{M}, \mathrm{M}-1} \quad \ldots \quad \mathrm{~T}_{2}^{+-} \mathrm{T}_{2,1} \mathrm{~T}_{1}^{+-}\binom{\psi_{1}^{-}}{\phi_{1}^{-}} \tag{4.55}
\end{equation*}
$$

Finally, eqs. (4.55), (4.50) and (4.51) form an operator system of 4 equations with 4 unknown functions, which can be solved on a computer after representing each operator by a matrix, as described in section 6 .

### 4.4.2. Profiles without interpenetration

This case is simpler than the situation in sec.4.4.1 because it is possible to use the plane-wave expansion between the grating profiles and thus to decouple the integral presentation used in eq.(4.44). The idea is illustrated in figure 4.3.

If it is possible to introduce a plane layer between the profiles, the plane wave expansion is valid inside this layer. The advantage is that the plane waves (propagating and evanescent) that represent each diffraction order m can be easily separated into to sets: (i) upgoing waves having amplitudes of the $y$-component of the field equal to $r_{j, m}$ that are
generated by the lower surface $S_{j}$, (ii) downgoing waves with amplitudes $t_{j, m}$ generated by the upper grating surface $\mathrm{S}_{\mathrm{j}+1}$.


Figure.4.3. Layer with two profiles that are separable by a plane layer
Let us first rewrite eqs.(4.184) at $\mathrm{z}>\mathrm{z}_{\mathrm{j}}^{+}$:

$$
\begin{equation*}
\mathrm{r}_{\mathrm{j}, \mathrm{~m}}=\int_{\mathrm{s}=0}^{\mathrm{s}_{\mathrm{j}, \mathrm{~d}}}\left[\mathrm{~N}_{\mathrm{j}, \mathrm{~m}}^{+}\left(\mathrm{s}_{\mathrm{j}}\right) \psi_{\mathrm{j}}^{+}\left(\mathrm{s}_{\mathrm{j}}\right)+\mathrm{G}_{\mathrm{j}, \mathrm{~m}}^{+}\left(\mathrm{s}_{\mathrm{j}}\right) \phi_{\mathrm{j}}^{+}\left(\mathrm{s}_{\mathrm{j}}\right)\right] \mathrm{ds}_{\mathrm{j}}, \tag{4.56}
\end{equation*}
$$

with

$$
\begin{gather*}
\mathrm{G}_{\mathrm{j}, \mathrm{~m}}^{+}\left(\mathrm{s}_{\mathrm{j}}\right)=\frac{1}{2 \mathrm{id} \gamma_{\mathrm{j}, \mathrm{~m}}^{+}} \exp \left[-\mathrm{imKx}\left(\mathrm{~s}_{\mathrm{j}}\right)-\mathrm{i} \gamma_{\mathrm{j}, \mathrm{~m}}^{+} \mathrm{z}\left(\mathrm{~s}_{\mathrm{j}}\right)\right],  \tag{4.57}\\
\mathrm{N}_{\mathrm{j}, \mathrm{~m}}^{+}\left(\mathrm{s}_{\mathrm{j}}\right)=\frac{1}{2 \mathrm{~d}}\left[\frac{\mathrm{dx}\left(\mathrm{~s}_{\mathrm{j}}\right)}{\mathrm{ds}_{\mathrm{j}}}-\frac{\alpha_{\mathrm{m}}}{\gamma_{\mathrm{j}, \mathrm{~m}}^{+}} \frac{\mathrm{dz}\left(\mathrm{~s}_{\mathrm{j}}\right)}{\mathrm{ds}_{\mathrm{j}}}\right] \exp \left[-\mathrm{imKx}\left(\mathrm{~s}_{\mathrm{j}}\right)-\mathrm{i} \gamma_{\mathrm{j}, \mathrm{~m}}^{+} \mathrm{z}\left(\mathrm{~s}_{\mathrm{j}}\right)\right] .
\end{gather*}
$$

and with $s_{j}$ denoting the curvilinear abscissa on the $j^{\text {th }}$ profile, $s_{j, d}$ being the curvilinear abscissa of the point of $\mathrm{S}_{\mathrm{j}}$ of abscissa $\mathrm{x}=\mathrm{d}$.

We then can directly use eq.(4.48) to express the field on the interface $\mathrm{z}=\mathrm{z}\left(\mathrm{s}_{\mathrm{j}+1}\right)$ :
$\frac{\psi_{\mathrm{j}+1}^{-}\left(\mathrm{s}_{\mathrm{j}+1}\right)}{2}=\sum_{\mathrm{m}} \exp \left[\operatorname{imKx}\left(\mathrm{s}_{\mathrm{j}+1}\right)+\mathrm{i} \gamma_{\mathrm{j}, \mathrm{m}}^{+} \mathrm{z}\left(\mathrm{s}_{\mathrm{j}+1}\right)\right] \mathrm{r}_{\mathrm{j}, \mathrm{m}}^{+}-\boldsymbol{\mathcal { G }}_{\mathrm{j}+1}^{-} \phi_{\mathrm{j}+1}^{-}-\mathcal{N}_{\mathrm{j}+1}^{-} \psi_{\mathrm{j}+1}^{-}$.
Let us consider the sum in eq. (4.58):

$$
\begin{equation*}
\zeta_{\mathrm{j}+1}=\sum_{\mathrm{m}} \exp \left[\operatorname{imKx}\left(\mathrm{~s}_{\mathrm{j}+1}\right)+\mathrm{i} \gamma_{\mathrm{j}, \mathrm{~m}}^{+} \mathrm{z}\left(\mathrm{~s}_{\mathrm{j}+1}\right)\right] \mathrm{r}_{\mathrm{j}, \mathrm{~m}} . \tag{4.59}
\end{equation*}
$$

Inserting in this equation the value of $\mathrm{r}_{\mathrm{j}, \mathrm{m}}$ given by eq. (4.56) yields:

$$
\begin{align*}
& \zeta_{\mathrm{j}+1}\left(\mathrm{~s}_{\mathrm{j}+1}\right)= \\
& \quad \sum_{\mathrm{m}} \exp \left[\operatorname{imKx}\left(\mathrm{~s}_{\mathrm{j}+1}\right)+\mathrm{i} \gamma_{\mathrm{j}, \mathrm{~m}}^{+} \mathrm{z}\left(\mathrm{~s}_{\mathrm{j}+1}\right)\right] \int_{\mathrm{s}_{\mathrm{j}}=0}^{\mathrm{s}_{\mathrm{j}, \mathrm{~d}}}\left[\mathrm{~N}_{\mathrm{j}, \mathrm{~m}}^{+}\left(\mathrm{s}_{\mathrm{j}}\right) \psi_{\mathrm{j}}^{+}\left(\mathrm{s}_{\mathrm{j}}\right)+\mathrm{G}_{\mathrm{j}, \mathrm{~m}}^{+}\left(\mathrm{s}_{\mathrm{j}}\right) \phi_{\mathrm{j}}^{+}\left(\mathrm{s}_{\mathrm{j}}\right)\right] \mathrm{ds}_{\mathrm{j}}, \tag{4.60}
\end{align*}
$$

which can be written in operator form:

$$
\begin{equation*}
\zeta_{\mathrm{j}+1}=\mathbf{N}_{\mathrm{j}, \mathrm{j}+1}^{+} \psi_{\mathrm{j}}^{+}+\mathbf{G}_{\mathrm{j}, \mathrm{j}+1}^{+} \phi_{\mathrm{j}}^{+}, \tag{4.61}
\end{equation*}
$$

where the operators $\mathbf{N}_{\mathrm{j}, \mathrm{j}+1}^{+}$and $\mathbf{G}_{\mathrm{j}, \mathrm{j}+1}^{+}$are obtained from a summation in m and an integral in $s_{j}$. Thus we can write eq. (4.58) in the operator form:

$$
\begin{equation*}
\frac{\psi_{\mathrm{j}+1}^{-}}{2}=\mathbf{N}_{\mathrm{j}, \mathrm{j}+1}^{+} \psi_{\mathrm{j}}^{+}+\mathbf{G}_{\mathrm{j}, \mathrm{j}+1}^{+} \phi_{\mathrm{j}}^{+}-\boldsymbol{\mathcal { G }}_{\mathrm{j}+1}^{-} \phi_{\mathrm{j}+1}^{-}-\boldsymbol{\mathcal { N }}_{\mathrm{j}+1}^{-} \psi_{\mathrm{j}+1}^{-} . \tag{4.62}
\end{equation*}
$$

The second relation comes from eq. (4.49) by using the amplitudes of the diffraction orders coming down from the upper interfaces and valid below $z=z_{j+1}^{-}$. Following the same lines as in the previous paragraph yields:

$$
\begin{equation*}
\frac{\psi_{\mathrm{j}}^{-}}{2}=\mathbf{N}_{\mathrm{j}+1, \mathrm{j}}^{-} \psi_{\mathrm{j}+1}^{-}+\mathbf{G}_{\mathrm{j}+1, \mathrm{j}}^{-} \phi_{\mathrm{j}+1}^{-}+\boldsymbol{\mathcal { G }}_{\mathrm{j}}^{+} \phi_{\mathrm{j}}^{+}+\mathcal{N}_{\mathrm{j}}^{+} \psi_{\mathrm{j}}^{+} . \tag{4.63}
\end{equation*}
$$

The transmission matrix between the $\mathrm{j}^{\text {th }}$ and the $\mathrm{j}+1^{\text {st }}$ interface takes a form similar to the case of interpenetrating layers, eq.(4.53). However, the difference is essential, because each series used in (4.63) is evaluated on a single interface:

$$
\mathrm{T}_{\mathrm{j}+1, \mathrm{j}}=\left(\begin{array}{cc}
\mathbf{N}_{\mathrm{j}+1, \mathrm{j}}^{-} & \mathbf{G}_{\mathrm{j}+1, \mathrm{j}}^{-}  \tag{4.64}\\
\boldsymbol{\mathcal { N }}_{\mathrm{j}+1}^{-}+\frac{\mathbb{I}}{2} & \boldsymbol{\mathcal { G }}_{\mathrm{j}+1}^{-}
\end{array}\right)^{-1}\left(\begin{array}{cc}
\boldsymbol{\mathcal { N }}_{\mathrm{j}}^{+}-\frac{\mathbb{I}}{2} & \boldsymbol{\mathcal { G }}_{\mathrm{j}}^{+} \\
\mathbf{N}_{\mathrm{j}, \mathrm{j}+1}^{+} & \mathbf{G}_{\mathrm{j}, \mathrm{j}+1}^{+}
\end{array}\right)
$$

The second difference is that the exponential terms are explicitly given in the $\mathbf{N}_{\mathrm{j}, \mathrm{j}+1}^{+}$, $\mathbf{N}_{\mathrm{j}+1, \mathrm{j}}^{-}, \quad \mathbf{G}_{\mathrm{j}, \mathrm{j}+1}^{+}$and $\mathbf{G}_{\mathrm{j}+1, \mathrm{j}}^{-}$through the functions $\exp \left[\operatorname{imKx}\left(\mathrm{s}_{\mathrm{j}+1}\right)+\mathrm{i} \gamma_{\mathrm{j}, \mathrm{m}}^{+} \mathrm{z}\left(\mathrm{s}_{\mathrm{j}+1}\right)\right]$ and $\exp \left[\operatorname{imKx}\left(\mathrm{s}_{\mathrm{j}}\right)-\mathrm{i} \gamma_{\mathrm{j}, \mathrm{m}}^{+} \mathrm{z}\left(\mathrm{s}_{\mathrm{j}}\right)\right]$ in such a way that it can be extracted from each of these operators a part containing all the growing and decreasing exponential terms, which allow a much better stability of the numerical implementation through adequate treatments, for example the S-matrix algorithm described in appendix A and B of Chapter 7.

### 4.5. Gratings in conical mounting

When classical gratings with one-dimensional periodicity are used with incidence plane perpendicular to the grooves, the diffraction orders lie in the same plane. Off-plane incidence brings the diffraction orders out of the plane and their directions lie on a cone, which explains the term of conical diffraction. One of the first experimental works can be found in [24, 25]. The use of conical mount is typical for concave gratings and in some spectrographs aiming to separate off-plane the incident and the diffracted directions.

The first theoretical studies were made in 1971 using the integral [26] and the differential [30] methods. An interesting theoretical development came in 1972 when Maystre and Petit demonstrated that under special conditions, mechanically ruled perfectly conducting gratings can have very high and constant efficiency over a large spectral domain [31]. They also established the theorem of invariance [32] that gives an expression of the diffraction efficiency of a perfectly conducting gratings in conical mounting, expressed as a linear combination of efficiencies in an in-plane (classical) mount for the two fundamental polarizations. Since the theorem is not valid for finitely conducting metals, later development of the integral method allowed studies of diffraction gratings behaviour in conical mount when working in the UV and visible [33, 34]. An interested reader can find the demonstration of the invariance theorem in [7, 33].


Figure 4.4. Parameters of the incident wave in conical mount.The angle $\varphi$ denotes the angle between the incident wavevector $\overrightarrow{\mathrm{k}}^{\mathrm{i}}$ and its projection $\overrightarrow{\mathrm{k}}_{\mathrm{xz}}^{\mathrm{i}}$ on the $x z$ plane. The angle $\theta^{\mathrm{i}}$ is the angle between the $z$ axis and $\overrightarrow{\mathrm{k}}_{\mathrm{xz}}^{\mathrm{i}}$. In order to define the polarization of the incident field, we construct the circle $M N M^{\prime} N^{\prime}$ in the plane perpendicular to the incident wavevector $\overrightarrow{\mathrm{k}}{ }^{\mathrm{i}}$, with the continuation of $N N^{\prime}$ intersecting the $z$ axis and $M M^{\prime}$ being perpendicular to $N N^{\prime}$. The polarization angle $\delta$ is the angle between M'M and the direction of the incident wavevector $\overrightarrow{\mathrm{k}}^{\mathrm{i}}$.

The notations are summarized in figure 4.4. The mathematical formulation of the invariance theorem takes the form of an equivalence between the conical case and an associated classical case:
(i) conical case:

With incident angles $\theta^{i}$ and $\varphi$, incident polarization angle $\delta$, incident wavelength $\lambda$, the efficiencies in the various orders are denoted by $\rho_{\mathrm{m}}\left(\theta^{\mathrm{i}}, \varphi, \delta, \lambda\right)$.
(ii) fictitious equivalent classical case:

The wavelength is increased to $\lambda^{\prime}=\lambda / \cos \varphi$, the angle $\varphi^{\prime}$ is now equal to 0 (in-plane incidence), the angle $\theta^{\prime i}=\theta^{i}$, then efficiencies in TE and TM polarizations are denoted by $\rho_{\mathrm{m}}^{\mathrm{TE}}\left(\theta^{\mathrm{i}}, \lambda^{\prime}\right)$ and $\rho_{\mathrm{m}}^{\mathrm{TM}}\left(\theta^{\mathrm{i}}, \lambda^{\prime}\right)$, respectively. It can be noticed that the projection $\overrightarrow{\mathrm{k}}_{\mathrm{xz}}^{\mathrm{i}}$ of the
wavevector of the incident wave on the $x z$ plane in conical mount identifies with the wavevector of the incident waves in the fictitious equivalent case.

The invariance theorem states that:

$$
\begin{equation*}
\rho_{\mathrm{m}}\left(\theta^{\mathrm{i}}, \varphi, \delta, \lambda\right)=(\cos \delta)^{2} \rho_{\mathrm{m}}^{\mathrm{TE}}\left(\theta^{\mathrm{i}}, \lambda^{\prime}\right)+(\sin \delta)^{2} \rho_{\mathrm{m}}^{\mathrm{TM}}\left(\theta^{\mathrm{i}}, \lambda^{\prime}\right) . \tag{4.65}
\end{equation*}
$$

It is to be noticed that, as for the incident wavevector, the projections of the wavevectors of the scattered waves on the xz plane in conical mount are identical to the wavevectors of the scattered waves in the fictitious equivalent case.

### 4.6. Numerical tools for an efficient numerical implementation

### 4.6.1. Integration schemes for the integral equation

All the integral equations in this chapter link the value of an unknown function $u(s)$ at a given point of $S$ to its value on its entire domain of definition:

$$
\begin{equation*}
\mathrm{cu}(\mathrm{~s})=\mathrm{v}(\mathrm{~s})+\int_{0}^{\mathrm{s}_{\mathrm{d}}} \mathrm{~W}\left(\mathrm{~s}, \mathrm{~s}^{\prime}\right) \mathrm{u}\left(\mathrm{~s}^{\prime}\right) \mathrm{d} \mathrm{~s}^{\prime} \tag{4.66}
\end{equation*}
$$

where all functions are periodical, with v and W being known functions, W being possibly singular but integrable. The constant c takes values 0 or 1 according to whether the integral equation is of the first or second kind.

There are many different ways to solve such an equation for the grating problem. Pavageau et al. proposed an iterative method [13] that does not require any matrix inversion, like the well known Born method for scattering problems. Unfortunately, it may diverge [35, 36].

A well known general methods is based on the periodicity of all functions of the equation, which allows a projection of these functions and of the equation on the Fourier space:

$$
\begin{equation*}
\mathrm{u}(\mathrm{~s})=\sum_{\mathrm{m}} \mathrm{u}_{\mathrm{m}} \exp \left(\mathrm{imK}_{\mathrm{s}} \mathrm{~s}\right), \quad \mathrm{K}_{\mathrm{s}}=\frac{2 \pi}{\mathrm{~s}_{\mathrm{d}}} \tag{4.67}
\end{equation*}
$$

and similar expressions for v and W . The integral equation is transformed into a linear system of algebraic equations:

$$
\begin{equation*}
\sum_{\mathrm{m}}\left(\mathrm{~W}_{\mathrm{nm}}-\mathrm{c} \delta_{\mathrm{nm}}\right) \mathrm{u}_{\mathrm{m}}=\mathrm{v}_{\mathrm{n}}, \quad \forall \mathrm{n} \tag{4.68}
\end{equation*}
$$

which can be solved numerically after truncation. However, this approach requires computing a double Fourier decomposition:

$$
\begin{equation*}
\mathrm{W}_{\mathrm{nm}}=\int_{0}^{\mathrm{s}_{\mathrm{d}} \mathrm{~s}_{\mathrm{d}}} \int_{0}^{\mathrm{W}}\left(\mathrm{~s}, \mathrm{~s}^{\prime}\right) \mathrm{u}\left(\mathrm{~s}^{\prime}\right) \exp \left(-\mathrm{inK} \mathrm{~s}_{\mathrm{s}} \mathrm{~s}-\mathrm{imK}_{\mathrm{s}^{\prime}} \mathrm{s}^{\prime}\right) \mathrm{ds}^{2} \mathrm{ds} \mathrm{~s}^{\prime} \tag{4.69}
\end{equation*}
$$

The method has been applied to gratings with profiles consisting of few straight segments, because in that case the double Fourier integral can be calculated in closed form. It is so for triangular profiles [11] or trapezoidal gratings [37].

The most widespread method is the so-called point-matching (or discretization) method. Instead ot using discrete Fourier components, the unknown function is discretized on the
grating profile and represented by its vales $u_{j}=u\left(s_{j}\right)$ inside the interval of integration. Ther integration process leads toan equation quite similar to eq.(4.68)

$$
\begin{equation*}
\sum_{\mathrm{p}}\left(\mathrm{~W}_{\mathrm{jp}}-\mathrm{c} \delta_{\mathrm{jp}}\right) \mathrm{u}_{\mathrm{p}}=\mathrm{v}_{\mathrm{j}}, \quad \forall \mathrm{j} \in[1, \mathrm{P}], \tag{4.70}
\end{equation*}
$$

with $P$ being the number of matching points. It is worth noting that the value of $W_{j p}$ may differ from the value of $\mathrm{W}\left(\mathrm{s}_{\mathrm{j}}, \mathrm{s}_{\mathrm{p}}\right)$ obtained through the rectangular rule of integration (multiplied by the weight of integration) and can require much more complicated treatments, specially if $W$ is singular. The sophistication of this treatment is one of the decisive keys for the precision of the solution of the integral solution. The second key is the analytical study of the kernel, which is described in the next two sections.

When W is regular, continuous, with a continuous first derivative, the rectangular rule of integration is quite precise since the function to be integrated is periodic and it can be noticed that the rectangular and trapezoidal rules are completely equivalent in that case. However, the derivative of $\mathrm{W}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$ is generally discontinuous when $\mathrm{s}=\mathrm{s}$ ' and a trapezoidal rule or higher order treatment provide a better precision [38, 39]. In what follows we assume that $u(s)$ is a continuous function. This is obviously the case when the grating profile has no edges. Several more detailed arguments in favour of the rectangular rule can be found in [6, 7]. Let us shortly repeat one of them. We consider an integral of a periodical continuous function a(s):

$$
\begin{equation*}
\mathrm{a}_{0}=\int_{0}^{\mathrm{s}_{\mathrm{d}}} \mathrm{a}(\mathrm{~s}) \mathrm{ds} . \tag{4.71}
\end{equation*}
$$

The exact integral is equal to the $0^{\text {th }}$ term in the Fourier expansion of $\mathrm{a}(\mathrm{s})$ :

$$
\begin{equation*}
\mathrm{a}(\mathrm{~s})=\sum_{\mathrm{m}} \mathrm{a}_{\mathrm{m}} \exp \left(\mathrm{imK}_{\mathrm{s}} \mathrm{~s}\right) \tag{4.72}
\end{equation*}
$$

Simple calculations show that when using the rectangular rule with P discretization points, the integration error is proportional to the $\mathrm{P}^{\text {th }}$ Fourier coefficient of $\mathrm{a}(\mathrm{x})$. Since it is continuous, the Fourier series converges like $1 / \mathrm{P}^{2}$ at least.

The implementation of the rectangular rule is very simple. We define the values of $\mathrm{W}_{\mathrm{jp}}$ in the following manner:

$$
\begin{equation*}
W_{j, p}=\frac{s_{d}}{P} W\left(\frac{j}{P} s_{d}, \frac{p}{P} s_{d}\right), \quad j, p=1, \ldots, P-1 \tag{4.73}
\end{equation*}
$$

Using this result, the product of the unknown functions $u(s)$ with the non-singular parts of the kernels can be integrated in the form of a simple matrix product:

$$
\begin{equation*}
\int_{0}^{\mathrm{s}_{\mathrm{d}}} \mathrm{~W}\left(\mathrm{~s}_{\mathrm{j}}, \mathrm{xs}\right) \mathrm{u}\left(\mathrm{~s}^{\prime}\right) \mathrm{ds}^{\prime} \approx \sum_{\mathrm{p}} \mathrm{~W}_{\mathrm{j}, \mathrm{p}} \mathrm{u}_{\mathrm{p}} \tag{4.74}
\end{equation*}
$$

In the case of multilayer gratings with profiles interpenetration, there are two main situations that can complicate the numerical evaluation of the cross-layer functions which link the fields and normal derivatives on both sides of a layer:
(i) The wavelength $\lambda$ is much larger than the layer thickness $t$. In that case the profiles are located too close to each other, compared to $\lambda$, so that the functions $\mathcal{G}$ and $\mathcal{N}$ become large in
modulus for $\mathrm{s}_{\mathrm{j}+1} \rightarrow \mathrm{~s}_{\mathrm{j}}$. This behaviour has the same origin as the singularities of the kernels $\mathcal{\mathcal { G }}$ and $\mathcal{N}$ for a bare grating, which will be discussed in section 6.3 and can be eliminated in a similar manner. Numerical results show no problems as long as layer thicknesses exceed $\lambda / 20$. (ii) The period d is much greater than the layer thickness t . If the wavelength is not too much larger than $t$, then the kernels have no singularities, but their moduli present peaks when the distance $\left|\mathrm{P}_{\mathrm{j}+1, \mathrm{p}} \mathrm{P}_{\mathrm{j}, \mathrm{q}}\right|$ between two points located on the two different profiles is small with respect to the discretization distance between the points located on the same profile. The width of these maxima is of the order of magnitude of the layer thickness, thus the correct implementation the trapezoidal integration rule requires the distance between two consecutive points of the profile discretization $\Delta=\left|\mathrm{s}_{\mathrm{j}, \mathrm{p}}-\mathrm{s}_{\mathrm{j}, \mathrm{p}-1}\right|$ to be less than the width of the maxima. As a rule of thumb, if $\Delta \approx \mathrm{d} / \mathrm{N}_{\mathrm{p}}$, where $\mathrm{N}_{\mathrm{p}}$ is the number of integration points, then its lower limit is determined by the relation:

$$
\begin{equation*}
\mathrm{N}_{\mathrm{p}, \min } \propto \frac{\mathrm{~d}}{\mathrm{t}} \tag{4.75}
\end{equation*}
$$

Thus, for echelles, a values of $d$ of about $10 \mu \mathrm{~m}$ and t of 20 nm requires the number of integration points to exceed 500 . Note that $\mathrm{N}_{\mathrm{p}}$ directly determines the number of unknown values of $\phi_{\mathrm{j}}$ and $\psi_{\mathrm{j}}$ and thus the size of the matrices to be used. Practically, it is not worth nowadays for $\mathrm{N}_{\mathrm{p}}$ to exceed 10000 , because of the computation time, memory requirements, round-off errors and limited digits. As a consequence, it is necessary to find another way of integration instead of the trapezoidal rule.

There is another problem that can come from the matrix inversion in the construction of the transmission matrix between the layers, eq.(4.53). Contrary to the transmission matrix that contains growing and decreasing exponential terms in the plane wave expansion, (thus requires some type of recursive algorithm to contain the contamination of the growing exponentials, S-matrix algorithm, for example, see appendix 4.C), the distance between the profiles in the z-direction that appears in the kernels in eqs. (4.45) and (4.46) restricts the terms to only propagation or evanescently decreasing ones. However, the matrix inversion of these terms that is required in eq.(4.53) can create exponentially growing terms. Two situations can appear:

1. The matrix inversion in eq.(4.53) can be done without numerical problems. This happens when the layer thickness is not quite large. In that case it is possible to progress upwards in the stack of layers by following the S-matrix algorithm.
2. The matrix inversion does not work. This could happen if the distance between two consecutive interface is large. Two different geometries can be concerned:
(i) there is no interpenetration of these two profiles. In that case one can easily apply the technique described in the next section.
(ii) there is interpenetration of two very deep interfaces. It is possible to use directly eq.(4.52) in the S-matrix algorithm without inverting the matrix to calculate the entire Tmatrix in eq. (4.53). The formulation of the S-matrix algorithm to an equation having the form given in (4.52) is quite similar to the classical aplication, but it requires one additional iteration step. The advantage is that it avoids the inversion of small terms that can lead to singular matrices. This special technique is given in Appendix 7.B and is not quite popular, but can be used in other methods that apply for multilayer stack, for example, in the coordinate transformation method.

### 4.6.2. Summation of the kernels

There are several problems in the calculation of the functions included in eqs (4.19) and (4.23):

$$
\begin{align*}
& \boldsymbol{\mathcal { G }}^{ \pm}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{id}} \sum_{\mathrm{m}=-\infty}^{\infty} \frac{1}{\gamma_{\mathrm{m}}^{ \pm}} \mathrm{e}^{\mathrm{imK}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}^{\prime}\left(\mathrm{s}^{\prime}\right)\right]+\mathrm{i} \gamma_{\mathrm{m}}^{ \pm}\left|\mathrm{z}(\mathrm{~s})-\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right|},  \tag{4.76}\\
& \mathcal{N}^{ \pm}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{~d}} \sum_{\mathrm{m}=-\infty}^{\infty}\left[\frac{\mathrm{dx} \mathrm{x}^{\prime}}{\mathrm{ds}} \operatorname{sgn}\left(\mathrm{z}(\mathrm{~s})-\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right)-\frac{\alpha_{\mathrm{m}}}{\gamma_{\mathrm{m}}^{ \pm}} \frac{\mathrm{dz}}{\mathrm{~d} \mathrm{~s}^{\prime}}\right] \mathrm{e}^{\operatorname{imK}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}^{\prime}\left(\mathrm{s}^{\prime}\right)\right]+\mathrm{i} \gamma_{\mathrm{m}}^{ \pm}\left|\mathrm{z}(\mathrm{~s})-\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right|},  \tag{4.77}\\
& \mathscr{K}\left(s, s^{\prime}\right)=\frac{1}{2 d} \sum_{m=-\infty,+\infty}\left[\operatorname{sgn}\left(z-z^{\prime}\right) \frac{d x}{d s}-\frac{\alpha_{m}}{\gamma_{\mathrm{m}}^{+}} \frac{d z}{d s}\right] e^{\operatorname{imK}\left[x(s)-x^{\prime}\left(s^{\prime}\right)\right]+i \gamma_{m}^{ \pm}\left|z(s)-z^{\prime}\left(s^{\prime}\right)\right|}, \tag{4.78}
\end{align*}
$$

with:

$$
\begin{gather*}
\alpha_{\mathrm{m}}=\alpha_{0}+\mathrm{mK},  \tag{4.79}\\
\gamma_{\mathrm{m}}^{ \pm}=\sqrt{\left(\mathrm{kn}^{ \pm}\right)^{2}-\alpha_{\mathrm{m}}^{2}} \tag{4.80}
\end{gather*}
$$

For the sake of simplicity, we assume here that $\mathrm{n}^{ \pm}=1$ and we cancel the superscript $\pm$ in $\gamma_{\mathrm{m}}^{ \pm}, \mathcal{G}^{ \pm}$and $\boldsymbol{N}^{ \pm}$.

Let us at first evaluate the asymptotic values of $\gamma_{\mathrm{m}}$ and $\alpha_{\mathrm{m}}$ for large values of m :

$$
\begin{gather*}
\alpha_{\mathrm{m}} \underset{\mathrm{~m} \rightarrow \infty}{\rightarrow} \mathrm{mK}, \\
\gamma_{\mathrm{m}} \approx \mathrm{i}\left|\alpha_{\mathrm{m}}\right|-\frac{\mathrm{ik} \mathrm{k}^{2}}{2\left|\alpha_{\mathrm{m}}\right|} \underset{\mathrm{m} \rightarrow \infty}{\rightarrow} \mathrm{i}\left|\alpha_{0}+\mathrm{mK}\right|,  \tag{4.81}\\
\frac{1}{\gamma_{\mathrm{m}}} \approx \frac{1}{\mathrm{i}\left|\alpha_{\mathrm{m}}\right|}+\frac{\mathrm{k}^{2}}{2 \mathrm{i}\left|\alpha_{\mathrm{m}}\right|^{3}} \underset{\mathrm{~m} \rightarrow \infty}{\rightarrow} \frac{1}{\mathrm{i}|\mathrm{~m}| \mathrm{K}} .
\end{gather*}
$$

We consider the function

$$
\begin{equation*}
\boldsymbol{\mathcal { G }}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{id}} \sum_{\mathrm{m}=-\infty}^{\infty} \frac{1}{\gamma_{\mathrm{m}}} \exp \left\{\operatorname{imK}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)\right]+\mathrm{i} \gamma_{\mathrm{m}}\left|\mathrm{z}(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right|\right\} . \tag{4.82}
\end{equation*}
$$

At point $\mathrm{s}=\mathrm{s}^{\prime}$, it is obvious that the sum does not converge since the terms decrease in $1 /|\mathrm{m}|$. Of course, a very slow convergence can be expected when the two points are close to each other. Different techniques have been proposed to accelerate the convergence. Neureuther and Zaki [15] have employed a transformation technique based on the use of Mellin transforms. Other authors have proposed accelerating processes [40-43]. Here we describe a direct approach [7]. Let us determine at first the asymptotic expression of the kernel. If we replace (4.81) into eq.(4.82), we obtain the asymptotic term in the sum:

$$
\begin{align*}
\boldsymbol{\mathcal { G }}_{\infty}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) & =\mathrm{e}^{\alpha_{0}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)\right]} \sum_{\mathrm{m}=-1}^{-\infty} \frac{1}{4 \pi \mathrm{~m}} \mathrm{e}^{\mathrm{mK}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)+\mathrm{i}\left|z(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right|\right]}  \tag{4.83}\\
& -\mathrm{e}^{-\alpha_{0}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)\right]} \sum_{\mathrm{m}=1}^{\infty} \frac{1}{4 \pi \mathrm{~m}} \mathrm{e}^{-\mathrm{mK}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)-\mathrm{i}\left|z(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right|\right]},
\end{align*}
$$

which contains two sums of the form $\sum_{\mathrm{m}=1}^{\infty} \xi^{\mathrm{m}} / \mathrm{m}$ and can be summed in closed form:

$$
\begin{align*}
\boldsymbol{\mathcal { G }}_{\infty}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) & =\frac{1}{4 \pi} \mathrm{e}^{\alpha_{0}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)\right]} \log \left(1-\mathrm{e}^{-\mathrm{K}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)+\mathrm{i}\left|\mathrm{z}(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right|\right]}\right)  \tag{4.84}\\
& +\frac{1}{4 \pi} \mathrm{e}^{-\alpha_{0}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)\right]} \log \left(1-\mathrm{e}^{\left.-\mathrm{K}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)-\mathrm{i} \mathrm{i}(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right]\right]}\right) .
\end{align*}
$$

The calculation of the kernel in (4.82) is achieved by subtracting the asymptotic value:

$$
\begin{equation*}
\boldsymbol{\mathcal { G }}=\boldsymbol{\mathcal { G }}_{\infty}+\left(\boldsymbol{\mathcal { g }}-\boldsymbol{\mathcal { G }}_{\infty}\right) \tag{4.85}
\end{equation*}
$$

The first term in the right-and side is explicitely given in (4.84). It is singular for $\mathrm{s}=\mathrm{s}$ 'e. This singularity is integrable and is be treated in the next section.

The term between parenthesis in eq.(4.85) is obtained by combining the terms in the sums in eqs.(4.82) and (4.83). As far as the second one is the asymptotic value of the former, the series converges, whatever the values of s and s'. Furthermore, it is possible to show that by combining the terms m and -m in the sum, we finally obtain a rapidly converging series, whose terms decrease as $\mathrm{m}^{-3}$ when $\mathrm{s}=\mathrm{s}^{\prime}$, as shown later in eq.(4.87) Moreover, in this case this series is continuous and its value is simply given by:

$$
\begin{equation*}
\left(\boldsymbol{\mathcal { G }}-\boldsymbol{\mathcal { G }}_{\infty}\right)_{\mid \mathrm{s}=\mathrm{s}^{\prime}}=\frac{1}{2 \mathrm{id} \gamma_{0}}+\sum_{\mathrm{m} \neq 0}\left(\frac{1}{2 \mathrm{id} \gamma_{\mathrm{m}}}+\frac{1}{4 \pi|\mathrm{~m}|}\right) \tag{4.86}
\end{equation*}
$$

Using the third identity of eq.(4.81), for large values of m , the term in the sum is equal to:

$$
\begin{equation*}
\frac{1}{2 \mathrm{id} \gamma_{\mathrm{m}}}+\frac{1}{4 \pi|\mathrm{~m}|} \rightarrow-\frac{\mathrm{k}^{2}}{4 \mathrm{~d}\left|\alpha_{\mathrm{m}}\right|^{3}} \rightarrow-\frac{\mathrm{d}^{2}}{8 \pi \lambda^{2}} \frac{1}{|\mathrm{~m}|^{3}} \tag{4.87}
\end{equation*}
$$

Obviously, the singularity and the slow convergence of $\mathcal{G}$ have been carried out by introducing the series $\boldsymbol{\mathcal { G }}_{\infty}$. Fortunaley, this singularity is logarithmic and thus integrable, as shown in the next section.

Let us now deal with the second function defined in eq. (4.77):

$$
\begin{equation*}
\mathcal{N}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{~d}} \sum_{\mathrm{m}=-\infty}^{\infty}\left\{\frac{\mathrm{dx}\left(\mathrm{~s}^{\prime}\right)}{\mathrm{ds} s^{\prime}} \operatorname{sgn}\left[\mathrm{z}(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right]-\frac{\alpha_{\mathrm{m}}}{\gamma_{\mathrm{m}}} \frac{\mathrm{dz}\left(\mathrm{~s}^{\prime}\right)}{\mathrm{ds} s^{\prime}}\right\} \mathrm{e}^{\mathrm{imK}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)\right]+\mathrm{i} \gamma_{\mathrm{m}}\left|\mathrm{z}(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right|} \tag{4.88}
\end{equation*}
$$

At the first glance, the term $\operatorname{sgn}\left[z(s)-z\left(s^{\prime}\right)\right]$ suggests us that this function is not continuous for $\mathrm{s}^{\prime}=\mathrm{s}$, at least if $\frac{\mathrm{dz}^{\prime}}{\mathrm{ds}^{\prime}} \neq 0$ for $\mathrm{s}^{\prime}=\mathrm{s}$. To deal with this term, we proceed in the same way as in eq.(4.85), by introducing an asymptotic value $\boldsymbol{\mathcal { N }}_{\infty}$ and we set now:

$$
\begin{equation*}
\mathcal{N}=\mathcal{N}_{\infty}+\left(\mathcal{N}-\mathcal{N}_{\infty}\right) \tag{4.89}
\end{equation*}
$$

with

$$
\begin{align*}
& \mathcal{N}_{\infty}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{~d}} \operatorname{sgn}\left[\mathrm{z}(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right] \\
& +\frac{1}{2 \mathrm{~d}}\left\{\frac{\mathrm{dx}\left(\mathrm{~s}^{\prime}\right)}{\mathrm{ds} s^{\prime}} \operatorname{sgn}\left[\mathrm{z}(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right]-\mathrm{i} \frac{\mathrm{dz}\left(\mathrm{~s}^{\prime}\right)}{\mathrm{ds} s^{\prime}}\right\} \mathrm{e}^{-\alpha_{0}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)\right]} \sum_{\mathrm{m}=1}^{\infty} \mathrm{e}^{-\mathrm{mK}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)\right]+\mathrm{imK}\left[\mathrm{z}(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right]}  \tag{4.90}\\
& +\frac{1}{2 \mathrm{~d}}\left\{\frac{\mathrm{dx}\left(\mathrm{~s}^{\prime}\right)}{\mathrm{ds}} \operatorname{sgn}\left[\mathrm{z}(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right]+\mathrm{i} \frac{\mathrm{dz}\left(\mathrm{~s}^{\prime}\right)}{\mathrm{ds} s^{\prime}}\right\} \mathrm{e}^{\alpha_{0}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)\right]} \sum_{\mathrm{m}=-1}^{-\infty} \mathrm{e}^{\mathrm{mK}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)\right]+\mathrm{imK}\left[\mathrm{z}(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right]} .
\end{align*}
$$

The sums can be evaluated in a closed form:

$$
\begin{align*}
& \mathcal{N}_{\infty}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{~d}} \operatorname{sgn}\left[\mathrm{z}(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right] \\
& \quad+\frac{1}{2 \mathrm{~d}}\left\{\frac{\mathrm{dx}\left(\mathrm{~s}^{\prime}\right)}{\mathrm{ds} s^{\prime}} \operatorname{sgn}\left[\mathrm{z}(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right]-\mathrm{i} \frac{\mathrm{dz}\left(\mathrm{~s}^{\prime}\right)}{\mathrm{ds} s^{\prime}}\right\} \frac{\mathrm{e}^{-\alpha_{0}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)\right]}}{\mathrm{e}^{\mathrm{K}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)\right]-\mathrm{imK}\left[\mathrm{z}(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right]}-1}  \tag{4.91}\\
& \quad+\frac{1}{2 \mathrm{~d}}\left\{\frac{\mathrm{dx}\left(\mathrm{~s}^{\prime}\right)}{\mathrm{ds} s^{\prime}} \operatorname{sgn}\left[\mathrm{z}(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right]+\mathrm{i} \frac{\mathrm{dz}\left(\mathrm{~s}^{\prime}\right)}{\mathrm{ds} s^{\prime}}\right\} \frac{\mathrm{e}^{\alpha_{0}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)\right]}}{\mathrm{e}^{\left.\mathrm{K}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)\right]+\mathrm{imK}\left[\mathrm{z}(\mathrm{~s})-\mathrm{z(s}^{\prime}\right)\right]-1} .}
\end{align*}
$$

As noticed for $\mathcal{G}$, the term $\mathcal{N}-\mathcal{N}_{\infty}$ must be considered as a series each term of which is obtained by making the difference of the corresponding terms in the sums in eqs.(4.88) and (4.90). This series converges much more rapidly than the series in eq. (4.88) and it is continuous at $\mathrm{s}=\mathrm{s}$ '.

The limit of $\mathcal{N}$ when $\mathrm{s}^{\prime} \rightarrow \mathrm{s}$ can be determined calculating the limits of the two terms $\mathcal{N}_{\infty}$ and $\left(\mathcal{N}-\mathcal{N}_{\infty}\right)$. After tedious calculations, we can deduce that this limit exists and is given by:
$\mathcal{N}(\mathrm{s}, \mathrm{s})=\lim _{\mathrm{s}^{\prime} \rightarrow \mathrm{s}} \mathcal{N}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)=-\frac{\mathrm{dz}}{\mathrm{ds}}\left(\frac{\mathrm{i}}{2 \pi \alpha_{0}}+\frac{1}{2 \mathrm{~d}} \sum_{\mathrm{m}=-\infty}^{\infty} \frac{\alpha_{\mathrm{m}}}{\gamma_{\mathrm{m}}}\right)+\frac{1}{4 \pi} \frac{\frac{\mathrm{~d}^{2} \mathrm{z}}{\mathrm{ds}^{2}}}{\left(\frac{\mathrm{dx}}{\mathrm{ds}}\right)^{2}+\left(\frac{\mathrm{dz}}{\mathrm{ds}}\right)^{2}}$.
This interesting results established by Pavageau and Bousquet [44] is very important for the numerical applications, because it shows that $\mathcal{N}(\mathrm{s}, \mathrm{s})$ contains a series that canverges like $1 / \mathrm{m}^{3}$ (after adding the terms with negative and positive values of m ).

Let us notice that the second derivative of the profile function appears in eq.(4.92) and it clearly requires the continuity of the first derivative, i.e., the absence of edges.

The third kernel $\mathscr{K}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$, given by eq. (4.36):

$$
\begin{equation*}
\mathscr{K}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{~d}} \sum_{\mathrm{m}=-\infty,+\infty}\left[\operatorname{sgn}\left[\mathrm{z}(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right] \frac{\mathrm{dx}}{\mathrm{ds}}-\frac{\alpha_{\mathrm{m}}}{\gamma_{\mathrm{m}}^{+}} \frac{\mathrm{dz}}{\mathrm{ds}}\right] \mathrm{e}^{\operatorname{imK}\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)\right]+i \gamma_{\mathrm{m}}^{+}\left|\mathrm{z}(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right|} \tag{4.93}
\end{equation*}
$$

is deduced from $\mathcal{N}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$ by replacing the derivatives $\frac{\mathrm{dx}^{\prime}}{\mathrm{ds}^{\prime}}$ and $\frac{\mathrm{dz}{ }^{\prime}}{\mathrm{ds} '^{\prime}}$ by $\frac{\mathrm{dx}}{\mathrm{ds}}$ and $\frac{\mathrm{dz}}{\mathrm{ds}}$ and its study is quite similar. After after tedious calculations it can be shown that:

$$
\begin{equation*}
\mathscr{K}(\mathrm{s}, \mathrm{~s})=\lim _{\mathrm{s}^{\prime} \rightarrow \mathrm{s}} \mathcal{N}_{0}\left(\mathrm{~s}, \mathrm{~s}^{\prime}\right)=-\frac{\mathrm{dz}}{\mathrm{ds}}\left(\frac{\mathrm{i}}{2 \pi \alpha_{0}}+\frac{1}{2 \mathrm{~d}} \sum_{\mathrm{m}=-\infty}^{\infty} \frac{\alpha_{\mathrm{m}}}{\gamma_{\mathrm{m}}}\right)-\frac{1}{4 \pi} \frac{\frac{\mathrm{~d}^{2} \mathrm{z}}{\mathrm{ds}^{2}}}{\left(\frac{\mathrm{dx}}{\mathrm{ds}}\right)^{2}+\left(\frac{\mathrm{dz}}{\mathrm{ds}}\right)^{2}} . \tag{4.94}
\end{equation*}
$$

### 4.6.3. Integration of kernel singularities

Clearly, the asymptotic part of $\boldsymbol{\mathcal { G }}_{\infty}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$ in eq.(4.84) is singular when $\mathrm{s} \rightarrow \mathrm{s}^{\prime}$. After some calculations, it can be found that

$$
\begin{equation*}
\lim _{\mathrm{s} \rightarrow \mathrm{~s}^{\prime}} \boldsymbol{\mathcal { G }}_{\infty}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\frac{1}{4 \pi} \ln \left(\mathrm{~K}^{2}\left\{\left[\mathrm{x}(\mathrm{~s})-\mathrm{x}\left(\mathrm{~s}^{\prime}\right)\right]^{2}+\left[\mathrm{z}(\mathrm{~s})-\mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right]^{2}\right\}\right) \tag{4.95}
\end{equation*}
$$

Noting that $x(s)-x\left(s^{\prime}\right) \approx\left(s-s^{\prime}\right) \frac{d x}{d s}$ and $z(s)-z\left(s^{\prime}\right) \approx\left(s-s^{\prime}\right) \frac{d z}{d s}$, eq.(4.95) yields:

$$
\begin{equation*}
\lim _{\mathrm{s} \rightarrow \mathrm{~s}^{\prime}} \boldsymbol{\mathcal { O }}_{\infty}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \pi}\left\{\ln (2 \pi)+\frac{1}{2} \ln \left[\left(\frac{\mathrm{dx}}{\mathrm{ds}}\right)^{2}+\left(\frac{\mathrm{dz}}{\mathrm{ds}}\right)^{2}\right]+\ln \frac{\left|\mathrm{s}-\mathrm{s}^{\prime}\right|}{\mathrm{d}}\right\} . \tag{4.96}
\end{equation*}
$$

The first two terms represent regular parts that can be integrated by using the rectangular rule, as shown later. Unfortunately, $\ln \frac{\left|s-s^{\prime}\right|}{d}$ is not periodic and the rectangular rule is very poor when applied to nonperiodic functions. It is possible to overcome this difficulty [7] by considering another function defined on ( $0, \mathrm{~d}$ ):

$$
\begin{equation*}
\tilde{\mathcal{G}}_{\infty}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \pi}\left[\ln \frac{\left|\mathrm{~s}-\mathrm{s}^{\prime}\right|}{\mathrm{d}}+\ln \left(1-\frac{\left|\mathrm{s}-\mathrm{s}^{\prime}\right|}{\mathrm{d}}\right)\right], \quad \mathrm{s}, \mathrm{~s}^{\prime} \in(0, \mathrm{~d}) \tag{4.97}
\end{equation*}
$$

which has the same singularity as $\frac{1}{2 \pi} \ln \frac{\left|\mathrm{~s}-\mathrm{s}^{\prime}\right|}{\mathrm{d}}$ in the interval $(0, \mathrm{~d})$, because $\mathrm{d}-\left|\mathrm{s}-\mathrm{s}^{\prime}\right|$ never vanishes when $\mathrm{s}, \mathrm{s}^{\prime} \in(0, \mathrm{~d})$. The advantage of his new function is that it is continuous except on the singularity, and all its derivatives with respect to $s$ are the same on $s^{\prime}=0$ and $s^{\prime}=d$ and thus we can use the rectangular rule.

We perform the integration of $\boldsymbol{\mathcal { G }}_{\infty}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$ by setting:

$$
\begin{equation*}
\int_{0}^{\mathrm{s}_{\mathrm{d}}} \boldsymbol{\mathcal { G }}_{\infty}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \phi\left(\mathrm{s}^{\prime}\right) \mathrm{ds} \mathrm{~s}^{\prime}=\int_{0}^{\mathrm{s}_{\mathrm{d}}}\left[\boldsymbol{\mathcal { G }}_{\infty}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \phi\left(\mathrm{s}^{\prime}\right)-\tilde{\boldsymbol{\mathcal { F }}}_{\infty}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \phi(\mathrm{s})\right] \mathrm{ds}^{\prime}+\int_{0}^{\mathrm{s}_{\mathrm{d}}} \tilde{\boldsymbol{\mathcal { F }}}_{\infty}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \phi(\mathrm{s}) \mathrm{ds}^{\prime} . \tag{4.98}
\end{equation*}
$$

The term in the square bracket is non-singular and can be integrated using the rectangular rule, if we take into account that:
$\lim _{\mathrm{s}^{\prime} \rightarrow \mathrm{s}}\left[\boldsymbol{\mathcal { G }}_{\infty}\left(\mathrm{s}, \mathrm{s}^{\prime}\right) \phi\left(\mathrm{s}^{\prime}\right)-\tilde{\boldsymbol{\mathcal { F }}}_{\infty}\left(\mathrm{s}, \mathrm{s}^{\prime}\right) \phi(\mathrm{s})\right]=\frac{1}{2 \pi}\left\{\ln (2 \pi)+\frac{1}{2} \ln \left[\left(\frac{\mathrm{dx}}{\mathrm{ds}}\right)^{2}+\left(\frac{\mathrm{dz}}{\mathrm{ds}}\right)^{2}\right]\right\} \phi(\mathrm{s})$.

The second term in eq.(4.98) contains the singular part, but it can be integrated analytically:

$$
\begin{equation*}
\int_{0}^{\mathrm{s}_{\mathrm{d}}} \tilde{\boldsymbol{\mathcal { F }}}_{\infty}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \phi(\mathrm{s}) \mathrm{ds}^{\prime}=\phi(\mathrm{s}) \int_{0}^{\mathrm{s}_{\mathrm{d}}} \tilde{\boldsymbol{\mathcal { F }}}_{\infty}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \mathrm{ds}^{\prime}=-\frac{\mathrm{d}}{\pi} \phi(\mathrm{~s}) . \tag{4.100}
\end{equation*}
$$

In conclusion, the integration of the singular kernel is made by introduction at first $\boldsymbol{\mathcal { G }}_{\infty}$, which permits to define a series $\boldsymbol{\mathcal { G }}-\boldsymbol{\mathcal { G }}_{\infty}$ that is continuous and rapidly converging hence easily integrable by the use of the rectangular rule. Second, the integration of the term containing $\boldsymbol{\mathcal { F }}_{\infty}$ is performed by defining a new function $\tilde{\mathcal{G}}_{\infty}$, which has the same singularity as $\boldsymbol{\mathcal { G }}_{\infty}$, has the property of a periodic function, and can be analytically integrated.

### 4.6.4. Kernel singularities for highly conducting metals

When the conductivity of a metallic grating tends to infinity, the Green function in the metal tends to a delta function. This property is rather obvious: for very large conductivities, the field generated by a line current (delta function) placed in the metal decreases very rapidly since it is absorbed on very short distances. This behaviour have drastic consequences on the kernels of the integral equations dealing with metallic gratings, which are directly derived from the Green function: the two variable functions relative to the metallic part of the grating tend to delta functions as well. The integration of such functions through a point matching method requires more and more points of discretization around $\mathrm{s}^{\prime}=\mathrm{s}$ and, since s can take any value in the interval $\left(0, s_{d}\right)$ the integration and the inversion of the final linear system of equations (bearing in mind that its size is the total number of discretization points) becomes impossible. This remarks explains why the first attempts at implementing the integral equations on computers were not able to give any result for metallic gratings in the visible and infrared regions.

A very efficient way to overcome this difficulty is to apply an approach called local summation [7], using another form of the Green function [45]:

$$
\begin{equation*}
\mathcal{G}\left(\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}\right)=\frac{1}{4 \mathrm{i}} \sum_{\mathrm{m}} \mathrm{e}^{\mathrm{imd} \alpha_{0}} \mathrm{H}_{0}^{+}\left(\mathrm{k}\left|\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}\right|-\mathrm{md} \hat{\mathrm{x}}\right), \quad \overrightarrow{\mathrm{r}}=(\mathrm{x}, \mathrm{z}) \tag{4.101}
\end{equation*}
$$

with $\hat{\mathrm{x}}$ being the unit vector of the x axis. This form is the direct consequence of the fact that $\frac{1}{4 \mathrm{i}} \mathrm{H}_{0}^{+}(\mathrm{k}|\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}|)$ is the Green function of the Helmholtz equation:

$$
\begin{equation*}
\nabla^{2}\left[\frac{1}{4 \mathrm{i}} \mathrm{H}_{0}^{+}\left(\mathrm{k}\left|\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}\right|\right)\right]+\mathrm{k}^{2}\left[\frac{1}{4 \mathrm{i}} \mathrm{H}_{0}^{+}\left(\mathrm{k}\left|\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}\right|\right)\right]=\delta\left(\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}\right) . \tag{4.102}
\end{equation*}
$$

Since the pseudo-periodic Green function $\mathcal{G}\left(\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}\right)$ in vacuum is defined by:

$$
\begin{equation*}
\nabla^{2} \tilde{\mathcal{G}}\left(\overrightarrow{\mathrm{r}}, \overrightarrow{\mathrm{r}}^{\prime}\right)+\mathrm{k}^{2} \tilde{\boldsymbol{\mathcal { G }}}\left(\overrightarrow{\mathrm{r}}, \overrightarrow{\mathrm{r}}^{\prime}\right)=\sum_{\mathrm{m}=-\infty,+\infty} \mathrm{e}^{\mathrm{i} \alpha_{\mathrm{m}} \mathrm{~d}} \delta\left(\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}-\mathrm{md} \hat{\mathrm{x}}\right) \tag{4.103}
\end{equation*}
$$

it follows that $\mathcal{G}\left(\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}\right)$ is a sum of Green functions $\frac{1}{4 \mathrm{i}} \mathrm{H}_{0}^{+}\left(\mathrm{k}\left|\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}\right|-\mathrm{md} \hat{\mathrm{x}}\right)$ satisfying:

$$
\begin{gather*}
\nabla^{2}\left[\frac{1}{4 \mathrm{i}} \mathrm{e}^{\mathrm{i} \alpha_{\mathrm{n}} \mathrm{~d}} \mathrm{H}_{0}^{+}\left(\mathrm{k}\left|\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}-\mathrm{md} \hat{\mathrm{x}}\right|\right)\right]+\mathrm{k}^{2}\left[\frac{1}{4 \mathrm{i}} \mathrm{e}^{\mathrm{i} \alpha_{\mathrm{n}} \mathrm{~d}} \mathrm{H}_{0}^{+}\left(\mathrm{k}\left|\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}-\mathrm{md} \hat{\mathrm{x}}\right|\right)\right]=  \tag{4.104}\\
=\mathrm{e}^{\mathrm{i} \alpha_{\mathrm{n}} \mathrm{~d}} \delta\left(\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}-\operatorname{md} \hat{\mathrm{x}}\right) .
\end{gather*}
$$

Inserting the value of $\mathcal{\mathcal { G }}\left(\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}\right)$ given by eq. (4.101) inside the integral $\int_{0}^{\mathrm{s}_{\mathrm{d}}} \boldsymbol{\mathcal { G }}\left(\mathrm{s}, \mathrm{s}^{\prime}\right) \phi\left(\mathrm{s}^{\prime}\right) \mathrm{ds}^{\prime}$, then making the change of variable $\overrightarrow{\mathrm{r}}^{\prime}-\mathrm{md} \hat{\mathrm{x}}=\overrightarrow{\mathrm{r}}^{\prime \prime}$, and finally gathering the infinite sum of integrals on one period into a single integral from $-\infty$ to $+\infty$ yields:

$$
\begin{equation*}
\int_{0}^{\mathrm{s}_{\mathrm{d}}} \mathcal{F}\left(\mathrm{~s}, \mathrm{~s}^{\prime}\right) \phi\left(\mathrm{s}^{\prime}\right) \mathrm{ds} \mathrm{~s}^{\prime}=\frac{1}{4 \mathrm{i}} \int_{-\infty}^{\infty} \mathrm{H}_{0}^{+}\left(\mathrm{k}\left|\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}\right|\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}^{\prime}-\mathrm{x}\right)} \phi\left(\mathrm{s}^{\prime}\right) \mathrm{ds} \mathrm{~s}^{\prime} \tag{4.105}
\end{equation*}
$$

In the same way it can be shown that

$$
\begin{equation*}
\int_{0}^{\mathrm{s}_{\mathrm{d}}} \mathcal{N}\left(\mathrm{~s}, \mathrm{~s}^{\prime}\right) \psi\left(\mathrm{s}^{\prime}\right) \mathrm{ds}^{\prime}=\frac{\mathrm{ik}}{4} \int_{-\infty}^{\infty} \mathrm{H}_{1}^{+}\left(\mathrm{k}\left|\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}\right|\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}^{\prime}-\mathrm{x}\right)} \frac{\mathrm{z}^{\prime}-\mathrm{z}-\frac{\mathrm{dz}}{\mathrm{ds}^{\prime}}\left(\mathrm{x}^{\prime}-\mathrm{x}\right)}{\left|\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}\right|} \psi\left(\mathrm{s}^{\prime}\right) \mathrm{ds}^{\prime} . \tag{4.106}
\end{equation*}
$$

When the permittivity becomes very large in modulus, the functions $\mathcal{G}^{-}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$ and $\mathcal{N}^{-}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$ are obtained by replacing k by $\mathrm{kn}^{-}$in $\mathcal{\mathcal { E }}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$ and $\mathcal{N}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$. Since the value of $\mathrm{n}^{-}$is close to an imaginary number in the visible and infrared regions for usual metals (for example, $\mathrm{n}^{-}=$ $1.3+\mathrm{i} 7.11$ for aluminum at 650 nm$)$, the Hankel functions $\mathrm{H}_{0}^{+}(\mathrm{k}|\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}|)$ and $\mathrm{H}_{1}^{+}\left(\mathrm{k}\left|\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}\right|\right)$ become very close to the modified Bessel functions $K_{0}\left[k\left|n^{-}\left(\vec{r}-\vec{r}^{\prime}\right)\right|\right]$ and $K_{1}\left[\mathrm{k}\left|\mathrm{n}^{-}\left(\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}\right)\right|\right]$ (see [45]) and tend to delta functions when $\left|\mathrm{n}^{-}\right| \rightarrow \infty$. Thus, these functions vary much more rapidly than the unknown function $\phi$, which can be considered as a constant. Thus, remarking that when $\mathrm{x}^{\prime} \approx \mathrm{x}$

$$
\begin{gather*}
\left|\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}\right| \approx\left|\mathrm{x}-\mathrm{x}^{\prime}\right| \sqrt{1+\left(\frac{\mathrm{dz}}{\mathrm{dx}}\right)^{2}} \\
\mathrm{z}^{\prime}-\mathrm{z}-\left(\mathrm{x}^{\prime}-\mathrm{x}\right) \frac{\mathrm{dz}}{\mathrm{dx}} \approx-\frac{1}{2} \frac{\mathrm{~d}^{2} \mathrm{z}}{\mathrm{dx}^{2}}\left(\mathrm{x}^{\prime}-\mathrm{x}\right)^{2}  \tag{4.107}\\
\exp \left[\mathrm{i} \alpha_{0}\left(\mathrm{x}^{\prime}-\mathrm{x}\right)\right] \approx 1
\end{gather*}
$$

yields finally:

$$
\begin{equation*}
\int_{0}^{\mathrm{s}_{\mathrm{d}}} \mathcal{G}\left(\mathrm{~s}, \mathrm{~s}^{\prime}\right) \phi\left(\mathrm{s}^{\prime}\right) \mathrm{ds}^{\prime} \approx \frac{\phi(\mathrm{s})}{4 \mathrm{i}} \int_{-\infty}^{\infty} \mathrm{H}_{0}^{+}\left(\mathrm{k}\left|\mathrm{n}^{-}\left(\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}\right)\right|\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}^{\prime}-\mathrm{x}\right)} \mathrm{ds}^{\prime} \approx \frac{\phi(\mathrm{s})}{2 \mathrm{ik} \sqrt{1+\left(\frac{\mathrm{dz}}{\mathrm{dx}}\right)^{2}}} \tag{4.108}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{0}^{\mathrm{s}_{\mathrm{d}}} \mathcal{N}\left(\mathrm{~s}, \mathrm{~s}^{\prime}\right) \psi\left(\mathrm{s}^{\prime}\right) \mathrm{d} \mathrm{~s}^{\prime} \approx \frac{\frac{\mathrm{d}^{2} \mathrm{z}}{\mathrm{dx}^{2}} \psi(\mathrm{~s})}{4 \mathrm{ik}\left[1+\left(\frac{\mathrm{dz}}{\mathrm{dx}}\right)^{2}\right]^{3 / 2}} \tag{4.109}
\end{equation*}
$$

It is not surprising to note that in this approximation, the calculations of the integrals require neither summation of the kernels, nor integration of the singularities, nor matrix multiplication. As the kernels tend toward delta-distributions with amplitudes determining the coefficients in eqs.(4.108) and (4.109), their matrix representations tend to diagonal matrices.

Numerical results have shown that this simpler formulation not only successfully applies in the domain where the summation and integration processes defined in the previous sections fail, but also remains valid with a good accuracy (about $10^{-3}$ in relative value) in a large domain of metals and wavelengths. For example, with aluminum, this approach works in the visible, a domain in which the classical method of integration can be used as well (but with a greater computation time), wheras the local summation is necessary for metals in the far-infrared and microwaves domain. It is very important to notice that, using the local summation and assuming that $\left|\mathrm{n}^{-}\right| \rightarrow \infty$, it can be shown that the integral equation for metallic gratings described in appendix 4.B tends towards the integral equations obtained for perfectly conducting metals.

### 4.6.5. Problems of edges and non-analytical profiles

When the grating profile has edges or corners (in 2D case), fundamental difficulties appear. First, the uniqueness of the solution of the electromagnetic field is not ensured. The hypothesis of the integrability of the unknown function $u(x)$ in the integral equation is equivalent to the Meixner condition of integrability [46], although it is singular.

The second problem lies in the validity of the boundary condition of electromagnetic field on the grating profile. Indeed, on the edges, the normal and tangential direction on the profile are not defined in a unique manner. Moreover, when establishing these boundary conditions, the demonstration does not work if the surface has edges; However, these warnings are not dramatic. The demonstration of Archimede theorem is also questionable if the object presents edges. Does it exist any doubt about the validity of this theorem?

The third problem in the integral method lies in the process of integration in the vicinity of the edges. The kernels become discontinuous or even meaningless, depending on whether the point of calculation of the integral coincides with the edge point or not, see eq.(4.90) (4.92). Moreover, the integration can fail due to the eventual singularity of the unknown function on the edge [44].

To overcome the edge problem, there exist several approaches. The most direct one consists in replacing the actual profile $\mathrm{z}=\mathrm{f}(\mathrm{x})$ with its truncated Fourier representation:

$$
\begin{equation*}
\tilde{\mathrm{f}}_{\mathrm{M}}(\mathrm{x})=\sum_{\mathrm{m}=-\mathrm{M}}^{\mathrm{M}} \mathrm{f}_{\mathrm{m}} \mathrm{e}^{\mathrm{imKx}} \tag{4.110}
\end{equation*}
$$

The new profile has no edges and in most cases of ruled or holographic gratings, it mimics quite well the true profile. Numerical tests have shown a very good convergence of results with respect to the number $2 \mathrm{M}+1$ of Fourier terms and the number P of discretization points, provided the empirical rule $\mathrm{P}>4(2 \mathrm{M}+1)$. This method can be used to describe some profiles that are not represented by continuous functions, for example lamellar gratings, provided that

M is larger than 10. Unfortanately, there are several important cases of classical gratings and classes of relatively new types of gratings and periodic systems that cannot be treated using this simple approach.

The first problem covers the case of echelle gratings, working in grazing incidence in high and very high orders (see Chapter 1). The fact that the period is some tens or hundreds times larger than the wavelength $\lambda$, and that the working facet is quite steep (sometimes going up to $86^{\circ}$ groove angle), requires that its geometry is represented in the method with an error smaller than $\lambda / 20$. Simple but incorrect estimations show that this rule of thumb would be acceptable for numerical treatment: if the period is close to 50 wavelengths, we need about 1000 points of discretization and thus 250 Fourier harmonics, according to the rule $\mathrm{P}>4(2 \mathrm{M}$ +1 ). However, Gibbs phenomenon will significantly modify the form of the working facet, which could be almost vertical. In order to correctly describe the field on this facet, we need to have at least 5 to 10 points per wavelength along the the working facet, rather than along its projection on the x-axis. With a $85^{\circ}$ groove angle, the length of this facet is about 11 times its x-projection, in such a way that the number of discretization points must be multiplied by a factor 10. The number of Fourier components in the profile follows the same rule.


Figure 4.5. Schematic representation of an etched grating profile with non-analytical function description.


Figure 4.6. Grating made of inclined rectangular cross-section rods.
Another class of problems consists of unconventional geometries, like inverted slope grooves, obtained during groove etching technologies, as seen in Figure.4.5, or rod gratings, shown in Figure.4.6. We have noticed that the problem of vertical segments adds difficulties. For example, the problem of edges cannot be solved any longer by a Fourier expansion of the profile. It exists a possibility to simultaneously solve the two problems by introducing a curvilinear coordinate that follows the grating profile. As far as the integration is made along the profile, this curvilinear integration comes as a natural way of calculating the integrals in the integral equations. Moreover, adaptive meshing can be used to reduce the influence of edges.

In the real life, edges do not exist. On each edge there is at least one atom that has no edges, even though the light has a wavelength much larger than the atom dimensions. Edges of nuclear particles are not discussed even in the most exotic theories. Thus the idea is to replace the edges by arcs, where the partial derivatives can be well defined till the second order, which is sufficient for integral equations as remarked in section 4.6.2. An adaptive
density of discretization points gives the possibility to significantly increase the density of the points close to the initial edges, and not elsewhere. Let us, for example, consider a rod grating having a straight rectangular cross section. The segments 1 and 3 are parallel to Ox, the segments 2 and 4, to Oz. Let us assume the origin of the curvilinear coordinate at the bottomleft corner (figure 4.7).


Figure 4.7. Left: rounding of corners of a rectangular cross-section rod with side lengths a and $b$, together with $x, y$, and s coordinate lines. Right: schematic representation of different straight and arc segments to obtain the links between the Cartesian and the curvilinear coordinates.

The coordinate s starts at $\mathrm{x}=\mathrm{R}$ and $\mathrm{z}=0$, and follows the interface along its different parts:

1) the first horizontal segment, $0<\mathrm{s}<\mathrm{a}-2 \mathrm{R}$ :

$$
\begin{gather*}
\mathrm{x}(\mathrm{~s})=\mathrm{R}+\mathrm{s}, \\
\mathrm{z}(\mathrm{~s})=0 . \tag{4.111}
\end{gather*}
$$

Here the final value of $s$ is equal to $s_{1, \max }=a-2 R$ :
2) the circular rounding of the bottom-right corner, defined by the equation:

$$
\begin{equation*}
[x-(a-R)]^{2}+(z-R)^{2}=R^{2} \tag{4.112}
\end{equation*}
$$

Here, $0<\mathrm{s}-\mathrm{s}_{1, \max }<\mathrm{R} \frac{\pi}{2}$ :

$$
\begin{gather*}
s=R \varphi+s_{1, \max } \\
x=a-R+R \sin (\varphi)=a-R+R \sin \frac{s-s_{1, \max }}{R},  \tag{4.113}\\
z=R-R \cos (\varphi)=R-R \cos \frac{s-s_{1, \max }}{R} .
\end{gather*}
$$

Here, $s_{2, \max }=R \frac{\pi}{2}+s_{1, \text { max }}$
3) the next vertical segment at $x=a, 0<s-s_{2, \max }<b-2 R$ :

$$
\begin{gather*}
x=a \\
z=s-s_{2, \max }+R, \tag{4.114}
\end{gather*}
$$

and the same for the rest of the profile. The process is straightforward and needs an adapted application for each class of profiles. The advantage is that the derivatives $\frac{\mathrm{dx}}{\mathrm{ds}}$ and $\frac{\mathrm{dz}}{\mathrm{ds}}$ exist and are continuous everywhere on the profile. Thus, the second-order derivatives, which are required for the explicit summation and integration of the kernels, exist at least piesewisely, too. This could easily be checked at the point $(a-R, 0)$, for example. For $s<s_{1, \max }$, we use eqs.(4.111):

$$
\begin{equation*}
\frac{\mathrm{x}(\mathrm{~s})}{\mathrm{ds}}=1 ; \quad \frac{\mathrm{z}(\mathrm{~s})}{\mathrm{ds}}=0 . \tag{4.115}
\end{equation*}
$$

For $\mathrm{s}_{1, \text { max }}<\mathrm{s}<\mathrm{s}_{2, \text { max }}$ it is necessary to use eqs.(4.113):

$$
\begin{array}{ll}
\frac{d x(s)}{d s}=\cos \frac{s-s_{1, \max }}{R}=1 & \text { for } s=s_{1, \max }  \tag{4.116}\\
\frac{d z(s)}{d s}=\sin \frac{s-s_{1, \max }}{R}=0 & \text { for } s=s_{1, \max }
\end{array}
$$

The comparison of (4.115) and (4.116) points out the existence and continuity of the first derivatives.

It is worth noting that, in contrast with the adaptive spatial resolution used in several other methods (FEM, FDTD, RCW, coordinate transformation methods), here it is more convenient to call it adaptive profile resolution method, as far as it represents a 1D curvilinear coordinate adaptation.

Let us impose the requirement that the arc segments require $\mathrm{N}_{\text {arc }}$ times larger density points than on the straight segments. The entire length along s of the profile in figure.4.7 (left) is equal to $L_{\text {tot }}=2 a+2 b-8 R+2 \pi R$. The total number of discretization points is related to the length ot the segments, R , and to $\mathrm{N}_{\text {arc }}$. If the distance between the points along the straight segments is $\Delta$, it will be equal to $\Delta / N_{\text {arc }}$ on the arcs. Thus the total number of points for a single-rod per period is equal to:

$$
\begin{equation*}
\mathrm{P}=\frac{2 \mathrm{a}+2 \mathrm{~b}-8 \mathrm{R}+2 \pi \mathrm{RN}_{\mathrm{arc}}}{\Delta} \tag{4.117}
\end{equation*}
$$

In practice, unless some automatic scheme of determining the technical parameters of the computation is used, this equation determines $\Delta$ when the total number of integration points is chosen.

Starting from $\mathrm{s}=0$ in figure 4.7 (right), the abscissa values along s of the points on the interface are given by consecutively adding the values of the point number $j$ to the end values of the previous segment:

1) first horizontal segment, $0<\mathrm{s}<\mathrm{s}_{1, \max }: \mathrm{s}_{\mathrm{j}}=\mathrm{s}_{\mathrm{j}-1}+\Delta$,
2) circular rounding of the bottom-right corner, $\mathrm{s}_{1, \max }<\mathrm{s}<\mathrm{s}_{2, \max }: \mathrm{s}_{\mathrm{j}}=\mathrm{s}_{\mathrm{j}-1}+\frac{\Delta}{\mathrm{N}_{\operatorname{arc}}}$,
3) next vertical segment at $x=a, s_{2, \max }<s<b-2 R+s_{2, \max }: s_{j}=s_{j-1}+\Delta$,
4) so on to close the rod, then eventually going to another object inside the same grating period.

An additional improvement can be performed by making a smooth transfer from $\Delta$ to $\Delta / \mathrm{N}_{\text {arc }}$, i.e., to smoothly go from the density defined on the straight segments and on the arcs. This could be important for large-period systems with respect to the wavelength, if there
is a restriction in the total number of points in the profile discretization. If so, we can lay on the fact that the smaller the curvature of the segment (i.e., the larger its length), the smoother the behaviour of the kernels and of the unknown functions. Thus in the middle of a straight segment we will place points that are more distant to each other than close to the extremeties of the segments. Maystre has adapted such approach in his code Grating 2000, by defining a specific distance along the large segments starting from their edges, so that the density of the discretization points increases when approaching the ends of the segments. This approach gave the possibility to model echelle gratings working in very high diffraction orders ( 600 or more) in the ' 90 s , when no alternative approach was able to provide reliable results. The method was unbeatable for echelles, before Li and Chandezon [47] formulated an improvement of the coordinate transformation method to work for profiles with edges. However, the latter does not apply to rod gratings, or to profiles having the form as given in figure.4.5.

### 4.7. Examples of numerical results

All the results shown in this section are obtained using the code Grating 2000.

### 4.7.1. Sinusoidal perfectly conducting grating

Table 1 shows the efficiencies in the two non-evanescent orders ( $-1^{\text {st }}$ and $0^{\text {th }}$ ) of a sinusoidal perfectly conducting grating of period 600 nm and height 180 nm (from the bottom to the top of the groove) illuminated under incidence angle $30^{\circ}$ and wavelength 600 nm . In this case, the $-1^{\text {st }}$ order is scattered with an angle of scattering (measured anticlockwise, in contrast with incident angle) of $-30^{\circ}$, which entails that it propagates just in the direction which is the opposite to that of the incident wave (this is called Littrow mounting by specialists of gratings).

The number of discretisation points is P and the series included in the kernel are summed from -M to +M . The symbol $\sum \rho_{\mathrm{m}}$ denotes the sum of the two efficiencies, the energy balance being satisfied when $\sum \rho_{\mathrm{m}}=1$.

|  | TE polarization |  |  | TM polarization |  |  |
| :--- | :---: | :---: | :---: | :---: | :---: | :---: |
| $\mathrm{P}, \mathrm{M}$ | $\rho_{-1}$ | $\rho_{-0}$ | $\sum \rho_{\mathrm{m}}$ | $\rho_{-1}$ | $\rho_{-0}$ | $\sum \rho_{\mathrm{m}}$ |
| 4,2 | 0.4658 | 0.5437 | 1.0095 | 0.9466 | 0.0514 | 0.9980 |
| 6,3 | 0.4703 | 0.5288 | 0.9991 | 0.9581 | 0.0411 | 0.9992 |
| 25,10 | 0.4669 | 0.5336 | 1.0005 | 0.9579 | 0.0421 | 1.0000 |
| 50,20 | 0.4659 | 0.5334 | 0.9993 | 0.9579 | 0.0421 | 1.0000 |
| 110,50 | 0.4659 | 0.5341 | 1.0000 | 0.9579 | 0.0421 | 1.0000 |

Table 4.1. The sinusoidal perfectly conducting grating.
It is worth noting that a precision better than 0.01 is reached as soon as $\mathrm{P}>4$ and $\mathrm{M}>2$ !. A precision of $10^{-3}$ needs $\mathrm{P}>50$ and $\mathrm{M}>20$.

### 4.7.2. Echelette perfectly conducting grating

The echelette grating is a grating with triangular groove (figure 4.8). The blaze angle b (angle of the large facet with the x axis) is equal to $30^{\circ}$ and the apex angle A (angle between the two facets) to $90^{\circ}$. The other parameters are the same as in section 7.1. It must be noticed that the
incidence angle and the blaze angle are equal, which entails that the incident wavevector is orthogonal to the large facet. In these conditions, it can be shown that for TM polarization, the grating problem can be solved in closed form: the efficiency in the $-1^{\text {st }}$ order is equal to unity while in the $0^{\text {th }}$ order it vanishes [48].

The demonstration is straightforward: the sum of the incident wave and of a plane wave with unit amplitude propagating in the opposite direction satisfies all the conditions of the boundary value problem stated in section 4.3. The reader can notice that this sum satisfies the Helmholtz equation. In addition, this sum of two plane waves propagating in opposite directions constitute an interference system that presents white areas and dark lines. The maximum of white lines coincide with the large facets of the grating, and on these lines, the derivative of the field (thus the normal derivative) vanishes. The normal derivative with respect to the small facet vanishes as well since the field is invariant in the normal direction. At the first glance, this property is obvious since the field is "reflected" by the large facets, or in other words, the scattering phenomenon reduces in that case to a simple reflection phenomenon. This reasoning fails since the same phenomenon is not observed for TE polarization: it is dangerous to invoke reflection phenomena on the large facets when the width of these facets has the same order of magnitude as the wavelength of the light! The concentration of the incident energy in a single order is called 'blazing effect' and the theorem of Marechal and Stroke is the origin of the name 'blazed gratings' given sometimes to ruled gratings.


Figure 4.8: A ruled grating

|  | TE polarization |  |  | TM polarization |  |  |
| :--- | ---: | ---: | :---: | ---: | ---: | :---: |
| $\mathrm{P}, \mathrm{M}$ | $\rho_{-1}$ | $\rho_{-0}$ | $\sum \rho_{\mathrm{m}}$ | $\rho_{-1}$ | $\rho_{-0}$ | $\sum \rho_{\mathrm{m}}$ |
| 6,3 | 0.6531 | 0.4451 | 1.0982 | 1.4452 | 0.0012 | 1.4464 |
| 25,10 | 0.5838 | 0.4123 | 0.9961 | 0.9976 | 0.0001 | 0.9977 |
| 50,20 | 0.5838 | 0.4123 | 0.9961 | 0.9976 | 0.0001 | 0.9977 |
| 110,50 | 0.5929 | 0.4055 | 0.9984 | 0.9984 | 0.0000 | 0.9984 |
| 250,100 | 0.5932 | 0.4035 | 0.9967 | 0.9989 | 0.0000 | 0.9989 |

Table 4.2. The echelette perfectly conducting grating and the Marechal and Stroke theorem.
Table 4.2 shows that the convergence is significantly slower than for the sinusoidal grating, due to the edges. Nevertheless, a precision of about 0.01 is obtained when $\mathrm{P}>25$ and $\mathrm{M}>10$. With the same values, the Marechal and Stroke theorem is satisfied with a precision better than 0.003 . The computation time is always less than 1 second on a PC computer except for the last line, for which 2 seconds are required.

### 4.7.3. Lamellar perfectly conducting grating

The profile of a lamellar grating is shown in figure 4.9. The widths of the hole and of the bump are denoted by $t$ and $b$ and the height by $h$. In this example, $b=t=300 \mathrm{~nm}$, and $h=180$ nm.


Figure 4.9. A lamellar grating

|  | TE polarization |  |  | TM polarization |  |  |
| :--- | ---: | ---: | :--- | :--- | :--- | :--- |
| $\mathrm{P}, \mathrm{M}$ | $\rho_{-1}$ | $\rho_{-0}$ | $\sum \rho_{\mathrm{m}}$ | $\rho_{-1}$ | $\rho_{-0}$ | $\sum \rho_{\mathrm{m}}$ |
| 6,3 | 0.7770 | 0.6999 | 1.4769 | 23.66 | 5.09 | 28.75 |
| 25,10 | 0.3490 | 0.6394 | 0.9884 | 0.8293 | 0.1724 | 1.0016 |
| 50,20 | 0.3347 | 0.6569 | 0.9915 | 0.8191 | 0.1705 | 0.9896 |
| 150,70 | 0.3279 | 0.6659 | 0.9938 | 0.8201 | 0.1718 | 0.9919 |
| 250,100 | 0.3288 | 0.6733 | 1.0021 | 0.8214 | 0.1725 | 0.9939 |

Table 4.3. The echelette perfectly conducting grating and the Marechal and Stroke theorem
The main conclusion to draw from Table 4.3 is that the convergence for lamellar gratings is even slower than for echelette gratings. This is not surprising if we notice that the number of edges is multiplied by 2 . A precision of about 0.02 is obtained when $\mathrm{P}>50$ and M $>20$. The results for $\mathrm{P}=6$ and $\mathrm{M}=3$ are aberrant, specially for TM polarization.

### 4.7.4. Aluminum sinusoidal grating in the near infrared

We consider a sinusoidal aluminum grating with period $\mathrm{d}=400 \mathrm{~nm}$, a height $\mathrm{h}=100 \mathrm{~nm}$, illuminated with incidence angle $10^{\circ}$ and wavelength 300 nm . With these parameters, three plane waves $\left(-1^{\text {st }}, 0^{\text {th }}\right.$ and $\left.+1^{\text {st }}\right)$ are reflected. The optical index of aluminum at 300 nm is equal to $4.2+$ i 21.5 . We give in Table 4.4 the efficiency in the $-1^{\text {st }}$ order and the sum of the three efficiencies.

|  | TE polarization |  | TM polarization |  |
| :--- | :---: | :---: | :---: | :---: |
| $\mathrm{P}, \mathrm{M}$ | $\rho_{-1}$ | $\sum \rho_{\mathrm{m}}$ | $\rho_{-1}$ | $\sum \rho_{\mathrm{m}}$ |
| 6,3 | 0.5205 | 0.9582 | 0.4367 | 0.9618 |
| 10,4 | 0.5201 | 0.9649 | 0.4325 | 0.9521 |
| 30,13 | 0.5203 | 0.9655 | 0.4321 | 0.9518 |
| 100,45 | 0.5204 | 0.9655 | 0.4320 | 0.9518 |

Table 4.4. The aluminum sinusoidal grating

Table 4.4 shows a very good convergence of the results, similar to the convergence observed for sinusoidal perfectly conducting gratings, thanks to the local summation of the two variable functions of the kernel derived from the Green function in aluminum.

### 4.7.5. Buried echelette silver grating in the visible.

The buried silver grating is shown in figure 4.10. A symmetric echelette silver grating with period 900 nm has been covered by a dielectric of index 1.5, the maximum depth e of dielectric being equal to 800 nm . This grating is illuminated in normal incidence .by a plane wave of wavelength 600 nm . The index of siver for this wavelength is equal to $0.006+\mathrm{i} 3.75$. Table 4.5 gives the efficiency in the $0^{\text {th }}$ order and the total of efficiencies of the three reflected orders ( $-1^{\text {st }}, 0^{\text {th }}$ and $\left.+1^{\text {st }}\right)$.


Figure 4.10. A symmetric silver grating coated by dielectric

|  | TE polarization |  | TM polarization |  |
| :--- | :--- | :---: | :---: | :---: |
| $\mathrm{P}, \mathrm{M}$ | $\rho_{0}$ | $\sum \rho_{\mathrm{m}}$ | $\rho_{0}$ | $\sum \rho_{\mathrm{m}}$ |
| 10,4 | 0.4892 D | 0.9830 | 0.6316 | 1.0622 |
| 30,13 | 0.5164 | 0.9651 | 0.5665 | 0.9339 |
| 100,45 | 0.5153 | 0.9600 | 0.6062 | 0.9178 |
| 200,90 | 0.5153 | 0.9593 | 0.6168 | 0.9153 |

Table 4.5. The buried silver grating
The convergence is slower than in the preceding case and the results for $\mathrm{P}=10$ and $\mathrm{M}=$ 4 are not correct, specially for TM polarization. For TE polarization, a convergence of the results with a precision better than 0.006 is obtained for $\mathrm{P}=30$ and $\mathrm{M}=13$. This is not so for TM polarization, in which it is necessary to reach $\mathrm{P}=100$ and $\mathrm{M}=45$ to get a precision of the order of 0.003 . Here, the method of local summation is not used, but this fact does not explain the slower convergence since the modulus of the optical index is not large. The main reason can be found in the edges of the profile. The slower convergence for TM polarization is rather general for metallic gratings, due to the existence of plasmon resonances on the grating surface.

### 4.7.6. Dielectric rod grating.

The grating is made of dielectric elliptic rods with optical index 1.4, width $\mathrm{w}=600 \mathrm{~nm}$ and height $\mathrm{h}=400 \mathrm{~nm}$ (figure 4.11). The period is equal to 800 nm . It is illuminated with incidence $20^{\circ}$ by a plane wave with wavelength 600 nm . Two orders ( $-1^{\text {st }}$ and $0^{\text {th }}$ ) are reflected and transmitted. We give in Table 4.6 the efficiency in the $0^{\text {th }}$ transmitted order and
the sum of efficiencies of the 4 scattered orders, which should be equal to 1 for a perfect energy balance.


Figure 4.11. A dielectric rod grating

|  | TE polarization |  | TM polarization |  |
| :--- | :--- | :--- | :--- | :--- |
| $\mathrm{P}, \mathrm{M}$ | $\tau_{0}$ | $\sum \rho_{\mathrm{m}}+$ | $\tau_{0}$ | $\sum \rho_{\mathrm{m}}+$ |
|  |  | $\sum \tau_{\mathrm{m}}$ |  | $\sum \tau_{\mathrm{m}}$ |
| 10,4 | 0.1824 D | 0.9139 | 0.7893 | 1.2489 |
| 30,13 | 0.2163 | 0.9961 | 0.8161 | 1.0054 |
| 100,45 | 0.2073 | 1.0005 | 0.8187 | 1.0004 |
| 200,90 | 0.2070 | 1.0001 | 0.8187 | 1.0001 |

Table4. 6. The dielectric rod grating with elliptic rods
The precision for $\mathrm{P}=30$ and $\mathrm{M}=13$ is equal to 0.01 and for $\mathrm{P}=100$ and $\mathrm{M}=90$, it reaches 0.0004 .

### 4.7.7. Flat perfectly conducting rod grating

The grating is similar to that of figure 4.11 , but its height is very small ( 8 nm ). In order to obtain a significant reflection, the dielectric has been replaced by a perfectly conducting material

|  | TE polarization |  | TM polarization |  |
| :---: | :---: | :---: | :---: | :---: |
| $\mathrm{P}, \mathrm{M}$ | $\tau_{0}$ | $\sum \rho_{\mathrm{m}}+$ | $\tau_{0}$ | $\sum \rho_{\mathrm{m}}+$ |
|  |  | $\sum \tau_{\mathrm{m}}$ |  | $\sum \tau_{\mathrm{m}}$ |
| 100,45 | 0.0598 | 1.1414 | 0.8940 | 0.9759 |
| 200,90 | 0.0634 | 1.0086 | 0.0292 | 0.9887 |
| 300,140 | 0.0639 | 1.0006 | 0.1059 | 0.9996 |
| 400,190 | 0.0639 | 1.0000 | 0.1119 | 1.0000 |

Table4.7. The flat perfectly conducting rod grating.
The convergence is very slow and it is necessary to reach $\mathrm{P}=300$ and $\mathrm{M}=140$ to obtain a precision better than 0.006 , the energy balance being satisfied with a precision better than $10^{-3}$. The reason must be found in the small height of the rods, as explained in section 6.1. The functions included in the kernel of the integral equation become very large in
modulus and very small in width for two points located on both sides of the rod for the same abscissa, thus the integration needs a large density of discretization points.

## Appendix 4.A. Mathematical bases of the integral theory

## 4.A.1. Presentation of the mathematical problem

We consider (figure 4.1) a function $F(x, z)=\left\{\begin{array}{l}F^{+}(x, z) \text { in } V^{+}, \\ F^{-}(x, z) \text { in } V^{-},\end{array}\right.$which satisfies the following conditions:

- it is pseudo-periodic along the x axis:

$$
\begin{equation*}
F(x+d, z)=F(x, z) \exp \left(i \alpha_{0} d\right) \tag{4.118}
\end{equation*}
$$

- it satisfies a Helmholtz equation:

$$
\begin{equation*}
\nabla^{2} \mathrm{~F}^{ \pm}+\mathrm{k}^{2}\left(\mathrm{n}^{ \pm}\right)^{2} \mathrm{~F}^{ \pm}=0 \quad \text { in } \mathrm{V}^{ \pm} \tag{4.119}
\end{equation*}
$$

- it satisfies a radiation condition for $\mathrm{z} \rightarrow \pm \infty$.

The aim of this appendix is to use the second Green's identity and basic theorems on boundary value problems in order to find an integral expression of this function and to analyze the properties of this integral expression. We will deduce the basic keys for writing an integral equation from a boundary value problem.

## 4.A.2. Calculation of the Green function

The first step of the calculation is to find the pseudo-periodic elementary solutions $\mathcal{G}^{+}(\overrightarrow{\mathrm{r}})$ and $\mathcal{G}^{-}(\overrightarrow{\mathrm{r}})$ of the two Helmholtz equations condensed in eq. (4.119) (with constants $\mathrm{k}^{2}\left(\mathrm{n}^{+}\right)^{2}$ for $\boldsymbol{\mathcal { G }}^{+}$and with constant $\mathrm{k}^{2}\left(\mathrm{n}^{+}\right)^{2}$ for $\boldsymbol{\mathcal { G }}^{-}$, which satisfy the radiation conditions for $\mathrm{z} \rightarrow \pm \infty$ and the following equations :

$$
\begin{gather*}
\nabla^{2} \mathcal{G}^{+}(\overrightarrow{\mathrm{r}})+\mathrm{k}^{2}\left(\mathrm{n}^{+}\right)^{2} \mathcal{G}^{+}(\overrightarrow{\mathrm{r}})=\sum_{\mathrm{m}=-\infty,+\infty} \mathrm{e}^{\mathrm{i} \alpha_{\mathrm{m}} \mathrm{~d}} \delta(\overrightarrow{\mathrm{r}}-\mathrm{md} \hat{\mathrm{x}}), \text { in } \mathrm{V}^{+} \text {and } \mathrm{V}^{-}  \tag{4.120}\\
\nabla^{2} \mathcal{G}^{-}(\overrightarrow{\mathrm{r}})+\mathrm{k}^{2}\left(\mathrm{n}^{-}\right)^{2} \mathcal{G}^{-}(\overrightarrow{\mathrm{r}})=\sum_{\mathrm{m}=-\infty,+\infty} \mathrm{e}^{\mathrm{i} \alpha_{\mathrm{m}} \mathrm{~d}} \delta(\overrightarrow{\mathrm{r}}-\mathrm{md} \hat{\mathrm{x}}), \text { in } \mathrm{V}^{+} \text {and } \mathrm{V}^{-}  \tag{4.121}\\
\boldsymbol{\mathcal { G }}^{ \pm}(\overrightarrow{\mathrm{r}}+\mathrm{d} \hat{\mathrm{x}})=\boldsymbol{\mathcal { G }}^{ \pm}(\overrightarrow{\mathrm{r}}) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{~d}} \tag{4.122}
\end{gather*}
$$

with:

$$
\begin{equation*}
\alpha_{\mathrm{m}}=\alpha_{0}+\mathrm{mK}, \quad \mathrm{~K}=\frac{2 \pi}{\mathrm{~d}} \tag{4.123}
\end{equation*}
$$

and $\hat{x}$ being the unit vector of the x axis.
We must emphasize that, in contrast with eq. (4.119), in which $\mathrm{F}^{+}$and $\mathrm{F}^{-}$do not satisfy the same Helmholtz equation, each Green function satisfies a unique Helmholtz equation in the entire space, with constant $\mathrm{k}^{2}\left(\mathrm{n}^{+}\right)^{2}$ for $\mathcal{\mathcal { G }}^{+}$and with constant $\mathrm{k}^{2}\left(\mathrm{n}^{+}\right)^{2}$ for $\boldsymbol{\mathcal { G }}^{-}$.

After expanding the periodic function $\mathcal{E}^{ \pm} \mathrm{e}^{-\mathrm{i} \alpha_{0} \mathrm{x}}$ in Fourier series, then multiplying the Fourier series by $\mathrm{e}^{\mathrm{i} 0_{0} \mathrm{x}}$, it can be deduced that:

$$
\begin{equation*}
\mathcal{G}^{ \pm}(\mathrm{x}, \mathrm{z})=\sum_{\mathrm{m}=-\infty,+\infty} \mathrm{G}^{ \pm}{ }_{\mathrm{m}}(\mathrm{z}) \mathrm{e}^{\mathrm{i} \alpha_{\mathrm{m}} \mathrm{x}} \tag{4.124}
\end{equation*}
$$

On the other hand, the right-hand member of eq. (4.120), called Dirac comb, can also be expanded in series:

$$
\begin{equation*}
\sum_{\mathrm{m}=-\infty,+\infty} \mathrm{e}^{\mathrm{i} \alpha_{\mathrm{m}} \mathrm{~d}} \delta(\overrightarrow{\mathrm{r}}-\mathrm{nd} \hat{\mathrm{x}})=\frac{1}{\mathrm{~d}} \delta(\mathrm{z}) \sum_{\mathrm{m}} \mathrm{e}^{\mathrm{i} \alpha_{\mathrm{m}} \mathrm{x}} . \tag{4.125}
\end{equation*}
$$

Introducing equations (4.124) and (4.125) in equation (4.120), multiplying by $\mathrm{e}^{-\mathrm{i} \alpha_{0} \mathrm{x}}$ then identifying the coefficients of the Fourier series yield:

$$
\begin{equation*}
\mathrm{G}_{\mathrm{m}}^{ \pm} "(\mathrm{z})+\left(\gamma_{\mathrm{m}}^{ \pm}\right)^{2} \mathrm{G}_{\mathrm{m}}^{ \pm}(\mathrm{z})=\frac{1}{\mathrm{~d}} \delta(\mathrm{z}) \tag{4.126}
\end{equation*}
$$

with $\left(\gamma_{\mathrm{m}}^{ \pm}\right)^{2}=\mathrm{k}^{2}\left(\mathrm{n}^{ \pm}\right)^{2}-\alpha_{\mathrm{m}}^{2}$.
For $\mathrm{z} \neq 0$, eq.(4.126) becomes the well-known one-dimension propagation equation in a homogeneous media without sources and have for solutions exponentials. We are searching for plane waves that satisfy the radiation condition for $\mathrm{z} \rightarrow \pm \infty$, thus:

$$
G_{m}^{ \pm}(z)= \begin{cases}A_{m}^{ \pm} \exp \left[i \gamma_{m}^{ \pm} z\right], & \text { if } z>0  \tag{4.127}\\ B_{m}^{ \pm} \exp \left[-i \gamma_{m}^{ \pm} z\right], & \text { if } z<0\end{cases}
$$

Distribution theory proves that the solution $\mathrm{G}_{\mathrm{m}}^{ \pm}$of eq.(4.126) is a continuous function of z and that its derivative has a jump equal to $1 / \mathrm{d}$ at $\mathrm{z}=0$. These two conditions allow one to find the unknown amplitudes:

$$
\begin{gather*}
\mathrm{A}_{\mathrm{m}}^{ \pm}=\mathrm{B}_{\mathrm{m}}^{ \pm} \\
\mathrm{i} \gamma_{m}^{ \pm} \mathrm{A}_{\mathrm{m}}^{ \pm}-\left(-\mathrm{i} \gamma_{\mathrm{m}}^{ \pm} \mathrm{B}_{\mathrm{m}}^{ \pm}\right)=\frac{1}{\mathrm{~d}} \tag{4.128}
\end{gather*}
$$

and thus:

$$
\begin{equation*}
\mathrm{G}_{\mathrm{m}}^{ \pm}(\mathrm{z})=\frac{1}{2 \mathrm{i} \gamma_{\mathrm{m}}^{ \pm} \mathrm{d}} \exp \left(\mathrm{i} \gamma_{\mathrm{m}}^{ \pm}|\mathrm{z}|\right) \tag{4.129}
\end{equation*}
$$

The final form of $\mathcal{G}^{ \pm}$becomes:

$$
\begin{equation*}
\boldsymbol{\mathcal { G }}^{ \pm}(\overrightarrow{\mathrm{r}})=\frac{1}{2 \mathrm{id}} \sum_{\mathrm{m}} \frac{1}{\gamma_{\mathrm{m}}^{ \pm}} \exp \left[\mathrm{i} \alpha_{\mathrm{m}}^{\mathrm{x}}+\mathrm{i} \gamma_{\mathrm{m}}^{ \pm}|\mathrm{z}|\right] \tag{4.130}
\end{equation*}
$$

Once the source is not in the origin, the Green functions is the function $\mathcal{\mathcal { G }}^{ \pm}\left(\overrightarrow{\mathrm{r}}-\overrightarrow{\mathrm{r}}^{\prime}\right)$.
Notice that $\boldsymbol{\mathcal { G }}^{ \pm}$symbolizes the Green functions.

## 4.A.3. Integral expression

Now, we apply the second Green theorem in order to find the expression of the function $\mathrm{F}^{ \pm}$. First, we consider the expression of $\mathrm{F}^{-}$in $\mathrm{V}^{-}$:

$$
\begin{equation*}
\mathrm{F}^{-}(\mathrm{x}, \mathrm{z})=\int_{\mathrm{S}_{\mathrm{T}}} \boldsymbol{\mathcal { G }}^{-}\left(\mathrm{x}-\mathrm{x}^{\prime}, \mathrm{z}-\mathrm{z}^{\prime}\right) \frac{\mathrm{dF}^{-}\left(\mathrm{x}^{\prime}, \mathrm{z}^{\prime}\right)}{\mathrm{dN}_{\mathrm{S}}} \mathrm{ds}^{\prime}-\int_{\mathrm{S}_{\mathrm{T}}} \frac{\mathrm{~d} \boldsymbol{\mathcal { G }}^{-}\left(\mathrm{x}-\mathrm{x}^{\prime}, \mathrm{z}-\mathrm{z}^{\prime}\right)}{\mathrm{dN}_{\mathrm{S}}} \mathrm{~F}^{-}\left(\mathrm{x}^{\prime}, \mathrm{z}^{\prime}\right) \mathrm{ds} \mathrm{~s}^{\prime} \tag{4.131}
\end{equation*}
$$

with the normal $\vec{N}_{S}$ being oriented towards the exterior of $V^{-}, x=x(s), x^{\prime}=x\left(s^{\prime}\right)$, and similar expressions for z and $\mathrm{z}^{\prime}$. The curve $\mathrm{S}_{\mathrm{T}}$ includes four parts: the vertical lines $\mathrm{S}_{\mathrm{L}}$ at $\mathrm{x}=$ 0 and $S_{R}$ at $x=d$, the horizontal segment $S_{H}$ at $z=z_{H}<0$ (figure 4.12), and, finally, one period of $S$. The variable s' denotes the curvilinear abscissa on $\mathrm{S}_{\mathrm{T}}^{-}$, with origine being located at the origin of the Cartesian coordinates.


Figure 4.12. Application of the second Green theorem.
The pseudo-periodicity in $\mathrm{x}^{\prime}$ of $\mathrm{F}^{-}\left(\mathrm{x}^{\prime}, \mathrm{z}^{\prime}\right)$ (and of $\mathrm{dF}^{-}\left(\mathrm{x}^{\prime}, \mathrm{z}^{\prime}\right) / \mathrm{dN}_{\mathrm{S}}$ ) is opposite to that of $\quad \boldsymbol{\mathcal { G }}^{-}\left(\mathrm{x}-\mathrm{x}^{\prime}, \mathrm{z}-\mathrm{z}^{\prime}\right) \quad\left(\right.$ or of $\left.\quad \mathrm{d} \boldsymbol{\mathcal { G }}^{-}\left(\mathrm{x}-\mathrm{x}^{\prime}, \mathrm{z}-\mathrm{z}^{\prime}\right) / \mathrm{dN}_{\mathrm{S}}\right)$, which entails that $\mathcal{G}^{-}\left(\mathrm{x}-\mathrm{x}^{\prime}, \mathrm{z}-\mathrm{z}^{\prime}\right) \frac{\mathrm{dF}^{-}\left(\mathrm{x}^{\prime}, \mathrm{z}^{\prime}\right)}{\mathrm{dN}_{\mathrm{S}}}$ and $\frac{\mathrm{d} \boldsymbol{\mathcal { G }}^{-}\left(\mathrm{x}-\mathrm{x}^{\prime}, \mathrm{z}-\mathrm{z}^{\prime}\right)}{\mathrm{dN}_{\mathrm{S}}} \mathrm{F}^{-}\left(\mathrm{x}^{\prime}, \mathrm{z}^{\prime}\right)$ are periodic. Furthermore, taking into account the orientation of the normal on $S_{R}$ and $S_{L}, \frac{\mathrm{dF}^{-}\left(x^{\prime}, z^{\prime}\right)}{\mathrm{dN}_{\mathrm{s}}}=-\frac{\mathrm{dF}^{-}\left(\mathrm{x}^{\prime}, \mathrm{z}^{\prime}\right)}{\mathrm{dx}}$ and $\frac{d \boldsymbol{\rho}^{-}\left(x-x^{\prime}, z-z^{\prime}\right)}{\mathrm{dN}_{\mathrm{S}}}=-\frac{\mathrm{d} \boldsymbol{\mathcal { F }}^{-}\left(\mathrm{x}-\mathrm{x}^{\prime}, \mathrm{z}-\mathrm{z}^{\prime}\right)}{\mathrm{dx}}$ on $\mathrm{S}_{\mathrm{L}}$, while $\frac{\mathrm{dF}^{-}\left(\mathrm{x}^{\prime}, \mathrm{z}^{\prime}\right)}{\mathrm{dN}_{\mathrm{s}}}=+\frac{\mathrm{dF}^{-}\left(\mathrm{x}^{\prime}, \mathrm{z}^{\prime}\right)}{\mathrm{dx}}$ on $S_{R}$. Thus, the integrals on $S_{R}$ and $S_{L}$ in eq. (4.131) cancel each other. On $S_{H}$, ds' and dx' identify and the integral takes the form:

$$
\begin{gather*}
\int_{\mathrm{S}_{\mathrm{H}}}\left[\boldsymbol{\mathcal { G }}^{-}\left(\mathrm{x}-\mathrm{x}^{\prime}, \mathrm{z}-\mathrm{z}^{\prime}\right) \frac{\mathrm{dF}^{-}\left(\mathrm{x}^{\prime}, \mathrm{z}^{\prime}\right)}{\mathrm{dN}_{\mathrm{s}}}-\frac{\mathrm{d} \boldsymbol{\mathcal { G }}^{-}\left(\mathrm{x}-\mathrm{x}^{\prime}, \mathrm{z}-\mathrm{z}^{\prime}\right)}{\mathrm{dN}_{\mathrm{S}}} \mathrm{~F}^{-}\left(\mathrm{x}^{\prime}, \mathrm{z}^{\prime}\right)\right] \mathrm{ds} \\
=\frac{1}{2 \mathrm{~d}} \sum_{\mathrm{m}=-\infty}^{\infty} \mathrm{e}^{\mathrm{i} \gamma_{\mathrm{m}}^{-}\left(\mathrm{z}-\mathrm{z}_{\mathrm{H}}\right)} \int_{0}^{\mathrm{d}} \mathrm{e}^{\mathrm{imK}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)}\left[\frac{1}{\mathrm{i} \gamma_{\mathrm{m}}^{-}} \phi_{\mathrm{H}}^{-}\left(\mathrm{s}^{\prime}\right)+\psi_{\mathrm{H}}^{-}\left(\mathrm{s}^{\prime}\right)\right] \mathrm{dx} \mathrm{x}^{\prime}  \tag{4.132}\\
\quad=\frac{1}{2} \sum_{\mathrm{m}=-\infty}^{\infty} \mathrm{e}^{\mathrm{i} \gamma_{\mathrm{m}}^{-}\left(\mathrm{z}-\mathrm{z}_{\mathrm{H}}\right)+i m K x}\left(\frac{1}{\mathrm{i} \gamma_{\mathrm{m}}^{-}} \phi_{\mathrm{H}, \mathrm{~m}}^{-}+\psi_{\mathrm{H}, \mathrm{~m}}^{-}\right)
\end{gather*}
$$

where $\phi_{\mathrm{H}, \mathrm{m}}^{-}$and $\psi_{\mathrm{H}, \mathrm{m}}^{-}$are the $\mathrm{m}^{\text {th }}$ Fourier components of the periodic functions $\phi_{\mathrm{H}}^{-}\left(\mathrm{s}^{\prime}\right)=\frac{\mathrm{dF}^{-}\left(\mathrm{x}^{\prime}, z^{\prime}\right)}{d N_{\mathrm{s}}} \mathrm{e}^{-\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}$ and $\psi_{\mathrm{H}}^{-}\left(\mathrm{s}^{\prime}\right)=\mathrm{F}^{-}\left(\mathrm{x}^{\prime}, \mathrm{z}^{\prime}\right) \mathrm{e}^{-\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}$ defined on $\mathrm{S}_{\mathrm{H}}$. Here, we can take advantage of the fact that the segment is parallel to the x axis and lies outside the groove region, which entails that $\phi_{\mathrm{H}, \mathrm{m}}^{-}$and $\psi_{\mathrm{H}, \mathrm{m}}^{-}$are related through the plane wave expansion valid in the homogeneous region (see chapter 2):

$$
\begin{equation*}
F^{-}=\sum_{m=-\infty,+\infty} t_{m} e^{i \alpha_{m} x-i \gamma_{m}^{-} z} \quad \text { if } z<0 \tag{4.133}
\end{equation*}
$$

which yields:

$$
\begin{equation*}
\psi_{\mathrm{H}, \mathrm{~m}}^{-}=-\mathrm{i} \gamma_{\mathrm{m}}^{-} \phi_{\mathrm{H}, \mathrm{~m}}^{-} . \tag{4.134}
\end{equation*}
$$

This relation allows us to cancel the integral along $\mathrm{S}_{\mathrm{H}}$ and thus the values of $\mathrm{F}^{-}(\mathrm{x}, \mathrm{z})$ can be determined by an integral along a single groove of S :

$$
\begin{equation*}
\mathrm{F}^{-}(\mathrm{x}, \mathrm{z})=\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \boldsymbol{\mathcal { G }}^{-}\left(\mathrm{x}-\mathrm{x}^{\prime}, \mathrm{z}-\mathrm{z}^{\prime}\right) \frac{\mathrm{dF}^{-}\left(\mathrm{x}^{\prime}, \mathrm{z}^{\prime}\right)}{\mathrm{dN}_{\mathrm{s}}} \mathrm{ds}^{\prime}-\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \frac{\mathrm{~d} \boldsymbol{\mathcal { G }}^{-}\left(\mathrm{x}-\mathrm{x}^{\prime}, \mathrm{z}-\mathrm{z}^{\prime}\right)}{\mathrm{dN}_{\mathrm{S}}} \mathrm{~F}^{-}\left(\mathrm{x}^{\prime}, \mathrm{z}^{\prime}\right) \mathrm{ds} \mathrm{~s}^{\prime} \tag{4.135}
\end{equation*}
$$

with $s_{d}$ being the curvilinear abscissa of the point of $S$ of abscissa $d$, the origin of curvilinear abscissa being the origin of the Cartesian coordinates system.

Similar considerations apply to $\mathrm{F}^{+}(\mathrm{x}, \mathrm{z})$, so that, after elementary calculations:

$$
\begin{equation*}
\mathrm{F}^{ \pm}(\mathrm{x}, \mathrm{z})= \pm \int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{d}}\left[\mathcal{G}^{ \pm}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \phi^{ \pm}\left(\mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}+\boldsymbol{N}^{ \pm}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \psi^{ \pm}\left(\mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}\right] \mathrm{ds}^{\prime}( \tag{4.136}
\end{equation*}
$$

with:

$$
\begin{equation*}
\mathcal{G}^{ \pm}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right)=\boldsymbol{\mathcal { G }}^{ \pm}\left(\mathrm{x}, \mathrm{z}, \mathrm{x}\left(\mathrm{~s}^{\prime}\right), \mathrm{z}\left(\mathrm{~s}^{\prime}\right)\right)=\frac{1}{2 \mathrm{id}} \sum_{\mathrm{m}=-\infty}^{\infty} \frac{1}{\gamma_{\mathrm{m}}^{ \pm}} \mathrm{e}^{\mathrm{imK}\left(\mathrm{x}-\mathrm{x}^{\prime}\left(\mathrm{s}^{\prime}\right)\right)+\mathrm{i} \gamma_{\mathrm{m}}^{ \pm}\left|\mathrm{z}-\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right|} \tag{4.137}
\end{equation*}
$$

$$
\begin{gather*}
\mathcal{N}^{ \pm}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right)=-\frac{\partial \boldsymbol{\mathcal { G }}^{ \pm}}{\partial \mathrm{N}_{\mathrm{s}}}\left(\mathrm{x}, \mathrm{z}, \mathrm{x}^{\prime}\left(\mathrm{s}^{\prime}\right), \mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right)= \\
=\frac{1}{2 \mathrm{~d}} \sum_{\mathrm{m}=-\infty}^{\infty}\left[\frac{\mathrm{dx}}{} \mathrm{ds}^{\prime} \operatorname{sgn}\left(\mathrm{z}-\mathrm{z}^{\prime}\right)-\frac{\alpha_{\mathrm{m}}}{\gamma_{\mathrm{m}}^{ \pm}} \frac{\mathrm{dz}}{} \mathrm{ds}^{\prime}\right.  \tag{4.138}\\
\mathrm{e}^{\mathrm{imK}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)+\mathrm{i} \gamma_{\mathrm{m}}^{ \pm}\left|\mathrm{z}-\mathrm{z}^{\prime}\right|}
\end{gather*}
$$

with $\psi^{ \pm}\left(s^{\prime}\right)=F^{-}\left(x^{\prime}, z^{\prime}\right) e^{-i \alpha_{0} x^{\prime}}$ and $\phi^{ \pm}\left(s^{\prime}\right)=\frac{\mathrm{dF}^{-}\left(\mathrm{x}^{\prime}, \mathrm{z}^{\prime}\right)}{\mathrm{dN}_{\mathrm{s}}} \mathrm{e}^{-\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}$ being defined on the grating profile S. Defining the periodic function $\mathrm{U}^{ \pm}(\mathrm{x}, \mathrm{z})=\mathrm{F}^{ \pm}(\mathrm{x}, \mathrm{z}) \mathrm{e}^{-\mathrm{i} \alpha_{0} \mathrm{x}}$ yields:

$$
\begin{equation*}
\mathrm{U}^{ \pm}(\mathrm{x}, \mathrm{z})= \pm \int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{ \pm}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \phi^{ \pm}\left(\mathrm{s}^{\prime}\right)+\boldsymbol{N}^{ \pm}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \psi^{ \pm}\left(\mathrm{s}^{\prime}\right)\right] \mathrm{ds}^{\prime} \tag{4.139}
\end{equation*}
$$

## 4.A.4. Equation of compatibility

In this section, we establish a crucial property of the integral theory of gratings, which unfortunately is ignored in most of the reference books of Electromagnetics. With this aim, it is necessary to point out a fundamental property of the expression of $\mathrm{U}^{ \pm}(\mathrm{x}, \mathrm{z})$ given by eqs. (4.13) and (4.139). Let us suppose that we introduce in eq. (4.136) arbitrary periodic functions $\tilde{\phi}^{-}\left(s^{\prime}\right)$ and $\tilde{\psi}^{-}\left(s^{\prime}\right)$. Since we have not introduced the actual physical values $\phi^{-}\left(s^{\prime}\right)$ and $\psi^{-}\left(\mathrm{s}^{\prime}\right)$ of these functions, we cannot expect to obtain the actual value of $\mathrm{F}^{-}(\mathrm{x}, \mathrm{z})$, but another function $\tilde{\mathrm{F}}^{-}(\mathrm{x}, \mathrm{z})$. More important, if the point $\mathrm{P}^{-}(\mathrm{x}, \mathrm{z})$ of $\mathrm{V}^{-}$tends towards a point $P_{S}$ of curvilinear abscissa s located on the profile (figure 4.12), the limit of $\tilde{F}^{-}(x, z)$ below the profile, denoted $\lim _{-}\left\{\tilde{\mathrm{F}}^{-}(\mathrm{x}, \mathrm{z})\right\}$ is not equal to $\tilde{\Psi}^{-}\left(\mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}$. The same remark can be made for the normal derivative $\lim _{-}\left\{\frac{d \tilde{F}^{-}(x, z)}{d N_{s}}\right\}$, which is different from $\tilde{\phi}^{-}\left(s^{\prime}\right) e^{i \alpha_{0} x^{\prime}}$. In order to understand this surprising property, it is necessary to give two results which can be demonstrated using the theory of distributions. Let us give these two fundamental results without demonstration.

1. The integral expression of $\tilde{\mathrm{F}}^{-}(\mathrm{x}, \mathrm{z})$ inside $\mathrm{V}^{-}$satisfies the Helmholtz equation in $\mathrm{V}^{-}$like the actual field. What happens to the same integral expression of $\tilde{\mathrm{F}}^{-}(\mathrm{x}, \mathrm{z})$, but now calculated in region $\mathrm{V}^{+}$? Denoting by $\tilde{\mathrm{F}}_{\mathrm{ext}}^{-}(\mathrm{x}, \mathrm{z})$ the function equal to $\tilde{\mathrm{F}}^{-}(\mathrm{x}, \mathrm{z})$ in $\mathrm{V}^{-}$and to this integral extension in $\mathrm{V}^{+}$, it is easy to verify that $\tilde{\mathrm{F}}_{\text {ext }}^{-}(\mathrm{x}, \mathrm{z})$ satisfies in the entire space the Helmholtz equation satisfied by $\tilde{\mathrm{F}}^{-}(\mathrm{x}, \mathrm{z})$ in $\mathrm{V}^{-}$, with constant $\mathrm{k}^{2}\left(\mathrm{n}^{-}\right)^{2}$.
2. The jumps $\quad \lim _{+}\left\{\tilde{\mathrm{F}}_{\text {ext }}^{-}(\mathrm{x}, \mathrm{z})\right\}-\lim _{-}\left\{\tilde{\mathrm{F}}_{\mathrm{ext}}^{-}(\mathrm{x}, \mathrm{z})\right\} \quad$ and $\quad \lim _{+}\left\{\frac{\mathrm{d} \tilde{\mathrm{F}}_{\text {ext }}^{-}(\mathrm{x}, \mathrm{z})}{\mathrm{dN}}\right\}-$ $\lim _{-}\left\{\frac{d \tilde{F}_{\text {ext }}^{-}(x, z)}{d N_{s}}\right\}$ of $\tilde{F}_{\text {ext }}^{-}(x, z)$ and of its normal derivative across $S$ (difference between the values in $\mathrm{V}^{+}$and in $\mathrm{V}^{-}$) are respectively equal to $-\tilde{\psi}^{-}\left(\mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}$ and $-\tilde{\phi}^{-}\left(\mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}$. The conclusion to draw from these properties is that the limit of $\tilde{\mathrm{F}}^{-}(\mathrm{x}, \mathrm{z})$ and of its normal derivative on $S$ are equal to $\tilde{\psi}^{-}\left(s^{\prime}\right) e^{i \alpha_{0} x^{\prime}}$ and $\tilde{\phi}^{-}\left(s^{\prime}\right) e^{\mathrm{i} \alpha_{0} x^{\prime}}$ in one case only: if $\tilde{\mathrm{F}}_{\text {ext }}^{-}(\mathrm{x}, \mathrm{z})$ vanishes throughout $\mathrm{V}^{+}$. Indeed, if this property is satisfied, the jumps are nothing but $-\lim _{-}\left\{\tilde{\mathrm{F}}_{\text {ext }}^{-}(\mathrm{x}, \mathrm{z})\right\} \quad$ and $\quad-\lim _{-}\left\{\frac{\mathrm{dF}_{\text {ext }}^{-}(\mathrm{x}, \mathrm{z})}{\mathrm{dN}}\right\}$, thus $\quad \lim _{-}\left\{\tilde{\mathrm{F}}_{\text {ext }}^{-}(\mathrm{x}, \mathrm{z})\right\}=\tilde{\psi}^{-}(\mathrm{s}) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}} \quad$ and $\lim _{-}\left\{\frac{d \tilde{F}_{\mathrm{ext}}^{-}(\mathrm{x}, \mathrm{z})}{\mathrm{dN}}\right\}=\tilde{\phi}_{\mathrm{s}}(\mathrm{s}) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}}$. This case occurs if the values of $\tilde{\Psi}^{-}\left(\mathrm{s}^{\prime}\right)$ and $\tilde{\phi}^{-}\left(\mathrm{s}^{\prime}\right)$ introduced in the integral expression are equal to the actual physical values $\psi^{-}\left(\mathrm{s}^{\prime}\right)$ and $\phi^{-}\left(s^{\prime}\right)$.

This result is not surprising for the specialist of boundary value problems. Indeed, using the second Green theorem in $\mathrm{V}^{-}$, we introduce in the integral expression of the field the limit values of both the field and of its normal derivative below S . In the domain of boundary value problems, it is well known that if a function must satisfy a Helmholtz equation in $\mathrm{V}^{-}$and a radiation condition for $\mathrm{z} \rightarrow-\infty$, one cannot impose the limit values of both this function and its normal derivative on $S$. In fact, we can impose either the limit values of this function or that of its normal derivative on $S$ : in both cases, the solution of the boundary value problem exists and is unique. Unfortunately, it does not exist any tool of applied mathematics which enables one to express the field in an integral form including either the limit values of the field, or that of its normal derivative on S : a prior solution of an integral equation is required, which is much more difficult.

On the other hand, if both the actual values of the field and of its normal derivative are known, the field can directly be expressed in an integral form through the second Green theorem, without any integral equation, and this is why the second Green theorem is considered as a basic tool of the integral mathods of scattering. It must be emphasized that in that case, we know a priori that the the limit values of the field and of its normal derivative are the actual ones, thus that they are compatible, which is not the case for arbitrarily chosen limits.

This fundamental property shows that when the second Green theorem is used with unknown values of the limits of the field and of its normal derivative, we must impose to these limits an equation of compatibility. This compatibility is satisfied if we impose to the limit $\lim _{-}\left\{\tilde{\mathrm{F}}_{\text {ext }}^{-}(\mathrm{x}, \mathrm{z})\right\}$ of the integral expression of $\tilde{\mathrm{F}}_{\text {ext }}^{-}(\mathrm{x}, \mathrm{z})$ to be equal to $\tilde{\psi}^{-}(\mathrm{s}) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}}$ or if we impose to the limit $\lim _{+}\left\{\tilde{\mathrm{e}}_{\text {ext }}^{-}(\mathrm{x}, \mathrm{z})\right\}$ of the integral expression of $\tilde{\mathrm{F}}_{\text {ext }}^{-}(\mathrm{x}, \mathrm{z})$ to be equal to zero, these two conditions being equivalent. Indeed, if $\lim _{+}\left\{\tilde{\mathrm{F}}_{\text {ext }}^{-}(\mathrm{x}, \mathrm{z})\right\}=0$, the expression of $\tilde{\mathrm{F}}_{\text {ext }}^{-}(\mathrm{x}, \mathrm{z})$ in $\mathrm{V}^{+}$satisfies a Helmholtz equation, a radiation condition at infinity and its limit
on S vanishes. The obvious solution of this boundary value problem is $\tilde{\mathrm{F}}_{\mathrm{ext}}^{-}(\mathrm{x}, \mathrm{z})=0$ in $\mathrm{V}^{+}$, and we have seen that the solution of this boundary value problem is unique, thus it is the solution. The consequence is that $\lim _{+}\left\{\frac{d \tilde{F}_{\text {ext }}^{-}(x, z)}{d N_{s}}\right\}=0$, thus $\lim _{-}\left\{\frac{d \tilde{F}_{\text {ext }}^{-}(x, z)}{d N_{s}}\right\}=$ $\tilde{\phi}^{-}(s) e^{i \alpha_{0} x}$.

In order to implement this condition of compatibility, we have to express $\lim _{+}\left\{\frac{d \tilde{F}_{\text {ext }}^{-}(x, z)}{d N_{s}}\right\}$ or $\lim _{-}\left\{\frac{d \tilde{F}_{\text {ext }}^{-}(x, z)}{d N_{s}}\right\}$. We can use a first equation:

$$
\begin{equation*}
\lim _{+}\left\{\tilde{\mathrm{F}}_{\text {ext }}^{-}(\mathrm{x}, \mathrm{z})\right\}-\lim _{-}\left\{\tilde{\mathrm{F}}_{\mathrm{ext}}^{-}(\mathrm{x}, \mathrm{z})\right\}=-\tilde{\Psi}^{-}(\mathrm{s}) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}} \tag{4.140}
\end{equation*}
$$

or equivalently:

$$
\begin{equation*}
\lim _{+}\left\{\tilde{\mathrm{U}}_{\mathrm{ext}}^{-}(\mathrm{x}, \mathrm{z})\right\}-\lim _{-}\left\{\tilde{\mathrm{U}}_{\mathrm{ext}}^{-}(\mathrm{x}, \mathrm{z})\right\}=-\tilde{\Psi}^{-}(\mathrm{s}) \tag{4.141}
\end{equation*}
$$

In order to find a second equation, we can consider eq. (4.139) which gives the expression of $\mathrm{U}^{-}(\mathrm{x}, \mathrm{z})=\mathrm{F}^{-}(\mathrm{x}, \mathrm{z}) \mathrm{e}^{-\mathrm{i} \alpha_{0} \mathrm{x}}$. If z is fixed, this expression is a Fourier series in x , which is discontinuous on $S$. It is well known that the value on $S$ of this Fourier series is the average value of the limits on both sides of S , thus:

$$
\begin{equation*}
\lim _{+}\left\{\tilde{\mathrm{U}}_{\text {ext }}^{-}(\mathrm{x}, \mathrm{z})\right\}+\lim _{-}\left\{\tilde{\mathrm{U}}_{\text {ext }}^{-}(\mathrm{x}, \mathrm{z})\right\}=-2 \int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{-}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \tilde{\phi}^{-}\left(\mathrm{s}^{\prime}\right)+\boldsymbol{N}^{-}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \tilde{\psi}^{-}\left(\mathrm{s}^{\prime}\right)\right] \mathrm{ds} \mathrm{~s}^{\prime} \tag{4.142}
\end{equation*}
$$

with $\boldsymbol{\mathcal { G }}^{-}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$ and $\boldsymbol{\mathcal { N }}^{-}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$ being the integral expressions of $\boldsymbol{\mathcal { G }}^{-}\left(\mathrm{x}, \mathrm{z}, \mathrm{s}{ }^{\prime}\right)$ and $\mathcal{N}^{-}(\mathrm{x}, \mathrm{z}, \mathrm{s})$ when the point of coordinates $x, z$ becomes a point of $S$ of curvilinear abscissa s:

$$
\begin{align*}
& \mathcal{G}^{-}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\boldsymbol{\mathcal { G }}^{-}\left(\mathrm{x}(\mathrm{~s}), \mathrm{z}(\mathrm{~s}), \mathrm{s}^{\prime}\right)  \tag{4.143}\\
& \mathcal{N}^{-}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\mathcal{N}^{-}\left(\mathrm{x}(\mathrm{~s}), \mathrm{z}(\mathrm{~s}), \mathrm{s}^{\prime}\right) \tag{4.144}
\end{align*}
$$

From eqs. (4.141) and (4.142), we deduce the limits of $\tilde{U}^{-}$on both sides of S , and we derive the equation of compatibility in $\mathrm{V}^{-}$:

$$
\begin{equation*}
\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{-}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \tilde{\phi}^{-}\left(\mathrm{s}^{\prime}\right)+\boldsymbol{\mathcal { N }}^{-}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \tilde{\psi}^{-}\left(\mathrm{s}^{\prime}\right)\right] \mathrm{ds}^{\prime}+\frac{\tilde{\psi}^{-}(\mathrm{s})}{2}=0 \tag{4.145}
\end{equation*}
$$

Achieving a similar calculation for $\mathrm{V}^{+}$yields the general equation of compatibility in $\mathrm{V}^{ \pm}$:

$$
\begin{align*}
& \int_{s^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{ \pm}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \tilde{\phi}^{ \pm}\left(\mathrm{s}^{\prime}\right)+\mathcal{N}^{ \pm}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \tilde{\Psi}^{ \pm}\left(\mathrm{s}^{\prime}\right)\right] \mathrm{ds} \mathrm{~s}^{\prime} \mp \frac{\tilde{\Psi}^{ \pm}(\mathrm{s})}{2}=0  \tag{4.146}\\
& \boldsymbol{\mathcal { G }}^{ \pm}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{id}} \sum_{\mathrm{m}=-\infty}^{\infty} \frac{1}{\gamma_{\mathrm{m}}^{ \pm}} \mathrm{e}^{\mathrm{imK}\left(\mathrm{x}(\mathrm{~s})-\mathrm{x}^{\prime}\left(\mathrm{s}^{\prime}\right)\right)+\mathrm{i} \gamma_{\mathrm{m}}^{ \pm}\left|\mathrm{z}(\mathrm{~s})-\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right|} \tag{4.147}
\end{align*}
$$

$$
\begin{equation*}
\mathcal{N}^{ \pm}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{~d}} \sum_{\mathrm{m}=-\infty}^{\infty}\left[\frac{\mathrm{dx}{ }^{\prime}}{\mathrm{ds} \mathrm{~s}^{\prime}} \operatorname{sgn}\left(\mathrm{z}(\mathrm{~s})-\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right)-\frac{\alpha_{\mathrm{m}}}{\gamma_{\mathrm{m}}^{ \pm}} \frac{\mathrm{dz}}{\mathrm{ds}}\right] \quad \mathrm{e}^{\mathrm{im} K\left(\mathrm{x}(\mathrm{~s})-\mathrm{x}^{\prime}\left(\mathrm{s}^{\prime}\right)\right)+\mathrm{i} \mathrm{\gamma} \gamma_{\mathrm{m}}^{ \pm}\left|\mathrm{z}(\mathrm{~s})-\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right|} \tag{4.148}
\end{equation*}
$$

It can be shown [7] that $\mathcal{\mathcal { G }}^{ \pm}$( $\mathrm{s}, \mathrm{s}$ ') has an integrable logarithmic singularity (it behaves like $a_{0}+a_{1} \log \left|s-s^{\prime}\right|$ when $s^{\prime} \rightarrow s, a_{0}$ and $a_{1}$ complex numbers) which can be taken into account in the integral by removing the singularity $a_{1} \log \left|s-s^{\prime}\right|$ from $\mathcal{G}^{ \pm}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$ then by integrating it in closed form. At the first glance, $\boldsymbol{N}^{ \pm}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$ cannot be continuous, due to the discontinuity of $\operatorname{sgn}\left(z(s)-z^{\prime}\left(s^{\prime}\right)\right)$ for $s=s^{\prime}$. In fact, a careful analysis of this function around $\mathrm{s}=\mathrm{s}^{\prime}$ shows that it is continuous and that its limit when $\mathrm{s}^{\prime} \rightarrow \mathrm{s}$ can be expressed in closed form [44,7], as stated in section 4.6.3.

## 4.A.5. Generalized compatibility

In the physical problem, the total field in $\mathrm{V}^{+}$includes the incident field. Here, we give an extension of the compatibility equation to the total field which allows a significant simplification of the use of integral theory.

$$
\Gamma_{\Gamma_{V^{-}}^{-}=\phi^{\mathrm{i}}}^{\mathrm{V}^{+}=\mathrm{F}^{\mathrm{i}}} \mathrm{~F}^{\mathrm{V}^{+}+\mathrm{F}^{+}+\mathrm{F}^{\mathrm{i}}}
$$

Figure 4.13. generalized: generalized compatibility
We define (figure 4.13) a function $\Gamma=\left\{\begin{array}{c}\Gamma^{+}=\mathrm{F}^{\mathrm{T}+}=\mathrm{F}^{+}+\mathrm{F}^{\mathrm{i}} \text { in } \mathrm{V}^{+} \text {, } \\ \Gamma^{-}=\mathrm{F}^{\mathrm{i}} \text { in } \mathrm{V}^{-} \text {. }\end{array}\right.$
$\mathrm{F}^{+}$being the field scattered in $\mathrm{V}^{+}$by a grating illuminated by the incident field $\mathrm{F}^{\mathrm{i}} . \Gamma^{-}$is the expression of the incident field in $\mathrm{V}^{-}$. Thus, in contrast with the preceding section, the function considered in this section does not satisfy a radiation condition at infinity in $\mathrm{V}^{+}$. According to eq. (4.136), the value of $\mathrm{F}^{+}$is given by:

$$
\begin{align*}
& \mathrm{F}^{+}(\mathrm{x}, \mathrm{z})=\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\boldsymbol{\mathcal { G }}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \phi^{+}\left(\mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}\right.  \tag{4.149}\\
&\left.\quad+\mathcal{N}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \psi^{+}\left(\mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}\right] \mathrm{ds}^{\prime}
\end{align*}
$$

thus $\Gamma^{+}$can be written:

$$
\begin{align*}
& \Gamma^{+}(\mathrm{x}, \mathrm{z})=\mathrm{F}^{\mathrm{i}}(\mathrm{x}, \mathrm{y})+\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \phi^{+}\left(\mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} 0_{0} \mathrm{x}^{\prime}}\right.  \tag{4.150}\\
&\left.\quad+\mathcal{N}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \psi^{+}\left(\mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}\right] \mathrm{ds}^{\prime}
\end{align*}
$$

We denote by $\Psi^{+}=\psi^{+}+\psi^{\mathrm{i}}$ and $\Phi^{+}=\phi^{+}+\phi^{\mathrm{i}}$ the limits of $\Gamma^{+}$on S and its normal derivative respectively. Introducing these values in eq. (4.149) yields:

$$
\begin{align*}
\mathrm{F}^{+}(\mathrm{x}, \mathrm{z}) & =\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{d}}\left[\mathcal{G}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)}\left(\Phi^{+}\left(\mathrm{s}^{\prime}\right)-\phi^{\mathrm{i}}\left(\mathrm{~s}^{\prime}\right)\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}\right.  \tag{4.151}\\
& \left.+\mathcal{N}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)}\left(\Psi^{+}\left(\mathrm{s}^{\prime}\right)-\psi^{\mathrm{i}}\left(\mathrm{~s}^{\prime}\right)\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}\right] \mathrm{ds}^{\prime}
\end{align*}
$$

Thus, the integral expression of the total field $\Gamma^{+}$is given by:

$$
\begin{align*}
& \Gamma^{+}(\mathrm{x}, \mathrm{z})=\mathrm{F}^{\mathrm{i}}(\mathrm{x}, \mathrm{z})+\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\boldsymbol{\mathcal { G }}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)}\left(\Phi^{+}\left(\mathrm{s}^{\prime}\right)-\phi^{\mathrm{i}}\left(\mathrm{~s}^{\prime}\right)\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}\right. \\
& \left.\quad+\boldsymbol{N}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)}\left(\Psi^{+}\left(\mathrm{s}^{\prime}\right)-\psi^{\mathrm{i}}\left(\mathrm{~s}^{\prime}\right)\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}\right] \mathrm{ds}^{\prime} \tag{4.152}
\end{align*}
$$

or, gathering the terms containing the incident field:

$$
\begin{align*}
& \Gamma^{+}(\mathrm{x}, \mathrm{z})= \\
& \mathrm{F}^{\mathrm{i}}(\mathrm{x}, \mathrm{z})-\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{d}}\left[\mathcal{G}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \phi^{\mathrm{i}}\left(\mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}+\boldsymbol{N}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \psi^{\mathrm{i}}\left(\mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} 0_{0} \mathrm{x}^{\prime}}\right] \mathrm{ds} \mathrm{~s}^{\prime}  \tag{4.153}\\
& +\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \Phi^{+}\left(\mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}+\boldsymbol{\mathcal { N }}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \Psi^{+}\left(\mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}\right] \mathrm{ds} s^{\prime} .
\end{align*}
$$

In order to simplify this equation, we consider the part of the integral containing the incident field, the middle line in the equation above. This part is equal to 0 . Indeed, we know that $\phi^{\mathrm{i}}\left(\mathrm{s}^{\prime}\right)$ and $\psi^{\mathrm{i}}\left(\mathrm{s}^{\prime}\right)$ are compatible in $\mathrm{V}^{-}$since they are derived from the actual values of the incident field in $\mathrm{V}^{-}$(figure 4.13). We have shown in the preceding section that the integral expression of such a function vanishes in $\mathrm{V}^{+}$. Remarking that $\mathcal{\mathcal { G }}^{-}\left(\mathrm{x}, \mathrm{z}, \mathrm{s}^{\prime}\right)=\boldsymbol{\mathcal { G }}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{s}^{\prime}\right)$ and $\mathcal{N}^{-}\left(\mathrm{x}, \mathrm{z}, \mathrm{s}^{\prime}\right)=\boldsymbol{\mathcal { N }}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{s}^{\prime}\right)$ since $\Gamma^{+}$and $\Gamma^{-}$satisfy the same Helmholtz equation (with constant $\left.\mathrm{k}^{2}\left(\mathrm{n}^{+}\right)^{2}\right)$, we can write than in $\mathrm{V}^{+}$:

$$
\begin{equation*}
\int_{s^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \phi^{\mathrm{i}}\left(\mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}+\boldsymbol{N}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \psi^{\mathrm{i}}\left(\mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}\right] \mathrm{ds}^{\prime}=0 \tag{4.154}
\end{equation*}
$$

Thus finally the expression of $\Gamma^{+}$is given by:

$$
\begin{align*}
& \Gamma^{+}(\mathrm{x}, \mathrm{z})=\mathrm{F}^{\mathrm{i}}(\mathrm{x}, \mathrm{z})+\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \Phi^{+}\left(\mathrm{s} \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}\right.  \tag{4.155}\\
& \left.\quad+\mathcal{N}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \Psi^{+}\left(\mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}\right] \mathrm{ds}^{\prime}
\end{align*}
$$

The result is that the expression of the total field in $\mathrm{V}^{+}$can be obtained from its limit value $\Psi^{+}$on $S$ and from its normal derivative $\Phi^{+}$, by adding the incident field to the integral expression deduced from the Green theorem.

The generalized compatibility equation is derived from the compatibility equations for $\mathrm{F}^{+}$in $\mathrm{V}^{+}$and for $\mathrm{F}^{\mathrm{i}}$ in $\mathrm{V}^{-}$:

$$
\begin{align*}
& \int_{s^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \phi^{+}\left(\mathrm{s}^{\prime}\right)+\boldsymbol{N}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \psi^{+}\left(\mathrm{s}^{\prime}\right)\right] \mathrm{ds} \mathrm{~s}^{\prime}-\frac{\psi^{+}(\mathrm{s})}{2}=0  \tag{4.156}\\
& \int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{\mathcal { G }}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \phi^{\mathrm{i}}\left(\mathrm{~s}^{\prime}\right)+\mathcal{N}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \psi^{\mathrm{i}}\left(\mathrm{~s}^{\prime}\right)\right] \mathrm{ds}^{\prime}+\frac{\psi^{\mathrm{i}}(\mathrm{~s})}{2}=0 \tag{4.157}
\end{align*}
$$

By adding these two equations, we deduce that:

$$
\begin{equation*}
\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \Phi^{+}\left(\mathrm{s}^{\prime}\right)+\mathcal{N}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \Psi^{+}\left(\mathrm{s}^{\prime}\right)\right] \mathrm{ds}^{\prime}+\psi^{\mathrm{i}}=\frac{\Psi^{+}(\mathrm{s})}{2} \tag{4.158}
\end{equation*}
$$

The remarkable result is that the compatibility equation given in the preceding section can be extended to the total field: the left-hand side of eq. (4.158) represents the expression of the total field on S and the right-hand, one half of its limit on S . Thus the generalized compatibility condition can be stated in the following way:

The compatibility condition for a field in $\mathrm{V}^{+}$including incident and/or diffracted waves and satisfying the two conditions:

- it is pseudo-periodic along the x axis:

$$
\begin{equation*}
\mathrm{F}(\mathrm{x}+\mathrm{d}, \mathrm{z})=\mathrm{F}(\mathrm{x}, \mathrm{z}) \exp \left(\mathrm{i} \alpha_{0} \mathrm{~d}\right) \tag{4.159}
\end{equation*}
$$

- it satisfies a Helmholtz equation:

$$
\begin{equation*}
\nabla^{2} \mathrm{~F}^{ \pm}+\mathrm{k}^{2}\left(\mathrm{n}^{ \pm}\right)^{2} \mathrm{~F}^{ \pm}=0 \quad \text { in } \mathrm{V}^{ \pm} \tag{4.160}
\end{equation*}
$$

can be stated in the following way: The value on $S$ of the total field, obtained by adding the incident wave to the integral expression deduced from the second Green theorem (but in which the limits are those of the total field) is equal to the half of its limit value on $S$.

A similar generalized compatibility condition can be obtained for a total field in $\mathrm{V}^{-}$ when an incident wave propagating upward in $\mathrm{V}^{-}$(supposed to contain a lossless dielectric) illuminates the grating surface, but this case is not worth in the frame of this chapter.

## 4.A.6. Normal derivative of a field continuous on $S$.

The calculation of the normal derivative of $\mathrm{F}^{ \pm}$on the grating surface in the general case is difficult. However, this aim can be reached at least in one case: when it is possible to define both $\mathrm{F}^{+}$and $\mathrm{F}^{-}$which satisfy three conditions:

- F is continuous across S , or equivalently $\psi^{+}(\mathrm{s})=\psi^{-}(\mathrm{s})$,
- $\mathrm{F}^{+}$and $\mathrm{F}^{-}$satisfy the same Helmholtz equation, with constant $\mathrm{k}^{2}\left(\mathrm{n}^{+}\right)^{2}$, or equivalently $\mathrm{n}^{+}=\mathrm{n}^{-}$.
- F satisfies a radiation condition at infinity.

Due to the second condition, it seems that this case does not make sense: if $\mathrm{n}^{+}=\mathrm{n}^{-}$, one cannot expect any scattering phenomenon. However, the study of this purely mathematical problem is crucial, for example in the study of perfectly conducting gratings.

First, it is worth noting that, thanks to the first condition, the gradient of F can be calculated without any use of distributions, which is not the case if F is discontinuous on S . Since $\gamma_{\mathrm{m}}^{-}=\gamma_{\mathrm{m}}^{+}, \mathcal{G}^{-}\left(\mathrm{x}, \mathrm{z}, \mathrm{s}^{\prime}\right)=\boldsymbol{\mathcal { G }}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{s}^{\prime}\right)$ and $\mathcal{N}^{-}\left(\mathrm{x}, \mathrm{z}, \mathrm{s}^{\prime}\right)=\boldsymbol{\mathcal { N }}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{s}^{\prime}\right)$ it can be deduced from eq. (4.139) that:

$$
\begin{align*}
& \mathrm{F}^{+}(\mathrm{x}, \mathrm{z})=\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}} \phi^{+}\left(\mathrm{s}^{\prime}\right)+\boldsymbol{N}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}} \psi^{+}\left(\mathrm{s}^{\prime}\right)\right] \mathrm{ds} \mathrm{~s}^{\prime},  \tag{4.161}\\
& \mathrm{F}^{-}(\mathrm{x}, \mathrm{z})=-\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\boldsymbol{\mathcal { G }}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}} \phi^{-}\left(\mathrm{s}^{\prime}\right)+\boldsymbol{N}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}} \psi^{-}\left(\mathrm{s}^{\prime}\right)\right] \mathrm{ds}^{\prime} . \tag{4.162}
\end{align*}
$$

We have seen in this appendix that, if $\phi^{+}\left(\mathrm{s}^{\prime}\right)$ and $\psi^{+}\left(\mathrm{s}^{\prime}\right)$ are compatible, the expression of $\mathrm{F}^{+}(\mathrm{x}, \mathrm{z})$ given by eq. (4.161) vanishes in $\mathrm{V}^{-}$. The same property holds for the expression of $\mathrm{F}^{-}(\mathrm{x}, \mathrm{z})$ given by eq. (4.162), which vanishes in $\mathrm{V}^{+}$if $\phi^{-}\left(\mathrm{s}^{\prime}\right)$ and $\psi^{-}\left(\mathrm{s}^{\prime}\right)$ are compatible. Bearing in mind that $\psi^{+}(\mathrm{s})=\psi^{-}(\mathrm{s})$, the expression of F in the entire space is given by adding the right-hand sides of eqs. (4.161) and (4.162):

$$
\begin{equation*}
\mathrm{F}(\mathrm{x}, \mathrm{z})=\mathrm{F}^{+}(\mathrm{x}, \mathrm{z})+\mathrm{F}^{-}(\mathrm{x}, \mathrm{z})=\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \boldsymbol{\mathcal { G }}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}}\left(\phi^{+}\left(\mathrm{s}^{\prime}\right)-\phi^{-}\left(\mathrm{s}^{\prime}\right)\right) \mathrm{ds}^{\prime} . \tag{4.163}
\end{equation*}
$$

Thanks to the continuity of F on S , the expression of its value on the profile does not make problem and is given by:

$$
\begin{equation*}
F(s)=\int_{s^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \boldsymbol{\mathcal { G }}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}}\left(\phi^{+}\left(\mathrm{s}^{\prime}\right)-\phi^{-}\left(\mathrm{s}^{\prime}\right)\right) \mathrm{ds} \mathrm{~s}^{\prime} \tag{4.164}
\end{equation*}
$$

In order to obtain the normal derivative of F , let us calculate its gradient:

$$
\begin{gather*}
\nabla F(\mathrm{x}, \mathrm{z})=\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \nabla_{(\mathrm{x}, \mathrm{z})}\left[\mathcal{G}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}}\right]\left(\phi^{+}\left(\mathrm{s}{ }^{\prime}\right)-\phi^{-}\left(\mathrm{s}^{\prime}\right)\right) \mathrm{ds} \mathrm{~s}^{\prime},  \tag{4.165}\\
\nabla\left[\boldsymbol{\mathcal { G }}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}}\right]=\frac{1}{2 \mathrm{~d}}\binom{\sum_{\mathrm{m}=-\infty,+\infty} \frac{\alpha_{\mathrm{m}}}{\gamma_{\mathrm{m}}^{+}} \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}} \mathrm{e}^{\mathrm{imK}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)+\mathrm{i} \gamma_{\mathrm{m}}^{+}\left|\mathrm{z}-\mathrm{z}^{\prime}\right|}}{\sum_{\mathrm{m}=-\infty,+\infty} \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}} \mathrm{e}^{\mathrm{imK}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)+\mathrm{i} \gamma_{\mathrm{m}}^{+} \mathrm{z}-\mathrm{z}^{\prime}}} \tag{4.166}
\end{gather*}
$$

The components of the normal $\overrightarrow{\mathrm{N}}_{\mathrm{S}}$ are given by:

$$
\begin{equation*}
\overrightarrow{\mathrm{N}}_{\mathrm{S}}=\binom{-\frac{\mathrm{dy}}{\mathrm{ds}}}{\frac{\mathrm{dx}}{\mathrm{ds}}} \tag{4.167}
\end{equation*}
$$

and thus:

$$
\begin{gather*}
{\left[\frac{\mathrm{d} \boldsymbol{\mathcal { G }}^{+}(\mathrm{x}, \mathrm{z}, \mathrm{~s}) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}}}{\mathrm{dN}}\right]_{\mathrm{S}}^{ \pm}=}  \tag{4.168}\\
\frac{1}{2 \mathrm{~d}} \lim _{ \pm}\left\{\sum_{\mathrm{m}=-\infty,+\infty}\left[\operatorname{sgn}\left(\mathrm{z}-\mathrm{z}^{\prime}\right) \frac{\mathrm{dx}}{\mathrm{ds}}-\frac{\alpha_{\mathrm{m}}}{\gamma_{\mathrm{m}}^{+}} \frac{\mathrm{dy}}{\mathrm{ds}}\right] \mathrm{e}^{\mathrm{i} \alpha_{0} x} \mathrm{e}^{\operatorname{imK}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)+i \gamma_{\mathrm{m}}^{+} \mid \mathrm{z}-\mathrm{z}^{\prime}}\right\}
\end{gather*}
$$

Using eqs. (4.165) and (4.168) yields:

$$
\begin{equation*}
\frac{\mathrm{dF}^{ \pm}}{\mathrm{dN}_{\mathrm{S}}}=\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\frac{\mathrm{~d} \boldsymbol{\mathcal { F }}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}}}{\mathrm{dN}_{\mathrm{S}}}\right]^{ \pm}\left(\phi^{+}\left(\mathrm{s}^{\prime}\right)-\phi^{-}\left(\mathrm{s}^{\prime}\right)\right) \mathrm{ds}^{\prime} \tag{4.169}
\end{equation*}
$$

In order to eliminate the use of limits in the expression of $\left[\frac{\mathrm{d} \boldsymbol{\mathcal { G }}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}}}{\mathrm{dN}}\right]_{\mathrm{S}}^{ \pm}$given by eq. (4.168), it can be remembered that, by definition,

$$
\begin{equation*}
\frac{\mathrm{dF}^{+}}{\mathrm{dN}_{\mathrm{S}}}-\frac{\mathrm{dF}^{-}}{\mathrm{dN}_{\mathrm{S}}}=\left[\phi^{+}(\mathrm{s})-\phi^{-}(\mathrm{s})\right] \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}} \tag{4.170}
\end{equation*}
$$

Moreover, it is to be noticed that, when z is constant, the components of $\nabla\left[\mathcal{G}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}}\right]$ given by eq. (4.166) are Fourier series in $x$. Using again the property of discontinuous Fourier series on the discontinuity, it can be derived that:

$$
\begin{align*}
& \frac{\mathrm{dF}^{+}}{\mathrm{dN}_{\mathrm{S}}}+\frac{\mathrm{dF}^{-}}{\mathrm{dN}_{\mathrm{S}}}= \\
& \frac{1}{d} \int_{s^{\prime}=0}^{s_{d}} \sum_{m=-\infty}^{+\infty}\left[\operatorname{sgn}\left(z-z^{\prime}\right) \frac{d x}{d s}-\frac{\alpha_{m}}{\gamma_{m}^{+}} \frac{d y}{d s}\right] e^{i \alpha_{0} x} e^{i m K\left(x-x^{\prime}\right)+i \gamma_{m}^{+}\left|z-z^{\prime}\right|}\left[\phi^{+}\left(s^{\prime}\right)-\phi^{-}\left(s^{\prime}\right)\right] d s^{\prime} . \tag{4.171}
\end{align*}
$$

From equations (4.170) and (4.171), we deduce:

$$
\begin{equation*}
\frac{\mathrm{dF}^{ \pm}}{\mathrm{dN}_{\mathrm{S}}}= \pm \frac{\phi^{+}(\mathrm{s})-\phi^{-}(\mathrm{s})}{2} \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}}+\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \mathscr{K}\left(\mathrm{~s}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}}\left[\phi^{+}\left(\mathrm{s}^{\prime}\right)-\phi^{-}\left(\mathrm{s}^{\prime}\right)\right] \mathrm{ds}^{\prime} \tag{4.172}
\end{equation*}
$$

with:

$$
\begin{equation*}
\mathscr{K}\left(s, s^{\prime}\right)=\frac{1}{2 d} \sum_{m=-\infty,+\infty}\left[\operatorname{sgn}\left(z-z^{\prime}\right) \frac{d x}{d s}-\frac{\alpha_{m}}{\gamma_{m}^{+}} \frac{d y}{d s}\right] e^{i m K\left(x-x^{\prime}\right)+i \gamma_{m}^{+}\left|z-z^{\prime}\right|} \tag{4.173}
\end{equation*}
$$

It is interesting to notice that the expression of function $\mathscr{K}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$ is very close to that of $\mathcal{N}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$, the only difference being that $\frac{\mathrm{dx}^{\prime}}{\mathrm{ds}{ }^{\prime}}$ and $\frac{\mathrm{dy}^{\prime}}{\mathrm{ds} '^{\prime}}$ are replaced by $\frac{\mathrm{dx}}{\mathrm{ds}}$ and $\frac{\mathrm{dy}}{\mathrm{ds}}$.

Like $\mathcal{N}^{ \pm}\left(\mathrm{s}, \mathrm{s}^{\prime}\right), \mathscr{K}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$ is continuous and its limit when $\mathrm{s}^{\prime} \rightarrow \mathrm{s}$ can be expressed in closed form [44,7], as stated in section 4.6.3.

## 4.A.7. Limit values of a field with continuous normal derivative on $S$.

We consider a function F satisfying the following conditions

- F has the same normal derivative on both sides of S, or equivalently $\phi^{+}(s)=\phi^{-}(s)$,
- $\mathrm{F}^{+}$and $\mathrm{F}^{-}$satisfy the same Helmholtz equation, or equivalently $\mathrm{n}^{+}=\mathrm{n}^{-}$.
- F satisfies a radiation condition at infinity.

The aim of this section is to calculate the limits of F on both parts of S .
From the secong Green theorem (eq. (4.136)), F can be expressed from the values of $\phi^{+}(\mathrm{s}), \phi^{-}(\mathrm{s}), \psi^{+}(\mathrm{s})$ and $\psi^{-}(\mathrm{s})$ :

$$
\begin{align*}
& \mathrm{F}^{+}(\mathrm{x}, \mathrm{z})=\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \phi^{+}\left(\mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}\right.  \tag{4.174}\\
&\left.\quad+\mathcal{N}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \psi^{+}\left(\mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}\right] \mathrm{ds}^{\prime}
\end{align*}
$$

and, bearing in mind that $\phi^{-}(\mathrm{s})=\phi^{+}(\mathrm{s})$ and that $\mathrm{n}^{-}=\mathrm{n}^{+}, \mathcal{G}^{-}\left(\mathrm{x}, \mathrm{z}, \mathrm{s}^{\prime}\right)=\boldsymbol{\mathcal { G }}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{s}^{\prime}\right)$ and $\mathcal{N}^{-}\left(\mathrm{x}, \mathrm{z}, \mathrm{s}^{\prime}\right)=\mathcal{N}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{s}^{\prime}\right):$

$$
\begin{align*}
& \mathrm{F}^{-}(\mathrm{x}, \mathrm{z})=-\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\boldsymbol{\mathcal { G }}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \phi^{+}\left(\mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}\right.  \tag{4.175}\\
& \left.\quad+\mathcal{N}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)} \Psi^{-}\left(\mathrm{s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}\right] \mathrm{ds}^{\prime}
\end{align*}
$$

It has been shown in section 4.A. 4 that, if $\phi^{+}\left(s^{\prime}\right)$ and $\psi^{+}\left(s^{\prime}\right)$ are compatible, the expression of $\mathrm{F}^{+}(\mathrm{x}, \mathrm{z})$ given by eq. (4.174) vanishes in $\mathrm{V}^{-}$. The same property holds for the expression of $\mathrm{F}^{-}(\mathrm{x}, \mathrm{z})$ given by eq. (4.175), which vanishes in $\mathrm{V}^{+}$, thus we can write that, if the compatibility equations are satisfied:

$$
\begin{equation*}
\mathrm{F}(\mathrm{x}, \mathrm{z})=\mathrm{F}^{+}(\mathrm{x}, \mathrm{z})+\mathrm{F}^{-}(\mathrm{x}, \mathrm{z})=\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \mathcal{N}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}} \Psi\left(\mathrm{~s}^{\prime}\right) \mathrm{ds} \mathrm{~s}^{\prime} \tag{4.176}
\end{equation*}
$$

or

$$
\begin{equation*}
\mathrm{U}(\mathrm{x}, \mathrm{z})=\mathrm{F}(\mathrm{x}, \mathrm{z}) \mathrm{e}^{-\mathrm{i} \alpha_{0} \mathrm{x}}=\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \mathcal{N}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \Psi\left(\mathrm{s}^{\prime}\right) \mathrm{ds}{ }^{\prime}, \tag{4.177}
\end{equation*}
$$

with:

$$
\begin{equation*}
\Psi\left(\mathrm{s}^{\prime}\right)=\psi^{+}\left(\mathrm{s}^{\prime}\right)-\psi^{-}\left(\mathrm{s}^{\prime}\right) \tag{4.178}
\end{equation*}
$$

In order to express the limits $\lim \left\{\mathrm{U}^{+}(\mathrm{x}, \mathrm{z})\right\}$ or $\lim \left\{\mathrm{U}^{-}(\mathrm{x}, \mathrm{z})\right\}$ on both parts of S, . we can use a first equation:

$$
\begin{equation*}
\lim _{+}\left\{\mathrm{U}^{+}(\mathrm{x}, \mathrm{z})\right\}-\lim _{-}\left\{\mathrm{U}^{-}(\mathrm{x}, \mathrm{z})\right\}=\Psi(\mathrm{s}) \tag{4.179}
\end{equation*}
$$

To find a second equation, we can consider eqs. (4.176) and (4.177) which gives the expression of $U(x, z)$. If $z$ is fixed, the expression of $U(x, z)$ is a Fourier series in $x$, which is discontinuous on $S$. The value on $S$ of this Fourier series is the average value of the limits on both sides of S, thus:

$$
\begin{equation*}
\lim _{+}\left\{\mathrm{U}^{+}(\mathrm{x}, \mathrm{z})\right\}+\lim _{-}\left\{\mathrm{U}^{-}(\mathrm{x}, \mathrm{z})\right\}=2 \int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \mathcal{N}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \Psi\left(\mathrm{s}^{\prime}\right) \mathrm{ds} \mathrm{~s}^{\prime} \tag{4.180}
\end{equation*}
$$

From eqs. (4.179) and (4.180), we deduce the two limits:

$$
\begin{equation*}
\lim _{ \pm}\{\mathrm{U}(\mathrm{x}, \mathrm{y})\}= \pm \frac{\Psi(\mathrm{s})}{2}+\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \mathcal{N}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \Psi\left(\mathrm{s}^{\prime}\right) \mathrm{ds} \mathrm{~s}^{\prime} \tag{4.181}
\end{equation*}
$$

## 4.A.8. Calculation of the amplitudes of the plane wave expansions at infinity.

For $\mathrm{z} \rightarrow \pm \infty$, the expression of $\mathrm{F}^{ \pm}$given by eq. (4.136) can be simplified since $\operatorname{sgn}\left(z-z^{\prime}\right)= \pm 1$ and $\left|z-z^{\prime}\right|= \pm\left(z-z^{\prime}\right)$, in such a way that the expression of $F$ at infinity becomes a sum of plane waves:

$$
\begin{gather*}
F^{+}(x, y)=\sum_{m=-\infty}^{\infty} r_{m} \exp \left(i \alpha_{m} x+i \gamma_{m}^{+} z\right) \quad \text { if } z>z_{0},  \tag{4.182}\\
F^{-}(x, y)=\sum_{m=-\infty}^{\infty} t_{m} \exp \left(i \alpha_{m} x-i \gamma_{m}^{-} z\right) \quad \text { if } z<0,  \tag{4.183}\\
r_{m}=\frac{1}{2 d} \int_{s=0}^{s_{d}} e^{-i m K x(s)-i \gamma_{m}^{+} z(s)}\left[\frac{-i \phi^{+}(s)}{\gamma_{m}^{+}}+\left(\frac{d x(s)}{d s}-\frac{\alpha_{m}}{\gamma_{m}^{+}} \frac{d z(s)}{d s}\right) \psi^{+}(s)\right] d s,  \tag{4.184}\\
t_{m}=\frac{1}{2 d} \int_{s=0}^{s_{d}} e^{-i m K x(s)+i \gamma_{m}^{-} z(s)}\left[\frac{i \phi^{-}(s)}{\gamma_{m}^{-}}+\left(\frac{d x(s)}{d s}+\frac{\alpha_{m}}{\gamma_{m}^{+}} \frac{d z(s)}{d s}\right) \psi^{-}(s)\right] d s, \tag{4.185}
\end{gather*}
$$

with $z_{0}$ being the ordinate of the top of the grating profile. It must be noticed that a finite number of orders $m$, called propagating orders, are non-evanescent and propagate at infinity. They correspond to real values of $\gamma_{\mathrm{m}}^{+}$(for reflected orders) or $\gamma_{\mathrm{m}}^{-}$(for transmitted orders, if the optical index $\mathrm{n}^{-}$is real only).

Equation (4.184) can easily be generalized to the case in which we know the limit value of the total field on S (including incident waves) and its normal derivative. It suffices to analyze the behaviour at infinity of eq. (4.155) instead of eq. (4.136). The result is that it suffices to replace the values $\phi^{+}(\mathrm{s})$ and $\psi^{+}(\mathrm{s})$ relative to the scattered field by the values $\Phi^{+}(\mathrm{s})$ and $\Psi^{+}(\mathrm{s})$ relative to the total field:

$$
\begin{equation*}
\mathrm{r}_{\mathrm{m}}=\frac{1}{2 \mathrm{~d}} \int_{\mathrm{s}=0}^{\mathrm{s}_{\mathrm{d}}} \mathrm{e}^{-\mathrm{imKx}(\mathrm{~s})-\mathrm{i} \gamma_{\mathrm{m}}^{+} \mathrm{z}(\mathrm{~s})}\left[\frac{-\mathrm{i} \Phi^{+}(\mathrm{s})}{\gamma_{\mathrm{m}}^{+}}+\left(\frac{\mathrm{dx}(\mathrm{~s})}{\mathrm{ds}}-\frac{\alpha_{\mathrm{m}}}{\gamma_{\mathrm{m}}^{+}} \frac{\mathrm{dz}(\mathrm{~s})}{\mathrm{ds}}\right) \Psi^{+}(\mathrm{s})\right] \mathrm{ds} \tag{4.186}
\end{equation*}
$$

Diffraction efficiencies $\rho_{\mathrm{m}}$ in the reflected orders propagating above the grating can be obtained by using the Poynting theorem on segments of one period parallel to the x axis:

$$
\begin{equation*}
\rho_{\mathrm{m}}=\frac{\gamma_{\mathrm{m}}^{+}}{\gamma_{0}^{+}}\left|\mathrm{r}_{\mathrm{m}}\right|^{2} \tag{4.187}
\end{equation*}
$$

When the grating is made of a lossless dielectric, the transmitted efficiencies $\tau_{\mathrm{m}}$ are given by:

$$
\begin{equation*}
\tau_{\mathrm{m}}=\frac{\mathrm{q}^{-}}{\mathrm{q}^{+}} \frac{\gamma_{\mathrm{m}}^{+}}{\gamma_{0}^{+}}\left|\mathrm{t}_{\mathrm{m}}\right|^{2} \tag{4.188}
\end{equation*}
$$

## Appendix 4.B. Integral method leading to a single integral equation for bare, metallic or dielectric grating

Historically, the formalism presented in this Appendix was the first one to lead to a single integral equation for a dielectric or metallic grating. The steps of the method are summarized in figure 4.14.


Figure 4.14. Steps of the integral formalism leading to a single equation

## 4.B.1. Definition of the unknown function

The single unknown function $\Phi$ is defined from a function $\tilde{\Gamma}=\left\{\begin{array}{l}\tilde{\Gamma}^{+} \text {in } \mathrm{V}^{+}, \\ \tilde{\Gamma}^{-} \text {in } \mathrm{V}^{-},\end{array}\right.$, atisfying the following conditions:

- it satisfies the same Helmholtz equation in the entire space, except may be on the profile S ,:

$$
\begin{equation*}
\nabla^{2} \tilde{\Gamma}+\mathrm{k}^{2}\left(\mathrm{n}^{+}\right)^{2} \tilde{\Gamma}=0 \tag{4.189}
\end{equation*}
$$

- it has the same pseudo-periodicity as the actual solution of the grating problem:

$$
\begin{equation*}
\tilde{\Gamma}(\mathrm{x}+\mathrm{d}, \mathrm{z})=\tilde{\Gamma}(\mathrm{x}, \mathrm{z}) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{~d}} \tag{4.190}
\end{equation*}
$$

- it identifies to the actual physical solution of the grating problem in $\mathrm{V}^{+}$:

$$
\begin{equation*}
\tilde{\Gamma}^{+} \equiv \mathrm{F}^{+}, \tag{4.191}
\end{equation*}
$$

- it is continuous across $S$,
- it satisfies a radiation condition for $\mathrm{y} \rightarrow \pm \infty$.

We denote by $\psi^{ \pm \pm} \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}$ and $\phi^{ \pm \pm} \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}^{\prime}}$ the limit values of $\tilde{\Gamma}$ and of its normal derivative $\frac{\mathrm{d} \tilde{\Gamma}}{\mathrm{dN}_{\mathrm{S}}}$ on S, bearing in mind that by definition, $\psi^{\prime+} \equiv \psi^{+}$and $\phi^{\prime+} \equiv \phi^{+}$.

The question which arises is to know if $\tilde{\Gamma}^{-}$is well defined. The continuity of $\tilde{\Gamma}$ across S imposes the value of $\tilde{\Gamma}^{-}$on S: $\psi^{\prime^{-}} \equiv \psi^{\prime+} \equiv \psi^{+}$. In addition, $\tilde{\Gamma}^{-}$satisfies a Helmholtz equation and a radiation condition at infinity. It is worth noting that, in contrast with the function $\Gamma$ introduced in Appendix 4.A, here the function $\tilde{\Gamma}$ does not include the incident field and thus has no physical meaning below the profile. As mentioned in Appendix 4.A, the solution of this boundary value problem (since we impose that $\tilde{\Gamma}^{-}=\tilde{\Gamma}^{+}$) exists and is unique, thus $\tilde{\Gamma}$ is correctly defined in the entire space. The unknown function $\Phi$ of the integral equation is defined as the jump of the normal derivative of $\tilde{\Gamma}$ across S , more precisely:

$$
\begin{equation*}
\Phi=\phi^{\prime+}-\phi^{\prime^{-}} . \tag{4.192}
\end{equation*}
$$

As regards the physical interpretation of $\Phi$, it can be shown easily that $\Phi \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}}$ is the surface current density which, placed on S , generates in $\mathrm{V}^{+}$the actual scattered field. In other words, we have considered a fictitious structure consisting of an infinitely thin, perfectly conducting metallic sheet supporting a surface current density $j(s) \hat{y}$ placed on the grating surface, this surface separating two media having identical optical properties (refractive index $\mathrm{n}^{+}$). The unknown $\Phi$ is proportional to j .

## 4.B.2. Expression of the scattered field, its limit on $S$ and its normal derivative from $\Phi$.

It must be noticed that the function $\tilde{\Gamma}$ satisfies all the conditions of the function $F$ of section 4.A.5. Since $\tilde{\Gamma}$ is continuous on S , the calculation of $\tilde{\Gamma}$ from $\Phi$ can be achieved using eqs. (4.163) and (4.137):

$$
\begin{gather*}
\tilde{\Gamma}(\mathrm{x}, \mathrm{z})=\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \mathcal{G}^{+}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}} \Phi\left(\mathrm{~s}^{\prime}\right) \mathrm{ds}^{\prime}  \tag{4.193}\\
\mathcal{\mathcal { G }}^{+}\left(\mathrm{x}, \mathrm{y}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{id}} \sum_{\mathrm{m}=-\infty}^{\infty} \frac{1}{\gamma_{\mathrm{m}}^{+}} \exp \left[\operatorname{imK}\left(\mathrm{x}-\mathrm{x}^{\prime}\right)+\mathrm{i} \gamma_{\mathrm{m}}^{+}\left|\mathrm{z}-\mathrm{z}^{\prime}\right|\right] . \tag{4.194}
\end{gather*}
$$

The value of the limit $\psi^{\prime+}(\mathrm{s}) \mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}}$ of $\tilde{\Gamma}^{+}(\mathrm{x}, \mathrm{z})$ on S does not make problem, thanks to its continuity:

$$
\begin{equation*}
\psi^{\prime+}(\mathrm{s})=\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \boldsymbol{\mathcal { G }}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \Phi\left(\mathrm{s}^{\prime}\right) \mathrm{ds}^{\prime} \tag{4.195}
\end{equation*}
$$

with $\mathcal{G}^{+}\left(\mathrm{s}, \mathrm{s}^{\prime}\right)$ the value on S of $\boldsymbol{\mathcal { G }}^{+}\left(\mathrm{x}, \mathrm{y}, \mathrm{s}^{\prime}\right)$, given by eq. (4.147):

$$
\begin{equation*}
\boldsymbol{\mathcal { G }}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)=\frac{1}{2 \mathrm{id}} \sum_{\mathrm{m}=-\infty}^{\infty} \frac{1}{\gamma_{\mathrm{m}}^{+}} \mathrm{e}^{\mathrm{imK}\left(\mathrm{x}(\mathrm{~s})-\mathrm{x}^{\prime}\left(\mathrm{s}^{\prime}\right)\right)+\mathrm{i} \gamma_{\mathrm{m}}^{+}\left|\mathrm{z}(\mathrm{~s})-\mathrm{z}^{\prime}\left(\mathrm{s}^{\prime}\right)\right|} \tag{4.196}
\end{equation*}
$$

and finally its normal derivative can be derived from eq. (4.172) and (4.173):

$$
\begin{equation*}
\phi^{+}(\mathrm{s})=\frac{\Phi(\mathrm{s})}{2}+\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \mathscr{K}\left(\mathrm{~s}, \mathrm{~s}^{\prime}\right) \Phi\left(\mathrm{s}^{\prime}\right) \mathrm{ds}^{\prime} \tag{4.197}
\end{equation*}
$$

$$
\begin{equation*}
\mathscr{K}\left(s, s^{\prime}\right)=\frac{1}{2 d} \sum_{m=-\infty,+\infty}\left[\operatorname{sgn}\left(z-z^{\prime}\right) \frac{d x}{d s}-\frac{\alpha_{m}}{\gamma_{m}^{+}} \frac{d y}{d s}\right] e^{\operatorname{imK}\left(x-x^{\prime}\right)+i \gamma_{m}^{+}\left|z-z^{\prime}\right|} \tag{4.198}
\end{equation*}
$$

It is worth noting that the values of the limits of the field and its normal derivative on S that are the two unknown functions in the classical integral theory (section 2), are now expressed from the single function $\Phi(\mathrm{s})$. This is not surprising since the limits on S of the function $\tilde{\Gamma}$ given by eq. (4.193) satisfy automatically a relation of compatibility, whatever the function $\Phi\left(s^{\prime}\right)$ introduced in the integral may be: it is the field generated by a surface current on S . As a consequence, we have not to include in the theory a relation of compatibility in $\mathrm{V}^{+}$, which was the first integral equation in the classical formalism.

## 4.B.3. Integral equation

The single integral equation will be obtained by writing the relation of compatibility in $\mathrm{V}^{-}$, considered now to be filled by the actual grating material, i.e. a material of index $\mathrm{n}^{-}$. With this aim, we calculate the limits in $\mathrm{V}^{-}$of the field and its normal derivative on S . inserting in the continuity conditions of the field given by eqs (4.8) and (4.10) the the limit values given by eqs (4.195) and (4.197):

$$
\begin{gather*}
\psi^{-}(\mathrm{s})=\psi^{\mathrm{i}}(\mathrm{~s})+\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \mathcal{G}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \Phi\left(\mathrm{s}^{\prime}\right) \mathrm{ds}^{\prime}  \tag{4.199}\\
\phi^{-}(\mathrm{s})=\frac{\mathrm{q}^{+}}{\mathrm{q}^{-}}\left[\frac{\Phi(\mathrm{s})}{2}+\phi^{\mathrm{i}}(\mathrm{~s})+\int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \mathscr{K}\left(\mathrm{~s}, \mathrm{~s}^{\prime}\right) \Phi\left(\mathrm{s}^{\prime}\right) \mathrm{ds} \mathrm{~s}^{\prime}\right] \tag{4.200}
\end{gather*}
$$

with $\mathrm{q}^{ \pm}$given by eq. (4.12).
The field in $\mathrm{V}^{-}$can be deduced from eqs. (4.199) and (4.200) using eq. (4.136):

$$
\begin{equation*}
\mathrm{F}^{-}(\mathrm{x}, \mathrm{z})=-\mathrm{e}^{\mathrm{i} \alpha_{0} \mathrm{x}} \int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left[\mathcal{G}^{-}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \phi^{-}\left(\mathrm{s}^{\prime}\right)+\mathcal{N}^{-}\left(\mathrm{x}, \mathrm{z}, \mathrm{~s}^{\prime}\right) \psi^{-}\left(\mathrm{s}^{\prime}\right)\right] \mathrm{ds}^{\prime} \tag{4.201}
\end{equation*}
$$

and the equation of compatibility is given by eq. (4.146):

$$
\begin{align*}
& \int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}}\left\{\frac{\mathrm{q}^{+}}{\mathrm{q}^{-}} \boldsymbol{\mathcal { G }}^{-}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)\left[\frac{\Phi\left(\mathrm{s}^{\prime}\right)}{2}+\phi^{\mathrm{i}}\left(\mathrm{~s}^{\prime}\right)+\int_{\mathrm{s}^{\prime \prime}=0}^{\mathrm{s}_{\mathrm{d}}} \mathscr{K}\left(\mathrm{~s}^{\prime}, \mathrm{s}^{\prime \prime}\right) \Phi\left(\mathrm{s}^{\prime \prime}\right) \mathrm{ds} \mathrm{~s}^{\prime \prime}\right]\right. \\
& \left.+\boldsymbol{\mathcal { N }}^{-}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right)\left[\psi^{\mathrm{i}}\left(\mathrm{~s}^{\prime}\right)+\int_{\mathrm{s}^{\prime \prime}=0}^{\mathrm{s}_{d}} \mathcal{G}^{+}\left(\mathrm{s}^{\prime}, \mathrm{s}^{\prime \prime}\right) \Phi\left(\mathrm{s}^{\prime \prime}\right) \mathrm{ds}^{\prime \prime}\right]\right\} \mathrm{ds}^{\prime}+\frac{\psi^{\mathrm{i}}(\mathrm{~s})}{2}+\frac{1}{2} \int_{\mathrm{s}^{\prime}=0}^{\mathrm{s}_{\mathrm{d}}} \boldsymbol{\mathcal { G }}^{+}\left(\mathrm{s}, \mathrm{~s}^{\prime}\right) \Phi\left(\mathrm{s}^{\prime}\right) \mathrm{ds}^{\prime}=0 \tag{4.202}
\end{align*}
$$

which yields, after simplification:

$$
\begin{equation*}
\left[\left(\mathcal{N}^{-}+\frac{\mathbb{I}}{2}\right) \boldsymbol{\mathcal { G }}^{+}+\frac{\mathrm{q}^{+}}{\mathrm{q}^{-}} \boldsymbol{\mathcal { G }}^{-}\left(\mathscr{K}+\frac{\mathbb{I}}{2}\right)\right] \Phi=-\left(\boldsymbol{N}^{-}+\frac{\mathbb{I}}{2}\right) \psi^{\mathrm{i}}(\mathrm{~s})+\frac{\mathrm{q}^{+}}{\mathrm{q}^{-}} \boldsymbol{\mathcal { G }}^{-} \phi^{\mathrm{i}}, \tag{4.203}
\end{equation*}
$$

with the symbol $\mathcal{O} \eta$ denoting the function $\int_{s^{\prime}=0}^{s_{\mathrm{d}}} \mathcal{O}\left(\mathrm{s}, \mathrm{s}^{\prime}\right) \eta\left(\mathrm{s}^{\prime}\right) \mathrm{ds}$ ' in operator notation.
For $\mathrm{z} \rightarrow+\infty$, the expression of $\tilde{\Gamma}^{+} \equiv \mathrm{F}^{+}$given by eq. (4.193) can be simplified since $\operatorname{sgn}\left(z-z^{\prime}\right)= \pm 1$ and $\left|z-z^{\prime}\right|= \pm\left(z-z^{\prime}\right)$, in such a way that the expression of $F^{+}$at infinity becomes a sum of plane waves:

$$
\begin{equation*}
F^{+}(x, y)=\sum_{m=-\infty}^{\infty} r_{m} \exp \left(i \alpha_{m} x+i \gamma_{m}^{+} z\right) \quad \text { if } z>z_{0} \tag{4.204}
\end{equation*}
$$

with amplitudes given by:

$$
\begin{equation*}
\mathrm{r}_{\mathrm{m}}=\frac{1}{2 \mathrm{id} \gamma_{\mathrm{m}}^{+}} \int_{\mathrm{s}=0}^{\mathrm{s}_{\mathrm{d}}} \mathrm{e}^{-\mathrm{imKx}(\mathrm{~s})-\mathrm{i} \gamma_{\mathrm{m}}^{+} z(\mathrm{~s})} \Phi(\mathrm{s}) \mathrm{ds} \tag{4.205}
\end{equation*}
$$

The diffraction efficiencies of the reflected waves are then deduced by:

$$
\begin{equation*}
\rho_{\mathrm{m}}=\frac{\gamma_{\mathrm{m}}^{+}}{\gamma_{0}^{+}}\left|\mathrm{r}_{\mathrm{m}}\right|^{2} \tag{4.206}
\end{equation*}
$$

Similarly, it can be derived from eq. (4.201) that the transmitted field can be represented by a sum of plane waves below the profile:

$$
\begin{equation*}
\mathrm{F}^{-}(\mathrm{x}, \mathrm{y})=\sum_{\mathrm{m}=-\infty}^{\infty} \mathrm{t}_{\mathrm{m}} \exp \left(\mathrm{i} \alpha_{\mathrm{m}} \mathrm{x}+\mathrm{i} \gamma_{\mathrm{m}}^{-} \mathrm{z}\right) \quad \text { if } \mathrm{z}<0 \tag{4.207}
\end{equation*}
$$

The amplitudes of the transmitted plane waves are derived from eq. (4.185) after calculating the functions $\phi^{-}$and $\psi^{-}$from eqs. (4.199) and (4.200):

$$
\begin{equation*}
\mathrm{t}_{\mathrm{m}}=\frac{1}{2 \mathrm{~d}} \int_{\mathrm{s}=0}^{\mathrm{s}_{\mathrm{d}}} \mathrm{e}^{-\mathrm{imKx}(\mathrm{~s})+\mathrm{i} \gamma_{\mathrm{m}}^{-} \mathrm{z}(\mathrm{~s})}\left[\frac{\mathrm{i} \phi^{-}(\mathrm{s})}{\gamma_{\mathrm{m}}^{-}}+\left(\frac{\mathrm{dx}(\mathrm{~s})}{\mathrm{ds}}+\frac{\alpha_{\mathrm{m}}}{\gamma_{\mathrm{m}}^{+}} \frac{\mathrm{dz}(\mathrm{~s})}{\mathrm{ds}}\right) \psi^{-}(\mathrm{s})\right] \mathrm{ds} . \tag{4.208}
\end{equation*}
$$

It is interesting to notice that this method has been extended to other problems (including 3D problems of scattering) by many specialists of theoretical physics and applied mathematics [49-52].

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Chapter C :
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) $\Psi \mathrm{QP} \mathrm{UF}=\mathrm{R}(\mathrm{CD})$
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# Chapter 5 

# Finite Element Method 

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### 5.1 Introduction

Finite element methods (FEM) represent a very general set of techniques to approximate solutions of partial derivative equations. Their main advantage lies in their ability to handle arbitrary geometries via unstructured meshes of the domain of interest: The discretization of oblic geometry edges is natively built in. Finite Element Methods have been widely developed in many areas of physics and engineering: mechanics, thermodynamics...

But until the early 80 's, two major drawbacks prevented them from being used in electromagnetic problems. On the one hand, existing nodal element basis did not satisfy the physical (dis)continuity of the vector fields components and lead to spurious solutions [1]. On the other hand, there was no proper way to truncate unbounded regions in open wave problems.

These two major limitations were both overcome in the early 80's: Vector elements have been developed by Nédélec [2, 3], and Perfectly Matched Layers (PMLs) were discovered by Bérenger [4]. Since then, it has been shown that PMLs could be described in the general framework of transformation optics $[5,6,7,8]$.

All the mathematical and computational ingredients now exist and the goal of this chapter is to show how to combine them to implement a general 3D numerical scheme adapted to gratings using Finite Elements. In fact, we are now facing the physical difficulties inherent to the infinite spatial characteristics of the grating problem, whereas the computation domain has to be bounded in practice: (i) Both the superstrate and the substrate are infinite regions, (ii) there is an infinite number of periods and, last but not least, (iii) the sources of the incident field (a plane wave) are located in the superstrate at an infinite distance from the grating.

In this chapter, the infinite extension of the superstrate and substrate is addressed using cartesian PMLs. In the framework of transformation optics, we demonstrate that Bérenger's original PMLs can be extended to the challenging numerical cases of grazing incidence in order to deal with extreme oblic incidences or configurations near Wood's anomalies. The second issue of infinite number of period can be addressed via Bloch conditions. Finally, we are dealing with the distant plane wave sources through an equivalence of the diffraction problem with a radiation one whose sources are localized inside the diffractive element itself. The unknown field to be approximated using Finite Elements is a radiated field with sources inside the computation box and allows to retrieve easily the total field with the plane wave source.

In a first section, we derive and implement this approach in the so-called 2D non-conical, or scalar, case. We are dealing with the infinite issues rigorously in both TE and TM polarization cases. It results in a radiation problem with sources localized in the diffractive element itself. We mathematically split the whole problem into two parts. The first one consists in the classical calculation of the total field solution of a simple interface. The second one amounts to looking
for a radiated field with sources confined within the diffractive obstacles and deduced from the first elementary problem. From this viewpoint, the later radiated field can be interpreted as an exact perturbation of the total field. We show that our approach allows to tackle some kind of anisotropy without increasing the computational time or resource. Through a battery of examples, we illustrate its independence towards the geometry of the diffractive pattern. Finally, we present an Adaptative PML able to tackle grazing incidences or configurations near Wood's anomaly.

In a second section, we extend this approach to the most general configuration of vector diffraction by crossed gratings embedded in arbitrary multilayered stack. The main advantage of this method is, again, its complete independence towards the shape of the diffractive element, whereas other methods often require heavy adjustments depending on whether the geometry of the groove region presents oblique edges. This approach combined with the use of second order edge elements allows us to retrieve the few numerical academic examples found in the literature with an excellent accuracy. Furthermore, we provide a new reference case combining major difficulties: A non trivial toroidal geometry together with strong losses and a high permittivity contrast. Finally, we discuss computation time and convergence as a function of the mesh refinement as well as the choice of the direct solver.

### 5.2 Scalar diffraction by arbitrary mono-dimensional gratings : a Finite Element formulation

### 5.2.1 Set up of the problem and notations

We denote by $\mathbf{x}, \mathbf{y}$ and $\mathbf{z}$, the unit vectors of the axes of an orthogonal coordinate system Oxyz. We deal only with time-harmonic fields; consequently, the electric and magnetic fields are represented by the complex vector fields $\mathbf{E}$ and $\mathbf{H}$, with a time dependance in $\exp (-i \omega t)$.

Besides, in this chapter, we assume that the tensor fields of relative permittivity $\underline{\underline{\varepsilon}}$ and relative permeability $\mu$ can be written as follows:

$$
\underline{\underline{\varepsilon}}=\left(\begin{array}{ccc}
\varepsilon_{x x} & \bar{\varepsilon}_{a} & 0  \tag{5.1}\\
\varepsilon_{a} & \varepsilon_{y y} & 0 \\
0 & 0 & \varepsilon_{z z}
\end{array}\right) \quad \text { and } \quad \underline{\underline{\mu}}=\left(\begin{array}{ccc}
\mu_{x x} & \bar{\mu}_{a} & 0 \\
\mu_{a} & \mu_{y y} & 0 \\
0 & 0 & \mu_{z z}
\end{array}\right)
$$

where $\varepsilon_{x x}, \varepsilon_{a}, \ldots \mu_{z z}$ are possibly complex valued functions of the two variables $x$ and $y$ and where $\bar{\varepsilon}_{a}$ (resp. $\bar{\mu}_{a}$ ) represents the conjugate complex of $\varepsilon_{a}$ (resp. $\mu_{a}$ ). These kinds of materials are said to be $z$-anisotropic. It is of importance to note that with such tensor fields, lossy materials can be studied (the lossless materials correspond to tensors with real diagonal terms represented by Hermitian matrices) and that the problem is invariant along the $z$-axis but the tensor fields can vary continuously (gradient index gratings) or discontinuously (step index gratings). Moreover we define $k_{0}:=\omega / c$.

The gratings that we are dealing with are made of three regions (See Fig. 5.1).

- The superstratum $\left(y>h^{g}\right)$ which is supposed to be homogeneous, isotropic and lossless and characterized solely by its relative permittivity $\varepsilon^{+}$and its relative permeability $\mu^{+}$ and we denote $k^{+}:=k_{0} \sqrt{\varepsilon^{+} \mu^{+}}$
- The substratum $(y<0)$ which is supposed to be homogeneous and isotropic and therefore characterized by its relative permittivity $\varepsilon^{-}$and its relative permeability $\mu^{-}$and we denote $k^{-}:=k_{0} \sqrt{\varepsilon^{-} \mu^{-}}$
- The groove region $\left(0<y<h^{g}\right)$ which can be heterogeneous and $z$-anisotropic and thus characterized by the two tensor fields $\underline{\underline{\varepsilon}}^{g}(x, y)$ and $\underline{\underline{\mu}}^{g}(x, y)$. It is worth noting that the method does work irrespective of whether the tensor fields are piecewise constant. The groove periodicity along $x$-axis will be denoted $d$.


Fig. 5.1: Sketch and notations of the grating studied in this section.

This grating is illuminated by an incident plane wave of wave vector $\mathbf{k}_{p}^{+}=\alpha \mathbf{x}-\beta^{+} \mathbf{y}=$ $k^{+}\left(\sin \theta_{0} \mathbf{x}-\cos \theta_{0} \mathbf{y}\right)$, whose electric field (TM case) ( resp. magnetic field (TE case)) is linearly polarized along the $z$-axis:

$$
\begin{equation*}
\mathbf{E}_{e}^{0}=\mathbf{A}_{e}^{0} \exp \left(i \mathbf{k}_{p}^{+} \cdot \boldsymbol{r}\right) \mathbf{z} \quad\left(\text { resp. } \mathbf{H}_{m}^{0}=\mathbf{A}_{m}^{0} \exp \left(i \mathbf{k}_{p}^{+} \cdot \boldsymbol{r}\right) \mathbf{z}\right), \tag{5.2}
\end{equation*}
$$

where $\mathbf{A}_{e}^{0}\left(\right.$ resp. $\left.\mathbf{A}_{m}^{0}\right)$ is an arbitrary complex number and $\boldsymbol{r}=(x, y)^{\mathrm{T}}$. In this section, a plane wave is characterized by its wave-vector denoted $\mathbf{k}_{\{p, c\}}^{\{+,-\}}$. The subscript $p$ (resp. c) stands for "propagative" (resp. "counter-propagative"). The superscript + (resp. -) refers to the associated wavenumber $k^{+}$(resp. $k^{-}$), and indicates that we are dealing with a plane wave propagating in the superstrate (resp. substrate).The magnetic (resp. electric) field derived from $\mathbf{E}_{e}^{0}$ (resp. $\mathbf{H}_{m}^{0}$ ) is denoted $\mathbf{H}_{e}^{0}$ (resp. $\mathbf{E}_{m}^{0}$ ) and the electromagnetic field associated with the incident field is therefore denoted $\left(\mathbf{E}^{0}, \mathbf{H}^{0}\right)$ which is equal to $\left(\mathbf{E}_{e}^{0}, \mathbf{H}_{e}^{0}\right)\left(\right.$ resp. $\left(\mathbf{E}_{m}^{0}, \mathbf{H}_{m}^{0}\right)$ ).

The diffraction problem that we address consists in finding Maxwell equation solutions in harmonic regime i.e. the unique solution $(\mathbf{E}, \mathbf{H})$ of:

$$
\left\{\begin{array}{l}
\operatorname{curl} \mathbf{E}=i \omega \mu_{0} \underline{\underline{\mu} \mathbf{H}}  \tag{5.3a}\\
\operatorname{curl} \mathbf{H}=-i \omega \varepsilon_{0} \underline{\underline{\varepsilon}} \mathbf{E}
\end{array}\right.
$$

such that the diffracted field $\left(\mathbf{E}^{\mathrm{d}}, \mathbf{H}^{\mathrm{d}}\right):=\left(\mathbf{E}-\mathbf{E}_{e}^{0}, \mathbf{H}-\mathbf{H}_{m}^{0}\right)$ satisfies an Outgoing Waves Condition (O.W.C. [9]) and where $\mathbf{E}$ and $\mathbf{H}$ are quasi-periodic functions with respect to the $x$ coordinate.

### 5.2.2 Theoretical developments of the method

### 5.2.2.1 Decoupling of fields and z-anisotropy

We assume that $\underline{\underline{\delta}}(x, y)$ is a $z$-anisotropic tensor field ( $\delta_{x z}=\delta_{y z}=\delta_{z x}=\delta_{z y}=0$ ). Moreover, the left upper matrix extracted from $\underline{\underline{\delta}}$ is denoted $\underline{\underline{\delta}}$, namely:

$$
\underline{\underline{\tilde{\delta}}}=\left(\begin{array}{cc}
\delta_{x x} & \bar{\delta}_{a}  \tag{5.4}\\
\delta_{a} & \delta_{y y}
\end{array}\right)
$$

For $z$-anisotropic materials, in a non-conical case, the problem of diffraction can be split into two fundamental cases (TE case and TM case). This property results from the following equality which can be easily derived:

$$
\begin{equation*}
-\operatorname{curl}\left(\underline{\underline{\delta}}^{-1} \operatorname{curl}(u \mathbf{z})\right)=\operatorname{div}\left(\underline{\underline{\delta}}^{T} / \operatorname{det}(\underline{\underline{\tilde{\delta}}}) \nabla u\right) \mathbf{z} \tag{5.5}
\end{equation*}
$$

where $u$ is a function which does not depend on the $z$ variable. Relying on the previous equality, it appears that the problem of diffraction in a non conical mounting amounts to looking for an electric (resp. magnetic) field which is polarized along the $z$-axis; $\mathbf{E}=e(x, y) \mathbf{z}$ (resp. $\mathbf{H}=h(x, y) \mathbf{z}$ ). The functions $e$ and $h$ are therefore solutions of similar differential equations:

$$
\begin{equation*}
\mathscr{L}_{\underline{\underline{\xi}}, \chi}(u):=\operatorname{div}(\underline{\underline{\xi}} \nabla u)+k_{0}^{2} \chi u=0 \tag{5.6}
\end{equation*}
$$

with

$$
\begin{equation*}
u=e, \quad \underline{\underline{\xi}}=\underline{\underline{\tilde{\mu}^{T}}} / \operatorname{det}(\underline{\underline{\tilde{\mu}}}), \quad \chi=\varepsilon_{z z} \tag{5.7}
\end{equation*}
$$

in the TM case and

$$
\begin{equation*}
u=h, \quad \underline{\underline{\xi}}=\underline{\underline{\tilde{\varepsilon}}}^{T} / \operatorname{det}(\underline{\underline{\tilde{\varepsilon}}}), \quad \chi=\mu_{z z} \tag{5.8}
\end{equation*}
$$

in the TE case.

### 5.2.2.2 Boiling down the diffraction problem to a radiation one

In its initial form, the diffraction problem summed up by Eq. (5.6) is not well suited to the Finite Element Method. In order to overcome this difficulty, we propose to split the unknown function $u$ into a sum of two functions $u_{1}$ and $u_{2}^{d}$, the first term being known as a closed form and the latter being a solution of a problem of radiation whose sources are localized within the obstacles.

We have assumed that outside the groove region (cf. Fig. 5.1), the tensor field $\underline{\underline{\xi}}$ and the function $\chi$ are constant and equal respectively to $\underline{\underline{\xi}}^{-}$and $\chi^{-}$in the substratum $(y<0)$ and equal respectively to $\underline{\underline{\xi}}^{+}$and $\chi^{+}$in the superstratum $\left(y>h^{g}\right)$. Besides, for the sake of clarity, the superstratum is supposed to be made of an isotropic and lossless material and is therefore solely defined by its relative permittivity $\varepsilon^{+}$and its relative permeability $\mu^{+}$, which leads to:

$$
\begin{equation*}
\underline{\underline{\xi}}^{+}=\frac{1}{\mu^{+}} \mathrm{Id}_{2} \quad \text { and } \quad \chi^{+}=\varepsilon^{+} \quad \text { in TE case } \tag{5.9}
\end{equation*}
$$

or

$$
\begin{equation*}
\underline{\underline{\xi}}^{+}=\frac{1}{\varepsilon^{+}} \mathrm{Id}_{2} \quad \text { and } \quad \chi^{+}=\mu^{+} \quad \text { in TM case }, \tag{5.10}
\end{equation*}
$$

where $\mathrm{Id}_{2}$ is the $2 \times 2$ identity matrix. With such notations, $\underline{\underline{\underline{\xi}}}$ and $\chi$ are therefore defined as follows:

$$
\underline{\underline{\xi}}(x, y):=\left\{\begin{array}{lll}
\xi^{+} & \text {for } & y>h^{g}  \tag{5.11}\\
\bar{\xi}^{g} \\
\overline{\bar{\xi}}^{-} & (x, y) & \text { for } \\
h^{g}>y>0 \\
\underline{\underline{\xi}} & \text { for } & y<0
\end{array}, \chi(x, y):=\left\{\begin{array}{lll}
\chi^{+} & \text {for } & y>h^{g} \\
\chi^{g}(x, y) & \text { for } h^{g}>y>0 \\
\chi^{-} & \text {for } & y<0
\end{array}\right.\right.
$$

It is now apropos to introduce an auxiliary tensor field $\underline{\underline{\xi}}_{1}$ and an auxiliary function $\chi_{1}$ :

$$
\underline{\underline{\xi}}(x, y):=\left\{\begin{array}{ll}
\underline{\xi}^{+} & \text {for } y>0  \tag{5.12}\\
\overline{\bar{\xi}}^{-} & \text {for } y<0
\end{array}, \chi_{1}(x, y):=\left\{\begin{array}{lll}
\chi^{+} & \text {for } y>0 \\
\chi^{-} & \text {for } y<0
\end{array}\right.\right.
$$

these quantities corresponding, of course, to a simple plane interface. Besides, we introduce the constant tensor field $\underline{\underline{\xi}}_{0}$ which is equal to $\underline{\underline{\xi}}^{+}$everywhere and a constant scalar field $\chi_{0}$ which is equal to $\chi^{+}$everywhere. Finally, we denote $u_{0}$ the function which equals the incident field $u^{\text {inc }}$ in the superstratum and vanishes elsewhere (see Fig. 5.1):

$$
u_{0}(x, y):=\left\{\begin{array}{ccc}
u^{\text {inc }} & \text { for } & y>h^{g}  \tag{5.13}\\
0 & \text { for } & y<h^{g}
\end{array}\right.
$$

We are now in a position to define more precisely the diffraction problem that we are dealing with. The function $u$ is the unique solution of:

$$
\begin{equation*}
\mathscr{L}_{\underline{\underline{\xi}}, \chi}(u)=0, \text { such that } u^{d}:=u-u_{0} \text { satisfies an O.W.C. } \tag{5.14}
\end{equation*}
$$

In order to reduce this diffraction problem to a radiation problem, an intermediate function is necessary. This function, called $u_{1}$, is defined as the unique solution of the equation:

$$
\begin{equation*}
\mathscr{L}_{\underline{\underline{\xi}},}, \chi_{1}\left(u_{1}\right)=0, \text { such that } u_{1}^{d}:=u_{1}-u_{0} \text { satisfies an O.W.C. } \tag{5.15}
\end{equation*}
$$

The function $u_{1}$ corresponds thus to an annex problem associated to a simple interface and can be solved in closed form and from now on is considered as a known function. As written above, we need the function $u_{2}^{d}$ which is simply defined as the difference between $u$ and $u_{1}$ :

$$
\begin{equation*}
u_{2}^{d}:=u-u_{1}=u^{d}-u_{1}^{d} . \tag{5.16}
\end{equation*}
$$

The presence of the superscript $d$ is, of course, not irrelevant: As the difference of two diffracted fields, the O.W.C. of $u_{2}^{d}$ is guaranteed (which is of prime importance when dealing with PML cf. 5.2.2.4). As a result, the Eq. (5.14) becomes:

$$
\begin{equation*}
\mathscr{L}_{\underline{\underline{\xi}}, \chi}\left(u_{2}^{d}\right)=-\mathscr{L}_{\underline{\underline{\xi}, \chi}}\left(u_{1}\right), \tag{5.17}
\end{equation*}
$$

where the right hand member is a scalar function which may be interpreted as a known source term $-\mathscr{S}_{1}(x, y)$ and the support of this source is localized only within the groove region. To prove it, all we have to do is to use Eq. (5.15):

$$
\begin{equation*}
\mathscr{S}_{1}:=\mathscr{L}_{\underline{\underline{\xi}}, \chi}\left(u_{1}\right)=\mathscr{L}_{\underline{\underline{\xi}}, \chi}\left(u_{1}\right)-\underbrace{\mathscr{L}_{\underline{\xi}}^{\underline{\xi}}, \chi_{1}}_{=0}\left(u_{1}\right)=\mathscr{L}_{\underline{\underline{\xi}}-\underline{\xi_{1}}, \chi-\chi_{1}}\left(u_{1}\right) . \tag{5.18}
\end{equation*}
$$

Now, let us point out that the tensor fields $\underline{\underline{\xi}}$ and $\underline{\underline{\xi}}$ are identical outside the groove region and the same holds for $\chi$ and $\chi_{1}$. The support of $\mathscr{S}_{1}$ is thus localized within the groove region as expected. It remains to compute more explicitly the source term $\mathscr{S}_{1}$. Making use of the linearity of the operator $\mathscr{L}$ and the equality $u_{1}=u_{1}^{d}+u_{0}$, the source term can be split into two terms

$$
\begin{equation*}
\mathscr{S}_{1}=\mathscr{S}_{1}^{0}+\mathscr{S}_{1}^{d} \tag{5.19}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathscr{S}_{1}^{0}=\mathscr{L}_{\underline{\underline{\xi}}-\underline{\underline{\xi}}}^{1}, \chi-\chi-\chi_{1}\left(u_{0}\right) \tag{5.20}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathscr{S}_{1}^{d}=\mathscr{L}_{\underline{\underline{\xi}}-\underline{\underline{\xi}}_{1}, \chi-\chi_{1}}\left(u_{1}^{d}\right) . \tag{5.21}
\end{equation*}
$$

Now, bearing in mind that $u_{0}$ is nothing but a plane wave $u_{0}=\exp \left(i \mathbf{k}_{p}^{+} \cdot \mathbf{r}\right)$ (with $\mathbf{k}_{p}^{+}=\alpha \mathbf{x}-$ $\beta^{+} \mathbf{y}$, it is sufficient to give $\nabla u_{0}=i \mathbf{k}_{p}^{+} u_{0}$ for the weak formulation associated with Eq. (5.17):

$$
\begin{equation*}
\mathscr{S}_{1}^{0}=\left\{i \operatorname{div}\left[\left(\underline{\underline{\xi}}^{+}-\underline{\underline{\xi}}\right) \mathbf{k}_{p}^{+} \exp \left(i \mathbf{k}_{p}^{+} \cdot \mathbf{r}\right)\right]+k_{0}^{2}\left(\chi^{+}-\chi\right) \exp \left(i \mathbf{k}_{p}^{+} \cdot \mathbf{r}\right)\right\} . \tag{5.22}
\end{equation*}
$$

The same holds for the term associated with the diffracted field. Since, in the superstrate, we have of course $u_{1}^{d}=\rho \exp \left(i \mathbf{k}_{c}^{+} \cdot \mathbf{r}\right)$ with $\mathbf{k}_{c}^{+}=\alpha \mathbf{x}+\beta^{+} \mathbf{y}$,

$$
\begin{equation*}
\mathscr{S}_{1}^{d}=\rho\left\{i \operatorname{div}\left[\left(\underline{\underline{\xi^{+}}}-\underline{\underline{\xi}}\right) \mathbf{k}_{c}^{+} \exp \left(i \mathbf{k}_{c}^{+} \cdot \mathbf{r}\right)\right]+k_{0}^{2}\left(\chi^{+}-\chi\right) \exp \left(i \mathbf{k}_{c}^{+} \cdot \mathbf{r}\right)\right\} \tag{5.23}
\end{equation*}
$$

where $\rho$ is simply the complex reflection coefficient associated with the simple interface:

$$
\rho=\frac{p^{+}-p^{-}}{p^{+}-p^{-}} \text {with } p^{ \pm}=\left\{\begin{array}{l}
\beta^{ \pm} \text {in the TM case }  \tag{5.24}\\
\frac{\beta^{ \pm}}{\varepsilon^{ \pm}} \text {in the TE case }
\end{array}\right.
$$

### 5.2.2.3 Quasi-periodicity and weak formulation

The weak formulation follows the classical lines and is based on the construction of a weighted residual of Eq. (5.6), which is multiplied by the complex conjugate of a weight function $u^{\prime}$ and integrated by part to obtain:

$$
\begin{equation*}
\mathscr{R}_{\underline{\underline{\xi}}, \chi}\left(u, u^{\prime}\right)=-\int_{\Omega}(\underline{\underline{\xi}} \nabla u) \cdot \nabla \overline{u^{\prime}}+k_{0}^{2} \chi u \overline{u^{\prime}} \mathrm{d} \Omega+\int_{\partial \Omega} \overline{u^{\prime}}(\underline{\underline{\xi}} \nabla u) \cdot \mathbf{n} \mathrm{d} S \tag{5.25}
\end{equation*}
$$

The solution $u$ of the weak formulation can therefore be defined as the element of the space $L^{2}(\operatorname{curl}, d, \alpha)$ of quasiperiodic functions (i.e. such that $u(x, y)=u_{\#}(x, y) e^{i \alpha x}$ with $u_{\#}(x, y)=$ $u_{\#}(x+d, y)$, a $d$-periodic function) of $L^{2}($ curl $)$ on $\Omega$ such that:

$$
\begin{equation*}
\mathscr{R}_{\underline{\underline{\xi}, ~},}\left(u, u^{\prime}\right)=0 \forall u^{\prime} \in L^{2}(\text { curl }, d, \alpha) . \tag{5.26}
\end{equation*}
$$

As for the boundary term introduced by the integration by part, it can be classically set to zero by imposing Dirichlet conditions on a part of the boundary (the value of $u$ is imposed and the weight function $u^{\prime}$ can be chosen equal to zero on this part of the boundary) or by imposing homogeneous Neumann conditions $(\underline{\underline{\xi}} \nabla u) \cdot \mathbf{n}=0$ on another part of the boundary (and $u$ is


Fig. 5.2: Quasi-periodicity of the field and sample of a d-periodic mesh.
therefore an unknown to be determined on the boundary). A third possibility are the so-called quasi-periodicity conditions of particular importance in the modeling of gratings.

Denote by $\Gamma_{l}$ and $\Gamma_{r}$ the lines parallel to the $y$-axis delimiting a cell of the grating (see Fig. 5.2) respectively from its left and right neighbor cell. Considering that both $u$ and $u^{\prime}$ are in $L^{2}(\operatorname{curl}, d, \alpha)$, the boundary term for $\Gamma_{l} \cup \Gamma_{r}$ is

$$
\begin{aligned}
& \int_{\Gamma_{l} \cup \Gamma_{r}} \overline{u^{\prime}}(\underline{\underline{\xi}} \nabla u) \cdot \mathbf{n} \mathrm{d} S=\int_{\Gamma_{l} \cup \Gamma_{r}} \overline{u_{\#}^{\prime}} e^{-i \alpha x}\left(\underline{\underline{\xi}} \nabla\left(u_{\#} e^{+i \alpha x}\right)\right) \cdot \mathbf{n} \mathrm{d} S= \\
& \int_{\Gamma_{l} \cup \Gamma_{r}} \overline{u_{\#}^{\prime}}\left(\underline{\underline{\xi}}\left(\nabla u_{\#}+i \alpha u_{\#} \mathbf{x}\right)\right) \cdot \mathbf{n} \mathrm{d} S=0,
\end{aligned}
$$

because the integrand $\overline{u_{\#}^{\prime}}\left(\underline{\underline{\xi}}\left(\nabla u_{\#}+i \alpha u_{\#} \mathbf{x}\right)\right) \cdot \mathbf{n}$ is periodic along $x$ and the normal $\mathbf{n}$ has opposite directions on $\Gamma_{l}$ and $\Gamma_{r}$ so that the contributions of these two boundaries have the same absolute value with opposite signs. The contribution of the boundary terms vanishes therefore naturally in the case of quasi-periodicity.

The finite element method is based on this weak formulation and both the solution and the weight functions are classically chosen in a discrete space made of linear or quadratic Lagrange elements, i.e. piecewise first or second order two variable polynomial interpolation built on a triangular mesh of the domain $\Omega$ (cf. Fig. 5.3a). Dirichlet and Neumann conditions may be used to truncate the PML domain in a region where the field (transformed by the PML) is negligible. The quasi-periodic boundary conditions are imposed by considering the $u$ as unknown on $\Gamma_{l}$ (in a way similar to the homogeneous Neumann condition case) while, on $\Gamma_{r}, u$ is forced equal to the value of the corresponding point on $\Gamma_{l}$ (i.e. shifted by a quantity $-d$ along $x$ ) up to the factor $e^{i \alpha d}$. The practical implementation in the finite element method is described in details in [10, 11]

### 5.2.2.4 Perfectly Matched Layer for z-anisotropic materials

The main drawback encountered in electromagnetism when tackling theory of gratings through the finite element method is the non-decreasing behaviour of the propagating modes in superstratum and substratum (if they are made of lossless materials): The PML has been introduced
by [4] in order to get round this obstacle. The computation of PML designed for $z$-anisotropic gratings is the topic of what follows.

In the framework of transformation optics, a PML may be seen as a change of coordinate corresponding to a complex stretch of the coordinate corresponding to the direction along which the field must decay [12, 13, 14]. Transformation optics have recently unified various techniques in computational electromagnetics such as the treatment of open problems, helicoidal geometries or the design of invisibility cloaks ([15]). These apparently different problems share the same concept of geometrical transformation, leading to equivalent material properties. A very simple and practical rule can be set up ([10]): when changing the coordinate system, all you have to do is to replace the initial materials properties $\underline{\underline{\varepsilon}}$ and $\underline{\underline{\mu}}$ by equivalent material properties $\underline{\underline{\varepsilon}}_{s}$ and $\underline{\underline{\mu}}_{s}$ given by the following rule:

$$
\begin{equation*}
\underline{\underline{\varepsilon}}_{s}=\mathbf{J}^{-1} \underline{\underline{\varepsilon}} \mathbf{J}^{-\mathrm{T}} \operatorname{det}(\mathbf{J}) \text { and } \underline{\underline{\mu_{s}}}=\mathbf{J}^{-1} \underline{\underline{\mu}} \mathbf{J}^{-\mathrm{T}} \operatorname{det}(\mathbf{J}) \tag{5.27}
\end{equation*}
$$

where $\mathbf{J}$ is the Jacobian matrix of the coordinate transformation consisting of the partial derivatives of the new coordinates with respect to the original ones $\left(\mathbf{J}^{-T}\right.$ is the transposed of its inverse).
In this framework, the most natural way to define PMLs is to consider them as maps on a complex space $\mathbb{C}^{3}$, which coordinate change leads to equivalent permittivity and permeability tensors. We detail here the different coordinates used in this section.

- $(x, y, z)$ are the cartesian original coordinates.
- $\left(x_{s}, y_{s}, z_{s}\right)$ are the complex stretched coordinates. A suitable subspace $\Gamma \subset \mathbb{C}^{3}$ is chosen (with three real dimensions) such that $\left(x_{s}, y_{s}, z_{s}\right)$ are the complex valued coordinates of a point on $\Gamma$ (e.g. $x=\mathfrak{R e}\left(x_{s}\right), y=\mathfrak{R e}\left(y_{s}\right), z=\mathfrak{R e}\left(z_{s}\right)$ ).
- $\left(x_{c}, y_{c}, z_{c}\right)$ are three real coordinates corresponding to a real valued parametrization of $\Gamma \subset \mathbb{C}^{3}$.

We use rectangular PMLs ([12]) absorbing in the $y$-direction and we choose a diagonal matrix $\mathbf{J}=\operatorname{diag}\left(1, s_{y}(y), 1\right)$, where $s_{y}(y)$ is a complex-valued function of the real variable $y$, defined by:

$$
\begin{equation*}
y_{s}(y)=\int_{0}^{y} s_{y}\left(y^{\prime}\right) \mathrm{dy}^{\prime} . \tag{5.28}
\end{equation*}
$$

The expression of the equivalent permittivity and permeability tensors are thus:

$$
\underline{\underline{\varepsilon}}_{s}=\left(\begin{array}{ccc}
s_{y} \varepsilon_{x x} & \overline{\varepsilon_{a}} & 0  \tag{5.29}\\
\varepsilon_{a} & s_{y}^{-1} \varepsilon_{y y} & 0 \\
0 & 0 & s_{y} \varepsilon_{z z}
\end{array}\right) \quad \text { and } \quad \underline{\underline{\mu_{s}}}=\left(\begin{array}{ccc}
s_{y} \mu_{x x} & \overline{\mu_{a}} & 0 \\
\mu_{a} & s_{y}^{-1} \mu_{y y} & 0 \\
0 & 0 & s_{y} \mu_{z z}
\end{array}\right) .
$$

Note that the equivalent medium has the same impedance than the original one as $\underline{\underline{\varepsilon}}$ an $\underline{\underline{\mu}}$ are transformed in the same way, which guarantees that the PML is perfectly reflectionless. $=$ Now, let us define the so-called substituted field $\boldsymbol{F}_{s}=\left(\mathbf{E}_{s}, \mathbf{H}_{s}\right)$, solution of Eqs. (5.3) with $\boldsymbol{\xi}=\xi_{s}$ and $\chi=\chi_{s}$. It turns out that $\boldsymbol{F}_{s}$ equals the field $\boldsymbol{F}$ in the region $y^{b}<y<y^{t}$ (with $y^{b}=-h^{-}$ $\overline{\bar{a}} \mathrm{a} d \vec{y}^{+}=h^{g}+h^{+}$, see Fig. 5.3a), provided that $s_{y}(y)=1$ in this region. The main feature of this latest field $\boldsymbol{F}_{s}$ is the remarkable correspondence with the first field $\boldsymbol{F}$; whatever the function $s_{y}$ provided that it equals 1 for $y^{t}<y<y^{b}$, the two fields $\boldsymbol{F}$ and $\boldsymbol{F}_{s}$ are identical in the region $y^{t}<y<y^{b}[8]$. In other words, the PML is completely reflection-less. In addition, for complex
valued functions $s_{y}\left(\mathfrak{I} m\left\{s_{y}\right\}\right.$ strictly positive in PML), the field $\boldsymbol{F}_{s}$ converges exponentially towards zero (as $y$ tends to $\pm \infty$, cf. Fig. 5.3c and 5.3d) although its physical counterpart $\boldsymbol{F}$ does not. Note that in Fig. 5.3d, the value of the computed radiated field $u_{2}^{d}$ on each extreme boundary of the PMLs is at least $10^{-8}$ weaker than in the region of interest. As a consequence, $\boldsymbol{F}_{s}$ is of finite energy and for this substituted field a weak formulation can be easily derived which is essential when dealing with Finite Element Method.

Still remains to give a suitable function $s_{y}$. Let us consider the complex coordinate mapping $y\left(y_{c}\right)$, which is simply defined as the derivative of the stretching coefficient $s_{y}(y)$ with respect to $y_{c}$. With simple stretching functions, we can obtain a reliable criterion upon proper fields decay. A classical choice is:

$$
s_{y}(y)= \begin{cases}\zeta^{-} & \text {if } y<y^{b}  \tag{5.30}\\ 1 & \text { if } y^{b}<y<y^{t} \\ \zeta^{+} & \text {if } y>y^{t}\end{cases}
$$

where $\zeta^{ \pm}=\zeta^{\prime}, \pm+i \zeta^{\prime \prime}, \pm$ are complex constants with $\zeta^{\prime \prime}, \pm>0$.
In that case, the complex valued function $y\left(y_{c}\right)$ defined by Eq. (5.28) is explicitly given by:

$$
y\left(y_{c}\right)= \begin{cases}y^{b}+\zeta^{-}\left(y_{c}-y^{b}\right) & \text { if } y_{c}<y^{b}  \tag{5.31}\\ y_{c} & \text { if } y^{b}<y_{c}<y^{t}, \\ y^{t}+\zeta^{+}\left(y_{c}-y^{t}\right) & \text { if } y_{c}>y^{t}\end{cases}
$$

Finally, let us consider a propagating plane wave in the substratum $u_{n}(x, y):=\exp (i(\alpha x-$ $\left.\beta_{n}^{-} y\right)$ ). Its expression can be rewritten as a function of the stretched coordinates in the PML as follows:

$$
\begin{equation*}
u_{n}^{\mathrm{sc}}\left(x_{c}, y_{c}\right):=u_{n}\left(x\left(x_{c}\right), y\left(y_{c}\right)\right)=e^{i \alpha x_{c}} e^{-i \beta_{n}^{-}\left(y^{b}+\zeta^{-}\left(y_{c}-y^{b}\right)\right)} \tag{5.32}
\end{equation*}
$$

The behavior of this latest function along the $y_{c}$ direction is governed by the function $U^{\mathrm{sc}}\left(y_{c}\right):=$ $e^{-i \beta_{n}^{-} \zeta^{-} y_{c}}$. Letting $\beta_{n}^{\prime,-}:=\mathfrak{R e}\left\{\beta_{n}^{-}\right\}, \beta_{n}^{\prime \prime,-}:=\mathfrak{I} m\left\{\beta_{n}^{-}\right\}, \zeta^{\prime,-}:=\mathfrak{R e}\left\{\zeta^{-}\right\}$and $\zeta^{\prime \prime,-}:=\mathfrak{I} m\left\{\zeta^{-}\right\}$, the non-oscillating part of the function $U^{\mathrm{sc}}\left(y_{c}\right)$ is given by $\exp \left(\left(\beta_{n}^{\prime,-} \zeta^{\prime \prime,-}+\beta_{n}^{\prime \prime,-} \zeta^{\prime,-}\right) y_{c}\right)$. Keeping in mind that $\beta_{n}^{\prime,-}$ and/or $\beta_{n}^{\prime \prime,-}$ are positive numbers, the function $U^{\text {sc }}$ decreases exponentially towards zero as $y_{c}$ tends to $-\infty$ (Fig. 5.3d) provided that $\zeta^{-}$belongs to $\mathbb{C}^{+}:=\{z \in$ $\mathbb{C}, \mathfrak{R e} e\{z\}>0$, and $\mathfrak{I} m\{z\}>0\}$. In the same way, it can be shown that $\zeta^{+}$belongs to $\mathbb{C}^{+}$.

Let us conclude this section with two important remarks:

1. Practical choice of PML parameters. As for the complex stretch parameters, setting $\zeta^{ \pm}=1+i$ is usually a safe choice. For computational needs, the PML has to be truncated and the other constitutive parameter of the PML is its thickness $\hat{h}$ (see Fig. 5.3a). Setting $\hat{h}^{ \pm}=\lambda_{0} / \sqrt{\varepsilon^{ \pm}}$leads to a PML thick enough to "absorb" all incident radiation. These specific values will be used in the sequel, unless otherwise specified.
2. Special cases. The reader will notice that a configuration where $\beta_{n}^{\prime,-}$ is a very weak positive number compared to $k_{0}$ with $\beta_{n}^{\prime \prime,-}$ (this is precisely the case of a plane wave at grazing incidence on the bottom PML) leads to a very slow exponential decay of $U^{\text {sc }}$. In such a case, close to so-called Wood's anomalies or at extreme grazing incidences, classical PML fail. We will address this tricky situation extensively in Section 5.2.4.


Fig. 5.3: Example of computation of the radiated field $u_{2}^{d}$ (TM case).

### 5.2.2.5 Synthesis of the method

In order to give a general view of the method, all information is collected here that is necessary to set up the practical Finite Element Model. First of all, the computation domain $\Omega$ (cf. Fig. 5.3a) corresponds to a truncated cell of the grating which is a finite rectangle divided into five horizontal layers. These layers are respectively from top to bottom upper PML, the superstratum, the groove region, the substratum, and the lower PML. The unknown field is the scalar function $u_{2}^{d}$ defined in Eq. (5.16). Its finite element approximation is based on the second Lagrange elements built on a triangle meshing of the cell (cf. Fig. 5.3b). A complex algebraic system of linear equations is constructed via the Galerkin weighted residual method, i.e. the set of weight functions $u^{\prime}$ is chosen as the set of shape functions of interpolation on the mesh [10].

- In region 1 (upper PML, see Fig. 5.3a),

$$
\begin{equation*}
\mathscr{R}_{\underline{\underline{\xi}}}=, \chi_{s}^{+}\left(u_{2}^{d}, u^{\prime}\right)=0, \tag{5.33}
\end{equation*}
$$

with $\underline{\underline{\xi}}_{s}^{+}$and $\chi_{s}^{+}$depending on the equivalent anisotropic properties of the PML given by Eq. $\left(\overline{\overline{5.7}}{ }^{\prime}\right.$, Eq. (5.8) and Eqs. (5.29).

- In region 2 (superstratum),

$$
\begin{equation*}
\mathscr{R}_{\underline{\underline{\xi}}^{+}}, \chi^{+}\left(u_{2}^{d}, u^{\prime}\right)=0, \tag{5.34}
\end{equation*}
$$

with $\underline{\xi}^{+}$and $\chi^{+}$depending on the homogeneous isotropic properties of the superstratum given by Eq. (5.7), Eq. (5.8), Eq. (5.9) and Eq. (5.10).

- In region 3 (groove region),

$$
\begin{equation*}
\mathscr{R}_{\underline{\xi^{g}}, \chi^{g}}\left(u_{2}^{d}, u^{\prime}\right)=-\mathscr{R}_{\underline{\xi}^{g}}^{\underline{g}}, \chi^{g}\left(\mathscr{S}_{1}, u^{\prime}\right), \tag{5.35}
\end{equation*}
$$

with $\underline{\xi}^{g}$ and $\chi^{g}$ depending on the heterogeneous possibly anisotropic properties given by Eq. ( $\overline{\overline{5}} .7$ ), Eq. (5.8), Eq. (5.11) and $\mathscr{S}_{1}$ given by Eq. (5.19), Eq. (5.22), Eq. (5.23) and Eq. (5.24).

- In region 4 (substratum),

$$
\begin{equation*}
\mathscr{R}_{\underline{\xi^{-}}, \chi^{-}}\left(u_{2}^{d}, u^{\prime}\right)=0, \tag{5.36}
\end{equation*}
$$

with $\underline{\underline{\xi}}^{-}$and $\chi^{-}$depending on the homogeneous isotropic properties of the substratum given by Eq. (5.7), Eq. (5.8), Eq. (5.9) and Eq. (5.10).

- In region 5 (lower PML),

$$
\begin{equation*}
\mathscr{R}_{\underline{\underline{\xi}}^{-}}, \chi_{s}^{-}\left(u_{2}^{d}, u^{\prime}\right)=0, \tag{5.37}
\end{equation*}
$$

with $\xi^{-}$and $\chi_{s}^{-}$depending on the equivalent anisotropic properties of the PML given by Eq. (5.7), Eq. (5.8) and Eqs. (5.29).

### 5.2.2.6 Energy balance: Diffraction efficiencies and absorption

The rough result of the FEM calculation is the complex radiated field $u_{2}^{d}$. Using Eq. (5.16), it is straightforward to obtain the complex diffracted field $u^{d}$ solution of Eq. (5.14) at each point of the bounded domain. We deduce from $u^{d}$ the diffraction efficiencies with the following method. The superscripts ${ }^{+}$(resp. ${ }^{-}$) correspond to quantities defined in the superstratum (resp. substratum) as previously.

On the one hand, since $u^{d}$ is quasi-periodic along the $x$-axis, it can be expanded as a Rayleigh expansion (see for instance [9]):

$$
\begin{equation*}
\text { for } y<0 \text { and } y>h^{g}, u^{d}(x, y)=\sum_{n \in \mathbb{Z}} u_{n}^{d}(y) e^{i \alpha_{n} x} \tag{5.38}
\end{equation*}
$$

where

$$
\begin{equation*}
u_{n}^{d}(y)=\frac{1}{d} \int_{-d / 2}^{d / 2} u^{d}(x, y) e^{-i \alpha_{n} x} \mathrm{~d} x \text { with } \alpha_{n}=\alpha+\frac{2 \pi}{d} n \tag{5.39}
\end{equation*}
$$

On the other hand, introducing Eq. (5.38) into Eq. (5.6) leads to the Rayleigh coefficients:

$$
u_{n}^{d}(y)=\left\{\begin{array}{ll}
u_{n}^{+}(y)=r_{n} e^{i \beta_{n}^{+} y}+a_{n} e^{-i \beta_{n}^{+} y} & \text { for } y>h^{g}  \tag{5.40}\\
u_{n}^{-}(y)=t_{n} e^{-i \beta_{n}^{-} y}+b_{n} e^{i \beta_{n}^{-} y} & \text { for } y<0
\end{array} \text { with } \beta_{n}^{ \pm^{2}=k^{ \pm^{2}}-\alpha_{n}^{2} .}\right.
$$

For a temporal dependance in $e^{-i \omega t}$, the O.W.C. imposes $a_{n}=b_{n}=0$. Combining Eq. (5.39) and (5.40) at a fixed $y_{0}$ altitude leads to:

$$
\left\{\begin{array}{l}
r_{n}=\frac{1}{d} \int_{-d / 2}^{d / 2} u^{d}\left(x, y_{0}\right) e^{-i\left(\alpha_{n} x+\beta_{n}^{+} y_{0}\right)} \mathrm{d} x \text { for } y_{0}>h^{g}  \tag{5.41}\\
t_{n}=\frac{1}{d} \int_{-d / 2}^{d / 2} u^{d}\left(x, y_{0}\right) e^{-i\left(\alpha_{n} x-\beta_{n}^{-} y_{0}\right)} \mathrm{d} x \text { for } y_{0}<0
\end{array}\right.
$$

We extract these two coefficients by trapezoidal numerical integration along $x$ from a cutting of the previously calculated field map at $y_{0}$. It is well known that the mere trapezoidal integration method is very efficient for smooth and periodic functions (integration on one period) [16]. Now the restriction on a horizontal straight line crossing the whole cell in homogeneous media (substratum and superstratum) is of $C^{\infty}$ class. From a numerical point of view, it appears that the interpolated approximation of the unknown function, namely $u_{2}^{d}$ preserves the good behaviour of the numerical computation of these integrals. From this we immediately deduce the reflected and transmitted diffracted efficiencies of propagative orders ( $T_{n}$ and $R_{n}$ ) defined by:

$$
\left\{\begin{array}{c}
R_{n}:=r_{n} \overline{r_{n}} \frac{\beta_{n}^{+}}{\beta^{+}} \quad \text { for } \quad y_{0}>h^{g}  \tag{5.42}\\
T_{n}:=t_{n} \overline{t_{n}} \frac{\beta_{n}^{-}}{\beta^{-}} \frac{\gamma^{+}}{\gamma^{+}}
\end{array} \text {for } \quad y_{0}<0 . ~ \text { with } \gamma^{ \pm}=\left\{\begin{array}{l}
1 \text { in the TM case } \\
\varepsilon^{ \pm} \text {in the TE case }
\end{array}\right.\right.
$$

This calculation is performed at several different $y_{0}$ altitudes in the superstratum and the substratum, and the mean value found for each propagative transmitted or reflected diffraction order is presented in the numerical experiments of the following section.

Normalized losses $Q$ can be obtained according to Poynting's theorem through the straightforward computation of the following ratio:

$$
\begin{equation*}
Q:=\frac{\int_{S} \omega \varepsilon_{0} \Im m\left(\varepsilon^{g^{\prime}}\right) \mathbf{E} \cdot \overline{\mathbf{E}} \mathrm{d} s}{\int_{L} \mathfrak{R} e\left\{\mathbf{E}_{0} \times \overline{\mathbf{H}_{0}}\right\} \cdot \mathbf{n} \mathrm{d} l}, \tag{5.43}
\end{equation*}
$$

The numerator in Eq. (5.43) clarifies losses in Watts by period of the considered grating and are computed by integrating the Joule effect losses density over the surface $S$ of the lossy element. The denominator normalizes these losses to the incident power, i.e. the time-averaged incident Poynting vector flux across one period (a straight line $L$ of length $d$ in the superstrate parallel to $O x$, whose normal oriented along decreasing values of $y$ is denoted $\mathbf{n}$ ).

Finally, combining Eqs. (5.42) and Eq. (5.43), a self consistency check of the whole numerical scheme consists in comparing the quantity $B$ :

$$
\begin{equation*}
B:=\sum_{n} T_{n}+\sum_{m} R_{m}+Q \tag{5.44}
\end{equation*}
$$

to unity. In Eq. (5.44), the summation indexed by $n$ (resp. $m$ ) corresponds to the sum over the efficiencies of all transmitted (resp. reflected) propagative diffraction orders in the substrate (resp. superstrate). We give interpretations and concrete examples of such numerical energy balances over non trivial grating profiles in sections 5.2.3.2 and 5.2.3.3.

### 5.2.3 Numerical experiments

### 5.2.3.1 Numerical validation of the method

We can refer to [17] in order to test the accuracy of our method. The studied grating is isotropic, since we lack numerical values in the literature in anisotropic cases. We compute the following problem (cf. Fig. 5.4), as described in [18] and [17]. The wavelength of the plane wave is set to $1 \mu m$ and is incoming with an angle of $\pi / 6$ with respect to the normal to the grating.

We present the $R_{0}$ efficiency (cf. Table 5.1) in both cases of polarization versus the mesh refinement. So we have a good agreement to the reference values, and the accuracy reached is independent from the polarization case.


Fig. 5.4: Rectangular groove grating: This pattern is repeatedly set up with a period $d=1 \mu \mathrm{~m}$. This grating has been studied by [17] and is one of our points of reference

| Maximum element size | $R_{0}^{\mathrm{TE}}$ | $R_{0}^{\mathrm{TM}}$ |
| :---: | :---: | :---: |
| $\lambda_{0} /(4 \sqrt{\varepsilon})$ | 0.7336765 | 0.8532342 |
| $\lambda_{0} /(6 \sqrt{\varepsilon})$ | 0.7371302 | 0.8456592 |
| $\lambda_{0} /(8 \sqrt{\varepsilon})$ | 0.7347466 | 0.8482817 |
| $\lambda_{0} /(10 \sqrt{\varepsilon})$ | 0.7333739 | 0.850071 |
| $\lambda_{0} /(12 \sqrt{\varepsilon})$ | 0.7346569 | 0.8494844 |
| $\lambda_{0} /(14 \sqrt{\varepsilon})$ | 0.7341944 | 0.8483238 |
| $\lambda_{0} /(16 \sqrt{\varepsilon})$ | 0.7342714 | 0.8484774 |
| Result given by $[17]$ | 0.7342789 | 0.8484781 |

Tab. 5.1: Reflected efficiencies versus mesh refinement. Note that the efficiencies are properly computed (two significant digits) even for a rather coarse mesh.

### 5.2.3.2 Experiment set up based on existing materials

The method proposed in this section is adapted to $z$-anisotropic materials, such as transparent $\mathrm{CaCO}_{3}$ [19], $\mathrm{LiNbO}_{3}$ [20] or Ni:YIG [21] and lossy CoPt or CoPd [22]. Let us now consider a trapezoidal (cf. Fig. 5.5) anisotropic grating made of aragonite $\left(\mathrm{CaCO}_{3}\right)$ deposited on an isotropic substratum $\left(\mathrm{SiO}_{2}, \varepsilon_{\mathrm{SiO}_{2}}=2.25\right)$. Along the anisotropic crystal axis, its dielectric tensor can be written as follows [19]:

$$
\underline{\underline{\varepsilon}}_{\mathrm{CaCO}_{3}}=\left(\begin{array}{ccc}
2.843 & 0 & 0  \tag{5.45}\\
0 & 2.341 & 0 \\
0 & 0 & 2.829
\end{array}\right) \quad \text { and } \quad \underline{\underline{\mu}} \mathrm{CaCO}_{3}=\left(\begin{array}{ccc}
\mu_{0} & 0 & 0 \\
0 & \mu_{0} & 0 \\
0 & 0 & \mu_{0}
\end{array}\right)
$$



Fig. 5.5: Diffractive element pattern. This element is made of aragonite for which the dielectric tensor is given by Eq. (5.46) and is deposited on a silica substrate with a period d $=600 \mathrm{~nm}$.

Now let's assume that the natural axis of our aragonite grating are rotated by $45^{\circ}$ around
the grating infinite dimension. The dielectric tensor becomes:

$$
\underline{\underline{\varepsilon}}_{\mathrm{CaCO}_{3}}^{45^{\circ}}=\left(\begin{array}{ccc}
2.592 & 0.251 & 0  \tag{5.46}\\
0.251 & 2.592 & 0 \\
0 & 0 & 2.829
\end{array}\right)
$$

We shall here remind that our method remains strictly the same whatever the diffractive element geometry is. The 2D computational domain is bounded along the $y$-axis by the PMLs and along the $x$ since we consider only one pseudo period. We propose to calculate the diffractive efficiencies at $\lambda_{0}=633 \mathrm{~nm}$ in both polarization cases TE and TM, and for different incoming incidences $\left(0^{\circ}, 20^{\circ}\right.$ and $\left.40^{\circ}\right)$. Since both $\underline{\underline{\mu}}$ and $\underline{\underline{\varepsilon}}$ are Hermitian, the whole incident energy is diffracted and the sum of theses efficiencies ought to be equal to the incident energy, which will stand for validation of our numerical calculation.

Finally, the resulting bounded domain is meshed with a maximum mesh element side size of $\lambda_{0} / 10 \sqrt{\varepsilon}$. Efficiencies are still post-processed in accordance with the calculation presented section 5.2.2.6.

| $\mathbf{T M}$ | $T_{-2}$ | $T_{-1}$ | $T_{0}$ | $T_{1}$ | $R_{-1}$ | $R_{0}$ | $R_{1}$ | total |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $0^{\circ}$ | - | 0.203133 | 0.585235 | 0.203138 | - | 0.008473 | - | 0.999978 |
| $20^{\circ}$ | - | 0.399719 | 0.575625 | 0.004643 | 0.004412 | 0.015630 | - | 1.000029 |
| $40^{\circ}$ | 0.025047 | 0.420714 | 0.493491 | - | 0.002541 | 0.058238 | - | 1.000031 |
| $\mathbf{T E}$ | $T_{-2}$ | $T_{-1}$ | $T_{0}$ | $T_{1}$ | $R_{-1}$ | $R_{0}$ | $R_{1}$ | total |
| $0^{\circ}$ | - | 0.322510 | 0.538165 | 0.124722 | - | 0.014683 | - | 1.000080 |
| $20^{\circ}$ | - | 0.538727 | 0.444403 | 0.000369 | 0.005372 | 0.011180 | - | 1.000051 |
| $40^{\circ}$ | 0.012058 | 0.434191 | 0.541090 | - | 0.005032 | 0.007686 | - | 1.000057 |

Tab. 5.2: Transmitted and reflected efficiencies of propagative orders deduced from field maps shown
Fig. 5.6

At normal incidence, the $h$ field in the TE case (cf. Fig. 5.6b) is non symmetric whereas the $e$ field in the TM case is (cf. Fig. 5.6a). This is illustrated by the obvious non-symmetry of $T_{-1}^{\mathrm{TE}}$ and $T_{1}^{\mathrm{TE}}$ (cf. Table 5.2: 0.322510 versus $0.124722!$ ), whereas $T_{-1}^{\mathrm{TM}}=T_{1}^{\mathrm{TM}}=0.20313$.


Fig. 5.6: Real part of the total calculated field depending on $\theta_{0}$ and the polarization case

### 5.2.3.3 A non trivial geometry

Since the beginning of this chapter, we have laid great stress upon the independence of the method towards the geometry of the pattern. But we have considered so far diffractive objects of simple trapezoidal section. Let us tackle a way more challenging case and see what this approach is made of.

We can obtain an quite winding shape by extracting the contrast contour of an arbitrary image (see Fig. 5.7a-5.7b). The contour is approximated by a set of splines, and the resulting domain is finely meshed (Fig. 5.7c). Finally, as shown in Fig. 5.7b, the formed pattern ( $h^{g} / \lambda_{0}=$ $1.68)$, breathing in free space $\left(\varepsilon_{\text {substrat }}=1\right)$, is supposed to be periodically repeated $d / \lambda_{0}=1.26$ on a plane ground of glass ( $\varepsilon_{\mathrm{SiO}_{2}}=2.25$ ). The element is considered to be "made of" a lossy material of high optical index ( $\varepsilon_{\text {marsu }}=40+0.1 i$ ). The response of this system to a incident $s$-polarized plane wave at oblic incidence $\left(\theta_{0}=-30^{\circ}\right)$ is finally calculated. The real part of the quasi-periodic total field is represented in Fig. 5.7d for three periods.

Indeed, we do not have any tabulated data available to check our results. But what we
do have is a pretty reliable consistency check through the computation of the energy balance described by Eqs. (5.42) and (5.43). As shown in Fig. 5.7e, we obtain at least 7 significative digits on the energetic values. The total balance of 1.00000019 is computed taking into account (i) values of the total field inside the diffractive elements, (ii) values of the diffracted field at altitudes spanning the whole (modeled) superstrate, (iii) values of the total field at altitudes spanning the entire (modeled) substrate. Finally, (iv) the calculated field $u_{2}^{d}$ also nicely decays exponentially inside both PML. These four points allow us to check a posteriori the validity of the field everywhere in the computation cell.

(e) ENERGY BALANCE

| $R_{0}$ | 0.13436549 |
| :---: | :---: |
| $R_{-1}$ | 0.10229352 |
| $T_{-2}$ | 0.00802646 |
| $T_{-1}$ | 0.00870891 |
| $T_{0}$ | 0.00666213 |
| $T_{+1}$ | 0.00359647 |
| $Q$ | 0.73634721 |
| TOTAL | $\mathbf{1 . 0 0 0 0 0 0 1 9}$ |

Fig. 5.7: (a) Initial contrasted image. (b) Proposed set up. (c) Sample mesh. (d) $\mathfrak{R e}\left\{E_{z}\right\}$ in $V / m$. (e) Energy balance of the problem.

### 5.2.4 Dealing with Wood anomalies using Adaptative PML

As we have noticed at the end of Section 5.2.2.4, PMLs based on "traditional coordinate stretching" are inefficient for periodic problems when dealing with grazing angles of diffracted orders, i.e. when the frequency is near a Wood's anomaly ( $[23,24]$ ), leading to spurious reflexions and thus numerical pollution of the results. An important question in designing absorbing layers is thus the choice of their parameters: The PML thickness and the absorption coefficient. To this aim, adaptative formulations have already been set up, most of them employing a posteriori error estimate [25, 18, 26]. In this section, we propose Adaptative PMLs (APMLs) with a suitable coordinate stretching, depending both on incidence and grating parameters, capable of efficiently absorbing propagating waves with nearly grazing angles. This section is dedicated to the mathematical formulation used to determine PML parameters adapted to any diffraction orders. We provide at the end a numerical example of a dielectric slit grating showing the relevance of our approach in comparison with classical PMLs.

### 5.2.4. $\quad$ Skin depth of the $P M L$



Fig. 5.8: The basic cell used for the FEM computation of the diffracted field $u_{2}^{d}$.
As explained in Section 5.2.2.6, the diffracted field $u^{d}$ can be expanded as a Rayleigh expansion, i.e. into an infinite sum of propagating and evanescent plane waves called diffraction orders. As detailed at the end of Section 5.2.2.4, we are now in position to rewrite easily the expression of, say, a transmitted diffraction order into the substrate. Similar considerations also apply to the reflected orders in the top PML. Combining Eq. (5.32) and (5.40) lead to the expression $u_{n, s}^{-}\left(y_{c}\right)$ of a transmitted propagative order inside the PML:

$$
u_{n, s}^{-}\left(y_{c}\right)=u_{n}^{-}\left(y\left(y_{c}\right)\right)=t_{n} \mathbf{e}^{-i \beta_{n}^{-}\left[y^{t}+\zeta^{-}\left(y_{c}-y^{t}\right)\right]} .
$$

The non oscillating part of this function is given by:

$$
U_{n}^{-}(y)=t_{n} \exp \left(\left(\beta_{n}^{\prime,--} \zeta^{\prime \prime,-}+\beta_{n}^{\prime \prime,-} \zeta^{\prime,-}\right) y_{c}\right),
$$

where $\beta_{n}^{-}=\beta_{n}^{\prime,-}+i \beta_{n}^{\prime \prime,-}$. For a propagating order we have $\beta_{n}^{\prime,-}>0$ and $\beta_{n}^{\prime,-}=0$, while for an evanescent order $\beta_{n}^{\prime,-}=0$ and $\beta_{n}^{\prime \prime,-}>0$. It is thus sufficient to take $\zeta^{\prime \prime,-}>0$ and $\zeta^{\prime \prime,-}>0$ to ensure the exponential decay to zero of the field inside the PML if it was of infinite extent. But, of course, for practical purposes, the thickness of the PML is finite and has to be suitably chosen. Two pitfalls must be avoided:

1. The PML thickness is chosen too small compared to the skin depth. As a consequence, the electromagnetic wave cannot be considered as vanishing: An incident electromagnetic "sees the bottom of the PML". In other words, this PML of finite thickness is no longer reflection-less.
2. The PML thickness is chosen much larger than the skin depth. In that case, a significant part of the PML is not useful, which gives rise to the resolution of linear systems of unnecessarily large dimensions.

Then remains to derive the skin depth, $l_{n}^{-}$, associated with the propagating order $n$. This characteristic length is defined as the depth below the PML at which the field falls to $1 / e$ of its value near the surface:

$$
U_{n}^{-}\left(y-l_{n}^{-}\right)=\frac{U_{n}^{-}(y)}{e} .
$$

Finally, we find $l_{n}^{-}=\left(\beta_{n}^{\prime,--} \zeta^{\prime \prime,-}+\beta_{n}^{\prime \prime,-} \zeta^{\prime,--}\right)^{-1}$ and we define $l^{-}$as the largest value among the $l_{n}^{-}$:

$$
l^{-}=\max _{n \in \mathbb{Z}} l_{n}^{-} .
$$

The height of the bottom PML region is set to $\hat{h}^{-}=10 l^{-}$.

### 5.2.4.2 Weakness of the classical PML for grazing diffracted angles

Let us consider the (bottom) PML adapted to the substrate. Similar conclusions will hold for the top PML. The efficiency of the classical PML fails for grazing diffracted angles, in other words when a given order appears/vanishes: this is the so-called Wood's anomaly, well known in the grating theory. In mathematical terms, there exists $n_{0}$ such that $\beta_{n_{0}}^{-} \simeq 0$. The skin depth of the PML then becomes very large. To compensate this, it is tempting to increase the value of $\zeta^{\prime \prime \prime}$, , but it would lead to spurious numerical reflections due to an overdamping. For a fixed value of $\hat{h}^{-}$, if $\zeta^{\prime \prime,-}$ is too weak, the absorption in the PMLs is insufficient and the wave is reflected on the outward boundary of the PML. To illustrate these typical behaviors (cf. Fig. 5.9), we compute the field diffracted by a grating with a rectangular cross section of height $h^{g}=1.5 \mu \mathrm{~m}$ and width $L^{g}=3 \mu \mathrm{~m}$ with $\varepsilon^{g}=11.7$, deposited on a substrate with permittivity $\varepsilon^{-}=2.25$. The structure is illuminated by a $p$-polarized plane wave of wavelength $\lambda_{0}=10 \mu \mathrm{~m}$ and of angle of incidence $\theta_{0}=10^{\circ}$ in the air $\left(\varepsilon^{+}=1\right)$. All materials are non magnetic ( $\mu_{r}=1$ ) and the periodicity of the grating is $d=4 \mu \mathrm{~m}$. We set $\hat{h}^{-}=10 l_{0}^{-}$and $\zeta^{\prime,-}=1$.


Fig. 5.9: Zero ${ }^{\text {th }}$ transmitted order by a grating with a rectangular cross section (see parameters in text, part 5.2.4.2) for different values of $\zeta^{\prime \prime},-$ : blue line, $\zeta^{\prime \prime,-}=1$, correct damping; green line, $\zeta^{\prime \prime,-}=0.1$, underdamping; red line, $\zeta^{\prime \prime,-}=20$, overdamping.

### 5.2.4.3 Construction of an adaptative PML

To overcome the problems pointed out in the previous section, we propose a coordinate stretching that rigorously treats the problem of Wood's anomalies. The wavelengths "seen" by the system are very different depending on the order at stake:

- if the diffracted angle $\theta_{n}$ is zero, the apparent wavelength $\lambda_{0} / \cos \theta_{n}$ is simply the incident wavelength,
- if the diffracted angle is near $\pm \pi / 2$ (grazing angle), the apparent wavelength $\lambda_{0} / \cos \theta_{n}$ is very large.

Thus if a classical PML is adapted to one diffracted order, it will not be for another, and vice versa. The idea behind the APML is to deal with each and every order when progressing in the absorbing medium.

Once again the development will be conducted only for the PML adapted to the substrate. We consider a real-valued coordinate mapping $y_{d}(y)$, the final complex-valued mapping is then $y_{c}(y)=\zeta^{-} y_{d}(y)$, with the complex constant $\zeta^{-}$, with $\zeta^{\prime,-}>0$ and $\zeta^{\prime \prime},->0$, accounting for the damping of the PML medium.

We begin with transforming the equation $\beta_{n}^{ \pm^{2}}=k^{ \pm 2}-\alpha_{n}^{ \pm^{2}}$, so that the function with integer argument $n \mapsto \beta_{n}^{-}$becomes a function with real argument continuously interpolated between the imposed integer values. Indeed, the geometric transformations associated to the PML has to be continuous and differentiable in order to compute its Jacobian. To that extent, we choose the parametrization:

$$
\begin{equation*}
\alpha\left(y_{d}\right)=\alpha_{0}+\frac{2 \pi}{d} \frac{y_{d}}{\lambda_{0}}, \tag{5.47}
\end{equation*}
$$

so that the application $\beta^{-}$defined by $\beta^{-}\left(y_{d}\right)^{2}=k_{0}^{2} \varepsilon^{-}-\alpha\left(y_{d}\right)^{2}$ is continuous. Thus, the propagation constant of the $n^{\text {th }}$ transmitted order is given by $\beta_{n}^{-}=\beta^{-}\left(n \lambda_{0}\right)$. The key idea is to combine the complex stretching with a real non uniform contraction (given by the continuous function $y\left(y_{d}\right)$, Eq. (5.49)). This contraction is chosen in such a way that for each order $n$ there is a depth $y_{d}^{n}$ such that, around this depth, the apparent wavelength corresponding to the order in play is contracted to a value close to $\lambda_{0}$. At that point of the PML, this order is perfectly absorbed thanks to the complex stretch. We thus eliminate first the orders with quasi normal diffracted angles at lowest depths up to grazing orders (near Wood's anomalies) which are absorbed at greater depths. In mathematical words, the translation of previous considerations on the real contraction can be expressed as:

$$
\begin{equation*}
\exp \left[-i \beta^{-}\left(y_{d}\right) y\left(y_{d}\right)\right]=\exp \left(-i k_{0} y_{d}\right) \tag{5.48}
\end{equation*}
$$

The contraction $y\left(y_{d}\right)$ is thus given by:

$$
\begin{equation*}
y\left(y_{d}\right)=\frac{k_{0} y_{d}}{\beta^{-}\left(y_{d}\right)}=\frac{y_{d}}{\sqrt{\varepsilon^{-}-\left(\sin \theta_{0}+y_{d} / d\right)^{2}}} \tag{5.49}
\end{equation*}
$$

The function $y\left(y_{d}\right)$ has two poles, denoted $y_{d, \pm}^{\star}=d\left( \pm \sqrt{\varepsilon^{-}}-\sin \theta_{0}\right)$. When $y_{d, \pm}^{\star}= \pm n \lambda_{0}$ with $n \in \mathbb{N}^{\star}, \beta^{-}\left(y_{d, \pm}^{\star}\right)=\beta^{-}\left( \pm n \lambda_{0}\right)=\beta_{ \pm}^{-}=0$, i.e. we are on a Wood's anomaly associated with the appearance/disappearance of the $\pm n^{\text {th }}$ transmitted order. We now search for the nearest point to $y_{d, \pm}^{*}$ associated with a Wood's anomaly, denoting:

$$
\begin{cases}n_{+}^{\star} / & D_{+}=\min _{n_{+}^{\star} \in \mathbb{N}^{\star}}\left|y_{d,+}^{\star}-n_{+}^{\star} \lambda_{0}\right| \\ n_{-}^{\star} / & D_{-}=\min _{n_{-}^{\star} \in \mathbb{N}^{\star}}\left|y_{d,-}^{\star}+n_{-}^{\star} \lambda_{0}\right|\end{cases}
$$

In a second step, we look for the point $y_{d}^{0}=n^{\star} \lambda_{0}$ such that:

$$
\begin{equation*}
n^{\star} / \quad D=\min _{n^{\star} \in\left\{n_{+}^{\star}, n_{-}^{\star}\right\}}\left(D^{+}, D^{-}\right) . \tag{5.50}
\end{equation*}
$$

To avoid the singular behaviour at $y_{d}=y_{d, \pm}^{\star}$, we continue the graph of the function $y_{d}(y)$ by a straight line tangent at $y_{d}^{0}$, which equation is $t_{0}\left(y_{d}\right)=s\left(y_{d}^{0}\right)\left(y_{d}-y_{d}^{0}\right)+y\left(y_{d}^{0}\right)$, where $s\left(y_{d}\right)=$ $\frac{\partial y}{\partial y_{d}}\left(y_{d}\right)$ is the so-called stretching coefficient. The final change of coordinate is then given by :

$$
\tilde{y}\left(y_{d}\right)=\left\{\begin{array}{l}
y\left(y_{d}\right) \text { for } y_{d} \leq y_{d}^{0}  \tag{5.51}\\
t_{0}\left(y_{d}\right) \text { for } y_{d}>y_{d}^{0}
\end{array}\right.
$$

Figure 5.10 shows an example of this coordinate mapping. Eventually, the complex stretch $s_{y}$ used in Eq. (5.29) is given by:

$$
\begin{equation*}
s_{y}\left(y_{d}\right)=\zeta^{-} \frac{\partial \tilde{y}}{\partial y_{d}}\left(y_{d}\right) \tag{5.52}
\end{equation*}
$$

Equipped with this mathematical formulation, we can tailor a layer that is doubly perfectly matched:

- to a given medium, which is the aim of the PML technique, through Eq. (5.27),
- to all diffraction orders, through the stretching coefficient $s_{y}$, which depends on the characteristics of the incident wave and on opto-geometric parameters of the grating.


Fig. 5.10: Example of a coordinate mapping $\tilde{y}\left(y_{d}\right)$ used for the APML (black solid line). The graph of $y_{d}(y)$ (blue solid line) is continued by a straight line $t_{0}\left(y_{d}\right)$ tangent at $y_{d}^{0}$ (red dashed line) to avoid the singular behaviour at $y_{d}=y_{d}^{\star}$.

### 5.2.4.4 Numerical example

We now apply the method described in the preceding parts to design an adapted bottom PML for the same example as in part 5.2.4.2. The parameters are the same, and we choose the wavelength of the incident plane wave close to the Wood's anomaly related to the +1 transmitted order $\left(\lambda_{0}=0.999 y_{d,+}^{\star}\right)$ Moreover, we set the length of the PML $\hat{h}^{-}=1.1 y_{d,+}^{\star}$ and choose absorption coefficients $\zeta^{+}=\zeta^{-}=1+i$. For both cases (PML and APML), parameters are alike, the only difference being the complex stretch $s_{y}$.

The field maps of the norm of $H_{z}, E_{x}$ and $E_{y}$ are plotted in logarithmic scale on Fig. 5.11, for the case of a classical PML and our APML. We can observe that the field $H_{z}$ that is effectively computed is clearly damped in the bottom APML (leftmost on Fig. 5.11(b)) whereas it is not in the standard case (leftmost on Fig. 5.11(a)), causing spurious reflections on the outer boundary. The fields $E_{x}$ and $E_{y}$ are deduced from $H_{z}$ thanks to Maxwell's equations. The high values of $E_{y}$ at the tip of the APML (rightmost on Fig. 5.11(b)) are due to very high values of the optical equivalent properties of the APML medium (due to high values of $s_{y}$ ), which does not affect the accuracy of the computed field within the domain of interest.
Another feature of our approach is that it efficiently absorbs the grazing diffraction order, as illustrated on Fig. 5.12: the +1 transmitted order does not decrease in the standard PML (blue solid line), and reaches a high value at $y=-\hat{h}^{-}$, whereas the same order tends to zero as $y \rightarrow-\hat{h}^{-}$in the case of the adapted PML (blue dashed line).
To further validate the accuracy of the method, we compare the diffraction efficiencies computed by our FEM formulation with PML and APML to those obtained by another method. We choose the Rigorous Coupled Wave Analysis (RCWA), also known as the Fourier Modal Method (FMM, [27]). For the chosen parameters, only the $0^{\text {th }}$ order is propagative in reflexion and the orders $-1,0$ and +1 are non evanescent in transmission. We can also check the energy balance $B=R_{0}+T_{-1}+T_{0}+T_{+1}$ since there is no lossy medium in our example. Results are reported in Table 5.3, and show a good agreement of the FEM with APML with the results from RCWA. On the contrary, if classical PML are used, the diffraction efficiencies are less accurate


Fig. 5.11: Field maps of the logarithm of the norm of $H_{z}, E_{x}$ and $E_{y}$ for the dielectric slit grating at $\lambda_{0}=0.999 y_{d,+}^{\star}$ (same parameters as in part 5.2.4.2). (a): classical PML with inefficient damping of $H_{z}$ in the bottom PML. (b): APML where the $H_{z}$ field is correctly damped in the bottom PML. For both cases the thickness of the PML is $\hat{h}^{-}=1.1 y_{d,+}^{\star}$.


Fig. 5.12: Modulus of the $u_{n}$ for the three propagating orders with adapted (dashed lines) and classical PMLs (solid lines). Note that the classical PMLs are efficient for all orders except for the grazing one $(n=1)$ as expected. This drawback is bypassed when using the adaptative PML.
compared to those computed with RCWA. Checking the energy balance leads the same conclusions: the numerical result is perturbed by the reflection of the waves at the end of the PML if it is not adapted to the situation of nearly grazing diffracted orders.

|  | $R_{0}$ | $T_{-1}$ | $T_{0}$ | $T_{+1}$ | $B$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
| RCWA | 0.1570 | 0.3966 | 0.1783 | 0.2680 | 0.9999 |
| FEM + APML | 0.1561 | 0.3959 | 0.1776 | 0.2703 | 0.9999 |
| FEM + PML | 0.1904 | 0.4118 | 0.1927 | 0.2481 | 1.0430 |

Tab. 5.3: Diffraction efficiencies $R_{0}, T_{-1}, T_{0}$ and $T_{+1}$ of the four propagating orders, and energy balance $B=R_{0}+T_{-1}+T_{0}+T_{+1}$, computed by three methods: RCWA (line 1), FEM formulation with APML (line 2), FEM formulation with classical PML (line 3).


Fig. 5.13: Mean value of the norm of $H_{z}$ along the outer boundary of the bottom PML $\gamma=$ $\langle | H_{z}\left(-\hat{h}^{-}\right)| \rangle_{x}$, for $\lambda_{0}$ approaching the Wood's anomaly $y_{d,+}^{\star}$ by inferior values $\left(\lambda_{0}=\left(1-10^{-n}\right) y_{d,+}^{\star}\right.$, red squares) and by superior value $\left(\lambda_{0}=\left(1+10^{-n}\right) y_{d,+}^{\star}\right.$, blue circles) as a function of $n$.

Eventually, to illustrate the behavior of the adaptative PML when the incident wavelength gets closer to a given Wood's anomaly, we computed the mean value of the norm of $H_{z}$ along the outer boundary of the bottom PML $\gamma=\langle | H_{z}\left(-\hat{h}^{-}\right)| \rangle_{x}$, when $\lambda_{0}=\left(1+10^{-n}\right) y_{d,+}^{\star}$ and $\lambda_{0}=\left(1-10^{-n}\right) y_{d,+}^{\star}$, for $n=1,2, \ldots 10$. The results are shown in Fig. 5.13. As the wavelength gets closer to $y_{d,+}^{\star}, \gamma$ first increases but for $n>3$, it decreases exponentially. However, in all cases, the value of $\gamma$ remains small enough to ensure the efficiency of the PMLs.

### 5.2.5 Concluding remarks

A novel FEM formulation was adapted to the analysis of $z$-anisotropic gratings relying on a rigorous treatment of the plane wave sources problem through an equivalent radiation problem with localized sources. The developed approach presents the advantage of being very general in the sense that it is applicable to every conceivable grating geometry.

Numerical experiments based on existing materials at normal and oblique incidences in both TE and TM cases showed the efficiency and the accuracy of our method. We demonstrated we could generate strongly imbalanced symmetric propagative orders in the TE polarization case and at normal incidence with an aragonite grating on a silica substratum.

We also introduced the adaptative PML for grazing incidences configurations. It based on a complex-valued coordinate stretching that deals with grazing diffracted orders, yielding an efficient absorption of the field inside the PML. We provided an example in the TM polarization case (but similar results hold for the TE case), illustrating the efficiency of our method. The value of the magnetic field on the outward boundary of the PML remains small enough to consider there is no spurious reflection. The formulation is used with the FEM but can be applied to others numerical methods. Moreover, the generalization to the vectorial three-dimensional case is straightforward: the recipes given in this last section do work irrespective of the dimension and whether the problem is vectorial.

In the next section, the scalar formulation adapted to mono-dimensional gratings is extended to the the most general case of bi-dimensional grating embedded in an arbitrary multilayered dielectric stack with arbitrary incidence.

### 5.3 Diffraction by arbitrary crossed-gratings : a vector Finite Element formulation

### 5.3.1 Introduction

In this section, we extend the method detailed in Sec. 5.2 to the most general case of vector diffraction by an arbitrary crossed gratings. The main advantage of the Finite Element Method lies in its native ability to handle unstructured meshes, resulting in a build-in accurate discretization of oblique edges. Consequently, our approach remains independent of the shape of the diffractive element, whereas other methods require heavy adjustments depending on whether the geometry of the groove region presents oblique edges (e.g. RCWA [28], FDTD...). In this section, for the sake of clarity, we recall again the rigorous procedure allowing to deal with the issue of the plane wave sources through an equivalence of the diffraction problem with a radiation one whose sources are localized inside the diffractive element itself, as already proposed in Sec. 5.2 [29, 30].

This approach combined with the use of second order edge elements allowed us to retrieve with a good accuracy the few numerical academic examples found in the literature. Furthermore, we provide a new reference case combining major difficulties such as a non trivial toroidal geometry together with strong losses and a high permittivity contrast. Finally, we discuss computation time and convergence as a function of the mesh refinement as well as the choice of the direct solver.

### 5.3.2 Theoretical developments

### 5.3.2.1 Set up of the problem and notations

We denote by $\mathbf{x}, \mathbf{y}$ and $\mathbf{z}$ the unit vectors of the axes of an orthogonal coordinate system $O x y z$. We only deal with time-harmonic fields; consequently, electric and magnetic fields are represented by the complex vector fields $\mathbf{E}$ and $\mathbf{H}$, with a time dependance in $\exp (-i \omega t)$. Note that incident light is now propagating along the $z$-axis, whereas $y$-axis was used in the 2D case.

Besides, in this section, for the sake of simplicity, the materials are assumed to be isotropic and therefore are optically characterized by their relative permittivity $\varepsilon$ and relative permeability $\mu$ (note that the inverse of relative permeabilities are denoted here $v$ ). It is of importance to note that lossy materials can be studied, the relative permittivity and relative permeability being represented by complex valued functions. The crossed-gratings we are dealing with can be split into the following regions as suggested in Fig. 5.14:

- The superstrate $\left(z>z_{0}\right)$ is supposed to be homogeneous, isotropic and lossless, and therefore characterized by its relative permittivity $\varepsilon^{+}$and its relative permeability $\mu^{+}\left(=1 / v^{+}\right)$ and we denote $k^{+}:=k_{0} \sqrt{\varepsilon^{+} \mu^{+}}$, where $k_{0}:=\omega / c$,
- The multilayered stack $\left(z_{N}<z<z_{0}\right)$ is made of $N$ layers which are supposed to be homogeneous and isotropic, and therefore characterized by their relative permittivity $\varepsilon^{n}$, their relative permeability $\mu^{n}\left(=1 / \nu^{n}\right)$ and their thickness $e_{n}$. We denote $k_{n}:=k_{0} \sqrt{\varepsilon^{n} \mu^{n}}$ for $n$ integer between 1 and $N$.
- The groove region ( $z_{g}<z<z_{g-1}$ ), which is embedded in the layer indexed $g\left(\varepsilon^{g}, \mu^{g}\right)$ of the previously described domain, is heterogeneous. Moreover the method does work irrespective of whether the diffractive elements are homogeneous: The permittivity and permeability can vary continuously (gradient index gratings) or discontinuously (step
index gratings). This region is thus characterized by the scalar fields $\varepsilon^{g^{\prime}}(x, y, z)$ and $\mu^{g^{\prime}}(x, y, z)\left(=1 / v^{g^{\prime}}(x, y, z)\right)$. The groove periodicity along the $x$-axis, respectively (resp.) $y$-axis, is denoted $d_{x}$, resp. $d_{y}$, in the sequel.
- The substrate $\left(z<z_{N}\right)$ is supposed to be homogeneous and isotropic and therefore characterized by its relative permittivity $\varepsilon^{-}$and its relative permeability $\mu^{-}\left(=1 / v^{-}\right)$and we denote $k^{-}:=k_{0} \sqrt{\varepsilon^{-} \mu^{-}}$,
Let us emphasize the fact that the method principles remain unchanged in the case of several diffractive patterns made of distinct geometry and/or material.


Fig. 5.14: Scheme and notations of the studied bi-gratings.
The incident field on this structure is denoted:

$$
\begin{equation*}
\mathbf{E}^{\mathrm{inc}}=\mathbf{A}_{0}^{e} \exp \left(i \mathbf{k}_{p}^{+} \cdot \mathbf{r}\right) \tag{5.53}
\end{equation*}
$$

with

$$
\mathbf{k}^{+}=\left[\begin{array}{l}
\alpha_{0}  \tag{5.54}\\
\beta_{0} \\
\gamma_{0}
\end{array}\right]=k^{+}\left[\begin{array}{l}
-\sin \theta_{0} \cos \varphi_{0} \\
-\sin \theta_{0} \sin \varphi_{0} \\
-\cos \theta_{0}
\end{array}\right]
$$

and

$$
\mathbf{A}_{0}^{e}=\left[\begin{array}{l}
E_{x}^{0}  \tag{5.55}\\
E_{y}^{0} \\
E_{z}^{0}
\end{array}\right]=A^{e}\left[\begin{array}{l}
\cos \psi_{0} \cos \theta_{0} \cos \varphi_{0}-\sin \psi_{0} \sin \varphi_{0} \\
\cos \psi_{0} \cos \theta_{0} \sin \varphi_{0}+\sin \psi_{0} \cos \varphi_{0} \\
-\cos \psi_{0} \sin \theta_{0}
\end{array}\right],
$$

where $\varphi_{0} \in[0,2 \pi], \theta_{0} \in[0, \pi / 2]$ and $\psi_{0} \in[0, \pi]$ (polarization angle).
We recall here the diffraction problem: finding the solution of Maxwell equations in harmonic regime i.e. the unique solution $(\mathbf{E}, \mathbf{H})$ of:

$$
\left\{\begin{array}{l}
\operatorname{curl} \mathbf{E}=i \omega \mu_{0} \mu \mathbf{H}  \tag{5.56a}\\
\operatorname{curl} \mathbf{H}=-i \omega \varepsilon_{0} \varepsilon \mathbf{E}
\end{array}\right.
$$

such that the diffracted field satisfies the so-called Outgoing Waves Condition (OWC [31] ) and where $\mathbf{E}$ and $\mathbf{H}$ are quasi-bi-periodic functions with respect to $x$ and $y$ coordinates.

One can choose to calculate arbitrarily $\mathbf{E}$, since $\mathbf{H}$ can be deduced from Eq. (5.56a). The diffraction problem amounts to looking for the unique solution $\mathbf{E}$ of the so-called vectorial Helmholtz propagation equation, deduced from Eqs. (5.56a,5.56b):

$$
\begin{equation*}
\mathscr{M}_{\varepsilon, v}:=-\operatorname{curl}(v \operatorname{curl} \mathbf{E})+k_{0}^{2} \varepsilon \mathbf{E}=\mathbf{0} \tag{5.57}
\end{equation*}
$$

such that the diffracted field satisfies an OWC and where $\mathbf{E}$ is a quasi-bi-periodic function with respect to $x$ and $y$ coordinates.

### 5.3.2.2 From a diffraction problem to a radiative one with localized sources

According to Fig. 5.14, the scalar relative permittivity $\varepsilon$ and inverse permeability $v$ fields associated to the studied diffractive structure can be written using complex-valued functions defined by part and taking into account the notations adopted in Sec. 5.3.2.1:

$$
v(x, y, z):=\left\{\begin{array}{llccc}
v^{+} & \text {for } & z>z_{0} &  \tag{5.58}\\
v^{n} & \text { for } & z_{n-1}>z>z_{n} & \text { with } & 1 \leq n<g \\
v^{g^{\prime}}(x, y, z) & \text { for } & z_{g-1}>z>z_{g} & & \\
v^{n} & \text { for } & z_{n-1}>z>z_{n} & \text { with } & g<n \leq N \\
v^{-} & \text {for } & z<z_{N}
\end{array}\right.
$$

with $v=\{\varepsilon, v\} \quad, \quad z_{0}=0 \quad$ and $\quad z_{n}=-\sum_{l=1}^{n} e_{l} \quad$ for $\quad 1 \leq n \leq N$.
It is now convenient to introduce two functions defined by part $\varepsilon_{1}$ and $v_{1}$ corresponding to the associated multilayered case (i.e. the same stack without any diffractive element) constant over $O x$ and $O y$ :

$$
v_{1}(x, y, z):=\left\{\begin{array}{llc}
v^{+} & \text {for } & z>0  \tag{5.59}\\
v^{n} & \text { for } & z_{n-1}>z>z_{n} \\
v^{-} & \text {for } & z<z_{N}
\end{array} \text { with } \quad 1 \leq n \leq N\right.
$$

with $v=\{\varepsilon, v\}$.
We denote by $\mathbf{E}_{0}$ the restriction of $\mathbf{E}^{\text {inc }}$ to the superstrate region:

$$
\mathbf{E}_{0}:=\left\{\begin{array}{ccc}
\mathbf{E}^{\mathrm{inc}} & \text { for } & z>z_{0}  \tag{5.60}\\
\mathbf{0} & \text { for } & z \leq z_{0}
\end{array}\right.
$$

We are now in a position to define more explicitly the vector diffraction problem that we are dealing with in this section. It amounts to looking for the unique vector field $\mathbf{E}$ solution of:

$$
\begin{equation*}
\mathscr{M}_{\varepsilon, v}(\mathbf{E})=\mathbf{0} \quad \text { such that } \mathbf{E}^{d}:=\mathbf{E}-\mathbf{E}_{0} \text { satisfies an OWC. } \tag{5.61}
\end{equation*}
$$

In order to reduce this diffraction problem to a radiation one, an intermediary vector field denoted $\mathbf{E}_{1}$ is necessary and is defined as the unique solution of:

$$
\begin{equation*}
\mathscr{M}_{\varepsilon_{1}, v_{1}}\left(\mathbf{E}_{1}\right)=\mathbf{0} \quad \text { such that } \mathbf{E}_{1}^{d}:=\mathbf{E}_{1}-\mathbf{E}_{0} \text { satisfies an OWC. } \tag{5.62}
\end{equation*}
$$

The vector field $\mathbf{E}_{1}$ corresponds to an ancillary problem associated to the general vectorial case of a multilayered stack which can be calculated independently. This general calculation is seldom treated in the literature, we present a development in Appendix. Thus $\mathbf{E}_{1}$ is from now on considered as a known vector field. It is now apropos to introduce the unknown vector field $\mathbf{E}_{2}^{d}$, simply defined as the difference between $\mathbf{E}$ and $\mathbf{E}_{1}$, which can finally be calculated thanks to the FEM and:

$$
\begin{equation*}
\mathbf{E}_{2}^{d}:=\mathbf{E}-\mathbf{E}_{1}=\mathbf{E}^{d}-\mathbf{E}_{1}^{d} \tag{5.63}
\end{equation*}
$$

It is of importance to note that the presence of the superscript $d$ is not fortuitous: As a difference between two diffracted fields (Eq. (5.63), $\mathbf{E}_{2}^{d}$ satisfies an OWC which is of prime importance in our formulation. By taking into account these new definitions, Eq. (5.61) can be written:

$$
\begin{equation*}
\mathscr{M}_{\varepsilon, v}\left(\mathbf{E}_{2}^{d}\right)=-\mathscr{M}_{\varepsilon, v}\left(\mathbf{E}_{1}\right), \tag{5.64}
\end{equation*}
$$

where the right-hand member is a vector field which can be interpreted as a known vectorial source term $-\mathscr{S}_{1}(x, y, z)$ whose support is localized inside the diffractive element itself. To prove it, let us introduce the null term defined in Eq. (5.62) and make the use of the linearity of $\mathscr{M}$, which leads to:

$$
\begin{equation*}
\mathscr{S}_{1}:=\mathscr{M}_{\varepsilon, v}\left(\mathbf{E}_{1}\right)=\mathscr{M}_{\varepsilon, v}\left(\mathbf{E}_{1}\right)-\underbrace{\mathscr{M}_{\varepsilon_{1}, v_{1}}\left(\mathbf{E}_{1}\right)}_{=\mathbf{0}}=\mathscr{M}_{\varepsilon-\varepsilon_{1}, v-v_{1}}\left(\mathbf{E}_{1}\right) . \tag{5.65}
\end{equation*}
$$

### 5.3.2.3 Quasi-periodicity and weak formulation

The weak form is obtained by multiplying scalarly Eq. (5.61) by weighted vectors $\mathbf{E}^{\prime}$ chosen among the ensemble of quasi-bi-periodic vector fields of $L^{2}(\mathbf{c u r l})\left(\right.$ denoted $L^{2}\left(\mathbf{c u r l},\left(d_{x}, d_{y}\right), \mathbf{k}\right)$ ) in $\Omega$ :

$$
\begin{equation*}
\mathscr{R}_{\varepsilon, v}\left(\mathbf{E}, \mathbf{E}^{\prime}\right)=\int_{\Omega}-\operatorname{curl}(v \operatorname{curl} \mathbf{E}) \cdot \overline{\mathbf{E}^{\prime}}+k_{0}^{2} \varepsilon \mathbf{E} \cdot \overline{\mathbf{E}^{\prime}} \mathrm{d} \Omega \tag{5.66}
\end{equation*}
$$

Integrating by part Eq. (5.66) and making the use of the Green-Ostrogradsky theorem lead to:

$$
\begin{equation*}
\mathscr{R}_{\varepsilon, v}\left(\mathbf{E}, \mathbf{E}^{\prime}\right)=\int_{\Omega}-v \operatorname{curl} \mathbf{E} \cdot \operatorname{curl} \overline{\mathbf{E}^{\prime}}+k_{0}^{2} \varepsilon \mathbf{E} \cdot \overline{\mathbf{E}^{\prime}} \mathrm{d} \Omega-\int_{\partial \Omega}(\mathbf{n} \times(v \operatorname{curl} \mathbf{E})) \cdot \overline{\mathbf{E}^{\prime}} \mathrm{d} S \tag{5.67}
\end{equation*}
$$

where $\mathbf{n}$ refers to the exterior unit vector normal to the surface $\partial \Omega$ enclosing $\Omega$.
The first term of this sum concerns the volume behavior of the unknown vector field whereas the right-hand term can be used to set boundary conditions (Dirichlet, Neumann or so called quasi-periodic Bloch-Floquet conditions).

The solution $\mathbf{E}_{2}^{d}$ of the weak form associated to the diffraction problem, expressed in its previously defined equivalent radiative form at Eq. (5.64), is the element of $L^{2}\left(\mathbf{c u r l},\left(d_{x}, d_{y}\right), \mathbf{k}\right)$ such that:

$$
\begin{equation*}
\forall \mathbf{E}^{\prime} \in L^{2}\left(\mathbf{c u r l}, d_{x}, d_{y}, \mathbf{k}\right), \mathscr{R}_{\varepsilon, v}\left(\mathbf{E}_{2}^{d}, \mathbf{E}^{\prime}\right)=-\mathscr{R}_{\varepsilon-\varepsilon_{1}, v-v_{1}}\left(\mathbf{E}_{1}, \mathbf{E}^{\prime}\right) \tag{5.68}
\end{equation*}
$$

In order to rigorously truncate the computation a set of Bloch boundary conditions are imposed on the pair of planes defined by $\left(y=-d_{y} / 2, y=d_{y} / 2\right)$ and $\left(x=-d_{x} / 2, x=d_{x} / 2\right)$. One can refer to [11] for a detailed implementation of Bloch conditions adapted to the FEM. A set of Perfectly Matched Layers are used in order to truncate the substrate and the superstrate along $z$ axis (see [32] for practical implementation of PML adapted to the FEM). Since the proposed unknown $\mathbf{E}_{2}^{d}$ is quasi-bi-periodic and satisfies an OWC, this set of boundary conditions is perfectly reasonable: $\mathbf{E}_{2}^{d}$ is radiated from the diffractive element towards the infinite regions of the problem and decays exponentially inside the PMLs along $z$ axis. The total field associated to the diffraction problem $\mathbf{E}$ is deduced at once from Eq. (5.63).

### 5.3.2.4 Edge or Whitney 1-form second order elements

In the vectorial case, edge elements (or Whitney forms) make a much more relevant choice [33] than nodal elements. Note that a lot of work (see for instance [34]) has been done on higher order edge elements since their introduction by Bossavit [35]. These elements are suitable to the representation of vector fields such as $\mathbf{E}_{2}^{d}$, by letting their normal component be discontinuous and imposing the continuity of their tangential components. Instead of linking the Degrees Of Freedom (DOF) of the final algebraic system to the nodes of the mesh, the DOF associated to edges (resp. faces) elements are the circulations (resp. flux) of the unknown vector field along (resp. across) its edges (resp. faces).

Let us consider the computation cell $\Omega$ together with its exterior boundary $\partial \Omega$. This volume is sampled in a finite number of tetrahedron according to the following rules: Two distinct tetrahedrons have to either share a node, an edge or a face or have no contact. Let us denote by $\mathscr{T}$ the set of tetrahedrons, $\mathscr{F}$ the set of faces, $\mathscr{E}$ the set of edges and $\mathscr{N}$ the set of nodes. In the sequel, one will refers to the node $n=\{i\}$, the edge $e=\{i, j\}$, the face $f=\{i, j, k\}$ and the tetrahedron $t=\{i, j, k, l\}$.


Fig. 5.15: Degrees of freedom of a second order tetrahedral element.

Twelve DOF (two for each of the six edges of a tetrahedron) are classically derived from line integral of weighted projection of the field $\mathbf{E}_{2}^{d}$ on each oriented edge $e=\{i, j\}$ :

$$
\left\{\begin{array}{l}
\vartheta_{i j}=\int_{i}^{j} \mathbf{E}_{2}^{d} \cdot \mathbf{t}_{i j} \lambda_{i} \mathrm{~d} l  \tag{5.69}\\
\vartheta_{j i}=\int_{j}^{i} \mathbf{E}_{2}^{d} \cdot \mathbf{t}_{j i} \lambda_{j} \mathrm{~d} l
\end{array},\right.
$$

where $\mathbf{t}_{i j}$ is the unit vector and $\lambda_{i}$, the barycentric coordinate of node $i$, is the chosen weight function.

According to Yioultsis et al. [36], a judicious choice for the remaining DOF is to make the use of a tangential projection of the 1-form $\mathbf{E}_{2}^{d}$ on the face $f=\{i, j, k\}$.

$$
\left\{\begin{array}{l}
\vartheta_{i j k}=\int_{f}\left(\mathbf{E}_{2}^{d} \times \mathbf{n}_{i j k}^{+}\right) \cdot \operatorname{grad} \lambda_{j} \mathrm{~d} s  \tag{5.70}\\
\vartheta_{i k j}=\int_{f}\left(\mathbf{E}_{2}^{d} \times \mathbf{n}_{i j k}^{-}\right) \cdot \operatorname{grad} \lambda_{k} \mathrm{~d} s
\end{array} .\right.
$$

The expressions for the shape functions, or basis vectors, of the second order 1-form Whitney element are given by:

$$
\left\{\begin{array}{c}
\mathbf{w}_{i j}=\left(8 \lambda_{i}^{2}-4 \lambda_{i}\right) \operatorname{grad} \lambda_{j}+\left(-8 \lambda_{i} \lambda_{j}+2 \lambda_{j}\right) \operatorname{grad} \lambda_{i}  \tag{5.71}\\
\mathbf{w}_{i j k}=16 \lambda_{i} \lambda_{j} \operatorname{grad} \lambda_{k}-8 \lambda_{j} \lambda_{k} \operatorname{grad} \lambda_{i}-8 \lambda_{k} \lambda_{i} \operatorname{grad} \lambda_{j}
\end{array} .\right.
$$

This choice of shape function ensures [37] the following fundamental property: every degree of freedom associated with a shape function should be zero for any other shape function. Finally, an approximation of the unknown $\mathbf{E}_{2}^{d}$ projected on the shape functions of the mesh $m\left(\mathbf{E}_{2}^{d, m}\right)$ can be derived:

$$
\begin{equation*}
\mathbf{E}_{2}^{d, m}=\sum_{e \in \mathscr{E}} \vartheta_{e} \mathbf{w}_{e}+\sum_{f \in \mathscr{F}} \vartheta_{f} \mathbf{w}_{f} \tag{5.72}
\end{equation*}
$$

Weight functions $\mathbf{E}^{\prime}$ (c.f. Eq. (5.68) are chosen in the same space than the unknown $\mathbf{E}_{2}^{d}$, $L^{2}\left(\mathbf{c u r l},\left(d_{x}, d_{y}\right), \mathbf{k}\right)$. According to the Galerkin formulation, this choice is made so that their restriction to one bi-period belongs to the set of shape functions mentioned above. Inserting the decomposition of $\mathbf{E}_{2}^{d}$ of Eq. (5.72) in Eq. (5.68) leads to the final algebraic system which is solved, in the following numerical examples, thanks to direct solvers.

### 5.3.3 Energetic considerations: Diffraction efficiencies and losses

Contrarily to modal methods based on the determination of Rayleigh coefficients, the rough results of the FEM are three complex components of the vector field $\mathbf{E}^{d}$ interpolated over the mesh of the computation cell. Diffraction efficiencies are deduced from this field maps as follows.

As a difference between two quasi-periodic vector fields (see Eq. (5.61)), $\mathbf{E}^{d}$ is quasi-biperiodic and its components can be expanded as a double Rayleigh sum:

$$
\begin{equation*}
E_{x}^{d}(x, y, z)=\sum_{(n, m) \in \mathbb{Z}^{2}} u_{n, m}^{d, x}(z) e^{i\left(\alpha_{n} x+\beta_{m} y\right)} \tag{5.73}
\end{equation*}
$$

with $\alpha_{n}=\alpha_{0}+\frac{2 \pi}{d_{x}} n, \beta_{m}=\beta_{0}+\frac{2 \pi}{d_{y}} m$ and

$$
\begin{equation*}
u_{n, m}^{d, x}(z)=\frac{1}{d_{x} d_{y}} \int_{-d_{x} / 2}^{d_{x} / 2} \int_{-d_{y} / 2}^{d_{y} / 2} E_{x}^{d}(x, y, z) e^{-i\left(\alpha_{n} x+\beta_{m} y\right)} \mathrm{d} x \mathrm{~d} y \tag{5.74}
\end{equation*}
$$

By inserting the decomposition of Eq. (5.73), which is satisfied by $E_{x}^{d}$ everywhere but in the groove region, into the Helmholtz propagation equation, one can express Rayleigh coefficients in the substrate and the superstrate as follows:

$$
\begin{equation*}
u_{n, m}^{d, x}(z)=e_{n, m}^{x, p} e^{-i \gamma_{n, m}^{+} z}+e_{n, m}^{x, c} e^{i \gamma_{n, m}^{+} z} \tag{5.75}
\end{equation*}
$$

with $\gamma_{n, m}^{ \pm^{2}}=k^{ \pm^{2}}-\alpha_{n}^{2}-\beta_{m}^{2}$, where $\gamma_{n, m}\left(\right.$ or $\left.-i \gamma_{n, m}\right)$ is positive. The quantity $u_{n, m}^{d, x}$ is the sum of a propagative plane wave (which propagates towards decreasing values of $z$, superscript $p$ ) and of a counterpropagative one (superscript $c$ ). The OWC verified by $\mathbf{E}^{d}$ imposes:

$$
\forall(n, m) \in \mathbb{Z}^{2}\left\{\begin{array}{lll}
e_{n, m}^{x, p}=0 & \text { for } & z>z_{0}  \tag{5.76}\\
e_{n, m}^{x, c}=0 & \text { for } & z<z_{N}
\end{array}\right.
$$

Eq. (5.74) allows to evaluate numerically $e_{n, m}^{x, c}$ (resp. $e_{n, m}^{x, p}$ ) by double trapezoidal integration of a slice of the complex component $E_{x}^{d}$ at an altitude $z_{c}$ fixed in the superstrate (resp. substrate).

It is well known that the mere trapezoidal integration method is very efficient for smooth and periodic functions (integration on one period). The same holds for $E_{y}^{d}$ and $E_{z}^{d}$ components as well as their coefficients $e_{n, m}^{y,\{c, p\}}$ and $e_{n, m}^{z,\{c, p\}}$.

The dimensionless expression of the efficiency of each reflected and transmitted $(n, m)$ order [38] is deduced from Eqs. $(5.75,5.76)$ :

$$
\left\{\begin{array}{c}
R_{n, m}=\frac{1}{\left|A_{e}\right|^{2}} \frac{\gamma_{n, m}^{+}}{\gamma_{0}} \mathbf{e}_{n, m}^{c}\left(z_{c}\right) \cdot \overline{\mathbf{e}_{n, m}^{c}\left(z_{c}\right)} \quad \text { for } \quad z_{c}>z_{0}  \tag{5.77}\\
T_{n, m}=\frac{1}{A_{e}^{2}} \frac{\gamma_{n, m}^{-}}{\gamma_{0}} \mathbf{e}_{n, m}^{p}\left(z_{c}\right) \cdot \overline{\mathbf{e}_{n, m}^{p}\left(z_{c}\right)} \quad \text { for } \quad z_{c}<z_{N}
\end{array}\right.
$$

with $\mathbf{e}_{n, m}^{\{c, p\}}=e_{n, m}^{x,\{c, p\}} \mathbf{x}+e_{n, m}^{y,\{c, p\}} \mathbf{y}+e_{n, m}^{z,\{c, p\}} \mathbf{z}$.
Furthermore, normalized losses Q can be obtained through the computation of the following ratio:

$$
\begin{equation*}
Q=\frac{\int_{V} \frac{1}{2} \omega \varepsilon_{0} \Im m\left(\varepsilon^{g^{\prime}}\right) \mathbf{E} \cdot \overline{\mathbf{E}} \mathrm{d} V}{\int_{S} \frac{1}{2} \Re e\left\{\mathbf{E}_{0} \times \overline{\mathbf{H}_{0}}\right\} \cdot \mathbf{n} \mathrm{d} S} . \tag{5.78}
\end{equation*}
$$

The numerator in Eq. (5.78) clarifies losses in watts by bi-period of the considered crossedgrating and are computed by integrating the Joule effect losses density over the volume $V$ of the lossy element. The denominator normalizes these losses to the incident power, i.e. the timeaveraged incident Poynting vector flux across one bi-period (a rectangular surface $S$ of area $d_{x} d_{y}$ in the superstrate parallel to $O x y$, whose normal oriented along decreasing values of $z$ is denoted $\mathbf{n}$ ). Since $\mathbf{E}_{0}$ is nothing but the plane wave defined at Eqs. (5.54,5.55), this last term is equal to $\left(A_{e}^{2} \sqrt{\varepsilon_{0} / \mu_{0}} d_{x} d_{y}\right) /\left(2 \cos \left(\theta_{0}\right)\right)$. Volumes and normal to surfaces being explicitly defined, normalized losses losses $Q$ are quickly computed once $\mathbf{E}$ determined and interpolated between mesh nodes.

Finally, the accuracy and self-consistency of the whole calculation can be evaluated by summing the real part of transmitted and reflected efficiencies $(n, m)$ to normalized losses:

$$
Q+\sum_{(n, m) \in \mathbb{Z}^{2}} \mathfrak{R} e\left\{R_{n, m}\right\}+\sum_{(n, m) \in \mathbb{Z}^{2}} \mathfrak{R e} e\left\{T_{n, m}\right\},
$$

quantity to be compared to 1 . The sole diffraction orders taken into account in this conservation criterium correspond to propagative orders whose efficiencies have a non-null real part. Indeed, diffraction efficiencies of evanescent orders, corresponding to pure imaginary values of $\gamma_{n, m}^{ \pm}$for higher values of $(n, m)$ (see Eq. (5.75)) are also pure imaginary values as it appears clearly in Eq. (5.77). Numerical illustrations of such global energy balances are presented in the next section.

### 5.3.4 Accuracy and convergence

### 5.3.4.1 Classical crossed gratings

There are only a few references in the literature containing numerical examples. For each of them, the problem only consists of three regions (superstrate, grooves and substrate) as summed up on Figure 5.16. For the four selected cases, among six found in the literature, published


Fig. 5.16: Configuration of the studied cases.
results are compared to ones given by our formulation of the FEM. Moreover, in each case, a satisfying global energy balance is detailed. Finally a new validation case combining all the difficulties encountered when modeling crossed-gratings is proposed: A non-trivial geometry for the diffractive pattern (a torus), made of an arbitrary lossy material leading to a large step of index and illuminated by a plane wave with an oblique incidence. Convergence of the FEM calculation as well as computation time will be discussed in Sec. 5.3.4.2.

Checkerboard grating In this example worked out by $\mathrm{L} . \mathrm{Li}$ [27], the diffractive element is a rectangular parallelepiped as shown Fig. 5.17a and the grating parameter highlighted in Fig. 5.16 are the following: $\varphi_{0}=\theta_{0}=0^{\circ}, \psi_{0}=45^{\circ}, d_{x}=d_{y}=5 \lambda_{0} \sqrt{2} / 4, h=\lambda_{0}, \varepsilon^{+}=\varepsilon^{g^{\prime}}=$ 2.25 and $\varepsilon^{-}=\varepsilon^{g}=1$.


Fig. 5.17: Diffractive element with vertical edges (a). $\mathfrak{R e}\left\{E_{x}\right\}$ in $V / m(b)$.

|  | FMM [27] | FEM |
| :--- | :---: | :---: |
| $T_{-1,-1}$ | 0.04308 | 0.04333 |
| $T_{-1,0}$ | 0.12860 | 0.12845 |
| $T_{-1,+1}$ | 0.06196 | 0.06176 |
| $T_{0,-1}$ | 0.12860 | 0.12838 |
| $T_{0,0}$ | 0.17486 | 0.17577 |
| $T_{0,+1}$ | 0.12860 | 0.12839 |
| $T_{+1,-1}$ | 0.06196 | 0.06177 |
| $T_{+1,0}$ | 0.12860 | 0.12843 |
| $T_{+1,+1}$ | 0.04308 | 0.04332 |
| $\sum_{(n, m) \in \mathbb{Z}} \Re\left\{R_{n, m}\right\}$ | - | 0.10040 |
| TOTAL | - | 1.00000 |

Tab. 5.4: Energy balance [27].
Our formulation of the FEM shows good agreement with the Fourier Modal Method developed by $\mathrm{L} . \mathrm{Li}([27], 1997)$ since the maximal relative difference between the array of values presented in Table 5.4 remains lower than $10^{-3}$. Moreover, the sum of the efficiencies of propagative orders given by the FEM is very close to 1 in spite of the addition of all errors of determination upon the efficiencies.

Pyramidal crossed-grating In this example firstly worked out by Derrick et al. [39], the diffractive element is a pyramid with rectangular basis as shown Fig. 5.18a and the grating parameters highlighted in Fig. 5.16 are the following: $\lambda_{0}=1.533, \varphi_{0}=45^{\circ}, \theta_{0}=30^{\circ}, \psi_{0}=0^{\circ}$, $d_{x}=1.5, d_{y}=1, h=0.25, \varepsilon^{+}=\varepsilon^{g}=1$ and $\varepsilon^{-}=\varepsilon^{g^{\prime}}=2.25$. Results given by the FEM show


Fig. 5.18: Diffractive element with oblique edges (a). $\mathfrak{R e}\left\{E_{y}\right\}$ in $V / m(b)$.

| Given in | $[39]$ | $[40]$ | $[41]$ | $[42]$ | FEM |
| :--- | :---: | :---: | :---: | :---: | :---: |
| $R_{-1,0}$ | 0.00254 | 0.00207 | 0.00246 | 0.00249 | 0.00251 |
| $R_{0,0}$ | 0.01984 | 0.01928 | 0.01951 | 0.01963 | 0.01938 |
| $T_{-1,-1}$ | 0.00092 | 0.00081 | 0.00086 | 0.00086 | 0.00087 |
| $T_{0,-1}$ | 0.00704 | 0.00767 | 0.00679 | 0.00677 | 0.00692 |
| $T_{-1,0}$ | 0.00303 | 0.00370 | 0.00294 | 0.00294 | 0.00299 |
| $T_{0,0}$ | 0.96219 | 0.96316 | 0.96472 | 0.96448 | 0.96447 |
| $T_{1,0}$ | 0.00299 | 0.00332 | 0.00280 | 0.00282 | 0.00290 |
| TOTAL | 0.99855 | 1.00001 | 1.00008 | 0.99999 | 1.00004 |

Tab. 5.5: Comparison with the results given in [39, 40, 41, 42].
good agreement with ones of the C method [39, 42], the Rayleigh method [40] and the RCWA [41]. Note that, in this case, some edges of the diffractive element are oblique.

Bi-sinusoidal grating In this example worked out by Bruno et al. [43], the surface of the grating is bi-sinusoidal (see Fig. 5.19a) and described by the function $f$ defined by:

$$
\begin{equation*}
f(x, y)=\frac{h}{4}\left[\cos \left(\frac{2 \pi x}{d}\right)+\cos \left(\frac{2 \pi y}{d}\right)\right] \tag{5.79}
\end{equation*}
$$

The grating parameters et al.highlighted in Fig. 5.16 are the following: $\lambda_{0}=0.83, \varphi_{0}=\theta_{0}=$ $\psi_{0}=0^{\circ}, d_{x}=d_{y}=1, h=0.2, \varepsilon^{+}=\varepsilon^{g}=1$ and $\varepsilon^{-}=\varepsilon^{g^{\prime}}=4$. Note that in order to define this


Fig. 5.19: Diffractive element with oblique edges (a). $\mathfrak{R e}\left\{E_{z}\right\}$ in $V / m(b)$.

|  | $[43]$ | FEM |
| :--- | :---: | :---: |
| $R_{-1,0}$ | 0.01044 | 0.01164 |
| $R_{0,-1}$ | 0.01183 | 0.01165 |
| $T_{-1,-1}$ | 0.06175 | 0.06299 |
| $\sum_{(n, m) \in \mathbb{Z}} \Re e\left\{R_{n, m}\right\}$ | - | 0.10685 |
| $\sum_{(n, m) \in \mathbb{Z}} R\left\{T_{n, m}\right\}$ | - | 0.89121 |
| TOTAL | - | 0.99806 |

Tab. 5.6: Energy balance [43].
surface, the bi-sinusoid was first sampled ( $15 \times 15$ points), then converted to a 3D file format. This sampling can account for the slight differences with the results obtained using the method of variation of boundaries developed by Bruno et al. (1993).

Circular apertures in a lossy layer In this example worked out by Schuster et al. [44], the diffractive element is a circular aperture in a lossy layer as shown Fig. 5.20a and the grating parameter highlighted in Fig. 5.16 are the following: $\lambda_{0}=500 \mathrm{~nm}, \varphi_{0}=\theta_{0}=0^{\circ}, \varepsilon^{+}=\varepsilon^{g}=1$, $\varepsilon^{g^{\prime}}=0.8125+5.2500 i$ and $\varepsilon^{-}=2.25$.


Fig. 5.20: Lossy diffractive element with vertical edges (a). $\mathfrak{R e}\left\{E_{y}\right\}$ in $V / m$ (b).

|  | $[45]$ | $[27]$ | $[44]$ | FEM |
| :--- | :---: | :---: | :---: | :---: |
| $R_{0,0}$ | 0.24657 | 0.24339 | 0.24420 | 0.24415 |
| $\sum_{(n, m) \in \mathbb{Z}} \Re e\left\{T_{n, m}\right\}$ | - | - | - | 0.29110 |
| $\sum_{(n, m) \in \mathbb{Z}} \Re e\left\{R_{n, m}\right\}$ | - | - | - | 0.26761 |
| $Q$ | - | - | - | 0.44148 |
| TOTAL | - | - | - | 1.00019 |

Tab. 5.7: Comparison with [45, 27, 44] and energy balance.
In this lossy case, results obtained with the FEM show good agreement with the ones obtained with the FMM [27], the differential method [44, 46] and the RCWA [45]. Joule losses inside the diffractive element can be easily calculated, which allows to provide a global energy balance for this configuration. Finally, the convergence of the value $R_{0,0}$ as a function of the mesh refinement will be examined.

Lossy tori grating We finally propose a new test case for crossed-grating numerical methods. The major difficulty of this case lies both in the non trivial geometry (see Fig. 5.21a) of the diffractive object and in the fact that it is made of a material chosen so that losses are optimal inside it. The grating parameters highlighted in Fig. 5.16 and Fig. 5.21a are the following: $\lambda_{0}=1, \varphi_{0}=\psi_{0}=0^{\circ}, d_{x}=d_{y}=0.3, a=0.1, b=0.05, R=0.15, h=500 \mathrm{~nm}, \varepsilon^{+}=\varepsilon^{g}=1$, $\varepsilon^{g^{\prime}}=-21+20 i$ and $\varepsilon^{-}=2.25$.


Fig. 5.21: Torus parameters (a). Coarse mesh of the computational domain (b).

| FEM 3D | $\theta=0^{\circ}$ | $\theta=40^{\circ}$ |
| :--- | :---: | :---: |
| $R_{0,0}$ | 0.36376 | 0.27331 |
| $T_{0,0}$ | 0.32992 | 0.38191 |
| $Q$ | 0.30639 | 0.34476 |
| TOTAL | 1.00007 | 0.99998 |

Tab. 5.8: Energy balances at normal and oblique incidence.
Tab. 5.8 illustrates the independence of our method towards the geometry of the diffractive element. $\boldsymbol{\varepsilon}^{g^{\prime}}$ is chosen so that the skin depth has the same order of magnitude as $b$, which maximizes losses. Note that energy balances remain very accurate at normal and oblique incidence, in spite of both the non-triviality of the geometry and the strong losses.

### 5.3.4.2 Convergence and computation time

Convergence as a function of mesh refinement When using modal methods such as the RCWA or the differential method, based on the calculation of Rayleigh coefficients, a number proportional to $N_{R}$ have to be be determined a priori. Then, the unknown diffracted field is expanded as a Fourier serie, injected under this form in Maxwell equations, which annihilates $x-$ and $y$-dependencies. This leads to a system of coupled partial differential equations whose coefficients can structured in a matrix formalism. The resulting matrix is sometimes directly invertible (RCWA) depending on whether the geometry allows to suppress the $z$-dependance, which makes this method adapted to diffractive elements with vertically (or decomposed in staircase functions) shaped edge. In some other cases, one has to make the use of integral methods in order to solve the system, as in the pyramidal case for instance, which leads to the socalled differential method. The diffracted field map can be deduced from these coefficients. If the grating configuration only calls for a few propagative orders and if the field inside the groove region is not the main information sought for, these two close methods allow to determine the repartition of the incident energy quickly. However, if the field inside the groove region is the main piece of information, it is advisable to calculate many Rayleigh coefficients corresponding to evanescent waves which increases the computation time as $\left(N_{R}\right)^{3}$ or even $\left(N_{R}\right)^{4}$.

FEM relies on the direct calculation of the vectorial components of the complex field. Rayleigh coefficients are determined a posteriori. The parameter limiting the computation time is the number of tetrahedral elements along which the computational domain is split up. We suppose that it is necessary to calculate at least two or three points (or mesh nodes) per period of the field $\left(\lambda_{0} / \sqrt{\mathfrak{R} e\{\varepsilon\}}\right)$. Figure 5.22 shows the convergence of the efficiency $R_{0,0}$ (circular apertures case, see Fig. 5.20a) as a function of the mesh refinement characterized by the parameter $N_{M}$ : The maximum size of each element is set to $\lambda_{0} /\left(N_{M} \sqrt{\Re e} e\{\varepsilon\}\right)$.


Fig. 5.22: Convergence of $R_{0,0}$ in function of $N_{m}$ (circular apertures crossed-grating).

It is of interest to note that even if $N_{M}<3$ the FEM still gives pertinent diffraction efficiencies: $R_{0,0}=0.2334$ for $N_{M}=1$ and $R_{0,0}=0.2331$ for $N_{M}=2$. The Galerkin method (see Eq. (5.67)) corresponds to a minimization of the error (between the exact solution and the approximation) with respect to a norm that can be physically interpreted in terms of energy-related quantities. Therefore, the finite element methods usually provide energy-related quantities that are more accurate than the local values of the fields themselves.

Computation time All the calculations were performed on a server equipped with 8 dual core Itanium1 processors and 256Go of RAM. Tetrahedral quadratic edge elements were used together with the direct solver PARDISO. Among different direct solvers adapted to sparse matrix algebra (UMFPACK, SPOOLES and PARDISO), PARDISO turned out to be the less time-consuming one as shown in Tab 5.9.

| Solver | Computation time for 41720 DOF | Computation time for 205198 DOF |
| :--- | :---: | :---: |
| SPOOLES | 15 mn 32 s | 14 h 44 mn |
| UMFPACK | 2 mn 07 s | 1 h 12 mn |
| PARDISO | 57 s | 16 mn |

Tab. 5.9: Computation time variations from solver to solver.


Fig. 5.23: Computation time and number of DOF as a function of $N_{M}$.
Figure 5.23 shows the computation time required to perform the whole FEM computational process for a system made of a number of DOF indicated on the right-hand ordinate. It is of importance to note that for values of $N_{M}$ lower than 3 , the problem can be solved in less than a minute on a standard laptop ( 4 Go RAM, $2 \times 2 \mathrm{GHz}$ ) with 3 significant digits on the diffraction efficiencies. This accuracy is more than sufficient in numerous experimental cases. Furthermore, as far as integrated values are at stake, relatively coarse meshes ( $N_{M} \approx 1$ ) can be used trustfully, authorizing fast geometric, spectral or polarization studies.

Nowadays, the efficiency of the numerical algorithms for sparse matrix algebra together with the available power of computers and the fact that the problem reduces to a basic cell with a
size of a small number of wavelengths make the finite element problem very tractable as proved here.

### 5.4 Concluding remarks

In this chapter, we demonstrate a general formulation of the FEM allowing to calculate the diffraction efficiencies from the electromagnetic field diffracted by arbitrarily shaped gratings embedded in a multilayered stack lightened by a plane wave of arbitrary incidence and polarization angle. It relies on a rigorous treatment of the plane wave sources problem through an equivalent radiation problem with localized sources. Bloch conditions and a new dedicated PML have been implemented in order to rigorously truncate the computational domain.

The principles of the method were discussed in detail for mono-dimensional gratings in TE/TM polarization cases (2D or scalar case) in a first part, and for the most general bidimensional or crossed gratings (3D or vector case) in a second part. Note that the very same concepts could be applied to the intermediate case of mono-dimensional gratings enlighten by an arbitrary incident plane wave (so-called conical case). The reader will find detail about the element basis relevant to this case in [11].

The main advantage of this formulation is its complete generality with respect to the studied geometries and the material properties, as illustrated with the lossy tori grating nontrivial case. Its principle remains independent of both the number of diffractive elements by period and number of stack layers. Its flexibility allowed us to retrieve with accuracy the few numerical academic examples found in the literature and established with independent methods.

The remarkable accuracy observed in the case of coarse meshes, makes it a fast tool for the design and optimization of diffractive optical components (e.g. reflection and transmission filters, polarizers, beam shapers, pulse compression gratings. . .). The complete independence of the presented approach towards both the geometry and the isotropic constituent materials of the diffractive elements makes it a handy and powerful tool for the study of metamaterials, finitesize photonic crystals, periodic plasmonic structures...The method described in this chapter has already been successfully applied to various problems, from homogenization theory [47] or transformation optics [48] to more applied concerns as the modeling of complex CMOS nanophotonic devices [49] or ultra-thin new generation solar cells [50].

## 5.A APPENDIX

This appendix is dedicated to the determination of the vector electric field in a dielectric stack enlightened by a plane wave of arbitrary polarization and incidence angle. This calculation, abundantly treated in the 2 D scalar case, is generally not presented in the literature since, as far as isotropic cases are concerned, it is possible to project the general vectorial case on the two reference TE and TM cases. However, the presented formulation can be extended to a fully anisotropic case for which this TE/TM decoupling is no longer valid and the three components of the field have to be calculated as follows.

Let us consider the ancillary problem mentioned in Sec. 5.3.2.2, i.e. a dielectric stack made of $N$ homogeneous, isotropic, lossy layers characterized by there relative permittivity denoted $\varepsilon^{j}$ and their thickness $e_{j}$. This stack is deposited on a homogeneous, isotropic, possibly lossy substrate characterized by its relative permittivity denoted $\varepsilon^{N+1}=\varepsilon^{-}$. The superstrate is air and its relative permittivity is denoted $\varepsilon^{+}=1$. Finally, we denote by $z_{j}$ the altitude of the interface between the $j^{\text {th }}$ and $j+1^{\text {th }}$ layers. The restriction of the incident field $\mathbf{E}^{\text {inc }}$ to the superstrate region is denoted $\mathbf{E}_{0}$. The problem amounts to looking for $\left(\mathbf{E}_{1}, \mathbf{H}_{1}\right)$ satisfying Maxwell equations in harmonic regime (see Eqs. (5.56a,5.56b)).

## Across the interface $z=z_{j}$

By projection on the main axis of the vectorial Helmholtz propagation equation (Eq. (5.57)), the total electric field inside the $j^{t h}$ layer can be written as the sum of a propagative and a counter-propagative plane waves:

$$
\mathbf{E}_{1}(x, y, z)=\left[\begin{array}{c}
E_{1}^{x, j,+}  \tag{5.80}\\
E_{1}^{y, j,+} \\
E_{1}^{z, j,+}
\end{array}\right] \exp \left(j\left(\alpha_{0} x+\beta_{0} y+\gamma_{j} z\right)\right)+\left[\begin{array}{c}
E_{1}^{x, j,-} \\
E_{1}^{y, j,-} \\
E_{1}^{z, j,-}
\end{array}\right] \exp \left(j\left(\alpha_{0} x+\beta_{0} y-\gamma_{j} z\right)\right)
$$

where

$$
\begin{equation*}
\gamma_{j}^{2}=k_{j}^{2}-\alpha_{0}^{2}-\beta_{0}^{2} \tag{5.81}
\end{equation*}
$$

What follows consists in writing the continuity of the tangential components of $\left(\mathbf{E}_{1}, \mathbf{H}_{1}\right)$ across the interface $z=z_{j}$, i.e. the continuity of the vector field $\Psi$ defined by:

$$
\Psi=\left[\begin{array}{l}
E_{1}^{x}  \tag{5.82}\\
E_{1}^{y} \\
i H_{1}^{x} \\
i H_{1}^{y}
\end{array}\right] .
$$

The continuity of $\Psi$ along $O z$ together with its analytical expression inside the $j^{t h}$ and $j+1^{t h}$ layers allows to establish a recurrence relation for the interface $z=z_{j}$.

Then, by projection of Eqs. $(5.56 \mathrm{a}, 5.56 \mathrm{~b})$ on $O x, O y$ and $O z$ :

$$
\left[\begin{array}{l}
i \beta_{0} H_{1}^{z}-\frac{\partial H_{1}^{y}}{d z}  \tag{5.83}\\
\frac{\partial H_{1}^{x}}{d z}-i \alpha_{0} H_{1}^{z} \\
i \alpha_{0} H_{1}^{y}-i \beta_{0} H_{1}^{x}
\end{array}\right]=-i \omega \varepsilon\left[\begin{array}{c}
E_{1}^{x} \\
E_{1}^{y} \\
E_{1}^{z}
\end{array}\right]
$$

and

$$
\left[\begin{array}{l}
i \beta_{0} E_{1}^{z}-\frac{\partial E_{1}^{y}}{\partial z}  \tag{5.84}\\
\frac{\partial E_{1}^{x}}{\partial z}-i \alpha_{0} E_{1}^{z} \\
i \alpha_{0} E_{1}^{y}-i \beta_{0} E_{1}^{x}
\end{array}\right]=i \omega \mu\left[\begin{array}{c}
H_{1}^{x} \\
H_{1}^{y} \\
H_{1}^{z}
\end{array}\right] .
$$

Consequently, tangential components of $\mathbf{H}_{1}$ can be expressed in function of tangential components of $\mathbf{E}_{1}$ :

$$
\underbrace{\left[\begin{array}{ccc}
\omega \mu & 0 & \beta_{0}  \tag{5.85}\\
0 & \omega \mu & -\alpha_{0} \\
-\beta_{0} & \alpha_{0} & -\omega \varepsilon
\end{array}\right]}_{B}\left[\begin{array}{l}
i H_{1}^{x} \\
i H_{1}^{y} \\
i H_{1}^{z}
\end{array}\right]=\left[\begin{array}{c}
\frac{\partial E_{1}^{y}}{d z} \\
-\frac{\partial E_{1}^{x}}{d z} \\
0
\end{array}\right] .
$$

By noticing the invariance and linearity of the problem along $O x$ and $O y$, the following notations are adopted:

$$
\left\{\begin{align*}
U_{x}^{j, \pm} & =E_{1}^{x, j, \pm} \exp \left( \pm i \gamma_{j} z\right)  \tag{5.86}\\
U_{y}^{j, \pm} & =E_{1}^{y, j, \pm} \exp \left( \pm i \gamma_{j} z\right)
\end{align*}\right.
$$

and

$$
\Phi_{j}=\left[\begin{array}{c}
U_{x}^{+, j}  \tag{5.87}\\
U_{x}^{-, j} \\
U_{y}^{+, j} \\
U_{y}^{-,, j}
\end{array}\right]
$$

Thanks to Eq. (5.80) and Eq. (5.84) and letting $M=B^{-1}$, it comes for the $j^{\text {th }}$ layer:

$$
\Psi(x, y, z)=\exp \left(i\left(\alpha_{0} x+\beta_{0} y\right)\right) \underbrace{\left[\begin{array}{cccc}
1 & 1 & 0 & 0  \tag{5.88}\\
0 & 0 & 1 & 1 \\
\gamma_{j} M_{12}^{j} & -\gamma_{j} M_{12}^{j} & -\gamma_{j} M_{11}^{j} & \gamma_{j} M_{11}^{j} \\
\gamma_{j} M_{22}^{j} & -\gamma_{j} M_{22}^{j} & -\gamma_{j} M_{21}^{j} & \gamma_{j} M_{21}^{j}
\end{array}\right]}_{\Pi_{j}}\left[\begin{array}{c}
U_{x}^{+, j} \\
U_{x}^{-, j} \\
U_{y}^{+, j} \\
U_{y}^{-,, j}
\end{array}\right] .
$$

Finally, the continuity of $\Psi$ at the interface $z=z_{j}$ leads to:

$$
\begin{equation*}
\Phi_{j+1}\left(z_{j}\right)=\Pi_{j+1}^{-1} \Pi_{j} \Phi_{j}\left(z_{j}\right) \tag{5.89}
\end{equation*}
$$

Normal components can be deduced using Eqs. $(5.83,5.84)$.

## Traveling inside the $j+1^{\text {th }}$ layer

Using Eq. (5.80), a simple phase shift allows to travel from $z=z_{j}$ to $z=z_{j+1}=z_{j}-e_{j+1}$ :
$\Phi_{j+1}\left(z_{j+1}\right)=\underbrace{\left[\begin{array}{cccc}\exp \left(-i \gamma_{j+1} e_{j+1}\right) & 0 & 0 & 0 \\ 0 & \exp \left(+i \gamma_{j+1} e_{j+1}\right) & 0 & 0 \\ 0 & 0 & \exp \left(-i \gamma_{j+1} e_{j+1}\right) & 0 \\ 0 & 0 & 0 & \exp \left(+i \gamma_{j+1} e_{j+1}\right)\end{array}\right]}_{T_{j+1}} \Phi_{j+1}\left(z_{j}\right)$
Thanks to Eq. (5.90) and Eq. (5.89), a recurrence relation can be formulated for the analytical expression of $\mathbf{E}_{1}$ in each layer:

$$
\begin{equation*}
\Phi_{j+1}\left(z_{j+1}\right)=T_{j+1} \Pi_{j+1}^{-1} \Pi_{j} \Phi_{j}\left(z_{j}\right) \tag{5.91}
\end{equation*}
$$

## Reflection and transmission coefficients

The last step consists in the determination of the first term $\Phi_{0}$, which is not entirely known, since the problem definition only specifies $U_{x}^{0,+}$ and $U_{y}^{0,+}$, imposed by the incident field $\mathbf{E}_{0}$. Let us make the use of the OWC hypothesis verified by $\mathbf{E}_{1}^{d}$ (see Eq. (5.62)). This hypothesis directly translates the fact that none of the components of $\mathbf{E}_{1}^{d}$ can either be traveling down in the superstrate or up in the substrate: $U_{y}^{N+1,-}=U_{x}^{N+1,-}=0$. Therefore, the four unknowns $U_{x}^{0,-}, U_{y}^{0,-}, U_{y}^{N+1,+}$ and $U_{x}^{N+1,+}$, i.e. transverse components of the vector fields reflected and transmitted by the stack, verify the following equation system:

$$
\begin{equation*}
\Phi_{N+1}\left(z_{N}\right)=\left(\Pi_{N+1}\right)^{-1} \Pi_{N} \prod_{j=0}^{N-1} T_{N-j}\left(\Pi_{N-j}\right)^{-1} \Pi_{N-j-1} \Phi_{0}\left(z_{0}\right) \tag{5.92}
\end{equation*}
$$

This allows to extend the definition of transmission and reflection widely used in the scalar case. Finally, $\Phi_{N+1}$ is entirely defined. Making the use of the recurrence relation of Eq. (5.91) and of Eq. (5.80) leads to an analytical expression for $\mathbf{E}_{1}^{d}$ in each layer.

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Chapter 6:
Spherical harmonic Lattice Sums for Gratings

Brian Stout

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## Chapter 6

# Spherical harmonic Lattice Sums for Gratings 

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### 6.1 Introduction to particulate gratings

Lattice sums of spherical harmonic functions are well suited for modeling gratings composed of periodic arrays of identical discrete particles, henceforth referred to as particulate gratings (cf. fig. (6.1)). By discrete particles, we mean that the particles have a physical boundary such that there exists a region between the individual particles that is governed by the host material's constitutive relations. This feature makes particulate gratings somewhat different from most of the other diffraction grating problems studied in this book which are usually characterized by a substrate and a superstrate with distinct constitutive parameters. The techniques of this chapter can be extended to include the effects of a nearby planar interface, [26, 27, 28] but such considerations complicate the problem somewhat and this chapter therefore concentrates on substrate-free particulate gratings.


Figure 6.1: Particulate grating with lattice vectors $\boldsymbol{a}$ and $\boldsymbol{b}$.
Theoretical analysis of the particulate grating problem can draw on both single-particle scattering theory and techniques originally developed for solid state physics. The solid state analogy is clear from the similarity of this problem to the scattering of waves by crystal lattices, particularly in the "muffin tin" approximation[25]. Summations of the spherical harmonic fields scattered by the (infinite) number of particles in the lattice involve semi-convergent series and will generally go under the name of "lattice sums". By lattice sum, we mean sums of the form
$\Sigma_{\Lambda} \Phi\left(\boldsymbol{r}_{j}\right)$ where $\Lambda$ refers to the ensemble of points, $\boldsymbol{r}_{j}$, in a periodic lattice, and $\Phi$ is a given function.

Lattice sums have applications in many fields and their study dates back to the 19th century treating conditionally convergent sums of solutions to the Laplace equations (most notably in the Madelund constant of ionic crystals). Nevertheless, they were not always recognized as a specific branch of study, and their derivations tended to be scattered throughout the literature. This situation is changing however with the appearance of extensive reviews in recent years [15, [16, 19, 20]. Furthermore, another monograph, dedicated entirely to lattice sums, is appearing at the same time as this one.[3]

As developed in detail in the aforementioned monograph, the study of the (scale invariant) Laplace equation lattice sums have generated a number of important analytic results. The grating problem on the other hand involves propagating waves and consequently requires lattice sums of Helmholtz-type solutions. Although there are fewer fully analytic results for the (scale dependent) Helmholtz lattice sums than for the Laplace equation case, analytic manipulations remain essential for regularizing and accelerating a numerical analysis.

In solid-state physics, Helmholtz (i.e. Schrödinger) equation lattice sums are a key aspect of the Korringa-Kohn-Rostoker (KKR) methods for band-structure calculations in crystals.[22, 14, 13] In KKR theory, lattice sums intervene in the calculation of the "structure constants" of the lattice Green function and their regularization generally goes under the name of Ewald sum techniques. The Ewald sum method is quite intricate, but its basic principle can be viewed as separating a semi-convergent sum into slowly and rapidly convergent parts and then to transform the slowly convergent part into reciprocal space via the Poisson sum formula where it becomes a rapidly convergent series.

Although Ewald sum methods are proven to be quite efficient for most of the problems encountered in solid state physics, there utility has been repeatedly criticized for gratingtype applications (requiring numerically unwieldy evaluations of incomplete Gamma functions with negative real arguments[32], poor numerical properties for high multipole orders or large wavenumber, $k$ ). A number of authors have consequently looked for alternative lattice sum techniques since the pioneering work of Kerker over 30 years ago. In this chapter, we simply discuss and compare some of our preferred methods in the appendices. Our emphasis will instead be placed on painting a complete gratings-picture analysis capable of describing both near and far-field phenomenon in particulate gratings.

The matrix elements of the $\Omega$ propagation matrix introduced in section 6.3 .3 correspond to the "structure constants" of a KKR theory. More precisely, due to the differences between the Schrödinger equation and Maxwell's propagation equation, the $\Omega$ matrix elements can be written as a superposition of the KKR structure constants. In both KKR and particulate grating theory, one desires to calculate the lattice Green function. A fundamental method choice in this chapter is to use the language of T-matrices. Notably, we will see that the quasi-periodic Green function can be expressed as a lattice sum of multiple-scattering T-matrices. The multiplescattering T-matrices themselves are calculated in terms of the single-particle T-matrices, and the $\Omega$ matrix.

The T-matrix manipulations are carried out on a basis set of solutions to the Helmholtz equation, which we generally refer to as partial waves, (PWs), also commonly referred to as spherical wave functions (SWFs). This T-matrix approach is also generally adopted in the KKR calculations[22], but in the light scattering community, the terminology "T-matrix method" is often considered to be synonymous with extended boundary condition technique (also called

Null-field methods), but the T-matrix is a general theoretical construct that relates the field incident on a particle to the field scattered by the particle. As such, it can be seen as providing a complete solution to the single-particle scattering problem. In practice, the T-matrix can be generated by a wide variety of techniques including DDA, method of moments, and fictitious source techniques.

The T-matrix of an individual particle depends on the shape and the constitutive parameters of the particles, both of which can be quite arbitrary as long as the particle response is linear (including anisotropic constitutive parameters, magnetic permeability contrast etc.). However, since the T-matrix can be viewed as being the complete solution of a 1-body problem, its determination can be viewed as being separate from the grating problem. In this chapter, we simplify the T-matrix part of the problem by considering only isotropic spherical scatterers. The T-matrix of such scatterers is diagonal in the partial wave basis with its elements being determined analytically from Mie theory. We insist however, that for particulate gratings composed of more exotic scatterers, it generally suffices to insert the appropriate T-matrix to obtain the response of lattices composed of such scatterers. We refer the interested to reader to reviews of the T-matrix methods. [17, 18]

The methods developed in this chapter can be adapted to the study gratings composed of periodic infinite cylinders. However, there are fundamental differences in the mathematics, since this problem is usually addressed by solving 2 -dimensional Helmholtz equations. We therefore neglect this problem in order to concentrate on the fully 3-dimensional problem of particulate gratings like those of figure 6.1.

The first five sections constitute the heart of this chapter since they describe the general mathematical analysis of gratings using spherical harmonic lattice sums. The last four sections treat numerical methods for calculating lattice sums and special functions. Support material, corrections, and erratum will be made available and updated during the year 2013 at www.fresnel.fr/perso/stout/index.htm

### 6.2 Waves and partial waves

A fundamental aspect of the particulate gratings is that they can be viewed as a multiple-scatting phenomenon with light propagating through the host medium between individual scatterings events. The wave equation for light in this homogeneous isotropic medium is:

$$
\begin{equation*}
\boldsymbol{\nabla} \times \boldsymbol{\nabla} \times \boldsymbol{E}+k^{2} \boldsymbol{E}=\mathbf{0}, \tag{6.1}
\end{equation*}
$$

where $k=\sqrt{\varepsilon_{r} \mu_{r}} \frac{\omega}{c}$ is the wavenumber of the host medium. Solutions of eq. 6.1) satisfy both the vector Helmholtz equation,

$$
\begin{equation*}
\Delta \boldsymbol{E}+k^{2} \boldsymbol{E}=\mathbf{0}, \tag{6.2}
\end{equation*}
$$

and the additional constraint that the longitudinal field components are null:

$$
\begin{equation*}
\boldsymbol{\nabla} \cdot \boldsymbol{E}=0 . \tag{6.3}
\end{equation*}
$$

A basis set for solutions to the vector Helmholtz equation of eq.(6.2) can be readily constructed starting from the scalar Helmholtz equation:

$$
\begin{equation*}
\Delta \varphi+k^{2} \varphi=0 . \tag{6.4}
\end{equation*}
$$

As well established in textbooks, eq.(6.4) can be solved by separation of variables in spherical coordinates with 'regular' solutions taking the form of spherical harmonics, $Y_{n, m}$ multiplied by spherical Bessel functions, $j_{n}(k r)$, that are regular for all values of $r$. There are also 'irregular' solutions spherical Neumann functions, $y_{n}(k r)$, that have essential singularities as $k r \rightarrow 0$. Details concerning the properties and calculation of the $Y_{n, m}(\theta, \phi)$ are given in section 6.8.1.

The spherical coordinate solutions to eq.(6.4) will henceforth be referred to as scalar partial waves and will be defined as:

$$
\begin{equation*}
\mathscr{J}_{n, m}(k \boldsymbol{r}) \equiv j_{n}(k r) Y_{n, m}(\hat{\boldsymbol{r}}), \quad \text { and } \quad \mathscr{Y}_{n, m}(k \boldsymbol{r}) \equiv y_{n}(k r) Y_{n, m}(\widehat{\boldsymbol{r}}) . \tag{6.5}
\end{equation*}
$$

The regular partial waves, $\mathscr{J}_{n, m}$, can serve as a basis set for any source free incident field solution to eq.(6.4). Outgoing partial waves solutions of the Helmholtz equation, denoted $\mathscr{H}_{n, m}$, will be of primary interest in grating theory since they will be used to describe fields scattered by the grating. They are defined as a superposition of the regular and irregular partial waves:

$$
\begin{equation*}
\mathscr{H}_{n, m}(k \boldsymbol{r}) \equiv h_{n}(k r) Y_{n, m}(\widehat{\boldsymbol{r}})=\mathscr{J}_{n, m}(k \boldsymbol{r})+i \mathscr{Y}_{n, m}(k \boldsymbol{r}) . \tag{6.6}
\end{equation*}
$$

Incident field solutions to the vector Helmholtz equation of eq. (6.2) can be expressed as scalar partial waves associated with unit vectors along each axis i.e.:

$$
\begin{equation*}
\boldsymbol{E}_{\mathrm{i} n c}(\boldsymbol{r})=\widehat{\boldsymbol{x}} \sum_{n, m} \alpha_{n, m}^{(x)} \mathscr{J}_{n, m}(k \boldsymbol{r})+\widehat{\boldsymbol{y}} \sum_{n^{\prime}, m^{\prime}} \alpha_{n^{\prime}, m^{\prime}}^{(y)} \mathscr{J}_{n^{\prime}, m^{\prime}}(k \boldsymbol{r})+\widehat{\boldsymbol{z}} \sum_{n^{\prime \prime}, m^{\prime \prime}} \alpha_{n^{\prime \prime}, m^{\prime \prime}}^{(z)} \mathscr{J}_{n^{\prime \prime}, m^{\prime \prime}}(k \boldsymbol{r}) . \tag{6.7}
\end{equation*}
$$

The field scattered field scattered by a particle in the context of the vector Helmholtz equation can likewise be developed in terms of the outgoing spherical waves:

$$
\begin{equation*}
E_{\text {scat }}(\boldsymbol{r})=\widehat{\boldsymbol{x}} \sum_{n, m} \beta_{n, m}^{(x)} \mathscr{H}_{n, m}(k \boldsymbol{r})+\widehat{\boldsymbol{y}} \sum_{n^{\prime}, m^{\prime}} \beta_{n^{\prime}, m^{\prime}}^{(y)} \mathscr{H}_{n^{\prime}, m^{\prime}}(k \boldsymbol{r})+\widehat{\boldsymbol{z}} \sum_{n^{\prime \prime}, m^{\prime \prime}} \beta_{n^{\prime \prime}, m^{\prime \prime}}^{(z)} \mathscr{H}_{n^{\prime \prime}, m^{\prime \prime}}(k \boldsymbol{r}) \tag{6.8}
\end{equation*}
$$

provided that the coordinate system origin is chosen to lie inside the particle and that the field description is applied only to regions lying outside the particle. The field in eq. (6.8) represents the field scattered by a single scatterer, so the grating problem in terms of partial waves must sum over the field scattered by all the particles in the lattice. Finding efficient ways for calculating the lattice sum will therefore figure prominently in the subsequent sections of this chapter.

Before studying T-matrices in the next section, we first address an important technical issue. The field expansions in eq. (6.7) and eq.(6.8) have both transverse and longitudinal components and therefore are not generally solutions of the light propagation problem of eq.(6.1). The transverse wave condition of eq. (6.3) can be satisfied by requiring that the Cartesian field coefficients, $\alpha_{n, m}^{(x, y, z)}$, (and respectively $\beta_{n, m}^{(x, y, z)}$ ) satisfy certain relations amongst themselves. The important point is to remark that the constraint conditions, although somewhat complex in spherical coordinates, only affect the partial wave coefficients, and not the partial waves themselves. Consequently, in the rest of this chapter, we can generally restrain our attention to lattice sums of scalar partial wave sums even though the end goal is describe electromagnetic field scattering.

Expressing the transverse vector partial waves in terms of the Cartesian axis partial waves of eq. (6.7) or eq. (6.8) is a relatively complex affair involving angular momentum coupling formalism, coordinate transformations, and recurrence relations.[21] Consequently, it is more
common to invoke one of the various methods that have been devised over the years to directly generate transverse partial waves: Debye potentials, Hertz potentials, the Boulenkamp-Casimir approach[10], pilot vector techniques[5], etc. Whatever one's "preferred" technique and notation, the two types of transverse partial waves, $\boldsymbol{\Psi}_{\mathscr{F}, q, p}$ (often denoted $\boldsymbol{M}_{\mathscr{J}, p}$ and $\boldsymbol{N}_{\mathscr{J}, p}$ in the literature), can be expressed:

$$
\begin{align*}
\boldsymbol{\Psi}_{\mathscr{J}, 1, p}(k \boldsymbol{r}) & \equiv j_{n}(k r) \boldsymbol{X}_{n, m}(\widehat{\boldsymbol{r}}) \\
\boldsymbol{\Psi}_{\mathscr{J}, 2, p}(k \boldsymbol{r}) & \equiv \frac{1}{k r}\left\{\sqrt{n(n+1)} j_{n}(k r) \boldsymbol{Y}_{n, m}(\widehat{\boldsymbol{r}})+\left[k r j_{n}(k r)\right]^{\prime} \boldsymbol{Z}_{n, m}(\widehat{\boldsymbol{r}})\right\}, \tag{6.9}
\end{align*}
$$

where $\boldsymbol{X}_{n, m}, \boldsymbol{Y}_{n, m}$, and $\boldsymbol{Z}_{n, m}$ are the vector spherical harmonics (VSHs), (described in section eq. 6.9 .3 ). The first subscript, $\mathscr{J}$, on $\Psi_{\mathscr{J}, q, p}$ serves to indicate that the radial dependence is governed by spherical Bessel functions. A value of $q=1$ in the second subscript indicates a "transverse electric" (TE) type wave (i.e. possessing no radial electric field component), while $q=2$ indicates transverse magnetic (TM) waves. In order to minimize the number of subscripts, we adopt the common procedure that the third subscript $p$ of $\Psi_{\mathscr{J}, q, p}$, replaces the $t w o$ multipole subscripts $n$ and $m$ by defining its value such that[31]:

$$
\begin{equation*}
p(n, m) \equiv n(n+1)-m \tag{6.10}
\end{equation*}
$$

With the notation of eq. 6.9, one can express any incident field satisfying equation 6.1) in terms of the transverse vector partial waves:

$$
\begin{equation*}
\boldsymbol{E}_{\mathrm{inc}}(\boldsymbol{r})=E \sum_{q=1,2} \sum_{p=1}^{\infty} \boldsymbol{\Psi}_{\mathscr{J}, q, p}(k \boldsymbol{r}) a_{q, p} \tag{6.11}
\end{equation*}
$$

where $a_{q, p}$ are (dimensionless) field coefficients, and $E$ is a constant with the dimension of electric field and which can be used to adjust the field strength. With this notation, the field scattered by a particle whose circumscribing sphere is centered at a position $\boldsymbol{x}_{j}$ can be written:

$$
\begin{equation*}
\boldsymbol{E}_{\mathrm{scat}}\left(\boldsymbol{r}_{j}\right)=E \sum_{q=1,2} \sum_{p=1}^{\infty} \boldsymbol{\Psi}_{\mathscr{H}, q, p}\left(k \boldsymbol{r}_{j}\right) f_{q, p}^{(j)} \tag{6.12}
\end{equation*}
$$

where $\boldsymbol{r}_{j} \equiv \boldsymbol{r}-\boldsymbol{x}_{j}$, and $f_{q, p}^{(j)}$ are the scattering coefficients of the particle $j$. The index, $\mathscr{H}$. on the $\Psi_{\mathscr{H}, q, p}$ indicates that the radial dependence should be governed by spherical Hankel functions, $h_{n}(k r)$, rather than the spherical Bessel functions, $j_{n}(k r)$, found in the $\Psi_{\mathscr{J}, q, p}$ functions of eq. 6.9).

### 6.3 T-matrix theory

### 6.3.1 Green functions and T-matrices

The fundamental object that one would like to calculate in a multiple-scattering system (like a particulate grating) is the system Green's function. However, the dyadic Green's function for a homogeneous medium has a strongly singular behavior and needs to be defined in the context of distributions.[5] The T-matrix formalism allows us to largely circumvent this singular behavior, and also to work directly in terms of fields which is often more manageable than the
relatively intricate dyadic Green's function formalism. For instance, the operator form of the Green function of a single object in a homogeneous medium can be written:

$$
\begin{equation*}
G=g+g t \boldsymbol{g} \tag{6.13}
\end{equation*}
$$

where $t$ is the isolated particle (or l-body) T-matrix operator, and $g$ is the Green function operator of the homogeneous medium (sometimes called a propagator). In this formalism, the singular behavior is relegated to the propagator, $\boldsymbol{g}$, leaving the (non-singular) scattering response due to the object being described by $\boldsymbol{t}$.

Furthermore, when considering excitations outside the scatterer, the homogeneous Green function, $\boldsymbol{g}$ to the right of the $\boldsymbol{t}$ operator acting on the sources generates the incident field, while the $\boldsymbol{g}$ to the left of it generates the scattered field.[24] In the partial wave basis, $\boldsymbol{t}$ then truly takes the form of a matrix, henceforth denoted, $t$, that relates the incident field coefficients to the scattered field coefficients:

$$
\begin{equation*}
f=t a, \tag{6.14}
\end{equation*}
$$

where $a$ and $f$ are column matrices composed respectively of the incident field and scattered field coefficients (cf. eqs. 6.11 and (6.12)).[30] The 1-body T-matrix, $t$, in this expression is now truly a matrix relating field coefficients of partial wave field decompositions.

This T-matrix formalism can be extended to include systems containing $N$ particles. The system Green function can be written,

$$
\begin{equation*}
\boldsymbol{G}=\boldsymbol{g}+\boldsymbol{g}\left(\sum_{j=1}^{N} \boldsymbol{T}^{(j)}\right) \boldsymbol{g} \tag{6.15}
\end{equation*}
$$

where the multiple-scattering (or $N$-body) T-matrix operators, $\boldsymbol{T}^{(j)}$, are associated with each particle and which incorporates all the multiple-scattering effects due to the presence of the $N-1$ other particles in the system. Passing once again to a partial wave field description, the multiple-scattering T-matrix, $T^{(j)}$, generates the field scattered by each particle in terms of the field incident on the system:

$$
\begin{equation*}
f^{(j)}=T^{(j)} a^{(j)} \tag{6.16}
\end{equation*}
$$

where $a^{(j)}$ indicates the incident field developed on a coordinate system centered on the $j^{t h}$ particle. All multiple-scattering phenomenon and some rather subtle technical difficulties have all been incorporated into the definition of $T^{(j)}$, but nowadays they can be calculated rather readily for systems with a finite number of particles starting from the 1-body T-matrices, $t^{(j)}$, of the individual particles.[30]

The number of particles in a grating problem is infinite (from the ideal mathematical standpoint), which given their physical content would render exact calculations of $T^{(j)} \mathrm{im}$ possible. Nevertheless, the fact that the system is identical when viewed from any lattice site allows the $T^{(j)}$ matrices to be calculated as a lattice sum as we shall see in section 6.3.3. We first rapidly review below our notation and terminology for lattices.

### 6.3.2 Direct and reciprocal lattices

A lattice, $\Lambda$, of dimension $d_{\Lambda}$, is invariant under a coordinate system translation along any vector, $\boldsymbol{r}_{\boldsymbol{j}}$, that can be expressed

$$
\begin{equation*}
r_{j}=\sum_{i=1}^{d_{\Lambda}} j_{i} a_{i} \tag{6.17}
\end{equation*}
$$

where $\boldsymbol{a}_{i}$ are the primitive lattice vectors, and $\boldsymbol{j} \equiv\left(j_{1}, \ldots, j_{d_{\Lambda}}\right)$ is a shorthand notation for a set of $d_{\Lambda}$ relative integers, $j_{i} \in \mathbb{Z}$. In order to diminish the number of subscripts, we will sometimes employ an alternative notation for the primitive lattice vectors: $\boldsymbol{a} \equiv \boldsymbol{a}_{1}, \boldsymbol{b} \equiv \boldsymbol{a}_{2}$, and $\boldsymbol{c} \equiv \boldsymbol{a}_{3}$. It also proves convenient to define the $x$ and $y$ axis of the system so that the primitive lattice vectors can be expressed: $\boldsymbol{a}=(a, 0,0), \boldsymbol{b}=\left(b_{x}, b_{y}, 0\right)$, and $\boldsymbol{c}=\left(c_{x}, c_{y}, c_{z}\right)$.

When $d_{\Lambda}=3$, the $\boldsymbol{r}_{\boldsymbol{j}}$ ensemble defines a crystalline type lattice, henceforth denoted ( $L$ ), as frequently encountered in photonic crystals and "meta-materials". A two-dimensional grating, or mono-layer lattice $(M L)$, like that of figure 6.1 , occurs when the system invariance only occurs for 2D displacements of $\boldsymbol{r}_{\boldsymbol{j}}=j_{a} \boldsymbol{a}+j_{b} \boldsymbol{b}$. Finally, linear chains $(C)$ are only invariant with respect to translations of $\boldsymbol{r}_{j}=j \boldsymbol{a}$.

The reciprocal lattice, $\Lambda^{*}$, is defined in terms of lattice 'wave-vectors', $\boldsymbol{p}_{\boldsymbol{g}}$, defined in terms of the primitive reciprocal lattice vectors, $\widetilde{\boldsymbol{a}}_{i}$ :

$$
\begin{equation*}
\boldsymbol{p}_{\boldsymbol{g}}=2 \pi \sum_{i=1}^{d_{\Lambda}} g_{i} \widetilde{\boldsymbol{a}}_{i} \tag{6.18}
\end{equation*}
$$

where $g_{i} \in \mathbb{Z}$. The primitive reciprocal vectors, $\widetilde{\boldsymbol{a}}_{j}$, are defined such that their scalar products with respect to $\boldsymbol{a}_{j}$ satisfy:

$$
\begin{equation*}
\boldsymbol{a}_{i} \cdot \widetilde{\boldsymbol{a}}_{j}=\delta_{i j} \quad i, j=1, \ldots, d_{\Lambda} \tag{6.19}
\end{equation*}
$$

From eqs. $6.17,\left(6.18\right.$ and 6.19 , one readily finds that the reciprocal lattice vectors, $\boldsymbol{p}_{\boldsymbol{g}}$, of eq. 6.18), have the property

$$
\begin{equation*}
\boldsymbol{r}_{\boldsymbol{j}} \cdot \boldsymbol{p}_{\boldsymbol{g}}=2 \pi N \tag{6.20}
\end{equation*}
$$

where $N$ is some integer (which results in $\exp \left(i \boldsymbol{r}_{\boldsymbol{j}} \cdot \boldsymbol{p}_{\boldsymbol{g}}\right)=1$ for all $\boldsymbol{r}_{\boldsymbol{j}}$ and $\boldsymbol{p}_{\boldsymbol{g}}$ ).
The unit cell "volume", $\mathscr{A}$, appears repeatedly in theories of particulate lattices. For lattice dimensions of $d_{\Lambda}=1,2$, and 3 , the corresponding $\mathscr{A}_{1,2,3}$ is given by:

$$
\begin{cases}\mathscr{A}_{1}=|\boldsymbol{a}| & d_{\Lambda}=1  \tag{6.21}\\ \mathscr{A}_{2}=|\boldsymbol{a} \times \boldsymbol{b}| & d_{\Lambda}=2 \\ \mathscr{A}_{3}=|(\boldsymbol{a} \times \boldsymbol{b}) \cdot \boldsymbol{c}| & \\ d_{\Lambda}=3\end{cases}
$$

with dimensions of "length" for $d_{\Lambda}=1$, "area" for $d_{\Lambda}=2$ and "volume" for $d_{\Lambda}=3$. The corresponding "volume" of the reciprocal space lattice sites are given by $\mathscr{A}^{-1}$.

### 6.3.3 Grating T-matrices

Each site of a lattice is identical to all the others so that the multiple scattering T-matrices of eq. 6.16 are all equal, i.e. $T^{(j)}=T$. The scattering coefficients $f^{(j)}$ are then given by:

$$
\begin{equation*}
f^{(j)}=T a^{(j)} \tag{6.22}
\end{equation*}
$$

The trouble with this equation is that the coefficients $f^{(j)}$ and $a^{(j)}$ are expressed on localized partial wave basis, but the long range nature of scattered fields would require the $T$-matrices to act on very high multipole orders in order to account for these long-rang interactions.

Since manipulating high multipole orders is numerically inefficient, one considerably simplifies this problem by only calculating the multiple-scattering $T$-matrices for incident fields
satisfying a quasi-periodic condition. Quasi-periodicity can be viewed as requiring the partialwave decomposition of the incident field on each lattice site, $\boldsymbol{r}_{\boldsymbol{j}}$, to satisfy,

$$
\begin{equation*}
a^{(j)}=e^{i \boldsymbol{\beta} \cdot \boldsymbol{r}_{j}} a \tag{6.23}
\end{equation*}
$$

where ' $a$ ' corresponds to the incident field coefficients at the origin, and $\boldsymbol{\beta}$, the 'on-shell' quasiperiodicity vector. The term 'on-shell' indicates that the quasi-periodic vector satisfy $|\boldsymbol{\beta}|=k$, and although this considerably restricts the type of incident fields (and modes) that one can treat, it does correspond to many situations of physical interest including both homogeneous and inhomogeneous plane-waves incident on diffraction gratings.

The quasi-periodic condition can be viewed as a partial Fourier transform description in that the overall field behavior is of an oscillatory nature, while the quasi-periodic T-matrix, $T_{\beta}$, describes local-field perturbations due to the presence of the particles. Consequently, one expects the $T_{\boldsymbol{\beta}}$ matrices to be well approximated on a truncated (i.e. finite) partial-wave basis (similar to the behavior of the isolated particle T-matrices[30]). The quasi-periodic condition allows the Foldy-Lax equations for the multiple-scattering T-matrices to take the form:

$$
\begin{equation*}
T_{\boldsymbol{\beta}}=t+t \boldsymbol{\Omega}_{\boldsymbol{\beta}} T_{\boldsymbol{\beta}}, \tag{6.24}
\end{equation*}
$$

where the $\Omega_{\beta}$ matrix designates a quasi-periodic lattice sum of the irregular translation-addition matrices:

$$
\begin{equation*}
\Omega_{\boldsymbol{\beta}}\left(k, \boldsymbol{a}_{i \leq d}\right)=\sum_{\substack{r_{j} \in \Lambda \\ r_{\boldsymbol{j}} \neq 0}} e^{i \boldsymbol{\beta} \cdot \boldsymbol{r}_{\boldsymbol{j}}} H\left(k \boldsymbol{r}_{\boldsymbol{j}}\right) . \tag{6.25}
\end{equation*}
$$

The analytical properties of the irregular translation-addition matrix, $H(\boldsymbol{x})$, are described in section 6.8.3 where one also gives expressions for its matrix elements.

The exclusion of the 'origin' lattice site, $\boldsymbol{r}_{\boldsymbol{j}}=\mathbf{0}$, from the sum in $\Omega_{\boldsymbol{\beta}}$ has a physical significance in that it accounts for propagation of the light scattered by all the other particles in the lattice onto the particle at the origin (the light 'scattered' by the particle onto itself has already been included in the individual T-matrix, $t$ ). One finds in section 6.8.3 that each matrix element of the translation-addition matrix, $H\left(k \boldsymbol{r}_{j}\right)$, can be written:

$$
\begin{equation*}
\left[H\left(k \boldsymbol{r}_{j}\right)\right]_{p, q ; p^{\prime}, q^{\prime}}=\sum_{l, m} C_{l, m}\left(p, q ; p^{\prime}, q^{\prime}\right) h_{l}\left(k r_{\boldsymbol{j}}\right) Y_{l, m}\left(\widehat{\boldsymbol{r}}_{\boldsymbol{j}}\right) \tag{6.26}
\end{equation*}
$$

where the sum over the multipole indices, $(l, m)$ is finite. Expressions for the $C_{l, m}\left(p, q ; p^{\prime}, q^{\prime}\right)$ coefficients [31, 5, 29] are given in the section 6.8, Inserting eq.(6.26) into eq. (6.25) and rearranging the summations, we find

$$
\begin{align*}
{\left[\Omega_{\boldsymbol{\beta}}\right]_{p, q ; p^{\prime}, q^{\prime}} } & =\sum_{l, m} C_{l, m}\left(p, q ; p^{\prime}, q^{\prime}\right) \sum_{\substack{r_{j} \in \Lambda \\
r_{j} \neq 0}} e^{i \boldsymbol{\beta} \cdot \boldsymbol{r}_{j}} h_{l}\left(k r_{j}\right) Y_{l, m}\left(\hat{\boldsymbol{r}}_{j}\right) \\
& \equiv \sum_{l, m} C_{l, m}\left(p, q ; p^{\prime}, q^{\prime}\right) S_{l, m}(k, \boldsymbol{\beta}) \tag{6.27}
\end{align*}
$$

where we have defined $S_{l, m}(k, \boldsymbol{\beta})$ as a Hankel function lattice sum such that:

$$
\begin{equation*}
S_{n, m}(k, \boldsymbol{\beta}) \equiv S_{n, m}^{\mathscr{H}}(k, \boldsymbol{\beta}) \equiv \sum_{\substack{r_{j} \in \Lambda \\ r_{\boldsymbol{j}} \neq 0}} e^{i \boldsymbol{\beta} \cdot \boldsymbol{r}_{\boldsymbol{j}}} \mathscr{H}_{n, m}\left(k \boldsymbol{r}_{\boldsymbol{j}}\right) . \tag{6.28}
\end{equation*}
$$

We recall that $\mathscr{H}_{n, m}(\boldsymbol{x})$ was defined in eq. 6.6) as a partial wave of the spherical Hankel function type.

It will occasionally prove useful to calculate the analogous lattice sums over the partial waves of the Bessel or Neumann types, denoted respectively, $S_{n, m}^{\mathscr{f}}(k, \boldsymbol{\beta})$ and $S_{n, m}^{\mathscr{V}}(k, \boldsymbol{\beta})$. Since we will principally be concerned with the partial wave lattice sums of the Hankel function type, $S_{n, m}$ without a superscript will always indicate a lattice sum of the Hankel function type. We also remark that the exclusion of the origin position from the lattice sum is important from a mathematical standpoint since the Hankel functions have an essential singularity at their origin.

The solution to eq. (6.24) for the multiple-scattering T-matrix is readily formulated in terms of matrix inversion:

$$
\begin{equation*}
T_{\boldsymbol{\beta}}=\left[t^{-1}-\Omega_{\boldsymbol{\beta}}\right]^{-1} . \tag{6.29}
\end{equation*}
$$

Once the $T_{\boldsymbol{\beta}}$ matrix is known, the scattering field coefficients for any particle, $j$, in the lattice is the same as the coefficients at the origin but multiplied by a $e^{i \boldsymbol{\beta} \cdot \boldsymbol{r}_{j}}$ phase factor. In the matrix notation, this is simply expressed:

$$
\begin{equation*}
f_{\boldsymbol{\beta}}^{(j)}=e^{i \boldsymbol{\beta} \cdot r_{j}} f_{\boldsymbol{\beta}}=e^{i \boldsymbol{\beta} \cdot \boldsymbol{r}_{j}} T_{\boldsymbol{\beta}} a \tag{6.30}
\end{equation*}
$$

where $a$ is the column matrix composed of the incident field coefficients developed around the origin.

### 6.3.3.1 Far-fields

The field 'scattered' by the grating (i.e. 'transmitted' and 'reflected' diffraction orders) can be determined by inserting eq. 6.30) into eq. (6.12) wherein the scattered field takes the form of a lattice sum of the transverse-outgoing-vector partial waves, $\boldsymbol{\Psi}_{\mathscr{H}, q, p}$, described in eq. 6.9 p section 6.2.

$$
\begin{align*}
\boldsymbol{E}_{\mathrm{s}, \Lambda}(\boldsymbol{r}) & =E \sum_{\boldsymbol{r}_{\boldsymbol{j}} \in \Lambda} e^{i \boldsymbol{\beta} \cdot \boldsymbol{r}_{\boldsymbol{j}}} \boldsymbol{\Psi}_{\mathscr{H}}\left(k \boldsymbol{r}_{\boldsymbol{j}}\right) f_{\boldsymbol{\beta}} \\
& \equiv E \sum_{q=1,2} \sum_{p=1}^{\infty}\left[\sum_{r_{j} \in \Lambda} \boldsymbol{\Psi}_{\mathscr{H}, q, p}\left(k \boldsymbol{r}_{\boldsymbol{j}}\right) e^{i \boldsymbol{\beta} \cdot \boldsymbol{r}_{\boldsymbol{j}}}\right]\left[f_{\boldsymbol{\beta}}\right]_{q, p} . \tag{6.31}
\end{align*}
$$

We will see in eq. 6.62) of section 6.4.4 that for each multipole order, $p=1, \ldots \infty$, and transverse wave type $q=1,2$; the term in brackets can be re-expressed as an infinite sum of plane waves. Only a finite subset of these waves are of the propagative type however (the rest are all of an evanescent nature). Consequently, the multipole summation of eq. 6.31) allows one to calculate the efficiency of each reflected or transmitted order in the far field.

### 6.3.3.2 Near-fields

Another quantity of physical interest is that of near fields in a particulate grating (non-linear effects, SERS, etc.). The plane wave expansion discussed above for far fields could be invoked in principal, but for near fields one must also calculate the (infinite) evanescent orders that one could neglect in the far field. The convergence of the plane wave expansion will generally be poor near the grating, which renders this approach unattractive.

As long as the incident field is quasi-periodic with respect to the grating, one needs only to determine the near fields in a single Brillouin zone around a given lattice site (the site at the origin being the most practical). In this case, it seems clear that the localized multipolar field developments are well adapted to the development of the local field in the Brillouin zone. In multiple scattering theory, the $f_{\boldsymbol{\beta}}$ coefficients give the field scattered by the particle at the origin, while the excitation field corresponds to the field at that was 'incident' on this particle, i.e. the superposition of the field incident on the grating and the field scattered by all the other particles in the system. This excitation field can be developed on the regular partial waves and its coefficients, $e_{\boldsymbol{\beta}}$, related to the scattering coefficients via the 1-body T-matrix via the relation:

$$
\begin{equation*}
e_{\boldsymbol{\beta}}=t^{-1} f_{\boldsymbol{\beta}} \tag{6.32}
\end{equation*}
$$

The total field in the Brillouin zone is simply a superposition of the scattered and excitation field:

$$
\begin{align*}
\boldsymbol{E}_{\mathrm{t}}^{(\mathrm{B} . \mathrm{z})}(\boldsymbol{r}) & =E\left(\boldsymbol{\Psi}_{\mathscr{H}}(k \boldsymbol{r}) f_{\boldsymbol{\beta}}+\boldsymbol{\Psi}_{\mathscr{J}}(k \boldsymbol{r}) t^{-1} f_{\boldsymbol{\beta}}\right) \\
& \equiv E \sum_{q=1,2} \sum_{p=1}^{\infty}\left[\boldsymbol{\Psi}_{\mathscr{H}, q, p}(k \boldsymbol{r}) f_{q, p}+\boldsymbol{\Psi}_{\mathscr{J}, q, p}(k \boldsymbol{r}) e_{q, p}\right] . \tag{6.33}
\end{align*}
$$

### 6.3.3.3 Propagating modes

The response of a grating or chain of particles will often be dominated by the excitation of nearly guided modes in the structure (commonly referred to as leaky modes). Unlike, lossless 3D lattices which have true propagating modes described by real values of $\boldsymbol{\beta}$, the leaky modes in a grating or chain of particles are so-named because they generally radiate energy to the far field even in the absence of intrinsic losses. Consequently, leaky modes can be described by a complex valued $\boldsymbol{\beta}$-vector (or complex frequency). The determination of a leaky mode thus involves searching for a complex pole in the determinant of the multiple-scattering T-matrix, $\left|T_{\beta}\right|$.

Since matrix inversions are numerically expensive, one may prefer to look for zero eigenvectors, $v_{\alpha}$, of the matrix $\left[t^{-1}-\Omega_{\boldsymbol{\beta}_{\alpha}}\right]$, i.e.

$$
\begin{equation*}
\left[t^{-1}-\Omega_{\boldsymbol{\beta}_{\alpha}}\right] v_{\alpha}=0 . \tag{6.34}
\end{equation*}
$$

However, the search for zero eigenvalues can limit the implantation of complex analysis methods that have proven useful in determining the position of poles in the complex plane.

The Floquet mode associated with the eigenvector, $v_{\alpha}$, can be constructed from eq.(6.34) coupled with the plane-wave development the terms in eigenvector, $v_{\alpha}$ :

$$
\begin{equation*}
\boldsymbol{E}_{\alpha}^{(\text {F.m. })}(\boldsymbol{r})=E \sum_{q=1,2} \sum_{p=1}^{\infty}\left[\sum_{\boldsymbol{r}_{j} \in \Lambda} \boldsymbol{\Psi}_{\mathscr{H}, q, p}\left(k \boldsymbol{r}_{\boldsymbol{j}}\right) e^{i \boldsymbol{\beta}_{\alpha} \cdot \boldsymbol{r}_{j}}\right]\left[v_{\alpha}\right]_{q, p} \tag{6.35}
\end{equation*}
$$

Before finishing this section, it should be pointed out that matrix inversion solutions to the multiple scattering problem (like that given in eq. (6.29)) were disregarded for a long time in favor of iterative solutions to the T-matrix or underlying linear system of equations. The reason for this is that the matrix $\left[t^{-1}-\Omega_{\boldsymbol{\beta}_{\alpha}}\right]$ is generally ill-conditioned. This difficulty can be generally overcome by analytical matrix balancing as described in the next section 6.3.4.

### 6.3.4 Matrix balancing

Although not necessary from a formal standpoint, analytical matrix balancing improves the conditioning of the matrices occurring in multiple-scattering calculations for both matrix inversion and eigenvalue resolution. [30] Analytical matrix balancing can be achieved by multiplying a matrix from both the right and left by diagonal matrices, $[\xi]$ and $[\psi]^{-1}$, whose matrix elements are given by:

$$
\begin{equation*}
[\psi]_{q, q^{\prime}, p, p^{\prime}}=\delta_{q, q^{\prime}} \delta_{p, p^{\prime}} \psi_{n}(k R), \quad[\xi]_{q, q^{\prime}, p, p^{\prime}}=\delta_{q, q^{\prime}} \delta_{p, p^{\prime}} \xi_{n}(k R), \tag{6.36}
\end{equation*}
$$

where $\psi_{n}(k R)$ and $\xi_{n}(k R)$ are respectively the regular and irregular spherical Ricatti-Bessel functions (cf. 6.168) and $R$ the radius of the minimal circumscribing sphere surrounding the scatterers.

Matrix balancing can be readily formulated by defining normalized incident and scattering coefficients, $\bar{a}$ and $\bar{f}$ respectively such that:

$$
\begin{equation*}
\bar{a} \equiv[\psi] a, \quad \bar{f}_{\boldsymbol{\beta}} \equiv[\xi] f_{\boldsymbol{\beta}} \tag{6.37}
\end{equation*}
$$

The associated normalized or 'balanced' matrices are defined[30]:

$$
\begin{equation*}
\bar{t} \equiv[\xi] t[\psi]^{-1}, \quad \bar{T}_{\boldsymbol{\beta}} \equiv[\xi] T_{\boldsymbol{\beta}}[\psi]^{-1}, \quad \bar{\Omega}_{\boldsymbol{\beta}} \equiv[\psi] \Omega_{\boldsymbol{\beta}}[\xi]^{-1} \tag{6.38}
\end{equation*}
$$

The above definitions were chosen such that eqs. 6.22 and 6.24 respectively take the form:

$$
\begin{equation*}
\bar{f}^{(j)} \equiv \bar{T}_{\boldsymbol{\beta}} \bar{a}^{(j)} \tag{6.39}
\end{equation*}
$$

and

$$
\begin{equation*}
\bar{T}_{\boldsymbol{\beta}} \bar{a}=\bar{t} \bar{a}+\bar{t} \bar{\Omega}_{\boldsymbol{\beta}} \bar{T}_{\boldsymbol{\beta}} \bar{a} . \tag{6.40}
\end{equation*}
$$

The normalized T-matrix, $\bar{T}$, is then obtained via the generally well-conditioned matrix inversion:

$$
\begin{equation*}
\bar{T}_{\boldsymbol{\beta}}=\left[\bar{t}^{-1}-\bar{\Omega}_{\boldsymbol{\beta}}\right]^{-1} \tag{6.41}
\end{equation*}
$$

Since we generally want the non-normalized T-matrix for applications, we reconstruct, $T_{\boldsymbol{\beta}}$ via a final multiplication by our diagonal matrices:

$$
\begin{equation*}
T_{\boldsymbol{\beta}}=[\xi]^{-1} \bar{T}_{\boldsymbol{\beta}}[\psi] \tag{6.42}
\end{equation*}
$$

### 6.4 Mathematical relations for lattice sums

This section is dedicated to reviewing the mathematical relations that allow one to treat lattice sums for lattices of dimensions $d_{\Lambda}=1,2,3$ (i.e. particulate chains, gratings, and crystals). They will notably allow us to evaluate the lattice sum in eq. 6.27) which is used to calculate far-field response from gratings. These relations were derived (and often rederived) in many places, and we refer the reader to refs. [15, 16, 19, 20, 8] for additional details and perspectives.

The most difficult mathematical problem to address will be the evaluation of Hankel function lattice sum, $S_{n, m}^{\mathscr{H}}$ that was introduced in eq. 6.28 for the calculation of the $\Omega_{\beta}$ matrix of eq. 6.27).

$$
\begin{equation*}
S_{n, m}^{\mathscr{H}} \equiv \sum_{\substack{\boldsymbol{r}_{j} \in \Lambda \\ r_{j} \neq 0}} e^{i \boldsymbol{\beta} \cdot \boldsymbol{r}_{j}} \mathscr{H}_{n, m}\left(k \boldsymbol{r}_{\boldsymbol{j}}\right) \tag{6.43}
\end{equation*}
$$

When not otherwise specified, partial waves lattice sums will always be assumed to be of the Hankel functions type. The underlying reason for this appears in the translation-addition where Hankel functions allow one to re-express waves scattered by a given lattice site as waves incident on a different lattice site.

We will occasionally consider Bessel and Neumann types of scalar partial waves:

$$
\begin{align*}
& S_{n, m}^{\mathscr{Y}} \equiv \sum_{\substack{\boldsymbol{r}_{j} \in \Lambda \\
r_{j} \neq 0}} e^{i \boldsymbol{\beta} \cdot \boldsymbol{r}_{j} \mathscr{Y}_{n, m}\left(k \boldsymbol{r}_{j}\right)} \\
& S_{n, m}^{\mathscr{Z}} \equiv \sum_{\substack{r_{j} \in \Lambda \\
r_{j} \neq 0}} e^{i \boldsymbol{\beta} \cdot \boldsymbol{r}_{j}} \mathscr{J}_{n, m}\left(k \boldsymbol{r}_{\boldsymbol{j}}\right), \tag{6.44}
\end{align*}
$$

The interest of these sums in part is due to the fact that $S_{n, m}^{\mathscr{H}}=S_{n, m}^{\mathscr{\mathscr { L }}}+i S_{n, m}^{\mathscr{Y}}$ but also because $S_{n, m}^{\mathscr{\mathscr { L }}}$ can potentially prove useful in certain applications. We will see that the relations developed in this chapter permit the $S_{n, m}^{\mathscr{f}}$ sum to be evaluated in closed form, but unfortunately the Neumann partial wave sum, $S_{n, m}^{\mathscr{Y}}$, appears to be as difficult to evaluate as the Hankel partial wave sum, and a closed form expression does not appear to be possible.

### 6.4.1 Lattice reduction

Although Ewald sums are a time honored technique in solid state physics, a considerable amount of effort has recently been devoted to what has come to be called Lattice reduction techniques. The basic idea turns around the fact that lattice sums tend to be more practical for $d_{\Lambda}=1$ and $d_{\Lambda}=3$ than for the grating dimension of $d_{\Lambda}=2$.

First one chooses a coordinate system such that a preferred axis (like the $z$ axis) will along a given lattice vector. For example $\boldsymbol{a}=(0,0, a)$, and $\boldsymbol{b}=\left(0, b_{2}, b_{1}\right)$. With this basis the 2D Mono-Layer lattice sum, $S_{n, m}^{M L}$ can then be written expressed as a superposition of a Chain sum, $S_{n, m}^{C}$, containing the origin ( $z$-axis), and a superposition of all the chain sums 'above' the central chain $(z>0)$, denoted $S_{n, m}^{M L+}$ or 'below' the central chain, $S_{n, m}^{M L-}(z<0)$. The central chain can be readily be evaluated using one of the techniques described in this chapter, while the integral expression for Hankel functions described in section 6.4.4 allows one to derive efficient expressions for $S_{n, m}^{M L+}$ and $S_{n, m}^{M L-}$.

Lattice reduction can also be applied in the reverse direction, with one expressing the 3D crystalline lattice sum, $S_{n, m}^{L}$, as the superposition of a monolayer sum in the $z=0$ plane, with sums of all monolayers with $z>0, S_{n, m}^{L+}$ and all monolayers with $z<0, S_{n, m}^{L-}$. There exists efficient techniques for calculating the crystalline lattice sum, $S_{n, m}^{L}$ while one can determine efficient expressions for $S_{n, m}^{L,+}$ and $S_{n, m}^{L,-}$ using again the integral expressions of section 6.4.4. In this reverse lattice sum method, the monolayer lattice sum is expressed:

$$
\begin{equation*}
S_{n, m}^{M L}=S_{n, m}^{L}-S_{n, m}^{M L,+}-S_{n, m}^{L,-} . \tag{6.45}
\end{equation*}
$$

One should that the choice of orientation of the coordinate axis will not be the same in general for different lattice sum techniques, but these differences can be be compensated for by using the rotation matrices of section 6.8.4

Lattice reduction is based on the idea that it can prove numerically efficient to carry out lattice sums for a lattice dimensions other than that desired. To construct a 3D periodic media,
we rotate the 2 D lattice of the preceding section back to an orientation in the $x O y$ plane with $\boldsymbol{a}=(a, 0,0)$, and $\boldsymbol{b}=\left(b_{1}, b_{2}, 0\right)$, and now $\boldsymbol{c}=\left(0, c_{2}, c_{3}\right)$. The quasi-periodic vector is given by $\boldsymbol{\beta}=\left(\beta_{1}, \beta_{2}, \beta_{3}\right)$. The 3D lattice sum can then be written:

$$
\begin{equation*}
S_{n, m}^{L}=S_{n, m}^{M L}+S_{n, m}^{L+}+S_{n, m}^{L-} \tag{6.46}
\end{equation*}
$$

The $S_{n, m}^{L+}$ denotes all the $z>0$ planes, while $S_{n, m}^{L-}$ sums all the $z<0$ planes.
The lattice reduction technique breaks the sum down into elements which tend to have a decreasing difficulties for divergence. An interest of the lattice reduction technique is that it can be adapted to partial lattices. For instance, large but finite chains, a finite number of infinite chains or finally a finite number of infinite planes.

### 6.4.2 Plane wave expansion

The expansion of a plane wave in terms of partial waves allows one to transform between partial wave and Fourier transforms. It reads:

$$
\begin{align*}
e^{i \boldsymbol{k} \cdot \boldsymbol{r}} & =4 \pi \sum_{v=0}^{\infty} \sum_{\mu=-v}^{\mu=v} i^{v} j_{v}(k r) Y_{v, \mu}^{*}(\widehat{\boldsymbol{k}}) Y_{v, \mu}(\widehat{\boldsymbol{r}}) \\
& =\sum_{v=0}^{\infty} \sum_{\mu=-v}^{\mu=v} p_{v, \mu} \Psi_{v, \mu}(\boldsymbol{r}), \tag{6.47}
\end{align*}
$$

where $\Psi_{v, \mu}(\boldsymbol{r})$ are the scalar partial wave functions discussed in section 6.2, and $p_{v, \mu}$ the coefficients in the development of a scalar plane wave on a partial wave basis i.e. :

$$
\begin{equation*}
\Psi_{v, \mu}(\boldsymbol{r}) \equiv j_{v}(k r) Y_{v, \mu}(\widehat{\boldsymbol{r}}), \quad p_{n, m}=4 \pi i^{n} Y_{n, m}^{*}(\widehat{\boldsymbol{k}}) . \tag{6.48}
\end{equation*}
$$

One can produce an integral expression of $j_{n}(k r) Y_{n, m}(\hat{\boldsymbol{r}})$ by multiplying both sides of eq. 6.47 by $Y_{n, m}(\widehat{\boldsymbol{k}})$ and integrating over all directions of $\widehat{\boldsymbol{k}}$.

$$
\begin{align*}
\int d \Omega_{\boldsymbol{k}} e^{i \boldsymbol{k} \cdot \boldsymbol{r}} Y_{n, m}(\widehat{\boldsymbol{k}}) & =4 \pi \int d \Omega_{k} \sum_{v=0}^{\infty} \sum_{\mu=-v}^{\mu=v} i^{v} j_{v}(k r) Y_{v, \mu}(\widehat{\boldsymbol{r}}) Y_{v, \mu}^{*}(\widehat{\boldsymbol{k}}) Y_{n, m}(\widehat{\boldsymbol{k}}) \\
& =4 \pi i^{n} j_{n}(k r) Y_{n, m}(\widehat{\boldsymbol{r}}) . \tag{6.49}
\end{align*}
$$

We have thus found that regular partial waves are an angular Fourier transform of the spherical harmonics:

$$
\begin{equation*}
\Psi_{\mathscr{J}, n, m} \equiv j_{n}(k r) Y_{n, m}(\widehat{\boldsymbol{r}})=\frac{1}{4 \pi i^{n}} \int d \Omega_{k} e^{i \boldsymbol{k} \cdot \boldsymbol{r}} Y_{n, m}(\widehat{\boldsymbol{k}}) \tag{6.50}
\end{equation*}
$$

Likewise, the transverse regular partial waves, $\Psi_{\mathscr{J}}$ can be expressed as an angular Fourier transform of the vector spherical harmonics:

$$
\begin{align*}
& \boldsymbol{\Psi}_{\mathscr{J}, q=1, n, m}(k \boldsymbol{r})=\frac{i^{-n}}{4 \pi} \int d \Omega_{\boldsymbol{k}} e^{i \boldsymbol{k} \cdot \boldsymbol{r}} \boldsymbol{X}_{n, m}(\widehat{\boldsymbol{k}}) \\
& \boldsymbol{\Psi}_{\mathscr{J}, q=2, n, m}(k \boldsymbol{r})=\frac{i^{1-n}}{4 \pi} \int d \Omega_{k} e^{i \boldsymbol{k} \cdot \boldsymbol{r}} \boldsymbol{Z}_{n, m}(\widehat{\boldsymbol{k}}) \tag{6.51}
\end{align*}
$$

### 6.4.3 Poisson summation formula

The Poisson summation formula is a crucial mathematical tool for evaluating lattice sums. It allows one to pass from a sum over the real lattice vectors to a sum over the reciprocal lattice vectors. Formally, it can be written:

$$
\begin{equation*}
\sum_{r_{j} \in \Lambda}^{\infty} e^{i \boldsymbol{k} \cdot r_{j}}=\frac{(2 \pi)^{d_{\Lambda}}}{\mathscr{A}_{d_{\Lambda}}} \sum_{p_{g} \in \Lambda^{*}} \delta\left(\boldsymbol{k}-\boldsymbol{p}_{g}\right) \tag{6.52}
\end{equation*}
$$

where $\mathscr{A}_{d_{\Lambda}}$ is the "volume" of the reciprocal lattice cell. Since long and short range interactions can both be strong for lattice problems the Poisson summation formula often does not directly accelerate the lattice sum, but it nevertheless proves invaluable for a number of useful that we will derive in the rest of this chapter.

For the 1-D sum in eq. 6.90, this can be written:

$$
\begin{equation*}
\sum_{j=-\infty}^{\infty} e^{i(\boldsymbol{k}+\boldsymbol{\beta}) \cdot(\hat{z} j a)}=\frac{2 \pi}{a} \sum_{g=-\infty}^{\infty} \delta\left(k_{z}+\beta_{z}-\frac{2 \pi}{a} g\right) \tag{6.53}
\end{equation*}
$$

We then write this relation in a dimensionless form:

$$
\begin{equation*}
\sum_{j=-\infty}^{\infty} e^{i(k+\boldsymbol{\beta}) \cdot(\widehat{z} j a)}=\frac{2 \pi}{k a i^{n}} \sum_{g=-\infty}^{\infty} \delta\left(\frac{k_{z}}{k}+\frac{\beta_{z}}{k}-g \frac{2 \pi}{k a}\right) \tag{6.54}
\end{equation*}
$$

### 6.4.4 Integral expressions for outgoing partial waves

The Weyl identity expressed the Hankel function of order 0 as an integral of plane waves:

$$
\begin{align*}
h_{0}(k r) & =\frac{1}{2 \pi k} \int_{-\infty}^{\infty} \int_{-\infty} d k_{x} d k_{y} \frac{\exp ( \pm i \boldsymbol{k} \cdot \boldsymbol{r})}{k_{z}} \\
& =\frac{1}{2 \pi k} \int_{-\infty}^{\infty} \int_{-\infty} d k_{x} d k_{y} \frac{\exp \left[ \pm i\left(k_{x} x+i k_{y} y+i k_{z} z\right)\right]}{k_{z}} \quad z \gtrless 0 \tag{6.55}
\end{align*}
$$

where the plus sign is taken for $z>0$ and the minus sign is used when $z<0$. The $k_{z}$ component is fixed by the constraint that $k_{x}^{2}+k_{y}^{2}+k_{z}^{2}=k^{2}$, namely $k_{z}=\sqrt{k^{2}-k_{x}^{2}-k_{y}^{2}}$. It is interesting to remark that the spherical Bessel function is a superposition of plane waves that are constrained to satisfy $\|\boldsymbol{k}\|=k$. Since the reciprocal space integration in eq. 6.55) is carried out in the $x O y$ plane, it is convenient to define a specific symbol for the wavevector in the $x O y$ plane, $\boldsymbol{K}=k_{x} \widehat{\boldsymbol{x}}+k_{y} \widehat{y}$, and the full wavevector is then, $\boldsymbol{k}=\boldsymbol{K}+k_{z} \widehat{\boldsymbol{z}}$. It is also convenient to define dimensionless or normalized quantities:

$$
\begin{equation*}
\bar{K} \equiv|\boldsymbol{K}| / k \quad \gamma_{z} \equiv k_{z} / k=\sqrt{1-\bar{K}^{2}} \tag{6.56}
\end{equation*}
$$

If we take the position vector $\boldsymbol{r}$ in eq. (6.55) to lie along the $z$ axis, $\boldsymbol{r}=r \widehat{\boldsymbol{z}}$, then we can integrate over the azimuthal angle to obtain a single integral expression for Hankel functions that can be used in lattice sums:

$$
\begin{equation*}
h_{0}(k r)=\int_{0}^{\infty} d \overline{K K} \frac{\exp \left[ \pm i \gamma_{z} k r\right]}{\gamma_{z}} \quad z \gtrless 0 \tag{6.57}
\end{equation*}
$$

Wittmann pointed out that the above Weyl identity of eq.(6.55) can be generalized to all partial waves of the Hankel function type[33]:

$$
\begin{align*}
\Psi_{\mathscr{H}, n, m} & \equiv h_{n}(k r) Y_{n, m}(\hat{\boldsymbol{r}}) \\
& =\frac{i^{-n}}{2 \pi} \iint_{-\infty}^{\infty} d k_{x} d k_{y}\left(\bar{k}_{x}+i \bar{k}_{y}\right)^{m} \widetilde{P}_{n}^{m}\left(\gamma_{z}\right) \frac{\exp \left( \pm i\left(k_{x} x+k_{y} y+k_{z} z\right)\right)}{\gamma_{z}}
\end{align*}
$$

If we take $\boldsymbol{r}$ again to lie along the $z$ axis, we find an integral expression for spherical Hankel functions:

$$
\begin{equation*}
h_{n}(k r)=i^{-n} \int_{0}^{\infty} d \overline{K K} P_{n}\left(\gamma_{z}\right) \frac{\exp \left[i \gamma_{z} k r\right]}{\gamma_{z}} . \tag{6.59}
\end{equation*}
$$

The integral relation of eq. (6.58) can also be extended to the outgoing vector partial waves:

$$
\begin{align*}
& \boldsymbol{\Psi}_{\mathscr{H}, q=1, n, m}(k \boldsymbol{r})=\frac{i^{-n}}{2 \pi} \int_{-\infty}^{\infty} \int_{x} d k_{x} d k_{y} \frac{\exp \left( \pm i\left(k_{x} x+k_{y} y+k_{z} z\right)\right)}{\gamma_{z}} \boldsymbol{X}_{n, m}(\widehat{\boldsymbol{k}}) \\
& \boldsymbol{\Psi}_{\mathscr{H}, q=2, n, m}(k \boldsymbol{r})=\frac{i^{1-n}}{2 \pi} \iint_{-\infty}^{\infty} d k_{x} d k_{y} \frac{\exp \left( \pm i\left(k_{x} x+k_{y} y+k_{z} z\right)\right)}{\gamma_{z}} \boldsymbol{Z}_{n, m}(\widehat{\boldsymbol{k}}) \tag{6.60}
\end{align*}
$$

The Poisson sum rule allows one to express quasi-periodic 2D lattice sum in terms of 2D reciprocal lattice vectors. For the scalar partial waves, one has:

$$
\begin{equation*}
\sum_{r_{j} \in \Lambda} \exp \left(i \boldsymbol{\beta} \cdot \boldsymbol{r}_{j}\right) \Psi_{\mathscr{H}, n, m}\left(k r_{\boldsymbol{j}}\right)=\sum_{\boldsymbol{p}_{\boldsymbol{g}} \in \Lambda^{*}} \frac{2 \pi i^{-n}}{k k_{\boldsymbol{g}, z^{\prime}}^{+} \mathscr{\mathscr { C }}_{2}} Y_{n, m}\left(\widehat{\boldsymbol{k}}_{\boldsymbol{g}}^{ \pm}\right) \exp \left(i \boldsymbol{k}_{\boldsymbol{g}}^{ \pm} \cdot \boldsymbol{r}\right) \quad z \gtrless 0 \tag{6.61}
\end{equation*}
$$

while for the vector partial waves,

$$
\begin{align*}
& \sum_{\boldsymbol{r}_{\boldsymbol{j}} \in \Lambda} \exp \left(i \boldsymbol{\beta} \cdot \boldsymbol{r}_{\boldsymbol{j}}\right) \boldsymbol{\Psi}_{\mathscr{H}, 1, n, m}\left(k \boldsymbol{r}_{\boldsymbol{j}}\right)=\sum_{\boldsymbol{p}_{\boldsymbol{g}} \in \Lambda^{*}} \frac{2 \pi i^{-n}}{k k_{\boldsymbol{g}, z^{-}}^{+} \mathscr{\mathscr { L }}_{2}} \boldsymbol{X}_{n, m}\left(\widehat{\boldsymbol{k}}_{\boldsymbol{g}}^{ \pm}\right) \exp \left(i \boldsymbol{k}_{\boldsymbol{g}}^{ \pm} \cdot \boldsymbol{r}\right) \\
& \sum_{\boldsymbol{r}_{\boldsymbol{j}} \in \Lambda} \exp \left(i \boldsymbol{\beta} \cdot \boldsymbol{r}_{\boldsymbol{j}}\right) \boldsymbol{\Psi}_{\mathscr{H}, 2, n, m}\left(k \boldsymbol{r}_{\boldsymbol{j}}\right)=\sum_{\boldsymbol{p}_{\boldsymbol{g}} \in \Lambda^{*}} \frac{2 \pi i^{1-n}}{k k_{\boldsymbol{g}, z^{\prime}}^{+} \mathscr{\mathscr { A }}_{2}} \boldsymbol{Z}_{n, m}\left(\widehat{\boldsymbol{k}}_{\boldsymbol{g}}^{ \pm}\right) \exp \left(i \boldsymbol{k}_{\boldsymbol{g}}^{ \pm} \cdot \boldsymbol{r}\right) \quad z \gtrless 0 \tag{6.62}
\end{align*}
$$

In the partial wave lattice sums of eqs. 6.61 and 6.62 , the wavevector $\boldsymbol{k}_{\boldsymbol{g}}^{ \pm}$is given by:

$$
\begin{equation*}
\boldsymbol{k}_{g}^{ \pm} \equiv\left(\boldsymbol{\beta}_{\|}+\boldsymbol{p}_{g} \pm \widehat{z} \sqrt{k^{2}-\left(\beta_{\|}+p_{g}\right)^{2}}\right), \tag{6.63}
\end{equation*}
$$

and $k_{g_{z}}^{+}$is its $z$ component:

$$
\begin{equation*}
k_{\boldsymbol{g}, z}^{+} \equiv \boldsymbol{k}_{\boldsymbol{g}}^{ \pm} \cdot \widehat{z}=\sqrt{k^{2}-\left(\boldsymbol{\beta}_{\|}+\boldsymbol{p}_{\boldsymbol{g}}\right)^{2}} \tag{6.64}
\end{equation*}
$$

One remarks that for a real Bloch vector $\boldsymbol{\beta}_{\|}$, the wavevector $\boldsymbol{k}_{g}^{ \pm}$is real i.e. propagative in nature only for those lattice vectors for which

$$
\begin{equation*}
k>\left\|\boldsymbol{\beta}_{\|}+\boldsymbol{p}_{\boldsymbol{g}}\right\| \tag{6.65}
\end{equation*}
$$

### 6.4.5 Partial wave rotation

Let us consider once again a row matrix, $\Psi^{t}{ }_{\mathscr{G}, \mathscr{Y}, \mathscr{H}}$, composed of one of the types of partial waves in eq.(6.5). The fact that the choice of the orientation of the coordinate system is arbitrary imposes transformation relations amongst the partial waves. Let us consider a position $M$ given by the vector $\boldsymbol{r}$ in our chosen coordinate system. We next consider another coordinate system with the same origin, but rotated by the 3 Euler angles, $\alpha, \beta$, and $\gamma$ in which the same point $M$ is now designated by a vector $\boldsymbol{r}^{\prime}\left(\right.$ n.b. $\left.\left|\boldsymbol{r}^{\prime}\right|=|\boldsymbol{r}|=r\right)$. The linear relationship between the row matrix in these 2 coordinate systems is then:

$$
\begin{equation*}
\Psi^{t}(\boldsymbol{r})=\Psi^{t}\left(\boldsymbol{r}^{\prime}\right) \mathscr{D}(\alpha, \beta, \gamma) \tag{6.66}
\end{equation*}
$$

If the rotated coordinate systems is taken such that $\boldsymbol{r}^{\prime}$ lies along the $z$ axis in the rotated coordinate n this relation takes the form:

$$
\begin{equation*}
\Psi^{t}(\boldsymbol{r})=\Psi^{t}(\widehat{r}) \mathscr{D}(\phi, \theta, 0) . \tag{6.67}
\end{equation*}
$$

In component form this reads for Hankel function sums:

$$
\begin{align*}
h_{n}(k r) Y_{n, m}(\widehat{\boldsymbol{r}}) & =h_{n}(k r) Y_{n 0}(0,0) \mathscr{D}_{n, 0 ; n, m}\left(\phi_{\widehat{r}}, \theta_{\widehat{r}}, 0\right) \\
& =\sqrt{\frac{2 n+1}{4 \pi}} h_{n}(k r) \mathscr{D}_{n, 0 ; n, m}\left(\phi_{\widehat{r}}, \theta_{\widehat{r}}, 0\right), \tag{6.68}
\end{align*}
$$

where we used eq. 6.75) for an expression of the $Y_{n 0}(0,0)$.

### 6.5 Numerical Examples

This section will be expanded considerably in the second edition.

### 6.5.1 Far and near field response from gratings

As discussed in section 6.3, once the lattice sums have been determined for all the $\Omega_{\boldsymbol{\beta}}$ matrix elements, and the lattice T-matrix of eq. 6.29) obtained, one has ready access to both the far and near field response of the system. However, the quasi-periodic lattice for all the multipole orders must be calculated anew whenever one looks for response to a different quasi-periodicity vector, $\boldsymbol{\beta}$ or wavenumber $k$ (i.e. frequency).

### 6.5.2 Modes for particulate chains

There is considerable interest in calculating and characterizing the 'propagating' modes of periodic chains, gratings, and finite stacks of particulate gratings. A major difficulty is that these systems are open, so propagating 'modes' are necessarily lossy which entails that they don't exist for real values of frequency and quasi-periodicity vector $\boldsymbol{\beta}$. This situation was in fact already present in the T-matrices of the individual particles which have do have poles at any real frequencies, but one can find poles at complex frequencies. For planar surfaces, Greffet has recently argued[2] that the Leaky-modes can be described by letting either frequency or
wavevector be described by a complex number. This idea has recently been employed by several authors for calculating modes in infinite particulate chains where the component of $\boldsymbol{\beta}$ along the chain axis is allowed to be a complex number.[23, 4, 6, 11, 12] There are still some technical issues under discussion considering modes in gratings due to ambiguities in the best manner to define a complex propagation vector in the grating plane.[9]. This issue will be treated in more depth in the second edition.

Typically, one has looked for propagating modes in particulate arrays of sub-wavelength particles metallic particles. There is however an increased in interest in high index dielectrics. Due to the complexity of the full multipole approach, most works searches for modes in the complex plane have adopted what amounts to be a electric dipole approximation to eq.(6.29]. [23] We have recently argued that electric dipole is insufficient in the presence of strong interactions that are provoked by resonances.[23] These results and conclusions are reviewed here.

We adopt the same parameters for a plasmonic chain as Conforti and Guasoni.[6] Namely, we consider an infinite chain of identical 50 nm diameter silver particles separated by $d=75 \mathrm{~nm}$ (center-to-center). The system is immersed in a non-magnetic medium with relative permittivity $\varepsilon=2.25$ ( $n=1.5$ ).

The figures are plotted with normalized frequencies and wave-vectors:

$$
\begin{equation*}
\bar{\omega} \equiv \frac{\omega d}{2 \pi c}=\frac{d}{\lambda_{\mathrm{v}}} \quad \bar{\beta} \equiv \frac{\beta d}{2 \pi} \tag{6.69}
\end{equation*}
$$

where $\lambda_{\mathrm{v}}$ is the vacuum wavelength. The light line for these parameters is given by $\bar{\omega}=\frac{\bar{\beta}}{n_{\text {med }}}$. The dispersion relations of the principal propagating modes calculated in the electric dipole approximation are plotted in figure 6.2 (dashed curves). They are then compared with fully converged $n_{\max }=10$ calculations of these dispersion relations (solid line) in this same figure by solving eq. (6.34). The imaginary part of the dispersion relations for dipolar and converged multipole calculations are given in figure 6.2a).



Figure 6.2: Real and imaginary parts of the dispersion relations in the dipole approximation (dashed curves), and fully converged multipole calculations with $n_{\max }=10$ (full lines). The Longitudinal mode with positive imaginary part is in cyan(gray) the "T1" mode with positive imaginary part is in blue(black line). The "T2" transverse mode with negative imaginary part is in orange(gray). reproduced with permission : http://dx.doi.org/10.1364/JOSAB.29.001012

Figures 6.2 p ) and 6.2 b ) merit some commentary. It is immediately clear that the dipole approximation provides a moderately accurate prediction of dispersion relations only over a narrow range of frequencies for which the imaginary part of the propagating wavevector is
rather small, and the real part is near the light line. One should also recall that symmetry dictates that if a given value of $\beta$ corresponds to mode at a given frequency, then by symmetry, $-\beta$ is also a solution to these equations. For the sake of clarity, these symmetric modes are not presented in these figures.

Like Conforti and Guasoni[6], we find a transverse mode, labeled "T2" whose imaginary part of $\beta$ is opposite in sign with the real part of $\beta$. It may prove physically relevant to think of this T2 mode as a backscattering mode, or to interpret this in terms of negative effective index. It is interesting to remark that the T2 mode tends toward the edge of the Brillouin zone at low frequencies. It has recently been argued that these

Our dipole approximation predictions for the longitudinal mode are quite similar to that of ref.[6] wherein the dipole prediction is that the mode "folds back" before reaching the edge of the Brillouin zone. The full multipole calculations on the other hand predict that the longitudinal mode goes to the edge of the Brillouin zone, and that the "fold back" only occurs after it has gone "beyond" the edge of the Brillouin zone. In our calculations, the "T1" mode is quite close to the light line, and henceforth rather poorly confined by the plasmon chain so its importance in applications seems limited. In our calculations, the dipole approximation for the "T1" mode is quite similar to the multipole solution except that we only found that the full multipole solution predicted both extremities of the T 1 mode to lie on the light line.


Figure 6.3: Normalized extinction is a solid blue (line) and scattering cross section given by a dashed green line of a silver monomer in terms of frequency $(a=25 n m)$. reproduced with permission : http://dx.doi.org/10.1364/JOSAB.29.001012

Due to the system design (sub-wavelength resonance particles) one expected to find significant guiding of modes appear only in the frequency domains where the scattering cross section of the individual particles is non-negligible. To illustrate this point, we plot the extinction and scattering cross section for an individual particle in the chain in figure 6.3 . We remark in particular that near individual particle resonance maxima, all the guided modes of figure 6.2 lie near the light line, and as seen in figure 6.2p) it is here also that their imaginary parts are smallest. Furthermore, with the exception of the 'backscattering' mode T2, all guided modes apparently cease to exist when one moves sufficiently far away from the scattering resonance frequency.

The reader has probably remarked some strange behavior of the modes in the electric dipole approximation at high frequencies. For instance, at around $\widetilde{\omega}=0.225$ a "kink" appears in the longitudinal mode, and a spurious T2 solution emerges from the light line. We carried
out mode calculations with various multipole cutoffs and found that such kinks and spurious solutions were relatively commonplace (at high or low frequencies) when low numbers of multipoles are used in the simulations and such behavior disappears when higher multipole orders are used. It is also worth remarking that for high order simulations, the $\operatorname{Re}[\beta]$ of the modes terminate at either the light line, or the edge of the Brillouin zone, but modes can terminate at undiscriminating positions in $\beta$ space when calculations are carried out at low order.


Figure 6.4: Positive and imaginary parts of the dispersion relations in the $\operatorname{Re}[\beta]>0$ part of the Brillouin zone. Transverse modes with: $\operatorname{Im}[\beta]>0$ modes are solid blue(black) lines, while that with $\operatorname{Im}[\beta]<0$ is given by a dashed blue(black) line. Longitudinal modes with $\operatorname{Im}[\beta]>0$ are solid cyan(gray) lines, while those with $\operatorname{Im}[\beta]<0$ is a dashed cyan(gray) line. reproduced with permission : http://dx.doi.org/10.1364/JOSAB.29.001012

The mode diagrams of figures 6.2 p ) and 6.2 p ) were somewhat unconventional since they did not display symmetric, $-\beta$, modes, and allowed the dispersion relation of the longitudinal mode to move outside the Brillouin zone. A more conventional representation of the dispersion relations is given in figure 6.4 which includes the symmetric modes, but only displays modes when $\operatorname{Re}[\beta]$ has positive values lying within Brillouin zone (here we display only the converged multipole calculations). Transverse modes with $\operatorname{Im}[\beta]>0$ modes are given by solid blue(black) lines, while transverse modes with $\operatorname{Im}[\beta]<0$ are dashed dashed blue(black) lines. Longitudinal modes with $\operatorname{Im}[\beta]>0$ are given by solid cyan(gray) lines while longitudinal modes with negative $\operatorname{Im}[\beta]$ are in dashed cyan(gray). It is interesting to note that the longitudinal modes extend to the positive edge of the Brillouin zone and that the "fold back" only occurs when $\operatorname{Im}[\beta]$ of the longitudinal mode is negative. One can also remark that transverse modes, T2, with both positive and negative $\operatorname{Im}[\beta]$ exist above the light line, but that their imaginary parts are quite large. Longitudinal modes above the light line also exist at frequencies below the particle resonance maximum, but these modes remain quite close to the light line.

### 6.6 Chain sums

### 6.6.1 Hankel function chain sums

A periodic chain of wave scattering is defined by a lattice vector $\boldsymbol{a}$, such that there is an elementary 'scatterer' at all positions $\boldsymbol{r}_{j}$ i.e.:

$$
\begin{equation*}
\boldsymbol{r}_{j} \equiv j \boldsymbol{a} \quad j \in \mathbb{Z} \quad j=-\infty, \ldots-2,-1,0,1,2, \ldots, \infty . \tag{6.70}
\end{equation*}
$$

A chain sum for a quasi-periodicity vector $\boldsymbol{\beta}$ is defined:

$$
\begin{equation*}
S_{n, m}^{C}(k, a, \beta ; \widehat{\boldsymbol{a}}) \equiv S_{n, m}^{C, \mathscr{H}}(k, a, \beta ; \widehat{\boldsymbol{a}}) \equiv \sum_{\substack{j \neq 0 \\ j \in \mathbb{Z}}} \mathscr{H}_{n, m}(j \boldsymbol{a}) e^{i j a \boldsymbol{\beta} \cdot \widehat{\boldsymbol{a}}} \tag{6.71}
\end{equation*}
$$

One can remark that the chain sum, $S_{n, m}^{C}$, depends on the amplitude of the lattice vector $a=|\boldsymbol{a}|$, and its direction, and the scalar product between $\boldsymbol{a}$ and another vector $\boldsymbol{\beta}$ which we will call the 'incident' or 'quasi-periodicity' vector. The chain sum in fact only depends on the scalar product between $\boldsymbol{\beta}$ and the periodicity vector:

$$
\begin{equation*}
\beta \equiv \boldsymbol{\beta} \cdot \widehat{\boldsymbol{a}} \tag{6.72}
\end{equation*}
$$

One remarks that the direction of $\boldsymbol{a}$ depends on the orientation of the coordinate system. We can take advantage of this fact to define the $z$ axis as the direction of $\boldsymbol{a}$ such that:

$$
\begin{equation*}
\boldsymbol{r}_{j}=j a \widehat{z}, \tag{6.73}
\end{equation*}
$$

but one must keep in mind that the expression for $S_{n, m}^{C}$ is reference frame dependent. In this coordinate system, the chain sum, $S_{n, m}^{C}(k, a, \beta ; \widehat{z})$, takes the form :

$$
\begin{align*}
S_{n, m}^{C}(k, a, \beta ; \widehat{\boldsymbol{z}}) & \equiv \sum_{\substack{j \neq 0 \\
j \in \mathbb{Z}}} \mathscr{H}_{n, m}(j a \widehat{\boldsymbol{z}}) e^{i \boldsymbol{\beta} \cdot \widehat{\mathbf{z}} j a} \\
& =\sum_{\substack{j \neq 0 \\
j \in \mathbb{Z}}} h_{n}(k|j| a) Y_{n, m}\left(\frac{j}{|j|} \widehat{\boldsymbol{z}}\right) e^{i \beta a j} \\
& =\delta_{m, 0} \lambda_{n, 0} \sum_{j=1}^{\infty} h_{n}(j k a)\left[e^{i j \beta a}+(-1)^{n} e^{-i j \beta a}\right], \tag{6.74}
\end{align*}
$$

where we used the fact that only the $m=0$ scalar spherical harmonics are non null at $\theta=0, \pi$ :

$$
\begin{equation*}
Y_{n, m}(0,0)=\delta_{m, 0} \lambda_{n, 0}=\sqrt{\frac{2 n+1}{4 \pi}} \quad Y_{n, 0}(\pi, 0)=\delta_{m, 0}(-)^{n} Y_{n, m}(0,0) \tag{6.75}
\end{equation*}
$$

The analytical expressions for the first few Hankel function are:

$$
\begin{align*}
& h_{0}(x)=-\frac{i}{x} e^{i x} \\
& h_{1}(x)=e^{i x}\left(\frac{-1}{x}-\frac{i}{x^{2}}\right) \\
& h_{2}(x)=e^{i x}\left(\frac{i}{x}-\frac{3}{x^{2}}-\frac{3 i}{x^{3}}\right), \tag{6.76}
\end{align*}
$$

which are readily obtained from the general analytic expression for Hankel functions of arbitrary order:

$$
\begin{equation*}
h_{n}(x)=(-i)^{n+1} \sum_{s=0}^{n} \frac{i^{s}}{2^{s} s!} \frac{(n+s)!}{(n-s)!} \frac{e^{i x}}{x^{s+1}} . \tag{6.77}
\end{equation*}
$$

### 6.6.2 Integral technique for Hankel lattice sums

Generalizing the Weyl integral to produce an integral expression for spherical Hankel functions gives us the integral:

$$
\begin{equation*}
h_{n}(k r)=\frac{1}{i^{n}} \int_{0}^{\infty} d \bar{K} \bar{K} P_{n}\left(k_{z} / k\right) \frac{\exp \left(i \gamma_{z} k r\right)}{\gamma_{z}}, \tag{6.78}
\end{equation*}
$$

where we recall that $k_{z}=\sqrt{k^{2}-K^{2}}$ and $K=\sqrt{k_{x}^{2}-k_{y}^{2}}$ is the wavevector component in the $x O y$ plane. In the last line we defined the dimensionless quantities, $\bar{K} \equiv K / k$ and $\gamma_{z}=\sqrt{1-\bar{K}^{2}}$ so that this integral could be calculated using dimensionless variables.

Now if we actually try to evaluate this integral, we will have numerical problems when we go past the point where $k_{z}=0$. Since the singularity coming from the $k_{z}$ denominator lies just above the real axis, we can analytically continue the integration into the fourth quadrant of the complex plane. Any angle will do as long as the resulting line integral is sufficiently far from the positive real axis or the negative imaginary axis. We will generally take an angle of $45^{\circ}$ as a reasonable compromise. We will find that the integrand will decrease exponentially for large $|K|$ in the complex plane so that we don't have much problem with the integral extending to infinity.

Thanks to the integral expression for Hankel functions of eq. 6.78), we are now ready to treat an infinite chain sum for a chain oriented along the $z$ axis:

$$
\begin{align*}
S_{n, m}^{C}(k, a, \beta ; \widehat{\boldsymbol{z}})= & \sum_{j \in \mathbb{Z}^{*}} \exp (i \beta a j) h_{n}(k|j| a) Y_{n, m}\left(\frac{j}{|j|} \widehat{\boldsymbol{z}}\right) \\
= & \sum_{j=1}^{\infty} \exp (i \beta a j) h_{n}(k a j) Y_{n, m}(0,0) \\
& +\sum_{j=1}^{\infty} \exp (-i \beta a j) h_{n}(k a j) Y_{n, m}(\pi, 0) \\
= & \delta_{m, 0} \sqrt{\frac{2 n+1}{4 \pi}}\left[\sum_{j=1}^{\infty} \exp (i \beta a j) h_{n}(k a j)+(-)^{n} \sum_{j=1}^{\infty} \exp (-i \beta a j) h_{n}(k a j)\right] \tag{6.79}
\end{align*}
$$

where we used:

$$
\begin{equation*}
Y_{n, m}(0,0)=\delta_{m, 0} \sqrt{\frac{2 n+1}{4 \pi}} P_{n}(1), \quad Y_{n, m}(\pi, 0)=\delta_{m, 0} \sqrt{\frac{2 n+1}{4 \pi}} P_{n}(-1) \tag{6.80}
\end{equation*}
$$

and

$$
\begin{equation*}
P_{n}(1)=1, \quad P_{n}(-1)=(-1)^{n} . \tag{6.81}
\end{equation*}
$$

Using the integral relation of eq. 6.78), we have:

$$
\begin{align*}
S_{n, m}^{C}(k, a, \beta ; \widehat{\boldsymbol{z}}) & =\delta_{m, 0} \sqrt{\frac{2 n+1}{4 \pi}} \frac{1}{i^{n} k} \int_{0}^{\infty} d K K \frac{P_{n}\left(k_{z} / k\right)}{k_{z}} \\
& \times\left[\sum_{j=1}^{\infty} \exp \left[i\left(k_{z}+\beta\right) j d\right]+(-)^{n} \sum_{j=1}^{\infty} \exp \left[i\left(k_{z}-\beta\right) j a\right]\right] . \tag{6.82}
\end{align*}
$$

We have finally an integral expression for the chain sums:

$$
\begin{align*}
S_{n, m}^{C}(k, a, \beta ; \widehat{z}) & =\sum_{j \in \mathbb{Z}^{*}} \exp (i \beta j d) h_{n}\left(k r_{j}\right) Y_{n, m}\left(\widehat{\boldsymbol{r}}_{j}\right) \\
& =\delta_{m, 0} \sqrt{\frac{2 n+1}{4 \pi}} \frac{1}{i^{n}} \int_{0}^{\infty} d \bar{K} \bar{K} \frac{P_{n}\left(\gamma_{z}\right)}{\gamma_{z}} \\
& \times\left[\frac{1}{\exp \left[-i\left(\bar{k}_{z}+\bar{\beta}\right) k a\right]-1}+\frac{(-)^{n}}{\exp \left[-i\left(\bar{k}_{z}-\bar{\beta}\right) k a\right]-1}\right] \tag{6.83}
\end{align*}
$$

### 6.6.3 Polylog approach to Hankel chain sums

Inspection of eqs. 6.79) and (6.77) shows that all terms in the chain sum can be expressed in terms of polylogarithm functions which are defined by[1]:

$$
\begin{equation*}
\mathrm{Li}_{n}(z)=\sum_{j=1}^{\infty} \frac{z^{j}}{j^{n}} \tag{6.84}
\end{equation*}
$$

The chain sum expressed in terms of polylogarithms is then:

$$
\begin{align*}
S_{n, m}^{C}(k, a, \beta ; \widehat{z})=\delta_{m, 0} \sqrt{\frac{2 n+1}{4 \pi}} & \sum_{s=0}^{n}\left[\left((-i)^{n+1} \frac{i^{s}}{2^{s} s!} \frac{(n+s)!}{(n-s)!}\right)\right. \\
& \left.\times \frac{\left(\operatorname{Li}_{s+1} \exp [i(k+\beta) a]+(-)^{n} \operatorname{Li}_{s+1} \exp [i(k-\beta) a]\right)}{(k a)^{s+1}}\right] . \tag{6.85}
\end{align*}
$$

This was the chain sum for outgoing Hankel functions, but we will also sometimes be interested in incoming Hankel functions, or Bessel functions of the fourth kind. These are expressed:

$$
\begin{equation*}
h_{n}^{(4)}(x) \equiv h_{n}^{-}(x)=j_{n}(x)-i y_{n}(x), \tag{6.86}
\end{equation*}
$$

and their chain sums are:

$$
\begin{align*}
S_{n, m}^{C}(k, a, \beta ; \widehat{\boldsymbol{z}})=\delta_{m, 0} \sqrt{\frac{2 n+1}{4 \pi}} & \sum_{s=0}^{n}\left[\left((-i)^{n+1} \frac{i^{s}}{2^{s} s!} \frac{(n+s)!}{(n-s)!}\right)\right. \\
& \left.\times \frac{\left(\operatorname{Li}_{s+1} \exp [-i(k-\beta) a]+(-)^{n} \operatorname{Li}_{s+1} \exp [-i(k+\beta) a]\right)}{(k a)^{s+1}}\right] \tag{6.87}
\end{align*}
$$

### 6.6.4 Bessel function chain sums

Although fully analytic expressions for Hankel function chain and lattice sums do not seem to exist currently, the Bessel functions lattice and chain sums do have analytic expressions. These Bessel function sums are useful in their own right for certain applications:

$$
\begin{align*}
S_{n}^{C, \mathscr{g}}(k, a, \beta ; \widehat{z}) & =\sum_{j=-\infty, j \neq 0}^{\infty} Y_{n, m}\left(\widehat{\boldsymbol{r}}_{j}\right) j_{n}(j k a) e^{i j \beta a} \\
& =\sum_{j=-\infty}^{\infty} Y_{n, m}\left(\widehat{\boldsymbol{r}}_{j}\right) j_{n}(j k a) e^{i j \beta a}-\sum_{j=-\infty}^{\infty} Y_{0,0}\left(\widehat{\boldsymbol{r}}_{0}\right) j_{0}(j k a) \\
& =\sum_{j=-\infty}^{\infty} Y_{n, m}\left(\widehat{\boldsymbol{r}}_{j}\right) j_{n}(j k a) e^{i j \beta a}-\frac{1}{\sqrt{4 \pi}} \delta_{n, 0} \tag{6.88}
\end{align*}
$$

Using the integral expression for $Y_{n, m}\left(\widehat{\boldsymbol{r}}_{j}\right) j_{n}(j k a)$ as an integral over the directions a wavenumber $\widehat{\boldsymbol{k}}$ as derived in eq. 6.50 allows us to write:

$$
\begin{equation*}
j_{n}\left(k R_{j}\right) Y_{n, m}\left(\widehat{\boldsymbol{r}}_{j}\right) e^{i \boldsymbol{\beta} \cdot(\widehat{z} j a)}=\frac{1}{4 \pi i^{n}} \int Y_{n, m}(\widehat{\boldsymbol{k}}) e^{i \boldsymbol{\beta} \cdot \widehat{z} j a} e^{i \boldsymbol{k} \cdot \bar{z} j a} d \Omega_{k} \tag{6.89}
\end{equation*}
$$

The lattice sum of the Bessel type then can be written:

$$
\begin{equation*}
S_{n}^{C, \mathscr{\mathscr { V }}}(k, a, \beta ; \widehat{z})=\frac{1}{4 \pi i^{n}} \int Y_{n, m}(\widehat{k})\left[\sum_{j=-\infty}^{\infty} e^{i(\boldsymbol{k}+\boldsymbol{\beta}) \cdot(\widehat{z} j a)}\right] d \Omega_{k}-\frac{1}{\sqrt{4 \pi}} \delta_{n, 0} \tag{6.90}
\end{equation*}
$$

At this point, one invokes the Poisson summation formula which can be written formally as:

$$
\begin{equation*}
\sum_{r_{j} \in \Lambda}^{\infty} e^{i \boldsymbol{k} \cdot \boldsymbol{r}_{j}}=\frac{(2 \pi)^{d_{\Lambda}}}{\mathscr{A}} \sum_{\boldsymbol{p}_{g} \in \Lambda^{*}} \delta\left(\boldsymbol{k}-\boldsymbol{p}_{g}\right) \tag{6.91}
\end{equation*}
$$

where $\mathscr{A}$ is the "volume" of the reciprocal cell. For the 1-D sum in eq. 6.90, this can be written:

$$
\begin{equation*}
\sum_{j=-\infty}^{\infty} e^{i(\boldsymbol{k}+\boldsymbol{\beta}) \cdot(\hat{z} j a)}=\frac{2 \pi}{a} \sum_{g=-\infty}^{\infty} \delta\left(k_{z}+\beta_{z}-\frac{2 \pi}{a} g\right) \tag{6.92}
\end{equation*}
$$

We then write this relation in a dimensionless form:

$$
\begin{equation*}
\sum_{j=-\infty}^{\infty} e^{i(k+\boldsymbol{\beta}) \cdot(\bar{z} j a)}=\frac{2 \pi}{k a i^{n}} \sum_{g=-\infty}^{\infty} \delta\left(\frac{k_{z}}{k}+\frac{\beta_{z}}{k}-g \frac{2 \pi}{k a}\right) \tag{6.93}
\end{equation*}
$$

Putting this relation into the $\widehat{k}$ integral of eq. 6.90 , we then obtain a finite sum expression for $S_{n}^{C, \mathscr{J}}$ :

$$
\begin{equation*}
S_{n}^{C, \mathscr{I}}(k, a, \beta ; \widehat{z})=-\frac{1}{\sqrt{4 \pi}} \delta_{n, 0}+\frac{\pi i^{n}}{k a} \sum_{g=g_{\min }}^{g_{\max }} Y_{n 0}\left(\cos \beta_{z, q}\right) \tag{6.94}
\end{equation*}
$$

where since $-1<k_{z} / k<1$ we only sum over those values of $g$ for which

$$
\begin{equation*}
-1<\mathfrak{R}\left[\frac{\beta_{z} a+2 \pi g}{k a}\right]<1 . \tag{6.95}
\end{equation*}
$$

The values $g_{\text {min }}$ and $g_{\text {max }}$ are:

$$
\begin{equation*}
g_{\min }=\left(-\frac{\beta_{z} a-k a}{2 \pi}\right)+1 \quad g_{\max }=-\frac{\beta_{z} a+k a}{2 \pi} \tag{6.96}
\end{equation*}
$$

The angle $\cos \beta_{z, g}$ in eq. 6.94 is given by:

$$
\begin{equation*}
\cos \beta_{z, g} \equiv \frac{\beta_{z} a+2 \pi g}{k a}=\mathfrak{R}\left[\beta_{z} / k\right]+i \mathfrak{I}\left[\beta_{z} / k\right]+g \frac{2 \pi}{k a}, \tag{6.97}
\end{equation*}
$$

where we used the parity relation:

$$
\begin{equation*}
Y_{n, m}(-\widehat{\boldsymbol{r}})=(-1)^{n} Y_{n, m}(-\widehat{\boldsymbol{r}}) . \tag{6.98}
\end{equation*}
$$

We recall that the sin and cosines for a complex angle, $\theta_{k}=\theta^{\prime}+\theta^{\prime \prime}$, are given by:

$$
\begin{align*}
\cos \theta_{k} & =\frac{e^{i \theta^{\prime}} e^{-\theta^{\prime \prime}}+e^{-i \theta^{\prime}} e^{\theta^{\prime \prime}}}{2} \\
& =\cos \theta^{\prime} \cosh \theta^{\prime \prime}-i \sin \theta^{\prime} \sinh \theta^{\prime \prime} \tag{6.99}
\end{align*}
$$

and

$$
\begin{align*}
\sin \theta_{k} & =\frac{e^{i \theta^{\prime}} e^{-\theta^{\prime \prime}}-e^{-i \theta^{\prime}} e^{\theta^{\prime \prime}}}{2 i} \\
& =\sin \theta^{\prime} \cosh \theta^{\prime \prime}+i \cos \theta^{\prime} \sinh \theta^{\prime \prime} \tag{6.100}
\end{align*}
$$

### 6.6.5 Chain sum rotation

The chain sums expressions given in eqs.(6.83), (6.87), and (6.94) all took advantage of the facilities presented by orienting the chain of particles along the $z$ axis. When performing lattice reduction techniques, it is necessary to have chain sums in other orientations. The chain sum is obtained by applying:

$$
\begin{equation*}
Y_{n, m}(\widehat{\boldsymbol{r}})=Y_{n, 0}(\widehat{z}) \mathscr{D}_{0, m}^{(n)}\left(\theta_{\widehat{r}}, \phi_{\widehat{r}}\right)=\sqrt{\frac{2 n+1}{4 \pi}} \mathscr{D}_{0, m}^{(n)}\left(\theta_{\widehat{\boldsymbol{r}}}, \phi_{\widehat{r}}\right), \tag{6.101}
\end{equation*}
$$

which just translates the relation derived in Edmonds (eq.(4.1.25) page 59) that:

$$
\begin{equation*}
\mathscr{D}_{0, m}^{(n)}(\theta, \phi)=\sqrt{\frac{4 \pi}{2 n+1}} Y_{n, m}(\theta, \phi) . \tag{6.102}
\end{equation*}
$$

Thus is trivial to write chain sums of any type $(\mathscr{J}, \mathscr{H}, \mathscr{Y})$ in an arbitrary orientation, $\widehat{\boldsymbol{r}}$, in terms of the chain sum in the direction $\widehat{\boldsymbol{z}}$ by the simple relation:

$$
\begin{equation*}
S_{n, m}^{C}(\beta ; \widehat{\boldsymbol{r}})=S_{n}^{C}(\beta ; \widehat{\boldsymbol{z}}) \sqrt{\frac{4 \pi}{2 n+1}} Y_{n, m}\left(\theta_{\widehat{\boldsymbol{r}}}, \phi_{\widehat{\boldsymbol{r}}}\right) \tag{6.103}
\end{equation*}
$$

An orientation along the $x$ axis is for example:

$$
\begin{equation*}
S_{n, m}^{C}(\beta ; \widehat{x})=S_{n}^{C}(\beta ; \widehat{z}) \sqrt{\frac{4 \pi}{2 n+1}} Y_{n, m}\left(\frac{\pi}{2}, 0\right), \tag{6.104}
\end{equation*}
$$

is useful when carrying out a monolayer sum in the next section. An orientation along the $y$ axis is for example:

$$
\begin{equation*}
S_{n, m}^{C}(\beta ; \widehat{\boldsymbol{y}})=S_{n}^{C}(\beta ; \widehat{\boldsymbol{z}}) \sqrt{\frac{4 \pi}{2 n+1}} Y_{n, m}\left(\frac{\pi}{2}, \frac{\pi}{2}\right) . \tag{6.105}
\end{equation*}
$$

### 6.7 2D Grating lattice sums

A 2D periodic media is characterized by two basic lattice vectors $\boldsymbol{a}, \boldsymbol{b}$. Although, we want to describe a system with lattice vectors $\boldsymbol{a}=(a, 0,0)$, and $\boldsymbol{b}=\left(b_{1}, b_{2}, 0\right)$, we are going to work in a rotated coordinate systems in which the lattice will be placed in the $x O y$ plane.

### 6.7.1 Integral technique

For the integral technique, it is useful to adopt a coordinate system where the $\boldsymbol{a}$ lattice vector lies along the $y$ axis. In this coordinate system, the basis vectors are $\boldsymbol{a}=(0, a, 0), \boldsymbol{b}=\left(0, b_{2}, b_{1}\right)$, then lattice sites are given by:

$$
\begin{equation*}
\boldsymbol{r}_{\boldsymbol{j}=\left(j_{a}, j_{b}\right)}=j_{a} \boldsymbol{a}+j_{b} \boldsymbol{b}=\left(0, j_{a} a+j_{b} b_{2}, j_{b} b_{1}\right) . \tag{6.106}
\end{equation*}
$$

In this case, fixing $j_{b}=0$ and summing over $j_{a}$ corresponds to a chain sum along the $y$ axis, and $j_{b} \lessgtr 0$, corresponds to a term in the $z \lessgtr 0$ half plane respectively.

The Mono-Layer (ML) lattice sum, $S_{n, m}^{M L}$, can be written:

$$
\begin{equation*}
S_{n, m}^{M L}=S_{n, m}^{C}\left(\beta_{1} ; \widehat{\boldsymbol{y}}\right)+S_{n, m}^{M L+}+S_{n, m}^{M L-}, \tag{6.107}
\end{equation*}
$$

where $S_{n, m}^{C}\left(\beta_{1} ; \widehat{\boldsymbol{y}}\right)$ is the chain sum along the $y$ axis, and $S_{n, m}^{M L}$ is the sum of all the sites with $z \neq 0$ The lattice sum for all $j_{b}>0(z>0)$ sites can be expressed as an integral:

$$
\begin{align*}
S_{n, m}^{M L+}=- & \frac{(-)^{m}}{i^{n} k a} \sum_{g=-\infty}^{\infty} \int_{0}^{\infty} \frac{d \bar{k}_{x}}{\gamma_{z}} \widetilde{P}_{n}^{(m)}(\gamma)\left[\left(\bar{k}_{g}-i \bar{k}_{x}\right)^{m}+\left(\bar{k}_{g}+i \bar{k}_{x}\right)^{m}\right] \\
& \times \frac{1}{1-\exp \left\{-i\left[k b_{2}\left(\bar{\beta}_{2}-\bar{k}_{g}\right)+k b_{1}\left(\bar{\beta}_{1}+\gamma\right)\right]\right\}} \tag{6.108}
\end{align*}
$$

while the lattice sum for all $j_{b}<0$, can be expressed:

$$
\begin{align*}
S_{n, m}^{M L-} & =-\frac{1}{(-i)^{n} k a} \sum_{g=-\infty}^{\infty} \int_{0}^{\infty} \frac{d \bar{k}_{x}}{\gamma_{z}} \widetilde{P}_{n}^{(m)}(\gamma)\left[\left(\bar{k}_{g}-i \bar{k}_{x}\right)^{m}+\left(\bar{k}_{g}+i \bar{k}_{x}\right)^{m}\right] \\
& \times \frac{1}{1-\exp \left\{-i\left[k b_{2}\left(-\bar{\beta}_{2}+\bar{k}_{g}\right)+k b_{1}\left(-\bar{\beta}_{1}+\gamma\right)\right]\right\}} . \tag{6.109}
\end{align*}
$$

In both of these expressions, we defined $k_{g}$ as the reciprocal lattice vector along the $y$ axis:

$$
\begin{equation*}
k_{g} \equiv \beta_{2}+2 \pi g / a=k \bar{k}_{g}, \tag{6.110}
\end{equation*}
$$

and normalized the components of the quasi-periodic vector $\beta$ as:

$$
\begin{equation*}
\bar{\beta}_{1} \equiv \beta_{1} / k \quad \bar{\beta}_{2} \equiv \beta_{2} / k, \tag{6.111}
\end{equation*}
$$

and $\gamma$ is reciprocal lattice vector along the $z$ axis:

$$
\begin{equation*}
\gamma \equiv \sqrt{1-k_{g}^{2}-k_{x}^{2}} . \tag{6.112}
\end{equation*}
$$

At the end of this calculation, one should keep in mind that lattice sum was carried out in a system where the lattice was in the $y O z$ plane. One can obtain the expression for the lattice sum in the $x O y$ plane by rotating the lattice sums by $90^{\circ}$ around the $y$ axis in a clockwise manner, and then $90^{\circ}$ around the new $y^{\prime}$ axis, and finally $90^{\circ}$ around the new $z^{\prime \prime}$ axis.

### 6.7.2 Modified Bessel function sums

Although the integral technique is rather efficient, one may prefer to obtain do a little more analytic work and obtain the lattice sum in a manner which takes the from of a lattice sum in a 2 D host space. As for the integral it proves convenient to place the lattice in the $y O z$ plane, but this time, one places the $\boldsymbol{a}$ lattice vector along the $z$ axis, so that the lattice vectors can be written:

$$
\boldsymbol{a}=(0,0, a), \quad \text { and } \quad \boldsymbol{b}=\left(0, b_{2}, b_{1}\right) .
$$

The lattice sites in this system can be expressed:

$$
\begin{equation*}
\boldsymbol{r}_{\boldsymbol{j}=\left(j_{a}, j_{b}\right)}=j_{a} \boldsymbol{a}+j_{b} \boldsymbol{b}=\left(0, j_{b} b_{2}, j_{a} a+j_{b} b_{1}\right)=j_{b} b_{2} \widehat{\boldsymbol{y}}+\left(j_{a} a+j_{b} b_{1}\right) \widehat{\boldsymbol{z}} \tag{6.113}
\end{equation*}
$$

The reciprocal lattice is given by:

$$
\begin{equation*}
\boldsymbol{K}_{\boldsymbol{g}=\left(g_{a}, g_{b}\right)}=g_{a} \widetilde{\boldsymbol{a}}+g_{b} \widetilde{\boldsymbol{b}}, \tag{6.114}
\end{equation*}
$$

where the reciprocal lattice vectors are:

$$
\begin{equation*}
\widetilde{\boldsymbol{a}}=\frac{1}{a b_{2}}\left(0,-b_{1}, b_{2}\right) \quad \widetilde{\boldsymbol{b}}=\left(0, \frac{1}{b_{2}}, 0\right) . \tag{6.115}
\end{equation*}
$$

This time, lattice reduction is performed by carrying out the lattice sum on the $z$ axis of the working coordinate system which is to say that one sets $j_{b}=0$, and sums over all $j_{a}$. The lattice sum is then achieved by:

$$
\begin{equation*}
S_{n, m}^{M L}=S_{n, m}^{C}\left(\beta_{1} ; \widehat{\boldsymbol{z}}\right)+S_{n, m}^{M L,+}+S_{n, m}^{M L,-}, \tag{6.116}
\end{equation*}
$$

where $S_{n, m}^{M L, \pm}$ is the sum of all sites except those along the $z$ axis:

$$
\begin{equation*}
S_{n, m}^{M L, \pm}=\sum_{j_{b} \in \mathbb{Z}^{*}} e^{i j_{b}\left(\beta_{1} b_{1}+\beta_{2} b_{2}\right)} \sum_{j_{a}=-\infty}^{\infty} e^{i j_{a} \beta_{1} a} \mathscr{H}_{n, m}\left(k \boldsymbol{r}_{\boldsymbol{j}}\right) \tag{6.117}
\end{equation*}
$$

Since $j_{b} \neq 0$ in this sum and all lattice sites are in the $y O z$ plane, the azimuthal angle of $\boldsymbol{r}_{j}$ is either $\pi / 2$ or $-\pi / 2$ for $j_{b}>0$ or $j_{b}<0$ respectively. This allows us to conclude in this plane, $\mathscr{H}_{n, m}$ doesn't depend on the sign of $m$ :

$$
\begin{equation*}
\mathscr{H}_{n,-m}\left(k \boldsymbol{r}_{\boldsymbol{j}}\right)=\mathscr{H}_{n, m}\left(k \boldsymbol{r}_{\boldsymbol{j}}\right) . \tag{6.118}
\end{equation*}
$$

We then appeal to an integral representation:

$$
\begin{align*}
h_{n}(k r) \bar{P}_{n}^{m}(\cos \theta) & =\frac{(-i)^{n+1}}{\pi} \int_{-\infty}^{\infty} e^{i k z t} K_{m}(-i k \rho \gamma(t)) \bar{P}_{n}^{m}(t) d t  \tag{6.119}\\
& =\frac{(-i)^{n-m}}{\pi} \int_{-\infty}^{\infty} e^{i k z t} H_{m}(k \rho \gamma(t)) \bar{P}_{n}^{m}(t) d t \tag{6.120}
\end{align*}
$$

where $z=r \cos \theta, k \rho=\sqrt{(k r)^{2}-(k z)^{2}}>0$, and $K_{m}(z)$ is a modified Bessel function defined by:

$$
\begin{equation*}
K_{m}(z) \equiv i^{m+1} H_{m}(i z), \tag{6.121}
\end{equation*}
$$

and finally $\gamma(t)$ is defined such that:

$$
\gamma(t)=\left\{\begin{array}{cc}
i \sqrt{t^{2}-1} & |t| \geq 1  \tag{6.122}\\
\sqrt{1-t^{2}} & t<1
\end{array}\right.
$$

The modified Bessel functions, $K_{m}(z)$. The $j_{a}$ sum can then be written:

$$
\begin{align*}
& \sum_{j_{a}=-\infty}^{\infty} e^{i j_{a} \beta_{1} a} \mathscr{H}_{n, m}\left(k \boldsymbol{r}_{\boldsymbol{j}}\right) \\
& =\frac{(-i)^{n}}{\pi}(-)^{m}\left[\operatorname{sgn}\left(j_{b}\right)\right]^{m} \sum_{j_{a}=-\infty}^{\infty} e^{i j_{a} \beta_{1} a} \int_{-\infty}^{\infty} e^{i k\left(j_{a} a+j_{b} b_{1}\right) t} H_{m}(k \rho \gamma(t)) \bar{P}_{n}^{m}(t) d t  \tag{6.123}\\
& =\frac{(-i)^{n}}{\pi}(-)^{m}\left[\operatorname{sgn}\left(j_{b}\right)\right]^{m} \int_{-\infty}^{\infty} \sum_{j_{a}=-\infty}^{\infty} e^{i j_{a}\left(\beta_{1}+k t\right) a} e^{i k j_{b} b_{1} t} H_{m}(k \rho \gamma(t)) \bar{P}_{n}^{m}(t) d t . \tag{6.124}
\end{align*}
$$

Using the 1D Poisson sum formula, we have:

$$
\begin{gather*}
\sum_{j_{a}=-\infty}^{\infty} e^{i j_{a}\left(k t+\beta_{1}\right) a}=\frac{2 \pi}{a} \sum_{g=-\infty}^{\infty} \delta\left(k t+\beta_{1}+g \frac{2 \pi}{a}\right) \\
=\frac{2 \pi}{k a} \sum_{g=-\infty}^{\infty} \delta\left(t+\frac{\beta_{1}}{k}+g \frac{2 \pi}{k a}\right)  \tag{6.125}\\
\sum_{j_{a}=-\infty}^{\infty} e^{i j_{a} \beta_{1} a} \mathscr{H}_{n, m}\left(k \boldsymbol{r}_{j}\right) \\
=\frac{(-i)^{n}}{\pi}(-)^{m}\left[\operatorname{sgn}\left(j_{b}\right)\right]^{m} \int_{-\infty}^{\infty} e^{i k j_{b} b_{1} t} \sum_{j_{a}=-\infty}^{\infty} e^{i j_{a}\left(\beta_{1}+k t\right) a} H_{m}\left(k \rho \gamma_{z}(t)\right) \bar{P}_{n}^{m}(t) d t \\
=\frac{2(-i)^{n}}{k a}(-)^{m}\left[\operatorname{sgn}\left(j_{b}\right)\right]^{m} \sum_{g=-\infty}^{\infty} e^{-i \beta_{1 p} j_{b} b_{1}} H_{m}\left(k b_{2}\left|j_{b}\right| \gamma_{g}\right) \bar{P}_{n}^{m}\left(-\bar{\beta}_{1, g}\right) \\
=\frac{2 i^{n}}{k a}\left[\operatorname{sgn}\left(j_{b}\right)\right]^{m} \sum_{g=-\infty}^{\infty} e^{-i \beta_{1, g} j_{b} b_{1}} H_{m}\left(k b_{2}\left|j_{b}\right| \gamma_{g}\right) \bar{P}_{n}^{m}\left(\bar{\beta}_{1, g}\right), \tag{6.126}
\end{gather*}
$$

where $\beta_{1, g}$ and $\bar{\beta}_{1, g}$ are defined:

$$
\begin{equation*}
\beta_{1, g} \equiv \beta_{1}+g \frac{2 \pi}{a} \quad \bar{\beta}_{1, g} \equiv \frac{\beta_{1}}{k}+g \frac{2 \pi}{k a} \tag{6.127}
\end{equation*}
$$

and

$$
\begin{equation*}
\gamma_{g} \equiv \gamma\left(\bar{\beta}_{1, g}\right) \tag{6.128}
\end{equation*}
$$

We have therefore the monolayer sum:

$$
\begin{align*}
S_{n, m}^{M L, \pm} & =\frac{2 i^{n}}{k a} \sum_{g=-\infty}^{\infty} \bar{P}_{n}^{m}\left(\bar{\beta}_{1, g}\right) \sum_{j_{b} \in \mathbb{Z}^{*}} e^{i j_{b}\left(\beta_{1}-\beta_{1, g}\right) b_{1}} e^{i j_{b} \beta_{2} b_{2}}\left[\operatorname{sgn}\left(j_{b}\right)\right]^{m} H_{m}\left(k b_{2}\left|j_{b}\right| \gamma_{g}\right) \\
& =\frac{2 i^{n}}{k a} \sum_{g=-\infty}^{\infty} \bar{P}_{n}^{m}\left(\bar{\beta}_{1, g}\right) \sum_{j_{b}=1}^{\infty}\left[e^{i j_{b} w_{g}}+(-)^{m} e^{-i j_{b} w_{g}}\right] H_{m}\left(k b_{2} j_{b} \gamma_{g}\right) \\
& =\frac{2 i^{n}}{k a} \sum_{g=-\infty}^{\infty} \bar{P}_{n}^{m}\left(\bar{\beta}_{1, g}\right) S_{m}\left(w_{g}, k b_{2} \gamma_{g}\right) \tag{6.129}
\end{align*}
$$

where we defined:

$$
\begin{equation*}
w_{g} \equiv \beta_{2} b_{2}-\frac{2 p \pi b_{1}}{a} \tag{6.130}
\end{equation*}
$$

One can rotate the lattice sum to put it in the desired coordinate system in the $x O y$ plane. This can be obtained simply by a rotation of $90^{\circ}$ about the $y$ axis. The expression of eq. (6.129) is still in the form of the 2 double infinite sums like our initial expression. However, only a finite number of $g$ values will correspond to propagating modes, i.e. $\left|\bar{\beta}_{1, g}\right|<1$, and the other values for $g$ correspond to evanescent modes and are exponentially convergent. therefore infinite series sums. The $j_{b}$ sum in eq. 6.129 is known as a Schlömilch series, and it can be expressed as a finite sum of Bernoulli polynomials.

### 6.7.3 Schlömilch series

The Schlömilch series can be expressed

$$
\begin{equation*}
S_{m}(\lambda, \mu) \equiv \sum_{j=1}^{\infty}\left[e^{i \lambda j}+(-)^{m} e^{-i \lambda j}\right] H_{m}(\mu j) \tag{6.131}
\end{equation*}
$$

The zero order sum is

$$
\begin{equation*}
S_{0}(\lambda, \mu)=-1-\frac{2 i}{\pi}\left(C+\log \frac{\mu}{4 \pi}\right)+\frac{2}{\Theta_{0}}+\sum_{g \in \mathbb{Z}^{*}}\left(\frac{2}{\Theta_{g}}+\frac{i}{\pi|g|}\right) \tag{6.132}
\end{equation*}
$$

where $C \simeq 0.5772$ is Euler's constant and

$$
\begin{equation*}
\Theta_{g}=\left(\mu^{2}-\lambda_{g}^{2}\right)^{1 / 2} \quad \lambda_{g}=\lambda+2 g \pi \tag{6.133}
\end{equation*}
$$

### 6.8 Addition theorem and Rotation matrices

### 6.8.1 Scalar spherical harmonics

The scalar spherical harmonics, $Y_{n, m}(\theta, \phi)$, are expressed in terms of the associated Legendre functions $P_{n}^{m}(\cos \theta)$ [7] :

$$
\begin{align*}
Y_{n, m}(\theta, \phi) & =\left[\frac{2 n+1}{4 \pi} \frac{(n-m)!}{(n+m)!}\right]^{\frac{1}{2}} P_{n}^{m}(\cos \theta) \exp (i m \phi) \\
& \equiv \bar{P}_{n}^{m}(\cos \theta) \exp (i m \phi) \tag{6.134}
\end{align*}
$$

where in the second line we have introduced the normalized associated Legendre functions, $\bar{P}_{n}^{m}\left(\cos \theta_{k}\right) \equiv \lambda_{n, m} P_{n}^{m}\left(\cos \theta_{k}\right)$, where the $\lambda_{n, m}$ normalization factor is defined:

$$
\begin{equation*}
\lambda_{n, m} \equiv\left[\frac{2 n+1}{4 \pi} \frac{(n-m)!}{(n+m)!}\right]^{\frac{1}{2}} \tag{6.135}
\end{equation*}
$$

These scalar spherical harmonics are normalized with respect to an integration over the solid angles :

$$
\begin{align*}
\int_{0}^{4 \pi} d \Omega Y_{v, \mu}^{*}(\theta, \phi) Y_{n, m}(\theta, \phi) & \equiv(-1)^{\mu} \int_{0}^{\pi} \sin \theta d \theta \int_{0}^{2 \pi} d \phi Y_{v,-\mu}(\theta, \phi) Y_{n, m}(\theta, \phi) \\
& =\delta_{n, v} \delta_{m, \mu} \tag{6.136}
\end{align*}
$$

In principal, the Legendre polynomials, $P_{n}(x)=P_{n}^{0}(x)$, can be obtained from Rodrigues' formula:

$$
\begin{equation*}
P_{n}(x)=\frac{1}{2^{n} n!} \frac{d^{n}}{d x^{n}}\left(x^{2}-1\right)^{n} \tag{6.137}
\end{equation*}
$$

but in practice we will calculate them with recurrence relations. Likewise, the associated Legendre functions could be obtained for from the expression:

$$
\begin{equation*}
P_{n}^{m}(x)=(-1)^{m}\left(1-x^{2}\right)^{m / 2} \frac{d^{m}}{d x^{m}} P_{n}(x) \tag{6.138}
\end{equation*}
$$

Their calculation is simplified by noting that the normalized associated Legendre functions have the convenient parity property that:

$$
\begin{equation*}
\bar{P}_{n}^{-m}(x)=(-1)^{m} \bar{P}_{n}^{m}(x) . \tag{6.139}
\end{equation*}
$$

There are alternative ways of calculating the scalar spherical harmonics that are better for formulating lattice sums and reflections from a physical interface. In lattice sums and reflections from surfaces, the spherical harmonics will be evaluated in terms of the direction of the incident or reflected wavevectors, $\widehat{k}$ :

$$
\begin{equation*}
Y_{n, m}\left(\theta_{k}, \phi_{k}\right)=Y_{n, m}(\widehat{\boldsymbol{k}})=Y_{n, m}(k / k)=Y_{n, m}\left(k_{x} / k, k_{y} / k, k_{z} / k\right), \tag{6.140}
\end{equation*}
$$

where we recall that:

$$
\begin{align*}
\frac{k_{z}}{k} & =\cos \theta_{k} \\
k_{x} / k & =\sin \theta_{k} \cos \phi_{k} \\
k_{y} / k & =\sin \theta_{k} \sin \phi_{k} \tag{6.141}
\end{align*}
$$

and we keep in mind that $\bar{P}_{n}^{m}$ are functions of $\cos \theta_{k}=k_{z} / k$.
Since $x=\cos \theta$, and the $P_{n}(x)$ are polynomials in $x$, the $\frac{d^{m}}{d x^{m}} P_{n}(x)$ are functions of $\cos \theta$. The factor $\left(1-x^{2}\right)^{m / 2}$ corresponds to $\sin ^{m} \theta_{k}$ with no ambiguity in sign since $\operatorname{Re}\left\{\theta_{k}\right\}$ is $\in(0, \pi)$. One should remark that the $\left(1-x^{2}\right)^{m / 2}$ is non-polynomial so that is why one refers to them as associated Legendre functions. For applications involving reciprocal space and/or integrations in the complex plane it proves useful to explicitly extract this factor, and define associated Legendre polynomials, which we shall denote, $\widetilde{P}_{n}^{m}$ (not to be confused with the normalized associated Legendre functions).

For positive $m$ we have then:

$$
\begin{align*}
Y_{n, m}\left(k_{x} / k, k_{y} / k, k_{z} / k\right) & =\bar{P}_{n}^{m}\left(\cos \theta_{k}\right) \exp \left(i m \phi_{k}\right) \\
& =\lambda_{n, m}(-1)^{m} \sin ^{m} \theta_{k}\left(\cos \phi_{k}+i \sin \phi_{k}\right)^{m} \frac{d^{m}}{d x^{m}} P_{n}\left(\frac{k_{z}}{k}\right) \\
& =\left(\sin \theta_{k} \cos \phi_{k}+i \sin \theta_{k} \sin \phi_{k}\right)^{m}(-1)^{m} \lambda_{n, m} \frac{d^{m}}{d x^{m}} P_{n}\left(\frac{k_{z}}{k}\right) \\
& =\left(k_{x} / k+i k_{y} / k\right)^{m}(-1)^{m} \lambda_{n, m} \frac{d^{m}}{d x^{m}} P_{n}\left(\frac{k_{z}}{k}\right) \\
& =\left(k_{x} / k+i k_{y} / k\right)^{m} \widetilde{P}_{n}^{m}\left(\frac{k_{z}}{k}\right)=\left(\frac{K}{k}\right)^{|m|} \exp \left(i m \phi_{k}\right) \widetilde{P}_{n}^{m}\left(\frac{k_{z}}{k}\right) . \tag{6.142}
\end{align*}
$$

where we have defined the normalized associated Legendre polynomials, $\widetilde{P}_{n}^{m}$, such that:

$$
\begin{equation*}
\widetilde{P}_{n}^{m}\left(\frac{k_{z}}{k}\right) \equiv(-1)^{m} \lambda_{n, m} \frac{d^{m}}{d x^{m}} P_{n}\left(\frac{k_{z}}{k}\right) . \tag{6.143}
\end{equation*}
$$

The parameter,

$$
\begin{equation*}
\boldsymbol{K} \equiv k_{x} \widehat{\boldsymbol{x}}+k_{y} \widehat{\boldsymbol{y}}, \tag{6.144}
\end{equation*}
$$

corresponds to the momentum space vector in the $x-y$ plane.
The wave vector components $k_{x}, k_{y}$, and $k_{z}$ are related to the possibly complex angles, $\theta_{k}$ and $\phi_{k}$, via the relations:

$$
\begin{align*}
k_{x} & =k \sin \theta_{k} \cos \phi_{k} \\
k_{y} & =k \sin \theta_{k} \sin \phi_{k} \\
k_{z}^{2} & =k^{2}-K^{2}=k^{2} \cos ^{2} \theta_{k} \tag{6.145}
\end{align*}
$$

This relation of eq. 6.142 for $Y_{n, m}$ can be extended to negative $m$ by writing :

$$
\begin{equation*}
Y_{n, m}\left(k_{x} / k, k_{y} / k, k_{z} / k\right)=\left(\frac{K}{k}\right)^{|m|} \exp \left(i m \phi_{k}\right) \widetilde{P}_{n}^{m}\left(\frac{k_{z}}{k}\right) \quad m \gtrless 0 \tag{6.146}
\end{equation*}
$$

as long as we define $\widetilde{P}_{n}^{-m}$ such that :

$$
\begin{equation*}
\widetilde{P}_{n}^{-m} \equiv(-1)^{m} \widetilde{P}_{n}^{m} \tag{6.147}
\end{equation*}
$$

The objective of the above procedure was to define $\widetilde{P}_{n}^{m}(x)$ that are always polynomials of $x$ for both positive and negative $m$. This is in contrast to the associated Legendre functions $\bar{P}_{n}^{m}(x)$ which are not polynomials in terms of $x$.

### 6.8.2 Translation-addition theorem for scalar partial waves

Let us consider a point $M$ in a system using spherical coordinates. We consider a second system of spherical coordinates centered on the position $\boldsymbol{r}_{0}$. The position of $M$ in this second system centered on $\boldsymbol{r}_{0}$ is:

$$
\begin{equation*}
r^{\prime}=r-r_{0} \tag{6.148}
\end{equation*}
$$

We take the usual convention of outgoing scalar partial waves as products of spherical Hankel function and scalar spherical harmonics:

$$
\begin{equation*}
\Psi_{\mathscr{H}, n, m}(k \boldsymbol{r}) \equiv \Psi_{n, m}^{(3)}(k \boldsymbol{r}) \equiv h_{n}(k r) Y_{n, m}(\theta, \phi) \tag{6.149}
\end{equation*}
$$

while the regular scalar partial waves replace the spherical Hankel functions with spherical Bessel functions:

$$
\begin{equation*}
\Psi_{\mathscr{J}, n, m}(k \boldsymbol{r}) \equiv \Psi_{n, m}^{(1)}(k \boldsymbol{r}) \equiv j_{n}(k r) Y_{n, m}(\theta, \phi) \tag{6.150}
\end{equation*}
$$

One can construct a row 'matrices' composed of the $\Psi_{\mathscr{J}, n, m}(k \boldsymbol{r})$ or $\Psi_{\mathscr{H}, n, m}(k \boldsymbol{r})$ functions respectively then the translation-addition theorem for scalar partial waves can be compactly expressed in matrix form:

$$
\begin{array}{rlrl}
\Psi_{\mathscr{H}}^{t}(k \boldsymbol{r}) & =\Psi_{\mathscr{H}}^{t}\left(k \boldsymbol{r}^{\prime}\right) \cdot \alpha\left(k \boldsymbol{r}_{0}\right) & r^{\prime}>r_{0} \\
\Psi_{\mathscr{H}}^{t}(\boldsymbol{r} \boldsymbol{r}) & =\Psi_{\mathscr{J}}^{t}\left(k \boldsymbol{r}^{\prime}\right) \cdot \beta\left(k \boldsymbol{r}_{0}\right) & & r^{\prime}<r_{0} \\
\Psi_{\mathscr{H}}^{t}(\boldsymbol{r}) & =\Psi_{\mathscr{J}}^{t}\left(\boldsymbol{r}^{\prime}\right) \cdot \beta\left(k \boldsymbol{r}_{0}\right) & & \forall\left|\boldsymbol{r}_{j}\right|, \tag{6.151}
\end{array}
$$

where the elements of the irregular translation-addition matrix, $\alpha$ have extremely simple expressions in terms of the $3 Y$ coefficients :

$$
\begin{equation*}
\alpha_{v, \mu, n m}\left(k \boldsymbol{r}_{0}\right)=4 \pi i^{v-n} \sum_{p=|n-v|}^{n+v} i^{p} 3 Y(n, m ; v, \mu ; p) h_{p}\left(k r_{0}\right) Y_{p, m-\mu}\left(\theta_{0}, \phi_{0}\right), \tag{6.152}
\end{equation*}
$$

while the the elements of the regular translation-addition matrix, $\beta_{v m u ; n, m}$, are the same as the $\alpha_{\nu, \mu ; n, m}$ coefficients but with the $j_{p}\left(k r_{0}\right)$ replacing the $h_{p}\left(k r_{0}\right)$ function i.e.:

$$
\begin{equation*}
\beta_{v, \mu ; n, m}\left(k \boldsymbol{r}_{0}\right) \equiv 4 \pi i^{v-n} \sum_{p=|n-v|}^{n+v} i^{p} 3 Y(n, m ; v, \mu ; p) j_{p}\left(k r_{0}\right) Y_{p, m-\mu}\left(\theta_{0}, \phi_{0}\right), \tag{6.153}
\end{equation*}
$$

where $3 Y(n, m ; v, \mu ; p)$ are the 3 Y coefficients defined by the angular integral of three scalar spherical harmonics:

$$
\begin{align*}
3 Y(n, m ; v, \mu ; p) & \equiv \int_{0}^{\pi} \int_{0}^{2 \pi} Y_{n, m}(\theta, \phi) Y_{v, \mu}^{*}(\theta, \phi) Y_{p, m-\mu}^{*}(\theta, \phi) \sin \theta d \theta d \phi \\
& =(-)^{\mu}(-)^{m-\mu} \int_{0}^{\pi} \int_{0}^{2 \pi} Y_{n, m}(\theta, \phi) Y_{v-\mu}(\theta, \phi) Y_{p, \mu-m}(\theta, \phi) \sin \theta d \theta d \phi \\
& =(-)^{m}\left[\frac{(2 n+1)(2 v+1)(2 p+1)}{4 \pi}\right]^{1 / 2}\left(\begin{array}{ccc}
n & v & p \\
0 & 0 & 0
\end{array}\right)\left(\begin{array}{ccc}
n & v & p \\
-m & \mu & \mu-m
\end{array}\right) . \tag{6.154}
\end{align*}
$$

The symbol,

$$
\left(\begin{array}{ccc}
n & v & p  \tag{6.155}\\
m & \mu & M
\end{array}\right)
$$

stands for the Wigner 3J coefficients. It is worth remarking that the ' 3 Y ' coefficients of eq.(6.154) are closely related to the Gaunt coefficients developed in quantum mechanics for treating the helium atom (mostly differing on account of different normalization conditions).

### 6.8.3 Vector translation-addition theorem

The vector translation-addition theorem are the vector analogue of the scalar translation theorem discussed above in section6.8.2. This would be an almost trivial extension of the scalar addition theorem if we were working with solutions of the vector Helmholtz equation in a Cartesian basis like those of eqs. (6.7) and (6.8). The additional complication is due to the fact we want to the vector translation-addition theorem to act on the purely transverse waves like those of eq. 6.9. Defining column matrices, $\Psi$, composed of the transverse vector partial waves, the vector translation-addition theorem is written:

$$
\begin{align*}
\boldsymbol{\Psi}_{\mathscr{H}}^{t}(k \boldsymbol{r}) & =\boldsymbol{\Psi}_{\mathscr{H}}^{t}\left(k \boldsymbol{r}^{\prime}\right) J\left(k \boldsymbol{r}_{0}\right) & & r^{\prime}>r_{0} \\
\boldsymbol{\Psi}_{\mathscr{H}}(k \boldsymbol{r}) & =\boldsymbol{\Psi}_{\mathscr{J}}^{t}\left(k \boldsymbol{r}^{\prime}\right) H\left(k \boldsymbol{r}_{0}\right) & & r^{\prime}<r_{0} \\
\boldsymbol{\Psi}_{\mathscr{H}}^{t}(\boldsymbol{r}) & =\boldsymbol{\Psi}_{\mathscr{J}}^{t}\left(\boldsymbol{r}^{\prime}\right) J\left(k \boldsymbol{r}_{0}\right) & & \forall\left|r_{0}\right|, \tag{6.156}
\end{align*}
$$

where the matrix $J\left(k \boldsymbol{r}_{0}\right)$ matrix can be expressed in terms of spherical scalar $\beta\left(k \boldsymbol{r}_{0}\right)$ matrices of eq. 6.153 (expressed in terms of spherical Hankel functions) while the $H\left(k \boldsymbol{r}_{0}\right)$ matrices can be expressed in terms of the $\alpha\left(k \boldsymbol{r}_{0}\right)$ matrices of eq. 6.152).

Explicitly, the $H\left(k \boldsymbol{r}_{0}\right)$ matrix can be expressed:

$$
H\left(k \boldsymbol{r}_{0}\right)=\left[\begin{array}{ll}
A_{v, \mu ; n, m}\left(k r_{0}, \theta_{0}, \phi_{0}\right) & B_{v, \mu ; n, m}\left(k r_{0}, \theta_{0}, \phi_{0}\right)  \tag{6.157}\\
B_{v, \mu ; n, m}\left(k r_{0}, \theta_{0}, \phi_{0}\right) & A_{v, \mu ; n, m}\left(k r_{0}, \theta_{0}, \phi_{0}\right)
\end{array}\right]
$$

The vector coefficients $A_{\nu, \mu ; n, m}$ are then calculated using :

$$
\begin{align*}
A_{v, \mu ; n, m} & =\frac{1}{2} \sqrt{\frac{1}{v(v+1) n(n+1)}}\left[2 \mu m \alpha_{v, \mu ; n, m}+\right. \\
& +\sqrt{(n-m)(n+m+1)} \sqrt{(v-\mu)(v+\mu+1)} \alpha_{v, \mu+1 ; n, m+1} \\
& \left.+\sqrt{(n+m)(n-m+1)} \sqrt{(v+\mu)(v-\mu+1)} \alpha_{v, \mu-1 ; n, m-1}\right] . \tag{6.158}
\end{align*}
$$

When filling up a matrix, with the $A_{\nu \mu, n m}$ coefficient, we should fill them up column by column (calculate all the $v, \mu$ ) elements for a fixed $n, m$. Then, for each $n, m$, the $B_{v \mu, n m}$ coefficients can be calculated from the previous (i.e. $v-1$ ) scalar coefficients :

$$
\begin{align*}
B_{v, \mu, n, m} & =-i \frac{1}{2} \sqrt{\frac{2 v+1}{2 v-1} \frac{1}{v(v+1) n(n+1)}}\left[2 m \sqrt{(v-\mu)(v+\mu)} \alpha_{v-1, \mu ; n, m}\right. \\
& +\sqrt{(n-m)(n+m+1)} \sqrt{(v-\mu)(v-\mu-1)} \alpha_{v-1, \mu+1 ; n, m+1} \\
& \left.-\sqrt{(n+m)} \sqrt{(n-m+1)} \sqrt{(v+\mu)(v+\mu-1)} \alpha_{v-1, \mu-1 ; n, m-1}\right] \tag{6.159}
\end{align*}
$$

### 6.8.4 Rotation matrices

Under rotation, each of the four blocks of a vector translation-addition matrix transform following the rotation matrix, $\mathscr{D}(\alpha, \beta, \gamma)$, which is expressed in terms of the 3 Euler angles, $\alpha, \beta$, and $\gamma$. The $\mathscr{D}(\alpha, \beta, \gamma)$ matrix elements are described in detail in ref.[7], and are block diagonal in the orbital (multipole) 'quantum' number, $n$ :

$$
\begin{equation*}
[\mathscr{D}(\alpha, \beta, \gamma)]_{v, \mu, n, m}=\delta_{n, v} \exp (i \mu \alpha) d_{\mu m}^{(n)}(\beta) \exp (i m \gamma) \tag{6.160}
\end{equation*}
$$

The elements $d_{\mu m}^{(n)}$ are standard, [7] and the $d_{\mu, m}^{(n)}$ term in the rotation matrices can be expressed in terms of the Jacobi polynomials [7] :

$$
\begin{equation*}
d_{\mu m}^{(n)}(\beta)=\left[\frac{(n+\mu)!(n-\mu)!}{(n+m)!(n-m)}\right]^{1 / 2}\left(\cos \frac{\beta}{2}\right)^{m+\mu} \times\left(\sin \frac{\beta}{2}\right)^{m-\mu} P_{n-\mu}^{(\mu-m, m+\mu)}(\cos \beta) \tag{6.161}
\end{equation*}
$$

### 6.9 Recurrence relations for special functions

Partial wave descriptions are composed of products of spherical harmonic and spherical Bessel types special functions. For a numerical analysis, it is important to calculate these functions rapidly and accurately. Recurrence relations prove to be a good manner to achieve this goal. Multipole expansions must be truncated to a given order $n_{\text {max }}$, which determines the the strength of spatial field variations. Inspection of the translation-addition theorem formulas show us that we will need to evaluate spherical harmonic and spherical Hankel functions up to order $2 n_{\max }$.

### 6.9.1 Recurrence relations for associated Legendre polynomials

The recurrence relations for the $\widetilde{P}_{n}^{m}$ polynomials are are initialized with:

$$
\begin{equation*}
\widetilde{P}_{0}^{0}(u)=\sqrt{\frac{1}{4 \pi}} . \tag{6.162}
\end{equation*}
$$

We can then calculate all the $\widetilde{P}_{n}^{n}$ up to $2 n_{\max }$ with the recurrence relation:

$$
\begin{equation*}
\widetilde{P}_{n}^{n}(u)=-\sqrt{\frac{2 n+1}{2 n}} \widetilde{P}_{n-1}^{n-1}(u) . \tag{6.163}
\end{equation*}
$$

The $\widetilde{P}_{n}^{m}$ with $m=n-1$ are calculated via:

$$
\begin{equation*}
\widetilde{P}_{n}^{n-1}(u)=u \sqrt{2 n+1} \widetilde{P}_{n-1}^{n-1} . \tag{6.164}
\end{equation*}
$$

All the remaining $\widetilde{P}_{n}^{m}$ with $m=1, \ldots, n-2$ can be calculated for each $n=3, \ldots, 2 n_{\max }$ using the relation :

$$
\begin{equation*}
\widetilde{P}_{n}^{m}(u)=\sqrt{\frac{(2 n+1)}{\left(n^{2}-m^{2}\right)}}\left[\sqrt{(2 n-1)} x \widetilde{P}_{n-1}^{m}-\sqrt{\frac{(n-1)^{2}-m^{2}}{(2 n-3)}} \widetilde{P}_{n-2}^{m}\right] . \tag{6.165}
\end{equation*}
$$

Although we obtain the $\widetilde{P}_{n}^{0}$ in the above scheme, it can sometimes prove useful to obtain the normalized Legendre Polynomials through the recurrence relation :

$$
\begin{equation*}
\widetilde{P}_{n}^{0}(u)=\frac{1}{n}\left(u \sqrt{4 n^{2}-1}\right) \widetilde{P}_{n-1}^{0}(u)-(n-1) \sqrt{\frac{2 n+1}{2 n-3}} \widetilde{P}_{n-2}^{0}(u) . \tag{6.166}
\end{equation*}
$$

The $\widetilde{P}_{n}^{m}$ with negative values of $m$ are calculated using:

$$
\begin{equation*}
\widetilde{P}_{n}^{-m}(u)=(-)^{m} \widetilde{P}_{n}^{m}(u) . \tag{6.167}
\end{equation*}
$$

### 6.9.2 Logarithmic Bessel functions

The spherical Bessel, Neumann and Hankel functions of the Ricatti form, are simply these functions multiplied by their argument. The advantage of this form is that they have better limit properties for small arguments. Their definitions are respectively:

$$
\begin{equation*}
\psi_{n}(z) \equiv z j_{n}(z), \quad \chi_{n}(z) \equiv z y_{n}(z), \quad \xi_{n}(z) \equiv z h_{n}(z) . \tag{6.168}
\end{equation*}
$$

Multiplying the logarithmic derivatives of these functions by their argument defines the functions, $\varphi^{(1)}, \varphi^{(2)}, \varphi^{(3)}$ :

$$
\begin{equation*}
\varphi_{n}^{(1)}(z) \equiv \frac{\psi_{n}^{\prime}(z)}{j_{n}(z)}, \quad \varphi_{n}^{(2)}(z) \equiv \frac{\chi_{n}^{\prime}(z)}{y_{n}(z)}, \quad \varphi_{n}^{(3)}(z) \equiv \frac{\xi_{n}^{\prime}(z)}{h_{n}(z)} . \tag{6.169}
\end{equation*}
$$

The $\varphi_{n}^{(i)}$ can also be generated by particularly efficient numerical recursion relations. They also provide particularly symmetric expression for the Mie coefficients of spherical scatterers and formulations of matrix balancing. The Wronskian relation for Ricatti-Bessel functions:

$$
\begin{equation*}
\psi_{n}(x) \xi_{n}^{\prime}(x)-\psi_{n}^{\prime}(x) \xi_{n}(x)=i \tag{6.170}
\end{equation*}
$$

takes a simple form in terms of the $\varphi_{n}^{(1,2,3)}$ functions as:

$$
\begin{equation*}
\varphi_{n}^{(3)}(z)-\varphi_{n}^{(1)}(z)=\frac{i}{x j_{n}(x) h_{n}(x)} \tag{6.171}
\end{equation*}
$$

The $\varphi^{(3)}$ can be reliably calculated numerically from the upward recurrence relation:

$$
\begin{equation*}
\varphi_{n}^{(3)}(z)=\frac{z^{2}}{n-\varphi_{n-1}^{(3)}(z)}-n, \tag{6.172}
\end{equation*}
$$

starting with an initialization of

$$
\begin{equation*}
\varphi_{0}^{(3)}(z)=i z \tag{6.173}
\end{equation*}
$$

The $\xi_{n}(z)$ functions in most situations can then be readily calculated numerically from the recurrence relation:

$$
\begin{equation*}
\xi_{n}(z)=\frac{\xi_{n-1}(z)}{z}\left(n-\varphi_{n-1}^{(3)}(z)\right) \tag{6.174}
\end{equation*}
$$

starting from the initial value $\xi_{0}(z)=-i e^{i z}$. One should note that analytical expressions exist for the spherical Ricatti-Hankel functions, $\xi_{n}(z)$, and these can be useful at low multipole order:

$$
\begin{align*}
& \xi_{0}(z)=-i e^{i z} \\
& \xi_{1}(z)=-e^{i z}\left(1+\frac{i}{z}\right) \\
& \xi_{2}(z)=e^{i z}\left(i-\frac{3}{z}-\frac{3 i}{z^{2}}\right) \tag{6.175}
\end{align*}
$$

When one deals with high multipole orders however, it usually is more practical to exploit the recurrence relations of eq. 6.172) and 6.174).

The regular Ricatti Bessel functions, $\varphi_{n}^{(1)}(z)$, obey the same recurrence relations as the $\varphi_{n}^{(3)}(z)$ functions. If one calculates them by via upward recurrence, things may work fine for the first few recurrence calculations, but at some point, the recurrence relations can go completely off course. The solution to this problem has been known for quite some time is that the $\varphi_{n}(z)$ functions should be calculated starting from high values of $n$ in a reverse recurrence relation. I start usually with $n$ equal to at least $n_{\max }+20$ where $n_{\max }$ is the largest value that I want to use in calculations, and start simply with $\varphi_{n_{\max }+20}(z)=0$. The $\varphi_{n}(z)$ functions so obtained have always been the correct ones up to machine precision. The reverse recurrence relation is:

$$
\begin{equation*}
\varphi_{n}(z)=n+1-\frac{z^{2}}{n+1+\varphi_{n+1}(z)} . \tag{6.176}
\end{equation*}
$$

One can check calculations by verifying that the $\varphi_{0}(z)$ obtained by backward recurrence is equal to the analytical result:

$$
\begin{equation*}
\varphi_{0}(z)=z \frac{\cos z}{\sin z} \tag{6.177}
\end{equation*}
$$

Once the $\varphi_{n}(z)$ functions have been calculated, one can readily generate the $\psi_{n}(z)$ functions with the upward recurrence relation which is the direct analogue of eq. 6.174

$$
\begin{equation*}
\psi_{n}(z)=\frac{\psi_{n-1}(z)}{z}\left(n-\varphi_{n-1}(z)\right) \tag{6.178}
\end{equation*}
$$

starting with the initial value $\psi_{0}(z)=\sin z$. Analytical expressions for the lowest $\psi_{n}(z)$ are:

$$
\begin{align*}
& \psi_{0}(z)=\sin z \\
& \psi_{1}(z)=\frac{\sin z}{z}-\cos z  \tag{6.179}\\
& \psi_{2}(z)=\left(\frac{3}{z^{2}}-1\right) \sin z-\frac{3}{z} \cos z \tag{6.180}
\end{align*}
$$

It is perhaps worth remarking that, there is one potential numerical problem with using eq. (6.174) to calculate spherical Hankel functions. For some values of $z$ the real and imaginary parts of the spherical Hankel functions can have extremely different absolute values. For concreteness, let us assume that $\left|\operatorname{Im}\left(h_{n}(z)\right)\right| \ll\left|\operatorname{Re}\left(h_{n}(z)\right)\right|$, then using eq. (6.174), the calculated value of $\operatorname{Im}\left(h_{n}(z)\right)$ will usually be quite inaccurate if its absolute value is less than last significant figure in the calculation of $\operatorname{Re}\left(h_{n}(z)\right)$. This problem can be circumvented (for real values of $z$ at least) by calculating the spherical Neumann functions (denoted here by $y_{n}(z)$ but some authors denote it $n_{n}(z)$ ). We recall that the Neumann functions are real-valued provided that $z$ is real valued.

The Ricatti Neumann functions are defined :

$$
\begin{equation*}
\chi_{n}(z) \equiv z y_{n}(z) \tag{6.181}
\end{equation*}
$$

The first few Ricatti Neumann functions, $\chi_{n}(z)$, are

$$
\begin{align*}
& \chi_{0}(z)=-\cos z \\
& \chi_{1}(z)=-\frac{\cos z}{z}-\sin z \\
& \chi_{2}(z)=-\left(\frac{3}{z^{2}}-1\right) \cos z-\frac{3}{z} \sin z \tag{6.182}
\end{align*}
$$

We define a $\varphi^{(2)}$ 'logarithmic derivative' Neumann function as :

$$
\begin{equation*}
\varphi_{n}^{(2)}(z) \equiv \frac{\chi_{n}^{\prime}(z)}{y_{n}(z)} \tag{6.183}
\end{equation*}
$$

We can calculate the $\varphi_{n}^{(2)}$ from the upward recurrence relation:

$$
\begin{equation*}
\varphi_{n}^{(2)}(z)=\frac{z^{2}}{n-\varphi_{n-1}^{(2)}(z)}-n, \tag{6.184}
\end{equation*}
$$

with an initialization of

$$
\begin{equation*}
\varphi_{0}^{(2)}(z)=-z \frac{\sin z}{\cos z} \tag{6.185}
\end{equation*}
$$

Once the $\varphi_{n}^{(2)}$ functions have been calculated, one can readily generate the $\chi_{n}(z)$ functions with the upward recurrence relation which is the direct analogue of eq. 6.174)

$$
\begin{equation*}
\chi_{n}(z)=\frac{\chi_{n-1}(z)}{z}\left(n-\varphi_{n-1}^{(2)}(z)\right) \tag{6.186}
\end{equation*}
$$

starting with the initial value $\chi_{0}(z)=-\cos z$.
Having calculated $\psi_{n}(z)$ via eq. 6.178) and $\chi_{n}(z)$ via eq. 6.186, one can finally construct $h_{n}(z)$ by

$$
\begin{equation*}
h_{n}(z) \equiv j_{n}(z)+i y_{n}(z) \tag{6.187}
\end{equation*}
$$

with the real and imaginary parts of $\xi_{n}(z)$ now both being calculated up to machine accuracy.
The ratio of spherical Bessel functions to spherical Hankel functions also occurs frequently in Mie theory, and I found it convenient and more accurate to calculate these ratios directly using upward recurrence relations, notably:

$$
\begin{equation*}
\frac{j_{n}(z)}{h_{n}(z)}=\frac{j_{n-1}(z)}{h_{n-1}(z)} \frac{n-\varphi_{n-1}(z)}{n-\varphi_{n-1}^{(3)}(z)} \tag{6.188}
\end{equation*}
$$

with the initialization

$$
\begin{equation*}
\frac{j_{1}(z)}{h_{1}(z)}=\frac{1}{2}\left(1-\frac{i-z}{i+z} \exp [-2 i z]\right) \tag{6.189}
\end{equation*}
$$

which has good properties for numerical calculations. If one needs $j_{0}(z) / h_{0}(z)$, one can usually calculate it as

$$
\begin{equation*}
\frac{j_{0}(z)}{h_{0}(z)}=-\frac{\sin z}{i e^{i z}}=\frac{1-e^{-2 i z}}{2} \tag{6.190}
\end{equation*}
$$

and $j_{1}(z) / h_{1}(z)$ satisfies the recurrence relation if we start with $j_{0}(z) / h_{0}(z)$ since

$$
\begin{align*}
\frac{j_{1}(z)}{h_{1}(z)} & =\frac{j_{0}(z)}{h_{0}(z)} \frac{1-\varphi_{0}(z)}{1-\varphi_{0}^{(3)}(z)}=\frac{1-e^{-2 i z}}{2} \frac{1-z \frac{\cos z}{\sin z}}{1-i z}=-\frac{1}{2 i} \frac{1-i z-e^{-2 i z}-i z e^{-2 i z}}{i+z} \\
& =\frac{1}{2}\left(1-\frac{i-z}{i+z} e^{-2 i z}\right) \tag{6.191}
\end{align*}
$$

For coated spheres, it is can also useful to use the analogous recurrence relation:

$$
\begin{equation*}
\frac{j_{n}(z)}{y_{n}(z)}=\frac{\psi_{n-1}(z)}{y_{n-1}(z)} \frac{n-\varphi_{n-1}(z)}{n-\varphi_{n-1}^{(2)}(z)} \tag{6.192}
\end{equation*}
$$

with the initialization of

$$
\begin{equation*}
\frac{j_{0}(z)}{y_{0}(z)}=-\frac{\sin z}{\cos z} \tag{6.193}
\end{equation*}
$$

with the first recurrence giving

$$
\begin{equation*}
\frac{j_{1}(z)}{y_{1}(z)}=\frac{z \cos z-\sin z}{z \sin z+\cos z} \tag{6.194}
\end{equation*}
$$

Other useful relations are obtained from the classic spherical Bessel function recurrence relations:

$$
\begin{align*}
& f_{n}(z)=\frac{z f_{n-1}(z)+z f_{n+1}(z)}{2 n+1} \\
& f_{n}^{\prime}(z)=\frac{n f_{n-1}(z)-(n+1) f_{n+1}(z)}{2 n+1} \tag{6.195}
\end{align*}
$$

with $f_{n}(z)=j_{n}(z), h_{n}(z)$ to obtain :

$$
\begin{equation*}
z j_{n}^{\prime}(z)+(n+1) j_{n}(z)=n \frac{\psi_{n-1}(z)}{2 n+1}+(n+1) \frac{\psi_{n-1}(z)}{2 n+1}=\psi_{n-1}(z) \tag{6.196}
\end{equation*}
$$

with

$$
\begin{equation*}
\psi_{n}^{\prime}(z) \equiv z j_{n}^{\prime}(z)+j_{n}(z) \tag{6.197}
\end{equation*}
$$

we obtain a convenient expression for the derivative of Ricatti-Bessel functions:

$$
\begin{equation*}
\psi_{n}^{\prime}(z)=\psi_{n-1}(z)-n j_{n}(z) \tag{6.198}
\end{equation*}
$$

or the expression:

$$
\begin{equation*}
\psi_{n}^{\prime}(z)=(n+1) j_{n}(z)-\psi_{n+1}(z) . \tag{6.199}
\end{equation*}
$$

### 6.9.3 Vector Spherical Harmonics

There is no universally accepted notation for the Vector Spherical Harmonics (VSHs). Our notation for their normalized forms is $\boldsymbol{X}_{n, m}, \boldsymbol{Y}_{n, m}$, and $\boldsymbol{Z}_{n, m}$ where they are respectively defined by

$$
\begin{align*}
\boldsymbol{X}_{n, m}(\boldsymbol{\theta}, \phi) & \equiv \boldsymbol{Z}_{n, m}(\theta, \phi) \times \widehat{\boldsymbol{r}} \\
\boldsymbol{Y}_{n, m}(\boldsymbol{\theta}, \phi) & \equiv \widehat{\boldsymbol{r}} Y_{n, m}(\boldsymbol{\theta}, \boldsymbol{\phi}) \\
\boldsymbol{Z}_{n, m}(\boldsymbol{\theta}, \phi) & \equiv \frac{r \boldsymbol{\nabla} Y_{n, m}(\boldsymbol{\theta}, \phi)}{\sqrt{n(n+1)}}=\widehat{\boldsymbol{r}} \times \boldsymbol{X}_{n, m}(\theta, \phi), \tag{6.200}
\end{align*}
$$

The scalar spherical harmonics, $Y_{n, m}(\theta, \phi)$, do have a nearly universal convention for their definitions [10] which we recalled in eq.(6.134].

For numerical calculations of the VSHs, it is convenient to introduce the normalized functions $\bar{u}_{n}^{m}$ and $\bar{s}_{n}^{m}$ defined as:

$$
\begin{align*}
\bar{u}_{n}^{m}(\cos \theta) & \equiv \gamma_{n, m} \frac{m}{\sin \theta} P_{n}^{m}(\cos \theta)  \tag{6.201}\\
\bar{s}_{n}^{m}(\cos \theta) & \equiv \gamma_{n, m} \frac{d}{d \theta} P_{n}^{m}(\cos \theta) \tag{6.202}
\end{align*}
$$

where the $P_{n}^{m}$ are the Legendre functions defined in eqs. 6.137, (6.138), and 6.139), and $\gamma_{n, m}$ a normalization factor given by

$$
\begin{equation*}
\gamma_{n, m} \equiv \frac{\lambda_{n, m}}{\sqrt{n(n+1)}}=\sqrt{\frac{(2 n+1)(n-m)!}{4 \pi n(n+1)(n+m)!}} . \tag{6.203}
\end{equation*}
$$

The transverse VSHs have compact expressions in terms of $\bar{u}_{n}^{m}$ and $\bar{s}_{n}^{m}$ :

$$
\begin{align*}
\boldsymbol{X}_{n, m}(\theta, \boldsymbol{\phi}) & =\left[i \bar{u}_{n}^{m}(\cos \theta) \widehat{\boldsymbol{\theta}}-\bar{s}_{n}^{m}(\cos \theta) \widehat{\boldsymbol{\phi}}\right] \exp (i m \phi) \\
\boldsymbol{Z}_{n, m}(\theta, \phi) & =\left[\bar{s}_{n}^{m}(\cos \theta) \widehat{\boldsymbol{\theta}}+i \bar{u}_{n}^{m}(\cos \theta) \widehat{\boldsymbol{\phi}}\right] \exp (i m \phi) \tag{6.204}
\end{align*}
$$

The normalized $\bar{u}_{n}^{m}$ functions can be readily calculated with recurrence relations:

$$
\begin{align*}
\bar{u}_{n}^{0}(x)= & 0, \quad \bar{u}_{1}^{1}(x)=-\frac{1}{4} \sqrt{\frac{3}{\pi}} \\
\bar{u}_{n}^{n}(x)= & -\sqrt{\frac{n(2 n+1)}{2(n+1)(n-1)}} \sqrt{1-x^{2}} \bar{u}_{n-1}^{n-1}(x) \\
\bar{u}_{n}^{m}(x)= & \sqrt{\frac{(n-1)\left(4 n^{2}-1\right)}{(n+1)\left(n^{2}-m^{2}\right)}} x \bar{u}_{n-1}^{m}(x) \\
& -\sqrt{\frac{(2 n+1)(n-1)(n-2)(n-m-1)(n+m-1)}{(2 n-3) n(n+1)\left(n^{2}-m^{2}\right)}} \bar{u}_{n-2}^{m}(x) \\
\bar{u}_{n}^{n-1}(x)= & \sqrt{\frac{(2 n+1)(n-1)}{(n+1)}} x \bar{u}_{n-1}^{n-1}(x), \tag{6.205}
\end{align*}
$$

while the $\bar{s}_{n}^{m}$ can be determined from the $\bar{u}_{n}^{m}$ functions via the formula:

$$
\begin{equation*}
\bar{s}_{n}^{m}(\cos \theta)=\frac{1}{(m+1)} \sqrt{(n+m+1)(n-m)} \sin \theta \bar{u}_{n}^{m+1}(\cos \theta)+\cos \theta \bar{u}_{n}^{m}(\cos \theta) \tag{6.206}
\end{equation*}
$$

The respective parity properties of $\bar{u}_{n}^{m}$ and $\bar{s}_{n}^{m}$ are:

$$
\begin{align*}
\bar{u}_{n}^{-m}(x) & =(-1)^{m+1} \bar{u}_{n}^{m}(x) \\
\bar{s}_{n}^{-m}(x) & =(-1)^{m} \bar{s}_{n}^{m}(x) . \tag{6.207}
\end{align*}
$$

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Chapter 7:

# Differential Theory of Periodic Structures 

Evgeny Popov

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## Chapter 7

# Differential Theory of Periodic Structures 

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The basic idea of the differential methods consisits in projecting the electromagnetic field on a set of basic functions in order to reduce Maxwell partial differential equations into a set of ordinary differential equations. When working in a Cartesian coordinates, the natural basis consists of exponentials, using the periodicity of the optogeometrical parameters. Diffraction by a single aperture requires working in the basis of cylindrical Bessel functions [7.1], while diffraction by an arbitrary-shaped single object requires vector spherical functions [7.2] as a basis.

The first studies using the differential method [7.3] appeared in the late 1960s, initiated by the birth of the computers. These studies concerned the modeling of diffusion of particles in nuclear potential by using the separation of variables of the radial Schrödinger equation. The method was called "optical method" due to the similarity between the Schrödinger and the Helmholtz equations. The first applications to grating diffraction appear in 1969 [7.4], but accurate and converging results required combining the differential method with conformal mapping techniques [7.5]. The classical differential theory as known nowadays was formulated in $[7.6,7.7]$. One can find a detailed review on the classical differential method in [7.8]

It appeared that the classical differential theory suffered from severe numerical problems in transverse magnetic (TM) polarization, as well as for deep gratings. The first breakthrough was made in the first half of the 1990s, by introducing orthonormalization of the differential equations during their integration [7.9] and followed later by the so-called Rmatrix or S-matrix propagating algorithms [7.10]. The second breakthrough improved considerably the convergence in TM polarization for lamellar gratings, by introducing the correct factorization rules (see further on), at first by chance [7.11] and after that using theoretical arguments [7.12], closely followed by a generalization to arbitrary profiles [7.13]. A detailed review can be found in [7.14].

### 7.1. Maxwell equations in the truncated Fourier space

Let us consider a structure with two-dimensional periodicity along the x - and y -axis (Fig.7.1) with periods equal to $d_{x}$ and $d_{y}$. The modulated (grating) region extends in $z$ from $z_{\text {min }}$ to $z_{\text {max }}$. Inside that region, for a given value of the vertical coordinate $z$, the permittivity $\varepsilon$ and permeability $\mu$ are periodic functions in x and y that can be projected on exponential Fourier basis:

$$
\begin{align*}
& \varepsilon(\mathrm{x}, \mathrm{y}, \mathrm{z})=\sum_{\mathrm{m}, \mathrm{n}=-\infty}^{+\infty} \varepsilon_{\mathrm{m}, \mathrm{n}}(\mathrm{z}) \exp \left(\mathrm{imK}_{\mathrm{x}} \mathrm{x}+\mathrm{inK}_{\mathrm{y}} \mathrm{y}\right) \\
& \mu(\mathrm{x}, \mathrm{y}, \mathrm{z})=\sum_{\mathrm{m}, \mathrm{n}=-\infty}^{+\infty} \mu_{\mathrm{m}, \mathrm{n}}(\mathrm{z}) \exp \left(\mathrm{imK}_{\mathrm{x}} \mathrm{x}+\mathrm{inK} \mathrm{y}_{\mathrm{y}} \mathrm{y}\right) \tag{7.1}
\end{align*}
$$



Fig.7.1. Schematical representation of a structure having two-dimensional periodicity in $x$ and $y$ directions, consisting of truncated pyramids with height $h$.
where $\mathrm{K}_{\mathrm{x}}=2 \pi / \mathrm{d}_{\mathrm{x}}$ and $\mathrm{K}_{\mathrm{y}}=2 \pi / \mathrm{d}_{\mathrm{y}}$. We shall deal with a monochromatic (wavelength $\lambda$ ) plane wave incident on the structure with a wavevector:

$$
\begin{equation*}
\overrightarrow{\mathrm{k}}_{\text {inc }}=\left(\alpha_{0}, \beta_{0},-\gamma_{0}\right) \tag{7.2}
\end{equation*}
$$

with components related to the incident polar angle $\theta$ (between the incident direction and the grating normal) and azimuthal angle $\varphi$ (between the plane of incidence and the xOz -plane):

$$
\begin{align*}
& \alpha_{0}=\mathrm{k}_{0} \sin \theta \cos \varphi, \quad \beta_{0}=\mathrm{k}_{0} \sin \theta \cos \varphi, \\
& \gamma_{0}=\sqrt{\mathrm{k}_{0}^{2} \mathrm{n}_{\mathrm{inc}}^{2}-\alpha_{0}^{2}-\beta_{0}^{2}}, \quad \mathrm{k}_{0}=2 \pi / \lambda \tag{7.3}
\end{align*}
$$

where $\mathrm{n}_{\text {inc }}$ is the refractive index of the cladding.
The existence and uniqueness of the solution of the diffraction problem is an interesting problem that is not discussed here. The reader can refer to several basic works (see for example [7.15, 7.16]. What is important to conclude is that the electromagnetic field is pseudo-periodic, so that similarly to eq.(7.1), the electric $\overrightarrow{\mathrm{E}}$ and magnetic $\overrightarrow{\mathrm{H}}$ field vectors can be represented in pseudo-Fourier series:

$$
\begin{align*}
& \overrightarrow{\mathrm{E}}(\mathrm{x}, \mathrm{y}, \mathrm{z})=\sum_{\mathrm{m}, \mathrm{n}=-\infty}^{+\infty} \overrightarrow{\mathrm{E}}_{\mathrm{m}, \mathrm{n}}(\mathrm{z}) \exp \left[\mathrm{i}\left(\alpha_{0}+\mathrm{mK}_{\mathrm{x}}\right) \mathrm{x}+\mathrm{i}\left(\beta_{0}+\mathrm{nK}_{\mathrm{y}}\right) \mathrm{y}\right]  \tag{7.4}\\
& \overrightarrow{\mathrm{H}}(\mathrm{x}, \mathrm{y}, \mathrm{z})=\sum_{\mathrm{m}, \mathrm{n}=-\infty}^{+\infty} \overrightarrow{\mathrm{H}}_{\mathrm{m}, \mathrm{n}}(\mathrm{z}) \exp \left[\mathrm{i}\left(\alpha_{0}+\mathrm{mK}_{\mathrm{x}}\right) \mathrm{x}+\mathrm{i}\left(\beta_{0}+\mathrm{nK} \mathrm{~K}_{\mathrm{y}}\right) \mathrm{y}\right]
\end{align*}
$$

In what follows, we use the notations:

$$
\begin{equation*}
\alpha_{\mathrm{m}}=\alpha_{0}+\mathrm{mK}_{\mathrm{x}}, \quad \beta_{\mathrm{n}}=\beta_{0}+\mathrm{nK}_{\mathrm{y}} . \tag{7.5}
\end{equation*}
$$

From a numerical point of view, it is necessary to truncate the series in eqs.(7.1) and (7.4), introducing truncation parameters $\mathrm{N}_{\mathrm{x}}$ and $\mathrm{N}_{\mathrm{y}}$, which limit the lower and the upper boundaries in the series.

Maxwell equations written in Fourier space take the form, assuming $\exp (-i \omega t)$ time dependence with circular frequency $\omega$ :

$$
\begin{align*}
& i \beta_{m, n} E_{z, m, n}(z)-\frac{d}{d z} E_{y, m, n}(z)=i \omega B_{x, m, n}(z) \\
& \frac{d}{d z} E_{x, m, n}(z)-i \alpha_{m, n} E_{z, m, n}(z)=i \omega B_{y, m, n}(z) \\
& i \alpha_{m, n} E_{y, m, n}(z)-i \beta_{m, n} E_{x, m, n}(z)=i \omega B_{z, m, n}(z)  \tag{7.6}\\
& i \beta_{m, n} H_{z, m, n}(z)-\frac{d}{d z} H_{y, m, n}(z)=-i \omega D_{x, m, n}(z) \\
& \frac{d}{d z} H_{x, m, n}(z)-i \alpha_{m, n} H_{z, m, n}(z)=-i \omega D_{y, m, n}(z) \\
& i \alpha_{m, n} H_{y, m, n}(z)-i \beta_{m, n} H_{x, m, n}(z)=-i \omega D_{z, m, n}(z) .
\end{align*}
$$

As can be observed, the third and the sixth equations are not differential equations, and they are used to eliminate the z-components of the fields, as shown further on. It has to be stressed out that the equations with different ( $\mathrm{m}, \mathrm{n}$ ) numbers are coupled through the Fourier components of $\vec{D}=\varepsilon \vec{E}$ and $\vec{B}=\mu \vec{H}$.

The next step is to factorize the products $\vec{D}=\varepsilon \vec{E}$ and $\vec{B}=\mu \vec{H}$. In this chapter we assume media with linear dielectric and magnetic properties and without spontaneous polarizations. The problem of Fourier transform of the product of two functions

$$
\begin{equation*}
\overrightarrow{\mathrm{D}}(\mathrm{x}, \mathrm{y}, \mathrm{z})=\sum_{\mathrm{m}, \mathrm{n}=-\infty}^{+\infty} \overrightarrow{\mathrm{D}}_{\mathrm{m}, \mathrm{n}}(\mathrm{z}) \exp \left[\mathrm{i}\left(\alpha_{0}+\mathrm{mK}_{\mathrm{x}}\right) \mathrm{x}+\mathrm{i}\left(\beta_{0}+\mathrm{nK} \mathrm{~K}_{\mathrm{y}}\right) \mathrm{y}\right] \tag{7.7}
\end{equation*}
$$

is, in generally, solved theoretically by convolution of the Fourier transformers of the two functions, using the so-called Laurent's rule:

$$
\begin{equation*}
\overrightarrow{\mathrm{D}}_{\mathrm{m}, \mathrm{n}}(\mathrm{z})=\sum_{\mathrm{m}^{\prime}, \mathrm{n}^{\prime}=-\infty}^{+\infty} \varepsilon_{\mathrm{m}-\mathrm{m}^{\prime}, \mathrm{n}-\mathrm{n}^{\prime}}(\mathrm{z}) \overrightarrow{\mathrm{E}}_{\mathrm{m}^{\prime}, \mathrm{n}^{\prime}}(\mathrm{z}) \tag{7.8}
\end{equation*}
$$

However, there are several problems in the numerical application of this rule:

First, numerical applications are simplified when using matrix notations. However, most of the standard routines use single-rank vectors and rectangular (2-ranks) matrices, while the vectors D and E in eq.(7.8) have two indexes, and the matrix $\varepsilon$ depend on four indexes. In the case of classical grating with one-dimensional periodicity, this problem does not exist. Fortunately, for structures having 2D periodicity, a reduction to standard arrays is possible by introduction of a single index instead of the double for the vectors, by the following substitution:

$$
\begin{equation*}
\mathrm{p}=\left(\mathrm{m}+\mathrm{N}_{\mathrm{x}}\right)\left(2 \mathrm{~N}_{\mathrm{y}}+1\right)+\left(\mathrm{n}+\mathrm{N}_{\mathrm{y}}+1\right) \tag{7.9}
\end{equation*}
$$

so that when $m$ varies between $-N_{x}$ and $+N_{x}$ and $n$ varies between $-N_{y}$ and $+N_{y}, p$ varies between 1 and $\mathrm{P}_{\text {max }}=\left(2 \mathrm{~N}_{\mathrm{x}}+1\right)\left(2 \mathrm{~N}_{\mathrm{y}}+1\right)$. Using these notations, we can introduce standard arrays in the following manner:

$$
\begin{align*}
& \overrightarrow{\mathrm{D}}_{\mathrm{p}}(\mathrm{z})=\overrightarrow{\mathrm{D}}_{\mathrm{m}, \mathrm{n}}(\mathrm{z}), \quad \overrightarrow{\mathrm{E}}_{\mathrm{p}}(\mathrm{z})=\overrightarrow{\mathrm{E}}_{\mathrm{m}, \mathrm{n}}(\mathrm{z}), \quad \text { etc. for } \overrightarrow{\mathrm{H}} \text { and } \overrightarrow{\mathrm{B}},  \tag{7.10}\\
& \varepsilon_{\mathrm{p}-\mathrm{p}^{\prime}}(\mathrm{z})=\varepsilon_{\mathrm{m}-\mathrm{m}^{\prime}, \mathrm{n}-\mathrm{n}^{\prime}}(\mathrm{z})
\end{align*}
$$

so that eq.(7.8) takes the standard truncated form

$$
\begin{equation*}
\overrightarrow{\mathrm{D}}_{\mathrm{p}}(\mathrm{z})=\sum_{\mathrm{p}^{\prime}=1}^{\mathrm{P}_{\max }} \varepsilon_{\mathrm{p}-\mathrm{p}^{\prime}}(\mathrm{z}) \overrightarrow{\mathrm{E}}_{\mathrm{p}^{\prime}}(\mathrm{z}) . \tag{7.11}
\end{equation*}
$$

That can be written in matrix notations in the form

$$
\begin{equation*}
[\overrightarrow{\mathrm{D}}(\mathrm{z})]=\llbracket \varepsilon(\mathrm{z}) \rrbracket[\overrightarrow{\mathrm{E}}(\mathrm{z})], \tag{7.12}
\end{equation*}
$$

where double square brackets stand for the Toepliz matrix.
In addition, two diagonal matrices are useful:

$$
\begin{align*}
& \alpha_{\mathrm{p}, \mathrm{p}^{\prime}}=\delta_{\mathrm{p}, \mathrm{p}^{\prime}} \alpha_{\mathrm{m}} \\
& \beta_{\mathrm{p}, \mathrm{p}^{\prime}}=\delta_{\mathrm{p}, \mathrm{p}^{\prime}} \beta_{\mathrm{n}} \tag{7.13}
\end{align*}
$$

with $\delta_{p, p^{\prime}}$ being the Kronecker's symbol.
Second, due to the vectorial character of the fields, the matrix form in eq.(7.12) has to be interpreted in a block form:

$$
\begin{align*}
& \left(\begin{array}{l}
{\left[D_{x}(z)\right]} \\
{\left[D_{y}(z)\right]} \\
{\left[D_{z}(z)\right]}
\end{array}\right)=\llbracket \varepsilon(z) \rrbracket\left(\begin{array}{l}
{\left[E_{x}(z)\right]} \\
{\left[E_{y}(z)\right]} \\
{\left[E_{z}(z)\right]}
\end{array}\right) \text {, isotropic media }  \tag{7.14}\\
& \left(\begin{array}{l}
{\left[\mathrm{D}_{\mathrm{x}}(\mathrm{z})\right]} \\
{\left[\mathrm{D}_{\mathrm{y}}(\mathrm{z})\right]} \\
{\left[\mathrm{D}_{\mathrm{z}}(\mathrm{z})\right]}
\end{array}\right)=\left(\begin{array}{lll}
\llbracket \varepsilon_{\mathrm{xx}}(\mathrm{z}) \rrbracket & \llbracket \varepsilon_{\mathrm{xy}}(\mathrm{z}) \rrbracket & \llbracket \varepsilon_{\mathrm{xz}}(\mathrm{z}) \rrbracket \\
\llbracket \varepsilon_{\mathrm{yx}}(\mathrm{z}) \rrbracket & \llbracket \varepsilon_{\mathrm{yy}}(\mathrm{z}) \rrbracket & \llbracket \varepsilon_{\mathrm{yz}}(\mathrm{z}) \rrbracket \\
\llbracket \varepsilon_{\mathrm{zx}}(\mathrm{z}) \rrbracket & \llbracket \varepsilon_{\mathrm{zy}}(\mathrm{z}) \rrbracket & \llbracket \varepsilon_{\mathrm{zz}}(\mathrm{z}) \rrbracket
\end{array}\right]\left(\begin{array}{l}
{\left[\mathrm{E}_{\mathrm{x}}(\mathrm{z})\right]} \\
{\left[\mathrm{E}_{\mathrm{y}}(\mathrm{z})\right]} \\
{\left[\mathrm{E}_{\mathrm{z}}(\mathrm{z})\right]}
\end{array}\right), \text { anisotropic media. } \tag{7.15}
\end{align*}
$$

The third problem linked with the truncation of eq. (7.8) has limited the use of the differential methods (including RCW method) for more than 30 years, and has been solved for lamellar gratings in the late 90s [7.11, 7.12], and for arbitrary-profile gratings in the start of the 2000s [7.13]. The problem is due to the very slow convergence with respect to the number of Fourier components in the truncated sum of eq. (7.8), when the two functions in the product are discontinuous. As demonstrated by Li [7.12], four different cases can be distinguished with respect to eq.(7.12):

1. Both $\varepsilon$ and $E$ are continuous functions of $x$ and $y$.
2. $\varepsilon$ is discontinuous, but E is continuous. This is the case of the tangential component of E.
3. Both $\varepsilon$ and E are discontinuous, but their product D is continuous, as it happens for the normal component of $D$.
4. All three functions are discontinuous.

In the first and second case, Laurent's rule assures relatively rapid convergence. In the third case, more rapidly converging scheme can be obtained through the following considerations for isotropic media.

If $\mathrm{D}=\varepsilon \mathrm{E}$ is continuous, then it is possible to factorize the product between D (continuous) and $1 / \varepsilon$ (discontinuous) using the Laurent's rule (called by Li direct rule):

$$
\begin{equation*}
[\overrightarrow{\mathrm{E}}(\mathrm{z})]=\llbracket \frac{1}{\varepsilon(\mathrm{z})} \rrbracket[\overrightarrow{\mathrm{D}}(\mathrm{z})], \tag{7.16}
\end{equation*}
$$

wherefrom the so called inverse rule is formulated:

$$
\begin{equation*}
[\overrightarrow{\mathrm{D}}(\mathrm{z})]=\llbracket \frac{1}{\varepsilon(\mathrm{z})} \|^{-1}[\overrightarrow{\mathrm{E}}(\mathrm{z})], \tag{7.17}
\end{equation*}
$$

which can be applied if the matrix $\llbracket \frac{1}{\varepsilon(\mathrm{z})} \rrbracket$ is not singular, a requirement that can create numerical problem for highly conducting gratings having small imaginary part of $\varepsilon$.


Fig.7.2.Convergence of the classical and the FFF version of the differential theory in the case of a dielectric sinusoidal grating with high contrast. Squares, old version of the differential theory for TM polarization; open triangles, new version, TM polarization; solid triangles, TE polarization (after [7.13]).

When infinite series are considered, eq.(7.17) is identical with eq.(7.12). However, as shown in Fig.7.2, the correct use of the direct or the inverse rules improves drastically the convergence of the differential methods with respect to the truncation parameter. Similarly to the abbreviation FFT, standing for Fast Fourier transformation, we have introduced the term Fast Fourier factorization (FFF) to name the correct use of the direct and the inverse rules, when applied numerically in the truncated Fourier space.

In the fourth case, neither the direct, nor the inverse rule result in acceptable convergence, so that this case must be avoided. Fortunately, this can be done by considering separately the electromagnetic field components, tangential or normal to the grating profile and taking into account that the electric field components tangential to the surface separating two different permittivities are continuous, in the same way as the normal components of the displacement D.

### 7.2. Differential theory for crossed gratings made of isotropic materials

In the isotropic case, the displacement vector $\overrightarrow{\mathrm{D}}$ can easily be separated into a continuous part $\overrightarrow{\mathrm{D}}_{\mathrm{N}}=\varepsilon \overrightarrow{\mathrm{E}}_{\mathrm{N}}$, normal to the profile surface, and $\overrightarrow{\mathrm{D}}_{\mathrm{T}}=\varepsilon \overrightarrow{\mathrm{E}}_{\mathrm{T}}$ that contains the continuous function $\overrightarrow{\mathrm{E}}_{\mathrm{T}}$. Let us define a unit vector $\overrightarrow{\mathrm{N}}$, normal to the grating profile. Although it is well defined on the profile (except edges), it is necessary to generalize its definition all over the grating region, which cannot be done in a unique manner. Different choices are shown further on for specific gratings having 1D or 2D periodicity. Using this generalized vector, the relations between $\vec{E}$ and $\vec{D}$ can be decomposed into two terms, for each of which we can apply the direct or the inverse factorization rules, skipping the explicit writing of the z -dependence:

$$
\begin{equation*}
\overrightarrow{\mathrm{D}}=\varepsilon \overrightarrow{\mathrm{E}}_{\mathrm{N}}+\varepsilon \overrightarrow{\mathrm{E}}_{\mathrm{T}}=\varepsilon \overrightarrow{\mathrm{N}}(\overrightarrow{\mathrm{~N}} \cdot \overrightarrow{\mathrm{E}})+\varepsilon[\overrightarrow{\mathrm{E}}-\overrightarrow{\mathrm{N}}(\overrightarrow{\mathrm{~N}} \cdot \overrightarrow{\mathrm{E}})] . \tag{7.18}
\end{equation*}
$$

The first term is a product of type 2 and requires the direct rule. The second term is of type 3 , demanding the inverse rule, so that:

$$
\begin{align*}
{[\overrightarrow{\mathrm{D}}] } & =\llbracket \varepsilon \rrbracket\left[\overrightarrow{\mathrm{E}}_{\mathrm{T}}\right]+\llbracket \frac{1}{\varepsilon} \rrbracket^{-1}\left[\overrightarrow{\mathrm{E}}_{\mathrm{N}}\right] \\
\quad= & \llbracket \varepsilon \rrbracket[\overrightarrow{\mathrm{E}}-\overrightarrow{\mathrm{N}}(\overrightarrow{\mathrm{~N}} \cdot \overrightarrow{\mathrm{E}})]+\llbracket \frac{1}{\varepsilon} \rrbracket^{-1}[\overrightarrow{\mathrm{~N}}(\overrightarrow{\mathrm{~N}} \cdot \overrightarrow{\mathrm{E}})] . \tag{7.19}
\end{align*}
$$

Introducing a square matrix representing a tensor product denoted $(\vec{N} \vec{N})$ with elements given by $\mathrm{N}_{\mathrm{i}} \mathrm{N}_{\mathrm{j}}$, we obtain:

$$
\begin{equation*}
\left.[\overrightarrow{\mathrm{D}}]=\llbracket \varepsilon \rrbracket[\overrightarrow{\mathrm{E}}]+\left(\llbracket \frac{1}{\varepsilon}\right]^{-1}-\llbracket \varepsilon \rrbracket\right) \llbracket \overrightarrow{\mathrm{N}} \overrightarrow{\mathrm{~N}} \rrbracket[\overrightarrow{\mathrm{E}}]=\mathrm{Q}_{\varepsilon}[\overrightarrow{\mathrm{E}}] \tag{7.20}
\end{equation*}
$$

where the matrix $\mathrm{Q}_{\varepsilon}$ has the form:

$$
\begin{equation*}
\left.\mathrm{Q}_{\varepsilon}=\llbracket \varepsilon \rrbracket+\left(\llbracket \frac{1}{\varepsilon}\right]^{-1}-\llbracket \varepsilon \rrbracket\right) \llbracket \overrightarrow{\mathrm{N}} \overrightarrow{\mathrm{~N}} \rrbracket . \tag{7.21}
\end{equation*}
$$

In a similar manner for magnetic materials, we can find the link between magnetic field and induction in the truncated Fourier space:

$$
\begin{gather*}
{[\overrightarrow{\mathrm{B}}]=\mathrm{Q}_{\mu}[\overrightarrow{\mathrm{H}}]} \\
\text { with } \quad \mathrm{Q}_{\mu}=\mathrm{Q}_{\varepsilon}=\llbracket \mu \rrbracket+\left(\llbracket \frac{1}{\mu} \rrbracket^{-1}-\llbracket \mu \rrbracket\right) \llbracket \overrightarrow{\mathrm{N}} \overrightarrow{\mathrm{~N}} \rrbracket \tag{7.22}
\end{gather*}
$$

Eq.(7.20) allows eliminating $\mathrm{E}_{\mathrm{z}}$ in the system (7.6):

$$
\begin{align*}
{\left[\mathrm{E}_{\mathrm{z}}\right] } & =\mathrm{Q}_{\varepsilon, z z}^{-1}\left(\left[\mathrm{D}_{\mathrm{z}}\right]-\mathrm{Q}_{\varepsilon, z \mathrm{x}}\left[\mathrm{E}_{\mathrm{x}}\right]-\mathrm{Q}_{\varepsilon, z \mathrm{y}}\left[\mathrm{E}_{\mathrm{y}}\right]\right) \\
& =-\mathrm{Q}_{\varepsilon, z z}^{-1}\left(\frac{\alpha\left[\mathrm{H}_{\mathrm{y}}\right]-\beta\left[\mathrm{H}_{\mathrm{x}}\right]}{\omega}+\mathrm{Q}_{\varepsilon, z \mathrm{x}}\left[\mathrm{E}_{\mathrm{x}}\right]+\mathrm{Q}_{\varepsilon, z \mathrm{z}}\left[\mathrm{E}_{\mathrm{y}}\right]\right) \tag{7.23}
\end{align*}
$$

where the matrices $\alpha$ and $\beta$ are defined in eq.(7.13).
Repeating the procedure for $\mathrm{H}_{\mathrm{z}}$ :

$$
\begin{equation*}
\left[\mathrm{H}_{\mathrm{z}}\right]=\mathrm{Q}_{\mu, z z}^{-1}\left(\frac{\alpha\left[\mathrm{E}_{\mathrm{y}}\right]-\beta\left[\mathrm{E}_{\mathrm{x}}\right]}{\omega}-\mathrm{Q}_{\mu, z \mathrm{x}}\left[\mathrm{H}_{\mathrm{x}}\right]-\mathrm{Q}_{\mu, \mathrm{zy}}\left[\mathrm{H}_{\mathrm{y}}\right]\right) \tag{7.24}
\end{equation*}
$$

it is also eliminated from eqs. (7.6).
For non-magnetic media, the last expression is further simplified:

$$
\begin{equation*}
\left[\mathrm{H}_{\mathrm{z}}\right]=\frac{\alpha\left[\mathrm{E}_{\mathrm{y}}\right]-\beta\left[\mathrm{E}_{\mathrm{x}}\right]}{\omega \mu_{0}} \tag{7.25}
\end{equation*}
$$

Thus the system (7.6) is replaced by a system of ordinary differential equations:

$$
\frac{\mathrm{d}}{\mathrm{dz}}\left(\begin{array}{l}
{\left[\mathrm{E}_{\mathrm{x}}\right]}  \tag{7.26}\\
{\left[\mathrm{E}_{\mathrm{y}}\right]} \\
{\left[\mathrm{H}_{\mathrm{x}}\right]} \\
{\left[\mathrm{H}_{\mathrm{y}}\right]}
\end{array}\right)=\mathrm{iM}\left(\begin{array}{l}
{\left[\mathrm{E}_{\mathrm{x}}\right]} \\
{\left[\mathrm{E}_{\mathrm{y}}\right]} \\
{\left[\mathrm{H}_{\mathrm{x}}\right]} \\
{\left[\mathrm{H}_{\mathrm{y}}\right]}
\end{array}\right) .
$$

This equation can be expressed in a compressed form:

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{dz}} \mathrm{~F}(\mathrm{z})=\mathrm{iM}(\mathrm{z}) \mathrm{F}(\mathrm{z}) \tag{7.27}
\end{equation*}
$$

Here the matrix $M$ has $4 \times 4$ blocks:

$$
\begin{align*}
& \mathrm{M}_{11}=-\alpha \mathrm{Q}_{\varepsilon, z z}^{-1} \mathrm{Q}_{\varepsilon, z x}-\mathrm{Q}_{\mu, \mathrm{yz}} \mathrm{Q}_{\mu, z z}^{-1} \beta \\
& \mathrm{M}_{12}=-\alpha \mathrm{Q}_{\varepsilon, z z}^{-1} \mathrm{Q}_{\varepsilon, z y}+\mathrm{Q}_{\mu, \mathrm{yz}} \mathrm{Q}_{\mu, z z}^{-1} \alpha \\
& \mathrm{M}_{13}=-\omega \mathrm{Q}_{\mu, \mathrm{yz}} \mathrm{Q}_{\mu, z z}^{-1} \mathrm{Q}_{\mu, \mathrm{zx}}+\frac{\alpha}{\omega} \mathrm{Q}_{\varepsilon, z z}^{-1} \beta+\omega \mathrm{Q}_{\mu, \mathrm{yx}} \\
& \mathrm{M}_{14}=-\omega \mathrm{Q}_{\mu, \mathrm{yz}} \mathrm{Q}_{\mu, z z}^{-1} \mathrm{Q}_{\mu, \mathrm{zy}}-\frac{\alpha}{\omega} \mathrm{Q}_{\varepsilon, z z}^{-1} \alpha+\omega \mathrm{Q}_{\mu, \mathrm{yy}} \\
& \mathrm{M}_{21}=-\beta \mathrm{Q}_{\varepsilon, z z}^{-1} \mathrm{Q}_{\varepsilon, z \mathrm{x}}+\mathrm{Q}_{\mu, \mathrm{xz}} \mathrm{Q}_{\mu, z z}^{-1} \beta \\
& M_{22}=-\beta Q_{\varepsilon, z z}^{-1} Q_{\varepsilon, z y}-Q_{\mu, x z} Q_{\mu, z z}^{-1} \alpha \\
& M_{23}=\omega Q_{\mu, x z} Q_{\mu, z z}^{-1} Q_{\mu, z x}+\frac{\beta}{\omega} Q_{\varepsilon, z z}^{-1} \beta-\omega Q_{\mu, x x} \\
& M_{24}=\omega Q_{\mu, x z} Q_{\mu, z z}^{-1} Q_{\mu, z y}-\frac{\beta}{\omega} Q_{\varepsilon, z z}^{-1} \alpha-\omega Q_{\mu, x y}  \tag{7.28}\\
& M_{31}=\omega \mathrm{Q}_{\varepsilon, y z} \mathrm{Q}_{\varepsilon, z z}^{-1} \mathrm{Q}_{\varepsilon, z \mathrm{x}}-\frac{\alpha}{\omega} \mathrm{Q}_{\mu, z z}^{-1} \beta-\omega \mathrm{Q}_{\varepsilon, y \mathrm{yx}} \\
& M_{32}=\omega Q_{\varepsilon, y z} Q_{\varepsilon, z z}^{-1} Q_{\varepsilon, z y}+\frac{\alpha}{\omega} Q_{\mu, z z}^{-1} \alpha-\omega Q_{\varepsilon, y y} \\
& \mathrm{M}_{33}=-\alpha \mathrm{Q}_{\mu, z z}^{-1} \mathrm{Q}_{\mu, z \mathrm{zx}}-\mathrm{Q}_{\varepsilon, \mathrm{yz}} \mathrm{Q}_{\varepsilon, z z}^{-1} \beta \\
& M_{34}=-\alpha Q_{\mu, z z}^{-1} Q_{\mu, z y}+Q_{\varepsilon, y z} Q_{\varepsilon, z z}^{-1} \alpha \\
& \mathrm{M}_{41}=-\omega \mathrm{Q}_{\varepsilon, \mathrm{xz}} \mathrm{Q}_{\varepsilon, z z}^{-1} \mathrm{Q}_{\varepsilon, z \mathrm{zx}}-\frac{\beta}{\omega} \mathrm{Q}_{\mu, z z}^{-1} \beta+\omega \mathrm{Q}_{\varepsilon, \mathrm{xx}} \\
& \mathrm{M}_{42}=-\omega \mathrm{Q}_{\varepsilon, \mathrm{xz}} \mathrm{Q}_{\varepsilon, z z}^{-1} \mathrm{Q}_{\varepsilon, z \mathrm{yy}}+\frac{\beta}{\omega} \mathrm{Q}_{\mu, z z}^{-1} \alpha+\omega \mathrm{Q}_{\varepsilon, \mathrm{xy}} \\
& M_{43}=-\beta Q_{\mu, z z}^{-1} Q_{\mu, z x}+Q_{\varepsilon, x z} Q_{\varepsilon, z z}^{-1} \beta \\
& M_{44}=-\beta \mathrm{Q}_{\mu, z z}^{-1} \mathrm{Q}_{\mu, z y}-\mathrm{Q}_{\varepsilon, \mathrm{xz}} \mathrm{Q}_{\varepsilon, z z}^{-1} \alpha \text {. }
\end{align*}
$$

This form looks like the form of the M-matrix obtained by Lifeng Li for crossed anisotropic (electrically and magnetically) gratings with profiles invariant with respect to z [7.17].

Whatever the form of the matrix M, eq.(7.26) represents a linear set of first-order ordinary differential equations. It can be solved numerically (with several problems, discussed further on), using well developed numerical schemes. In the case of vertical invariance of the optogeometrical parameters of the system inside the modulated region, the elements of the Mmatrix becomes constant in z , so that the solution of eq. (7.26) can be found through the eigenvectors and eigenvalues of M , a technique known under the name of Fourier modal method, or Rigorous coupled wave (RCW) method.

The solution of (7.26) gives a linear link between the field in the substrate and in the cladding

$$
\begin{equation*}
\mathrm{F}\left(\mathrm{z}_{\max }\right)=\mathrm{TF}\left(\mathrm{z}_{\min }\right), \tag{7.29}
\end{equation*}
$$

where T is called transmission matrix.
The advantage of this presentation comes from the fact that the field components participating in the calculations are tangential to the interfaces between the substrate and the modulated region, and between the cladding and the modulated region, so that they are continuous across these interfaces (in the absence of surface charges).

### 7.3. Electromagnetic field in the homogeneous regions - plane wave expansion

In most case, the substrate and cladding are homogeneous isotropic media. The electromagnetic field there can be expressed as a sum of plane waves. In particular, if the $x$ and $y$-dependencies are given as in eq.(7.7), the $z$-dependence is explicitly known, for example for the electric field it takes the form:

$$
\begin{equation*}
\overrightarrow{\mathrm{E}}_{\mathrm{p}}(\mathrm{z})=\overrightarrow{\mathrm{A}}_{\mathrm{p}}^{+} \exp \left(\mathrm{i} \gamma_{\mathrm{p}} \mathrm{z}\right)+\overrightarrow{\mathrm{A}}_{\mathrm{p}}^{-} \exp \left(-\mathrm{i} \gamma_{\mathrm{p}} \mathrm{z}\right) \tag{7.30}
\end{equation*}
$$

of two waves propagating upwards (sign + ) and downwards (sign - ) along the z -axis, with p given in eq.(7.9). Each diffraction order with a given p propagates independently of the others, the coupling is effective inside the grating region.

The z-propagation constant $\gamma$ depends on the medium properties:

$$
\begin{equation*}
\gamma_{\mathrm{p}}=\sqrt{\omega^{2} \varepsilon \mu-\alpha_{\mathrm{p}}^{2}-\beta_{\mathrm{p}}^{2}} \tag{7.31}
\end{equation*}
$$

Equations (7.6) enable us to express the magnetic field components through the electric ones:

$$
\begin{align*}
& \mathrm{H}_{\mathrm{x}, \mathrm{p}}=-\frac{1}{ \pm \gamma_{\mathrm{p}}}\left(\frac{\alpha_{\mathrm{p}} \beta_{\mathrm{p}}}{\omega \mu} \mathrm{E}_{\mathrm{x}, \mathrm{p}}+\frac{\beta_{\mathrm{p}}^{2}+\gamma_{\mathrm{p}}^{2}}{\omega \mu} \mathrm{E}_{\mathrm{y}, \mathrm{p}}\right) \\
& \mathrm{H}_{\mathrm{y}, \mathrm{p}}=\frac{1}{ \pm \gamma_{\mathrm{p}}}\left(\frac{\alpha_{\mathrm{p}}^{2}+\gamma_{\mathrm{p}}^{2}}{\omega \mu} \mathrm{E}_{\mathrm{x}, \mathrm{p}}+\frac{\alpha_{\mathrm{p}} \beta_{\mathrm{p}}}{\omega \mu} \mathrm{E}_{\mathrm{y}, \mathrm{p}}\right) \tag{7.32}
\end{align*}
$$

where the sign of $\gamma$ determines the direction of propagation in along z -axis.
With this link in mind, the column vector $F$ in eq.(7.27) takes the form:

$$
F \equiv\left(\begin{array}{l}
{\left[\mathrm{E}_{x}\right]}  \tag{7.33}\\
{\left[\mathrm{E}_{\mathrm{y}}\right]} \\
{\left[\mathrm{H}_{\mathrm{x}}\right]} \\
{\left[\mathrm{H}_{y}\right]}
\end{array}\right)=\Psi^{+} \mathrm{A}^{+}+\Psi^{-} \mathrm{A}^{-},
$$

where the column vectors

$$
\begin{equation*}
\mathrm{A}^{ \pm}=\binom{\left[\mathrm{A}_{\mathrm{x}}^{ \pm}\right]}{\left[\mathrm{A}_{\mathrm{y}}^{ \pm}\right]} \tag{7.34}
\end{equation*}
$$

contains the amplitudes of $\mathrm{E}_{\mathrm{x}}$ and $\mathrm{E}_{\mathrm{y}}$ propagating in positive or negative direction of the z axis, matrices $\Psi^{ \pm}$are block-diagonal:

$$
\Psi^{ \pm}=\left(\begin{array}{cc}
\mathbb{I} &  \tag{7.35}\\
& \mathbb{I} \\
\Psi_{\mathrm{xx}}^{ \pm} & \Psi_{\mathrm{xy}}^{ \pm} \\
\Psi_{\mathrm{yx}}^{ \pm} & \Psi_{\mathrm{yy}}^{ \pm}
\end{array}\right)
$$

with diagonal blocks

$$
\begin{align*}
& \mathbb{I}_{\mathrm{pp}}=1, \\
& \Psi_{\mathrm{xx}, \mathrm{pp}}^{ \pm}=\mp \frac{\alpha_{\mathrm{p}} \beta_{\mathrm{p}}}{\gamma_{\mathrm{p}} \omega \mu}, \quad \Psi_{\mathrm{xy}, \mathrm{pp}}^{ \pm}=\mp \frac{\beta_{\mathrm{p}}^{2}+\gamma_{\mathrm{p}}^{2}}{\gamma_{\mathrm{p}} \omega \mu}  \tag{7.36}\\
& \Psi_{\mathrm{yx}, \mathrm{pp}}^{ \pm}= \pm \frac{\alpha_{\mathrm{p}}^{2}+\gamma_{\mathrm{p}}^{2}}{\gamma_{\mathrm{p}} \omega \mu}, \quad \Psi_{\mathrm{yy}, \mathrm{pp}}^{ \pm}= \pm \frac{\alpha_{\mathrm{p}} \beta_{\mathrm{p}}}{\gamma_{\mathrm{p}} \omega \mu}
\end{align*}
$$

found from eq.(7.32)
Let us consider the case of a single incident wave from the cladding. The grating generates different diffraction order that propagate upwards in the cladding and downwards in the substrate. We attribute number 1 to the substrate and number 3 to the cladding. The total number of unknown diffracted field amplitudes will be equal to $4 \mathrm{P}_{\text {max }}$, two sets of $\mathrm{A}_{\mathrm{x}, \mathrm{p}}^{1-}$ and $A_{y, p}^{1-}$ transmitted in the substrate, and two sets of $A_{x, p}^{3+}$ and $A_{y, p}^{3+}$. These unknown amplitudes are subjected to $4 \mathrm{P}_{\text {max }}$ number of linear algebraic equation in (7.29).

In order to obtain the T-matrix, the numerical integration of eq.(7.26) is made by using the so-called shooting method, which consists of choosing $2 \mathrm{P}_{\text {max }}$ linearly independent representatives of the transmitted field. These representatives must correctly reflect the link between the electric and magnetic field components, as given by eqs.(7.32). A typical example for the shooting vectors starting from the substrate is that matrix $\Psi^{1+}$, which has $2 \mathrm{P}_{\text {max }}$ linearly independent columns. Here again, the number 1 indicates the substrate.

Thus the F column vector at $\mathrm{z}=\mathrm{z}_{\text {min }}$ can be formally written as a linear combination of the unknown amplitudes $\mathrm{A}^{1-}$ :

$$
\begin{equation*}
\tilde{\mathrm{F}}\left(\mathrm{z}_{\min }\right)=\Psi^{1-} \mathrm{A}^{1-}, \tag{7.37}
\end{equation*}
$$

Assuming that there is no incidence from the substrate side. Here the tilde indicates that the vector $F$ is not yet the true solution of the diffraction problem.

The result of the numerical integration from $z_{\min }$ to $z_{\max }$ will provide the values of $\tilde{F}$ at $z=z_{\text {max }}$, which are also a linear combination $\tilde{F}\left(z_{\max }\right) A^{1-}$ of $A^{1-}$, due to the linearity of the problem. On the other side, the column vector F at the upper interface is equal to $\psi^{3+} \mathrm{A}^{3+}+\psi^{3-} \mathrm{A}^{3-}$, according eq.(7.33), thus a linear set of algebraic equations for the unknown amplitudes $A^{1-}$ and $A^{3+}$ is obtained, with the free part determined by the wave incident from the cladding side:

$$
\begin{equation*}
\psi^{3+} \mathrm{A}^{3+}+\psi^{3-} \mathrm{A}^{3-}=\tilde{\mathrm{F}}\left(\mathrm{z}_{\max }\right) \mathrm{A}^{1+} . \tag{7.38}
\end{equation*}
$$

Once this system is solved, all field components can be calculated.

Unfortunately, this simple procedure creates enormous numerical problems that can be explained by using two different arguments:

First, it is known (but no quite well) in the theory of systems of ordinary differential equations, that numerical integration could become instable after a specific integration length, due to the fact that the set of shooting vectors can lose its linear independence during the integration. In other words, if the initial choice covers a vector space of $2 \mathrm{P}_{\max }$ dimensions, this space could shrink during the numerical integration to reduce its dimensions, so that the final algebraic system (7.38) could become singular. A solution of the problem based on this understanding was proposed in 1990 by G. Tayeb by using intermediate orthonormalization procedures during the numerical integration.

The second argument is based on the fact that inside the modulated region, as well as in the homogeneous regions, electromagnetic field contains components that propagate both in the positive and in the negative z-direction. During the integration, they both are treated in the same manner. As far as the solution requires taking into account the evanescent orders in addition to the propagating ones, a part of the former grows exponentially in z-direction, while the other part decreases exponentially. Due to the limited length of computer words, the ones that decrease substantially will be lost with respect to the ones that grow rapidly, even if the former could bring physical information. During the 90 s , several different algorithms were proposed for solving the problem, based on a different treatment of the diffraction orders propagating upwards and downwards [7.9, 7.10]. Among them, to so called S-matrix propagation algorithm [7.10c] is probably the easiest to implement. Moreover, it can be used with methods other than the differential one in, for example, treating a stack of layers by the integral method, or by methods based on a transformation of the coordinate system. Interested reader can find in Appendix 7.1 a brief description of the S-matrix algorithm.

### 7.4. Several simpler isotropic cases

In practice most applications use non-magnetic materials, for which the form of M-matrix is considerably simplified, taking into account that then $Q_{\mu}$ is diagonal and equal to $\mu_{0}$. Furthermore, several specific cases are of great interest for application, and they lead to a further simplification of the M-matrix.

### 7.4.1. Classical grating with one-dimensional periodicity, example of a sinusoidal profile

Let us consider a classical grating with grooves parallel to the $y$-axis and surface profile given by the equation $\mathrm{z}=\mathrm{g}(\mathrm{x})$. The vector normal to the surface is given by

$$
\begin{align*}
& \overrightarrow{\mathrm{N}}=\frac{1}{\sqrt{1+\mathrm{g}^{\prime 2}(\mathrm{x})}}\left(-\mathrm{g}^{\prime}(\mathrm{x}), 0,1\right), \quad \text { if } \mathrm{g}^{\prime}(\mathrm{x}) \text { exists, }  \tag{7.39}\\
& \overrightarrow{\mathrm{N}}=(1,0,0), \quad \text { if not }
\end{align*}
$$

where the prime stands for a derivative with respect to $x$. In case of vertical walls $\overrightarrow{\mathrm{N}}=(1,0,0)$. Thus the easiest way to generalize the normal vector to the entire modulated region is just to make it equal to eq.(7.39) not only on the profile $z=g(x)$, but everywhere inside the grating region for $\min [g(x)] \leq z \leq \max [g(x)]$. The advantage of this choice is that $\overrightarrow{\mathrm{N}}$ does not depend on z , and the Fourier transformation of the tensor $\overrightarrow{\mathrm{N}} \overrightarrow{\mathrm{N}}$ is done only once.

If the derivative of the profile function does not exist, or if the function is a multivalued one (e.g., circular or elliptical rods), but the interface can be expressed as a two-variable function:

$$
\begin{equation*}
\mathrm{g}(\mathrm{x}, \mathrm{y})=0 \tag{7.40}
\end{equation*}
$$

the normal vector is easily defined as the gradient of the profile function $\overrightarrow{\mathrm{N}}=\operatorname{grad}[\mathrm{g}(\mathrm{x}, \mathrm{z})] /[\operatorname{grad}[\mathrm{g}(\mathrm{x}, \mathrm{z})]$.

The $\mathrm{Q}_{\varepsilon}$ matrix takes the form:

$$
\mathrm{Q}_{\varepsilon}=\left(\begin{array}{ccc}
\llbracket \varepsilon \rrbracket \llbracket \mathrm{N}_{\mathrm{z}}^{2} \rrbracket+\llbracket \frac{1}{\varepsilon} \|^{-1} \llbracket \mathrm{~N}_{\mathrm{x}}^{2} \rrbracket & 0 & \left(\llbracket \frac{1}{\varepsilon} \|^{-1}-\llbracket \varepsilon \rrbracket\right) \llbracket \mathrm{N}_{\mathrm{x}} \mathrm{~N}_{\mathrm{z}} \rrbracket  \tag{7.41}\\
0 & \llbracket \varepsilon \rrbracket & 0 \\
\left(\llbracket \frac{1}{\varepsilon} \|^{-1}-\llbracket \varepsilon \rrbracket\right) \llbracket \mathrm{N}_{\mathrm{x}} \mathrm{~N}_{\mathrm{z}} \rrbracket & 0 & \llbracket \varepsilon \rrbracket \llbracket \mathrm{~N}_{\mathrm{x}}^{2} \rrbracket+\llbracket \frac{1}{\varepsilon} \rrbracket^{-1} \llbracket \mathrm{~N}_{\mathrm{z}}^{2} \rrbracket
\end{array}\right)
$$

where it is taken into account that $\mathrm{N}_{\mathrm{x}}^{2}+\mathrm{N}_{\mathrm{z}}^{2}=1$. The fact that the normal vector components participate in the form of products in couples is important, because it leads to the conclusion is that the choice of the sign of $\overrightarrow{\mathrm{N}}$ plays no role.

Further simplification of the M-matrix comes if limited to non-conical diffraction with $\beta_{0}=0$ :

$$
\mathrm{M}=\left(\begin{array}{cccc}
-\alpha \mathrm{Q}_{\varepsilon, z z}^{-1} \mathrm{Q}_{\varepsilon, z \mathrm{zx}} & 0 & 0 & -\frac{\alpha}{\omega} \mathrm{Q}_{\varepsilon, z z}^{-1} \alpha+\omega \mu_{0} \mathbb{I}  \tag{7.42}\\
0 & 0 & -\omega \mu_{0} \mathbb{I} & 0 \\
0 & \frac{\alpha \alpha}{\omega \mu_{0}}-\omega \llbracket \varepsilon \rrbracket & 0 & 0 \\
-\omega \mathrm{Q}_{\varepsilon, x z} \mathrm{Q}_{\varepsilon, z z}^{-1} \mathrm{Q}_{\varepsilon, z \mathrm{zx}}+\omega \mathrm{Q}_{\varepsilon, x x} & 0 & 0 & -\mathrm{Q}_{\varepsilon, x z} \mathrm{Q}_{\varepsilon, z z}^{-1} \alpha
\end{array}\right)
$$

This shows that the system to integrate decouples into two subsystems, corresponding to the two fundamental polarizations, transversal with respect to the plane of incidence, transverse electric (TE):

$$
\begin{align*}
& \frac{\mathrm{d}}{\mathrm{dz}}\left[\mathrm{E}_{\mathrm{y}}\right]=-\mathrm{i} \omega \mu_{0}\left[\mathrm{H}_{\mathrm{x}}\right] \\
& \frac{\mathrm{d}}{\mathrm{dz}}\left[\mathrm{H}_{\mathrm{x}}\right]=\mathrm{i}\left(\frac{\alpha^{2}}{\omega \mu_{0}}-\omega \llbracket \varepsilon \rrbracket\right)\left[\mathrm{E}_{\mathrm{y}}\right] \tag{7.43}
\end{align*}
$$

and transverse magnetic (TM):

$$
\begin{align*}
& \frac{\mathrm{d}}{\mathrm{dz}}\left[\mathrm{E}_{\mathrm{x}}\right]=-\mathrm{i} \alpha \mathrm{Q}_{\varepsilon, z z}^{-1} \mathrm{Q}_{\varepsilon, z \mathrm{z}}\left[\mathrm{E}_{\mathrm{x}}\right]-\mathrm{i}\left(\frac{\alpha}{\omega} \mathrm{Q}_{\varepsilon, z z}^{-1} \alpha-\omega \mu_{0}\right)\left[\mathrm{H}_{\mathrm{y}}\right]  \tag{7.44}\\
& \frac{\mathrm{d}}{\mathrm{dz}}\left[\mathrm{H}_{\mathrm{y}}\right]=\mathrm{i} \omega\left(\mathrm{Q}_{\varepsilon, \mathrm{xx}}-\mathrm{Q}_{\varepsilon, \mathrm{xz}} \mathrm{Q}_{\varepsilon, z z}^{-1} \mathrm{Q}_{\varepsilon, z \mathrm{x}}\right)\left[\mathrm{E}_{\mathrm{x}}\right]-\mathrm{iQ}_{\varepsilon, \mathrm{xz}} \mathrm{Q}_{\varepsilon, z z}^{-1} \alpha\left[\mathrm{H}_{\mathrm{y}}\right]
\end{align*}
$$

### 7.4.1.1. Fourier transformation of the permittivity

The set of ordinary differential equations to be integrated contains the Fourier transforms of $\varepsilon$, $1 / \varepsilon, \mu, 1 / \mu, N_{x}^{2}$, and $N_{z}^{2}$. In general, Fast Fourier transform (FFT) techniques can be easily applied. As already discussed with respect to eq.(7.39), the normal vector components must be transformed only once, if chosen to be independent on z . On the other hand, the permittivity and permeability depend on z and their Fourier components have to be calculated for each value of $z$ during the numerical integration. Fortunately, in the 1D case, it is possible and recommended to use analytical formulae for the Fourier transforms of $\varepsilon, 1 / \varepsilon, \mu, 1 / \mu$, which give faster more accurate results. This can be done because for a given value of $z$, they are piecewise constant functions of $y$. Fig. 7.3 presents schematically a grating with a period $d$ that separates two homogeneous media with permittivities $\varepsilon_{1}$ and $\varepsilon_{3}$. For a given value $\mathrm{z}_{0}$ of z , the Fourier transform of, for example, the permittivity intside the modulated region $0 \leq \mathrm{z} \leq \mathrm{h}$ is given by

$$
\begin{align*}
\varepsilon_{\mathrm{m}} & =\frac{\varepsilon_{1}}{\mathrm{~d}} \int_{\mathrm{x}_{1}}^{\mathrm{x}_{2}} \mathrm{e}^{-\mathrm{imK}_{\mathrm{x}}} \mathrm{dx}+\frac{\varepsilon_{3}}{\mathrm{~d}} \int_{\mathrm{x}_{2}}^{\mathrm{d}+\mathrm{x}_{1}} \mathrm{e}^{-\mathrm{imK}_{x} \mathrm{x}} \mathrm{dx} \\
& =\left(\varepsilon_{1}-\varepsilon_{3}\right) \frac{\sin \left(\mathrm{mK}_{\mathrm{x}} \frac{\mathrm{x}_{2}-\mathrm{x}_{1}}{2}\right)}{\pi \mathrm{m}} \mathrm{e}^{-\mathrm{imk}_{\mathrm{x}} \frac{\mathrm{x}_{1}+\mathrm{x}_{2}}{2}}+\varepsilon_{3} \delta_{\mathrm{m}, 0} \tag{7.45}
\end{align*}
$$

so that the two integrals can be solved analytically, once $x_{1}$ and $x_{2}$ are determined from the inverse of $g(x)$ :

$$
\begin{equation*}
\mathrm{x}_{1,2}=\mathrm{g}^{-1}\left(\mathrm{z}_{0}\right) . \tag{7.46}
\end{equation*}
$$

If the inverse of $\mathrm{g}(\mathrm{x})$ has more than two solutions, the sum of integrals (7.45) will contain several more terms. The same equations can be used to obtain the Fourier transforms of the inverse of the permittivity.


Fig.7.3.Piecewise constant representation of the permittivity for a one-dimensional grating

In the case of a sinusoidal profile:

$$
\begin{equation*}
\mathrm{z}=\frac{\mathrm{h}}{2}\left[1+\sin \left(\mathrm{K}_{\mathrm{x}} \mathrm{x}\right)\right] \Rightarrow \mathrm{x}_{1,2}=\arcsin \left(\frac{2 \mathrm{z}_{0}}{\mathrm{~h}}-1\right) . \tag{7.47}
\end{equation*}
$$

### 7.4.1.2. Fourier transformation of the normal vector

As already explained, the Fourier transformation of the normal vector requires its continuation all over the space. If the grating profile can be represented as a single-value function, we can use eq.(7.39) for $\overrightarrow{\mathrm{N}}$ and calculate the Fourier components of the tensor $\overrightarrow{\mathrm{N}} \overrightarrow{\mathrm{N}}^{\mathrm{T}}$ by use of the Fast Fourier transform (FFT) technique once for all z-values. For a sinusoidal grating having a profile defined in eq.(7.47), the normal vector takes the form:

$$
\begin{equation*}
\overrightarrow{\mathrm{N}}=\frac{1}{\sqrt{1+\mathrm{g}^{\prime 2}(\mathrm{x})}}\left(-\mathrm{g}^{\prime}(\mathrm{x}), 0,1\right)=\frac{\left(-\frac{\pi \mathrm{h}}{\mathrm{~d}} \cos \left(\mathrm{~K}_{\mathrm{x}} \mathrm{x}\right), 0,1\right)}{\sqrt{1+\left(\frac{\pi \mathrm{h}}{\mathrm{~d}}\right)^{2} \cos ^{2}\left(\mathrm{~K}_{\mathrm{x}} \mathrm{x}\right)}} \tag{7.48}
\end{equation*}
$$

### 7.4.2. Classical isotropic trapezoidal or triangular grating

A trapezoidal grating is shown schematically in Fig.7.4 with two flat regions $L$ at the top and the bottom of the groove and two different, in general, groove angles $\psi$. The Fourier transform of the permittivity and its inverse are calculated using eq.(7.45) with:

$$
\begin{align*}
& \mathrm{x}_{1}=\mathrm{z}_{0} \operatorname{cotg} \psi_{1}  \tag{7.49}\\
& \mathrm{x}_{2}=\mathrm{x}_{\mathrm{C}}-\mathrm{z}_{0} \operatorname{cotg} \psi_{2}
\end{align*}
$$

with $x_{C}=d-L_{2}$. For the normal vector, the period can be divided in four regions $A$ to $D$, as shown in the figure:

$$
\left.\begin{array}{l}
\left.\left.\begin{array}{l}
\mathrm{N}_{\mathrm{y}}=0 \\
\mathrm{~N}_{\mathrm{x}}=\sin \psi_{1} \\
\mathrm{~N}_{\mathrm{z}}=-\cos \psi_{1}
\end{array}\right\} \quad \text { in } \mathrm{A}, \quad \begin{array}{l}
\mathrm{N}_{\mathrm{x}}=\sin \psi_{2} \\
\mathrm{~N}_{\mathrm{z}}=\cos \psi_{2}
\end{array}\right\} \quad \text { in } \mathrm{C}, \\
\mathrm{~N}_{\mathrm{x}}=0
\end{array}\right\} \text { in B and D, } \quad \text {, }
$$

Their Fourier components do not depend on z and can be represented as a sum of several analytical terms, similar to eq.(7.45):


Fig.7.4. Trapezoidal profile with parameters. The normal vector direction is given in red arrows.

$$
\begin{align*}
& \left(N_{x}^{2}\right)_{m}=\frac{\sin ^{2} \psi_{1}}{d} \int_{0}^{h \operatorname{cotg} \psi_{1}} e^{-i m K_{x} x} d x+\frac{1}{d} \int_{h \operatorname{cotg} \psi_{1}}^{h \operatorname{cotg} \psi_{1}+L_{1}} e^{-i m K_{x} x} d x \\
& \quad+\frac{\sin ^{2} \psi_{2}}{d} \int_{h \operatorname{cotg} \psi_{1}+L_{1}}^{d-L_{2}} e^{-i m K_{x} x} d x+\frac{1}{d} \int_{d-L_{2}}^{d} e^{-i m K_{x} x} d x  \tag{7.51}\\
& \left(N_{y}^{2}\right)_{m}=\frac{\cos ^{2} \psi_{1}}{d} \int_{0}^{h \operatorname{cotg} \psi_{1}} e^{-i m K_{x} x} d x+\frac{\cos ^{2} \psi_{2}}{d} \int_{h \operatorname{cotg} \psi_{1}+L_{1}}^{d-L_{2}} e^{-i m K_{x} x} d x  \tag{7.52}\\
& \left(N_{x} N_{y}\right)_{m}=-\frac{\sin 2 \psi_{1}}{2 d} \int_{0}^{h \operatorname{cotg} \psi_{1}} e^{-i m K_{x} x} d x+\frac{\sin 2 \psi_{2}}{2 d} \int_{h \operatorname{cotg} \psi_{1}+L_{1}}^{d-L_{2}} e^{-i m K_{x} x} d x . \tag{7.53}
\end{align*}
$$

A triangular-groove grating can be considered as a particular case of a trapezoidal profile with no flat regions, $\mathrm{L}_{1}=\mathrm{L}_{2}=0, \mathrm{x}_{\mathrm{C}}=\mathrm{d}$. Moreover, the profile given in Fig.7.4 also includes the case with vertical facets, and some more exotic profiles with hanging back walls, Fig.7.5.


Fig.7.5. Two different profiles with slanted grooves

### 7.4.3. Classical lamellar grating

Lamellar profile with vertical walls is most easy to treat, because the normal to the profile vector has only one non-zero component, $\mathrm{N}_{\mathrm{x}}=1$. The $\mathrm{Q}_{\varepsilon}$ matrix takes the form:

$$
\left.\begin{array}{c}
\mathrm{Q}_{\varepsilon}=\left(\begin{array}{ccc}
\llbracket \frac{1}{\varepsilon} \|^{-1} & 0 & 0 \\
0 & \llbracket \varepsilon \rrbracket & 0 \\
0 & 0 & \llbracket \varepsilon \rrbracket
\end{array}\right) \\
\mathrm{M}=\left(\begin{array}{ccc}
0 & 0 & \frac{\alpha}{\omega} \llbracket \varepsilon \rrbracket^{-1} \beta_{0} \\
0 & -\frac{\alpha}{\omega} \llbracket \varepsilon \rrbracket^{-1} \alpha+\omega \mu_{0} \mathbb{I} \\
-\frac{\alpha}{\omega \mu_{0}} \beta_{0} & \frac{\beta^{2}}{\omega \mu_{0}}-\omega \llbracket \varepsilon \rrbracket^{-1}-\omega \mu_{0} \mathbb{I} & -\frac{\beta_{0}}{\omega} \llbracket \varepsilon \rrbracket^{-1} \alpha \\
\omega \llbracket \frac{1}{\omega} \|^{-1}-\frac{\beta_{0}^{2}}{\omega \mu_{0}} \mathbb{I} & \frac{\beta_{0}}{\omega \mu_{0}} \alpha & 0
\end{array}\right.  \tag{7.55}\\
0
\end{array}\right)
$$

In non-conical diffraction, when $\beta_{0}=0$, the two fundamental polarizations are decoupled and can be solved independently of each other. The M-matrix is simplified to obtain an antidiagonal block form:

$$
M=\left(\begin{array}{cccc}
0 & 0 & 0 & -\frac{\alpha}{\omega} \llbracket \varepsilon \rrbracket^{-1} \alpha+\omega \mu_{0} \mathbb{I}  \tag{7.56}\\
0 & 0 & -\omega \mu_{0} \mathbb{I} & 0 \\
0 & \frac{\alpha^{2}}{\omega \mu_{0}}-\omega \llbracket \varepsilon \rrbracket & 0 & 0 \\
\omega \llbracket \frac{1}{\varepsilon} \|^{-1} & 0 & 0 & 0
\end{array}\right)
$$

thus the two sets of differential equations for each polarization become:

$$
\begin{align*}
& \frac{\mathrm{d}}{\mathrm{dz}}\left[\mathrm{E}_{\mathrm{x}}\right]=\mathrm{i}\left(\omega \mu_{0} \mathbb{I}-\frac{\alpha}{\omega} \llbracket \varepsilon \rrbracket^{-1} \alpha\right)\left[\mathrm{H}_{\mathrm{y}}\right]  \tag{7.57}\\
& \frac{\mathrm{d}}{\mathrm{dz}}\left[\mathrm{H}_{\mathrm{y}}\right]=\mathrm{i} \omega \llbracket \frac{1}{\varepsilon} \|^{-1}\left[\mathrm{E}_{\mathrm{x}}\right]
\end{align*}
$$

and

$$
\begin{align*}
& \frac{d}{d z}\left[E_{y}\right]=-i \omega \mu_{0}\left[H_{x}\right] \\
& \frac{d}{d z}\left[H_{x}\right]=i\left(\frac{\alpha^{2}}{\omega \mu_{0}}-\omega \llbracket \varepsilon \rrbracket\right)\left[E_{y}\right] \tag{7.58}
\end{align*}
$$

Even in the case of conical diffraction, it is possible to define two other polarizations, for which the differential system decouples. These are the electric and magnetic polarizations that are transverse with respect to the $x$-axis. Let us denote the two polarization with superscripts (e), when $E_{x}=0$, and (h), when $H_{x}=0$. For (e) case, it is possible to express $H_{y}$ as a function of $\mathrm{H}_{\mathrm{x}}$ from eq.(7.26) and the first line of the M-matrix in eq.(7.55):

$$
\begin{equation*}
\left[\mathrm{H}_{\mathrm{y}}^{(\mathrm{e})}\right]=-\left(\omega \mu_{0} \mathbb{I}-\frac{\alpha}{\omega} \llbracket \varepsilon \rrbracket^{-1} \alpha\right)^{-1} \frac{\alpha}{\omega} \llbracket \varepsilon \rrbracket^{-1} \beta_{0}\left[\mathrm{H}_{\mathrm{x}}^{(\mathrm{e})}\right] \tag{7.59}
\end{equation*}
$$

which can be simplified into:

$$
\begin{equation*}
\left[\mathrm{H}_{\mathrm{y}}^{(\mathrm{e})}\right]=-\alpha\left(\omega^{2} \mu_{0} \llbracket \varepsilon \rrbracket-\alpha^{2}\right)^{-1} \beta_{0}\left[\mathrm{H}_{\mathrm{x}}^{(\mathrm{e})}\right] \tag{7.60}
\end{equation*}
$$

The second line of the matrix $M$ then results in:

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{dz}}\left[\mathrm{E}_{\mathrm{y}}^{(\mathrm{e})}\right]=\mathrm{i}\left[\frac{\beta_{0}^{2}}{\omega} \llbracket \varepsilon \rrbracket^{-1}-\omega \mu_{0} \mathbb{I}+\frac{\beta_{0}^{2}}{\omega} \llbracket \varepsilon \rrbracket^{-1} \alpha^{2}\left(\omega^{2} \mu_{0} \llbracket \varepsilon \rrbracket-\alpha^{2}\right)^{-1}\right]\left[\mathrm{H}_{\mathrm{x}}^{(\mathrm{e})}\right] \tag{7.61}
\end{equation*}
$$

This expression can be further simplified, and together with the third line of eq.(7.55) (when $E_{x}=0$ ) gives a set of equation for (e) polarization:

$$
\begin{align*}
& \frac{\mathrm{d}}{\mathrm{dz}}\left[\mathrm{E}_{\mathrm{y}}^{(\mathrm{e})}\right]=\mathrm{i} \omega \mu_{0}\left[\beta_{0}^{2}\left(\omega^{2} \mu_{0} \llbracket \varepsilon \rrbracket-\alpha^{2}\right)^{-1}-\mathbb{I}\right]\left[\mathrm{H}_{\mathrm{x}}^{(\mathrm{e})}\right]  \tag{7.62}\\
& \frac{\mathrm{d}}{\mathrm{dz}}\left[\mathrm{H}_{\mathrm{x}}^{(\mathrm{e})}\right]=\frac{\mathrm{i}}{\omega \mu_{0}}\left(\alpha^{2}-\omega^{2} \mu_{0} \llbracket \varepsilon \rrbracket\right)\left[\mathrm{E}_{\mathrm{y}}^{(\mathrm{e})}\right]
\end{align*}
$$

Similar procedure for (h) case when $\mathrm{H}_{\mathrm{x}}^{(\mathrm{h})}=0$, result in another system of differential equations, decoupled from the (e) case:

$$
\begin{align*}
& \left.\frac{\mathrm{d}}{\mathrm{dz}}\left[\mathrm{H}_{\mathrm{y}}^{(\mathrm{h})}\right]=\mathrm{i} \omega\left[\llbracket \frac{1}{\varepsilon}\right]^{-1}+\beta_{0}^{2}\left(\alpha \llbracket \varepsilon \rrbracket^{-1} \alpha-\omega^{2} \mu_{0} \mathbb{I}\right)^{-1}\right]\left[\mathrm{E}_{\mathrm{x}}^{(\mathrm{h})}\right]  \tag{7.63}\\
& \frac{\mathrm{d}}{\mathrm{dz}}\left[\mathrm{E}_{\mathrm{x}}^{(\mathrm{h})}\right]=\frac{\mathrm{i}}{\omega}\left(\omega^{2} \mu_{0} \mathbb{I}-\alpha \llbracket \varepsilon \rrbracket^{-1} \alpha\right)\left[\mathrm{H}_{\mathrm{y}}^{(\mathrm{h})}\right]
\end{align*}
$$

In non-conical mount, $\beta_{0}=0$ and eqs.(7.62) and (7.63) become equivalent to eqs.(7.58) and (7.57).

Both conical and nonconical cases of diffraction by lamellar gratings are solved by eigenvector technique, due to the fact that the coefficients of the differential equations are $z$ independent. Moreover, due to the separation of the two fundamental polarizations, it is possible to further reduce by half the size of the matrices, by dealing with second-order differential equations. For example, eq. (7.57) can be written in the form:

$$
\begin{align*}
& \frac{\mathrm{d}}{\mathrm{dz}}\left[\mathrm{E}_{\mathrm{x}}\right]=\mathrm{iM}_{14}\left[\mathrm{H}_{\mathrm{y}}\right]  \tag{7.64}\\
& \frac{\mathrm{d}}{\mathrm{dz}}\left[\mathrm{H}_{\mathrm{y}}\right]=\mathrm{iM}_{41}\left[\mathrm{E}_{\mathrm{x}}\right]
\end{align*}
$$

Thus

$$
\begin{align*}
& \frac{\mathrm{d}^{2}}{\mathrm{dz}^{2}}\left[\mathrm{E}_{\mathrm{x}}\right]=-\mathrm{M}_{14} \mathrm{M}_{41}\left[\mathrm{E}_{\mathrm{x}}\right] \\
& \frac{\mathrm{d}^{2}}{\mathrm{dz}^{2}}\left[\mathrm{H}_{\mathrm{y}}\right]=-\mathrm{M}_{41} \mathrm{M}_{14}\left[\mathrm{H}_{\mathrm{y}}\right] \tag{7.65}
\end{align*}
$$

Les us denote with $\rho_{p}^{2}$ the eigenvalues of the product $M_{14} M_{41}$ and with $V$ the matrix with its eigenvectors arranged in columns. The solution of the first eq.(7.65) can be written as:

$$
\begin{equation*}
\left[\mathrm{E}_{\mathrm{x}}(\mathrm{z})\right]=\mathrm{V} \Phi(\mathrm{z}) \mathrm{V}^{-1}\left[\mathrm{E}_{\mathrm{x}}(0)\right] \tag{7.66}
\end{equation*}
$$

with

$$
\begin{equation*}
\Phi_{\mathrm{pp}}(\mathrm{z})=\delta_{\mathrm{pp}}, \exp \left( \pm \mathrm{i}_{\mathrm{p}} \mathrm{z}\right) \tag{7.67}
\end{equation*}
$$

which shows that the elementary solutions along z (called modes, wherefrom the names Fourier modal method or Rigorous coupled waves method) exist in pairs that can propagate upwards or downwards with the same propagation constants.

By integrating the second eq.(7.65), we obtain that:

$$
\begin{equation*}
\left[\mathrm{H}_{\mathrm{y}}(\mathrm{z})\right]=\mathrm{W} \Phi(\mathrm{z}) \mathrm{W}^{-1}\left[\mathrm{H}_{\mathrm{y}}(0)\right] \tag{7.68}
\end{equation*}
$$

with W that can be written in different forms, because the eigenvectors are defined within an arbitrary factor. For example, if we take into account the second eq.(7.64), $\mathrm{W}=\mp \mathrm{i} \rho^{-1} \mathrm{M}_{41} \mathrm{~V}$, where the diagonal matrix $\rho$ has elements equal to $\rho_{\mathrm{p}}$. Another possibility, that is quite convenient in TM polarization described by eq.(7.57), is at first to calculate the eigenvectors of $\mathrm{M}_{41} \mathrm{M}_{14}$, instead of $\mathrm{M}_{14} \mathrm{M}_{41}$ (their eigenvalues are the same). Then the link between V and W contains the inverse of $\mathrm{M}_{41}$, which is just equal to $\frac{1}{\omega} \llbracket \frac{1}{\varepsilon} \rrbracket$, so that $\mathrm{W}= \pm \frac{1}{\omega} \llbracket \frac{1}{\varepsilon} \rrbracket \mathrm{~V} \rho$.

### 7.4.4. Crossed grating having vertical walls made of isotropic material

Most of the recent applications of the Fourier modal method are devoted to studies of light diffraction by structures with 2D periodicity and piecewise invariant in the third direction. This popularity has several reasons. First, extraordinary light transmission was found in the late 90s by Ebbesen [7.18] for such structures, namely metallic sheets with periodical hole arrays, and it attracted a lot of attention (see Chapter 1). Second, the technology of such structures has significantly advanced in the last 20 years. Third, the Fourier modal method is relatively simple to implement, and much faster than most of the other methods, because of the eigenvalue/vector technique of integration.

Detailed study of these structures will be described in a separate chapter. However, due to its importance, we are discussing different aspects of this theory, as it presents a particular case of the more general geometry, that is characterized by a constant value of the zcomponent of the normal vector on each cross-section having $\mathrm{z}=$ const. The prolongation of the normal vector within the entire grating cell is discussed in sections 7.6.2.2.

### 7.5. Differential theory for anisotropic media

If we consider anisotropic media that do not extend inside the grating structure, there is not necessary to reformulate the diffraction theory, only that in the general case it is not possible to separate the problem into two independent polarizations, and it is necessary to work with the complete 4Pmax vectors and matrices.

In the case of anisotropic medium that lies inside the grating, the equations linking the M-matrix with the $\mathrm{Q}_{\varepsilon}$ and $\mathrm{Q}_{\mu}$ matrices remain the same, eq.(7.28) for isotropic and anisotropic media. The difference comes from the fact that the Q-matrices take more complex form, because the link between the normal and tangential components of the couples E and D and H and B is made through the tensors of permittivity and permeability, which are not scalars. Let us establish this link in detail for E and D. As far as the continuous and discontinuous field components must be factorized differently, we construct a column vector $\mathrm{F}_{\varepsilon}$, which contains the continuous field components $\mathrm{E}_{\mathrm{T}}$ and $\mathrm{D}_{\mathrm{N}}$. There are two tangential components to the grating surface, and only a single normal:

$$
\mathrm{F}_{\varepsilon}=\left(\begin{array}{l}
\mathrm{D}_{\mathrm{N}}  \tag{7.69}\\
\mathrm{E}_{\mathrm{T}_{1}} \\
\mathrm{E}_{\mathrm{T}_{2}}
\end{array}\right)=\left(\begin{array}{c}
\overrightarrow{\mathrm{N}} \cdot(\overline{\bar{\varepsilon}} \overrightarrow{\mathrm{E}}) \\
\overrightarrow{\mathrm{T}_{1}} \cdot \overrightarrow{\mathrm{E}} \\
\overrightarrow{\mathrm{~T}}_{2} \cdot \overrightarrow{\mathrm{E}}
\end{array}\right)=\mathrm{U}_{\varepsilon}\left(\begin{array}{c}
\mathrm{E}_{\mathrm{x}} \\
\mathrm{E}_{\mathrm{y}} \\
\mathrm{E}_{\mathrm{z}}
\end{array}\right)
$$

where the double bar indicates a second-rank tensor with 3 dimensions, and the matrix $U_{\varepsilon}$ has the form:

$$
\mathrm{U}_{\varepsilon}=\left(\begin{array}{ccc}
(\overrightarrow{\mathrm{N}} \overline{\bar{\varepsilon}})_{\mathrm{x}} & (\overrightarrow{\mathrm{~N}} \overline{\bar{\varepsilon}})_{\mathrm{y}} & (\overrightarrow{\mathrm{~N}} \overline{\bar{\varepsilon}})_{\mathrm{z}}  \tag{7.70}\\
\mathrm{~T}_{1 \mathrm{x}} & \mathrm{~T}_{1 \mathrm{y}} & \mathrm{~T}_{1 \mathrm{z}} \\
\mathrm{~T}_{2 \mathrm{x}} & \mathrm{~T}_{2 \mathrm{y}} & \mathrm{~T}_{2 \mathrm{z}}
\end{array}\right)
$$

with $\overrightarrow{\mathrm{N}} \overline{\bar{\varepsilon}}$ representing a tensor product with contraction of indices, for example, $(\overrightarrow{\mathrm{N}} \overline{\bar{\varepsilon}})_{\mathrm{x}}=\mathrm{N}_{\mathrm{x}} \overline{\bar{\varepsilon}}_{\mathrm{xx}}+\mathrm{N}_{\mathrm{y}} \overline{\bar{y}}_{\mathrm{yx}}+\mathrm{N}_{\mathrm{z}} \overline{\bar{\varepsilon}}_{\mathrm{zx}}$, etc.

The vectors $\overrightarrow{\mathrm{N}}, \overrightarrow{\mathrm{T}}_{1}$, and $\overrightarrow{\mathrm{T}}_{2}$ are defined on the grating surface, but for their further Fourier transform, it is necessary to choose a suitable continuation. The necessary conditions are that (i) they are continuous on the surfaces where $\varepsilon$ and $\mu$ are discontinuous, and (ii) they form an orthonormal triad.

Since $\overline{\bar{\varepsilon}}$ never vanishes, the determinant of $U_{\varepsilon}$ represents a quadric non-null form, equal to:

$$
\begin{equation*}
\xi_{\varepsilon} \equiv \operatorname{det} \mathrm{U}_{\varepsilon}=(\overrightarrow{\mathrm{N}} \overline{\bar{\varepsilon}}) \cdot\left(\overrightarrow{\mathrm{T}}_{1} \times \overrightarrow{\mathrm{T}}_{2}\right)=\sum_{\mathrm{i}, \mathrm{j}=\mathrm{x}, \mathrm{y}, \mathrm{z}} \mathrm{~N}_{\mathrm{i}} \varepsilon_{\mathrm{ij}} \mathrm{~N}_{\mathrm{j}} \tag{7.71}
\end{equation*}
$$

since $\overrightarrow{\mathrm{N}}=\overrightarrow{\mathrm{T}}_{1} \times \overrightarrow{\mathrm{T}}_{2}$.
Thus $U_{\varepsilon}$ has an inverse $U_{\varepsilon}^{\text {inv }}$ in the form:

$$
\mathrm{U}_{\varepsilon}^{\mathrm{inv}}=\frac{1}{\xi_{\varepsilon}}\left(\begin{array}{ccc}
\mathrm{N}_{\mathrm{x}} & -\left[(\overrightarrow{\mathrm{N}} \overline{\bar{\varepsilon}}) \times \overrightarrow{\mathrm{T}}_{2}\right]_{\mathrm{x}} & {\left[(\overrightarrow{\mathrm{~N}} \overline{\bar{\varepsilon}}) \times \overrightarrow{\mathrm{T}}_{1}\right]_{\mathrm{x}}}  \tag{7.72}\\
\mathrm{~N}_{\mathrm{y}} & -\left[(\overrightarrow{\mathrm{N}} \overline{\bar{\varepsilon}}) \times \overrightarrow{\mathrm{T}}_{2}\right]_{\mathrm{y}} & {\left[(\overrightarrow{\mathrm{~N}} \overline{\bar{\varepsilon}}) \times \overrightarrow{\mathrm{T}}_{1}\right]_{\mathrm{y}}} \\
\mathrm{~N}_{\mathrm{z}} & -\left[(\overrightarrow{\mathrm{N}} \overline{\bar{\varepsilon}}) \times \overrightarrow{\mathrm{T}}_{2}\right]_{\mathrm{z}} & {\left[(\overrightarrow{\mathrm{~N}} \overline{\bar{\varepsilon}}) \times \overrightarrow{\mathrm{T}}_{1}\right]_{\mathrm{z}}}
\end{array}\right)
$$

It is not evident to derive this form, but it can easily be verified by using the equivalence $U_{\varepsilon} U_{\varepsilon}^{\text {inv }}=\mathbb{I}$ and the fact that $\vec{N}=\vec{T}_{1} \times \vec{T}_{2}$. For example, the product of the second line of $U_{\varepsilon}$ with the second column of $\mathrm{U}_{\varepsilon}^{\text {inv }}$ can be written in vectorial form:

$$
\begin{equation*}
\xi_{\varepsilon}\left(\mathrm{U}_{\varepsilon} \mathrm{U}_{\varepsilon}^{\text {inv }}\right)_{\mathrm{yy}}=-\overrightarrow{\mathrm{T}}_{1} \cdot\left[(\overrightarrow{\mathrm{~N}} \overline{\bar{\varepsilon}}) \times \overrightarrow{\mathrm{T}}_{2}\right]=-\left[(\overrightarrow{\mathrm{N}} \overline{\bar{\varepsilon}}) \times \overrightarrow{\mathrm{T}}_{2}\right] \cdot \overrightarrow{\mathrm{T}}_{1}=-(\overrightarrow{\mathrm{N}} \overline{\bar{\varepsilon}}) \cdot\left(\overrightarrow{\mathrm{T}}_{2} \times \overrightarrow{\mathrm{T}}_{1}\right)=(\overrightarrow{\mathrm{N}} \overline{\bar{\varepsilon}}) \cdot \overrightarrow{\mathrm{N}}=\xi_{\varepsilon}(7 \tag{7.73}
\end{equation*}
$$

Going back to the vector $\mathrm{F}_{\varepsilon}$, it is continuous across the grating surface, whereas the Cartesian components of the electric vector are, in general discontinuous, as well as the components of $\mathrm{U}_{\varepsilon}$. Thus for their Fourier transform, it is necessary to apply the inverse factorization rule:

$$
\begin{equation*}
\left[\mathrm{F}_{\varepsilon}\right]=\llbracket \mathrm{U}_{\varepsilon}^{\mathrm{inv}} \rrbracket^{-1}[\overrightarrow{\mathrm{E}}] \tag{7.74}
\end{equation*}
$$

At the other hand,

$$
\begin{equation*}
\overrightarrow{\mathrm{E}}=\mathrm{U}_{\varepsilon}^{\text {inv }} \mathrm{F}_{\varepsilon} \Rightarrow \overrightarrow{\mathrm{D}}=\overline{\bar{\varepsilon}} \mathrm{U}_{\varepsilon}^{\text {inv }} \mathrm{F}_{\varepsilon} \tag{7.75}
\end{equation*}
$$

with $\mathrm{F}_{\varepsilon}$ being continuous, so that the Fourier transform of $\overrightarrow{\mathrm{D}}$ requires the direct factorization rule:

$$
\begin{equation*}
[\overrightarrow{\mathrm{D}}]=\llbracket \overline{\bar{\varepsilon}} \mathrm{U}_{\varepsilon}^{\mathrm{inv}} \rrbracket \llbracket \mathrm{U}_{\varepsilon}^{\mathrm{inv}} \rrbracket^{-1}[\overrightarrow{\mathrm{E}}] \tag{7.76}
\end{equation*}
$$

i.e.,

$$
\begin{equation*}
\mathrm{Q}_{\varepsilon}=\llbracket \overline{\bar{\varepsilon}} \mathrm{U}_{\varepsilon}^{\mathrm{inv}} \rrbracket \llbracket \mathrm{U}_{\varepsilon}^{\mathrm{inv}} \rrbracket^{-1} \tag{7.77}
\end{equation*}
$$

For gratings having anisotropic magnetic properties, the corresponding $\mathrm{Q}_{\mu}$ matrix is obtained from eqs. (7.71), (7.72), and (7.77) by replacing $\mathrm{U}_{\varepsilon}^{\text {inv }}$ by $\mathrm{U}_{\mu}^{\text {inv }}$ and $\overline{\bar{\varepsilon}}$ by $\overline{\bar{\mu}}$.

### 7.5.1. Lamellar gratings made of anisotropic material

Such gratings are analyzed in the chapter devoted to the Fourir modal method by using more direct approaches, but here we want show how the corresponding equations can be obtained from the general eqs.(7.72). To this aim it is sufficient to realize that

$$
\begin{align*}
& \overrightarrow{\mathrm{N}}=(1,0,0) \\
& \overrightarrow{\mathrm{T}}_{1}=(0,1,0)  \tag{7.78}\\
& \overrightarrow{\mathrm{T}}_{2}=(0,0,1)
\end{align*}
$$

so that

$$
\begin{align*}
& (\overrightarrow{\mathrm{N}} \overline{\bar{\varepsilon}}) \times \overrightarrow{\mathrm{T}}_{2}=\left(\varepsilon_{\mathrm{xy}},-\varepsilon_{\mathrm{xx}}, 0\right) \\
& (\overrightarrow{\mathrm{N}} \overline{\bar{\varepsilon}}) \times \overrightarrow{\mathrm{T}}_{1}=\left(-\varepsilon_{\mathrm{xz}}, 0, \varepsilon_{\mathrm{xx}}\right)  \tag{7.79}\\
& (\overrightarrow{\mathrm{N}} \overline{\bar{\varepsilon}}) \cdot \overrightarrow{\mathrm{N}}=\varepsilon_{\mathrm{xx}}
\end{align*}
$$

and eq.(7.72) takes the form:

$$
\mathrm{U}_{\varepsilon}^{\mathrm{inv}}=\left(\begin{array}{ccc}
\frac{1}{\varepsilon_{\mathrm{xx}}} & -\frac{\varepsilon_{\mathrm{xy}}}{\varepsilon_{\mathrm{xx}}} & -\frac{\varepsilon_{\mathrm{xz}}}{\varepsilon_{\mathrm{xx}}}  \tag{7.80}\\
0 & 1 & 0 \\
0 & 0 & 1
\end{array}\right)
$$

with a determinant equal to $\llbracket \frac{1}{\varepsilon_{\mathrm{xx}}} \rrbracket$. Thus

$$
\llbracket \mathrm{U}_{\varepsilon}^{\mathrm{inv}} \rrbracket^{-1}=\left(\begin{array}{ccc}
\llbracket \frac{1}{\varepsilon_{\mathrm{xx}}} \rrbracket^{-1} & \llbracket \frac{1}{\varepsilon_{\mathrm{xx}}} \|^{-1} \llbracket \frac{\varepsilon_{\mathrm{xy}}}{\varepsilon_{\mathrm{xx}}} \rrbracket & \llbracket \frac{1}{\varepsilon_{\mathrm{xx}}} \rrbracket^{-1} \llbracket \frac{\varepsilon_{\mathrm{xz}}}{\varepsilon_{\mathrm{xx}}} \rrbracket  \tag{7.81}\\
0 & \mathbb{I} & 0 \\
0 & 0 & \mathbb{I}
\end{array}\right)
$$

The second matrix that is required takes the form:

$$
\overline{\bar{\varepsilon}} \mathrm{U}_{\varepsilon}^{\mathrm{inv}}=\left(\begin{array}{ccc}
1 & 0 & 0  \tag{7.82}\\
\frac{\varepsilon_{\mathrm{yx}}}{\varepsilon_{\mathrm{xx}}} & \varepsilon_{\mathrm{yy}}-\frac{\varepsilon_{\mathrm{yx}} \varepsilon_{\mathrm{xy}}}{\varepsilon_{\mathrm{xx}}} & \varepsilon_{\mathrm{yz}}-\frac{\varepsilon_{\mathrm{yx}} \varepsilon_{\mathrm{xz}}}{\varepsilon_{\mathrm{xx}}} \\
\frac{\varepsilon_{\mathrm{zx}}}{\varepsilon_{\mathrm{xx}}} & \varepsilon_{\mathrm{zy}}-\frac{\varepsilon_{\mathrm{zx}} \varepsilon_{\mathrm{xy}}}{\varepsilon_{\mathrm{xx}}} & \varepsilon_{\mathrm{zz}}-\frac{\varepsilon_{\mathrm{zx}} \varepsilon_{\mathrm{xz}}}{\varepsilon_{\mathrm{xx}}}
\end{array}\right)
$$

This form is valid even when the permittivity tensor is not symmetric, as happens in the modeling of magnetooptical effects.

The $\mathrm{Q}_{\varepsilon}$ matrix takes the form obtained in [7.19], using a completely different approach:

$$
\begin{aligned}
& \mathrm{Q}_{\varepsilon}=\llbracket \overline{\bar{\varepsilon}} \mathrm{U}_{\varepsilon}^{\mathrm{inv}} \rrbracket \llbracket \mathrm{U}_{\varepsilon}^{\mathrm{inv}} \rrbracket^{-1}
\end{aligned}
$$

### 7.6. Normal vector prolongation for 2D periodicity; Fourier transform

As observed, the proper use of the direct and the inverse factorization rules requires that the vector normal to the interfaces between different media is defined not only on these interfaces, but throughout the entire grating cell. In the case of classical gratings with onedimensional periodicity, the prolongation of the normal vector can be done quite easily, as shown in sec.7.4.1. For two-dimensional periodicity, the choice depends on the geometry, but also on its mathematical representation. Several different solutions have to be considered, without pretending to be exhaustive.

In general, the cross-section profile changes with z , so that the matrices $\mathrm{Q}_{\varepsilon}, \mathrm{Q}_{\mu}$, and M have to be recalculated for each value of $z$. If the geometry is $z$-invariant, this must be done only once. Concerning the Fourier components of the normal vector, there are two different classes of grating profiles that has to be treated separately. The first class consists of surfaces that can be expressed all over the unit cell (containing a single period in $x$ and $z$ ) as an analytical (at lease piecewisely) function $\mathrm{z}_{\mathrm{S}}=\mathrm{g}(\mathrm{x}, \mathrm{y})$, where S indicates that the point lies on the interface. In this case, it is possible to have a unique extension of $\overrightarrow{\mathrm{N}}$ whatever the values of $z$. In addition, it is not necessary to calculate the cross-section of the surface(s) with a plane perpendicular to the z -axis for each value of z . This case also includes multilayered homomorphous structure with constant layer(s) thickness in the z-direction. We consider this class of cases in sec.7.6.1.

The second class of surfaces includes surfaces that cannot be expressed through singlevalued functions, as the example given on the right-hand side of Fig.7.12. In that case, it is necessary for each fixed value of $z$ to know the cross-section function of the grating surface with the plane $\mathrm{z}=$ const. Subsection 7.6 .2 presents general analysis, some important specific cases are considered further in the following subsections.

### 7.6.1. General analytical surfaces

If the interface representing the structure can be expressed as a single-valued function, analytical over the entire unit cell (this is also valid if different analytical functions can be defined over different regions of the cell):

$$
\begin{equation*}
\mathrm{z}_{\mathrm{S}}=\mathrm{g}(\mathrm{x}, \mathrm{y}) \tag{7.84}
\end{equation*}
$$

then the components of the vector normal to the surface defined on the surface have the form:

$$
\begin{equation*}
\overrightarrow{\mathrm{N}}(\mathrm{x}, \mathrm{y})=\frac{\left(-\frac{\partial \mathrm{g}}{\partial \mathrm{x}},-\frac{\partial \mathrm{g}}{\partial \mathrm{y}}, 1\right)}{\sqrt{1+\operatorname{grad}^{2} \mathrm{~g}(\mathrm{x}, \mathrm{y})}} \tag{7.85}
\end{equation*}
$$

It can immediately be extended to whatever the values of z inside the modulated region. Moreover, its values do not depend on z , as it was a case of the classical one-dimensional grating already discussed in sec.7.4.1.

The permittivity and its inverse can easily be obtained on a mesh ( $\mathrm{x}, \mathrm{y}$ ) covering the grating cell for each z :

$$
\begin{align*}
& \mathrm{g}(\mathrm{x}, \mathrm{y})<\mathrm{z} \Rightarrow \varepsilon(\mathrm{x}, \mathrm{y})=\varepsilon_{3}  \tag{7.86}\\
& \mathrm{~g}(\mathrm{x}, \mathrm{y}) \geq \mathrm{z} \Rightarrow \varepsilon(\mathrm{x}, \mathrm{y})=\varepsilon_{1}
\end{align*}
$$

where 1 and 3 are the indexes representing respectively the inferior and the superior regions separated by the grating surface (as it was done in Fig.7.3). If the cross-section of the grating surface with the planes $\mathrm{z}=$ const are ellipses (or circles), and if $\mathrm{N}_{\mathrm{z}}$ does not depend on ( $\mathrm{x}, \mathrm{y}$ ) at each $z$, it is possible to replace the numerical Fourier transform by an analytical formulae. One important particular case is the z-invariant grating with elliptical cross section; another case includes the gratings having a rotational symmetry, as shown in Fig.7.12.

The same extension (7.85) for the normal vector is valid for a stack of layers having homogeneous thicknesses in the z-direction:

$$
\begin{equation*}
\mathrm{z}_{\mathrm{S}, \mathrm{j}}=\mathrm{g}(\mathrm{x}, \mathrm{y})+\mathrm{t}_{\mathrm{j}} \tag{7.87}
\end{equation*}
$$

The permittivity and its inverse inside the intermediate layers are simply given as:

$$
\begin{equation*}
\mathrm{z}_{\mathrm{S}, \mathrm{j}-1}(\mathrm{x}, \mathrm{y})<\mathrm{z} \leq \mathrm{z}_{\mathrm{S}, \mathrm{j}}(\mathrm{x}, \mathrm{y}) \Rightarrow \varepsilon(\mathrm{x}, \mathrm{y})=\varepsilon_{\mathrm{j}} . \tag{7.88}
\end{equation*}
$$

### 7.6.2. Irregular general surfaces

If the case does not fit into the preceding section, the interface is expressed through the more general function $\mathrm{u}\left(\mathrm{x}, \mathrm{y}, \mathrm{z}_{\mathrm{S}}\right)=0$, and the vector $\overrightarrow{\mathrm{N}}$ has to be determined for each inclusion:

$$
\begin{equation*}
\overrightarrow{\mathrm{N}}\left(\mathrm{x}, \mathrm{y}, \mathrm{z}_{\mathrm{S}}\right)=\frac{\left(\frac{\partial \mathrm{u}}{\partial \mathrm{x}}, \frac{\partial \mathrm{u}}{\partial \mathrm{y}} \frac{\partial \mathrm{u}}{\partial \mathrm{z}_{\mathrm{S}}}\right)}{\left|\operatorname{grad} \mathrm{u}\left(\mathrm{x}, \mathrm{y}, \mathrm{z}_{\mathrm{S}}\right)\right|} \tag{7.89}
\end{equation*}
$$

However, these values are well defined on the grating surface (except on its edges), and have to be extended over the entire cell. When considerintg a cross-section of the profile with a plane at $\mathrm{z}=$ const., several different cases exist:

### 7.6.2.1. Single-valued radial cross-section

At first, we shall consider that the cross section function $f\left(x_{S}, y_{S}\right)=0$ defines a single curve, and that curve can be expressed in cylindrical coordinates as

$$
\begin{equation*}
\rho_{\mathrm{S}}=\rho_{\mathrm{S}}(\varphi) \tag{7.90}
\end{equation*}
$$

where

$$
\begin{align*}
& \rho_{\mathrm{S}}=\sqrt{\left(\mathrm{x}_{\mathrm{S}}-\mathrm{x}_{\mathrm{C}}\right)^{2}+\left(\mathrm{y}_{\mathrm{S}}-\mathrm{y}_{\mathrm{C}}\right)^{2}}  \tag{7.91}\\
& \varphi=\arctan \left[\left(\mathrm{y}_{\mathrm{S}}-\mathrm{y}_{\mathrm{C}}\right) /\left(\mathrm{x}_{\mathrm{S}}-\mathrm{x}_{\mathrm{C}}\right)\right]
\end{align*}
$$



Fig.7.6. Single-curve cross-section of the grating surface at $z=$ const.
and $x_{C}$ and $y_{C}$ represent a "central" point of the curve, Fig.7.6. Here we assume that the values of $\rho_{\mathrm{S}}$ are unique for each $\varphi$. The other case is analyzed further on.

It is possible to extend to the entire cell the values of the normal vector, defined only on the curve, by assuming that it is constant for each fixed angle $\varphi$. This prolongation requires the following procedure:

1. Fixing the pair $(\mathrm{x}, \mathrm{y})$.
2. Calculating the angle $\varphi=\arctan \left[\left(y-y_{C}\right) /\left(x-x_{C}\right)\right]$.
3. Calculating $\rho_{\mathrm{S}}=\rho_{\mathrm{S}}(\varphi)$ from eq.(7.91).
4. Calculation of $x_{S}=\rho_{\mathrm{S}} \cos \varphi$, and $y_{\mathrm{S}}=\rho_{\mathrm{S}} \sin \varphi$.
5. Determining $\mathrm{N}_{\mathrm{z}}$, together with $\mathrm{N}_{\mathrm{x}}$ and $\mathrm{N}_{\mathrm{y}}$ from eq.(7.89).
6. Attributing these values of the components of $\overrightarrow{\mathrm{N}}\left(\mathrm{x}_{\mathrm{S}}, \mathrm{y}_{\mathrm{S}}\right)$ to the pair (x,y).
7. Fast Fourier transform after the normal vector components are determined for all the pairs ( $\mathrm{x}, \mathrm{y}$ ) on a mesh inside the grating cell.
The procedure can be simplified for most of the typical diffracting objects, as shown further for objects with elliptical or circular cross-section.

If the grating profile varies with z , the calculations of the Fourier components of the permittivity and its inverse has to be made at each value of $z$, both for the analytical profiles, for which the normal vector prolongation can be chosen z -invariant, or for the irregular surfaces. For each ( $\mathrm{x}, \mathrm{y}$ ) pair of the mesh used in the FFT method, it is possible to determine whether the point lies within or outside the cross-section part of Fig.7.6:

$$
\begin{array}{ll}
\rho(\varphi)<\rho_{\mathrm{S}}(\varphi), & \varepsilon(\mathrm{x}, \mathrm{y})=\varepsilon_{\text {inside }} \\
\rho(\varphi) \geq \rho_{\mathrm{S}}(\varphi), & \varepsilon(\mathrm{x}, \mathrm{y})=\varepsilon_{\text {outside }} \tag{7.92}
\end{array}
$$

with $\rho=\sqrt{\left(x-x_{C}\right)^{2}+\left(y-y_{C}\right)^{2}}$.

### 7.6.2.2. Objects with polygonal cross section

A typical example of such objects is presented in Fig.7.1. Its surface consists of different plates, and for their treatment the condition $N_{z}=$ const. for fixed $z$ is fulfilled, because $\vec{N}$ is constant at each plate. Another possible surface consists of plane ribbons with curvature in zdirection, Fig.7.7.


Fig.7.7. Surface made of plane ribbons

As shown in Fig.7.8, the cross-section represents a polygon. On each of its sides the modulus of the in-plane component of the normal vector is known:


Fig.7.8. Object with a polygonal cross-section.

$$
\begin{equation*}
\mathrm{N}_{/ /, i}=\sqrt{1-\mathrm{N}_{\mathrm{z}, i}^{2}} \tag{7.93}
\end{equation*}
$$

and its direction is perpendicular to the segment. If the $i$-th segment is located between the angles $\varphi_{i}$ and $\varphi_{i+1}$, we can extend the definition of the normal vector all over the unit cell situated within the range $\left(\varphi_{i}, \varphi_{i+1}\right)$, delimited by the bold dot-dashed lines in Fig.7.8, by
assuming that $\overrightarrow{\mathrm{N}}=\overrightarrow{\mathrm{N}}_{i}$. The normal vector extension will be continuous everywhere, except on the sector border lines (bold dot-dashed lines) and thus the only points where both permittivity (and/or permeability) and $\overrightarrow{\mathrm{N}}$ are simultaneously discontinuous are at the polygonal corners, where anyway $\overrightarrow{\mathrm{N}}$ is never continuous.

The procedure to follow requires that for each value of $z$ the polygon corners $\left(x_{i}, y_{i}\right)$ and the z-components of $\overrightarrow{\mathrm{N}}$ for each segment are determined, as well as fixing some "central" point $\left(\mathrm{x}_{\mathrm{C}}, \mathrm{y}_{\mathrm{C}}\right)$. Then the angular ranges of each segment with respect to that central point are calculated:

$$
\begin{equation*}
\varphi_{i}=\arctan \frac{\mathrm{y}_{i}-\mathrm{y}_{\mathrm{C}}}{\mathrm{x}_{i}-\mathrm{x}_{\mathrm{C}}} \tag{7.94}
\end{equation*}
$$

For each pair $(x, y)$, the azimuthal angle is given as $\varphi=\arctan \left[\left(y-y_{C}\right) /\left(x-x_{C}\right)\right]$, which value determines the number of the segment, say $j$, within the point lies. The unknown inplane part of the normal vector is perpendicular to the $j$-th segment:

$$
\begin{align*}
& N_{x, j}=\left(y_{j+1}-y_{j}\right) \sqrt{\frac{1-N_{z, j}^{2}}{\left(y_{j+1}-y_{j}\right)^{2}+\left(x_{j+1}-x_{j}\right)^{2}}} \\
& N_{y, j}=-\left(x_{j+1}-x_{j}\right) \sqrt{\frac{1-N_{z, j}^{2}}{\left(y_{j+1}-y_{j}\right)^{2}+\left(x_{j+1}-x_{j}\right)^{2}}} \tag{7.95}
\end{align*}
$$

The expression in the square root comes from the normalization of $\overrightarrow{\mathrm{N}}$.
The value of the permittivity depends on whether the point ( $\mathrm{x}, \mathrm{y}$ ) lies inside or outside the polygon. The calculations of $\varepsilon(\mathrm{x}, \mathrm{y})$ and $1 / \varepsilon$ are made simultaneously with the normal vector calculus. After the angular segment in which the point ( $\mathrm{x}, \mathrm{y}$ ) of the mesh in grating cell is determined (say the j -th one, as in eq.(7.95)), we can find the length of $\rho_{\mathrm{S}}$ between the central point and the polygon segment, shown in Fig.7.8. For this sake we show in Fig.7.9 the enlarged segment:


Fig.7.9. The j-segment of Fig. 7.8 together with notations

The sine theorem gives the possibility to determine the angle $\zeta$ :

$$
\begin{equation*}
\frac{\sqrt{\left(x_{j+1}-x_{j}\right)^{2}+\left(y_{j+1}-y_{j}\right)^{2}}}{\sin \left(\varphi_{j+1}-\varphi_{j}\right)}=\frac{\sqrt{\left(x_{j+1}-x_{C}\right)^{2}+\left(y_{j+1}-y_{C}\right)^{2}}}{\sin (\zeta)} \tag{7.96}
\end{equation*}
$$

wherefrom the radius $\rho_{\mathrm{S}}$ is given as:

$$
\begin{equation*}
\rho_{\mathrm{S}}=\sin (\zeta) \frac{\sqrt{\left(\mathrm{x}_{\mathrm{C}}-\mathrm{x}_{\mathrm{j}}\right)^{2}+\left(\mathrm{y}_{\mathrm{C}}-\mathrm{y}_{\mathrm{j}}\right)^{2}}}{\sin \left(\pi-\zeta-\varphi+\varphi_{\mathrm{j}}\right)} \tag{7.97}
\end{equation*}
$$

Eq.(7.92) enables us to obtain the values of the permittivity (and its inverse):

$$
\begin{array}{ll}
\rho(\varphi)<\rho_{\mathrm{S}}(\varphi), & \varepsilon(\mathrm{x}, \mathrm{y})=\varepsilon_{\text {inside }}  \tag{7.98}\\
\rho(\varphi) \geq \rho_{\mathrm{S}}(\varphi), & \varepsilon(\mathrm{x}, \mathrm{y})=\varepsilon_{\text {outside }}
\end{array}
$$

with $\rho=\sqrt{\left(x-x_{C}\right)^{2}+\left(y-y_{C}\right)^{2}}$.


Fig.7.10. Schematic presentation of corner rounding

Concerning the edges, in reality the surfaces never have such, as etching always ends by rounding the corners, as shown in Fig.7.10. Let us consider that the rounding between the segments numbered $i-1$ and $i$ is made preserving the values of $\mathrm{N}_{\mathrm{z}}$, and that in the cross-plane z $=$ const., the rounding can be considered as circular, having a center in $\left(\mathrm{x}_{\mathrm{C}_{i}}, \mathrm{y}_{\mathrm{C}_{i}}\right)$ and radius $\mathrm{r}_{i}$. The in-plane component of the normal vector follows the curvature radius and thus is given by equations, similar to eq.(7.95):

$$
\begin{align*}
& \mathrm{N}_{\mathrm{x}, \Delta_{i}}=\left(\mathrm{x}-\mathrm{x}_{\mathrm{C}_{i}}\right) \sqrt{\frac{1-\mathrm{N}_{\mathrm{z}, i}^{2}}{\left(\mathrm{y}-\mathrm{y}_{\mathrm{C}_{i}}\right)^{2}+\left(\mathrm{x}-\mathrm{x}_{\mathrm{C}_{i}}\right)^{2}}} \\
& \mathrm{~N}_{\mathrm{y}, \Delta_{i}}=\left(\mathrm{y}-\mathrm{y}_{\mathrm{C}_{i}}\right) \sqrt{\frac{1-\mathrm{N}_{\mathrm{z}, i}^{2}}{\left(\mathrm{y}-\mathrm{y}_{\mathrm{C}_{i}}\right)^{2}+\left(\mathrm{x}-\mathrm{x}_{\mathrm{C}_{i}}\right)^{2}}} \tag{7.99}
\end{align*}
$$

The prolongation is more complicated, if the consecutive values of $\mathrm{N}_{\mathrm{z}}$ at the two sides of the rounded corner differ significantly. In that case a linear interpolation of $\mathrm{N}_{\mathrm{z}}$ between $\varphi_{i}$ and $\varphi_{i}+\Delta_{i}$ can be applied.

### 7.6.2.3. Mutlivalued cross-sections

If the cross-section cannot be represented as a radial function, nother possibility arises if it is a piecewise analytical function in x (or y ), as shown in Fig.7.11, where we can use two different functions of x . We assume again that $\mathrm{N}_{\mathrm{z}}$ is known, as it happens for z -independent profiles, for which it is simply null. In the upper part of the figure, for each value of $x$ it is possible to calculate the normal vector on the profile:

$$
\begin{equation*}
\overrightarrow{\mathrm{N}}_{1}=\frac{\left(-\mathrm{f}_{1}^{\prime}(\mathrm{x}), 1, \mathrm{~N}_{\mathrm{z}}\right)}{\sqrt{1+\mathrm{f}_{1}^{\prime 2}(\mathrm{x})+\mathrm{N}_{\mathrm{z}}^{2}}} \tag{7.100}
\end{equation*}
$$

We can take this value to be the same for each $y$ in the upper region $A_{1}$, so that the numerical Fourier transform is made only once in $\mathrm{A}_{1}$ and once in $\mathrm{A}_{2}$.


Fig.7.11. Piecewise analytical c ross-section of the grating surface at $z=$ const.

The permittivity has to be calculated at each ( $\mathrm{x}, \mathrm{y}$ ) mesh point:

$$
\begin{array}{ll}
\mathrm{y}<\mathrm{y}_{\mathrm{S}}, & \varepsilon(\mathrm{x}, \mathrm{y})=\varepsilon_{\text {inside }}  \tag{7.101}\\
\mathrm{y} \geq \mathrm{y}_{\mathrm{S}}, & \varepsilon(\mathrm{x}, \mathrm{y})=\varepsilon_{\text {outside }} .
\end{array}
$$

### 7.6.4. Objects with cylindrical symmetry

Many periodic systems consist of inclusion having rotational cylindrical symmetry, like spheres, vertical cylinders, or ellipsoids with axe of rotation parallel to the $z$-axis, but also smooth surfaces, as presented in Fig.7.12.

These structures are characterized by a circular cross-section of the surface with the horizontal planes at $\mathrm{z}=$ const., but also with an independence on x and y of the values of $\mathrm{N}_{\mathrm{z}}$ on each horizontal plane. In addition, due to the circular cross-sections, the angular component $\mathrm{N}_{\varphi}=0$ everywhere. Once z is fixed, the variation of the interface in the vertical
direction fixes the value of $\mathrm{N}_{z}$, for example through eq.(7.89), wherefrom the radial normal vector component $N_{\rho}=\sqrt{1-N_{z}^{2}}$. For each pair of $(x, y)$ then:


Fig.7.12. Several profiles with cylindrical symmetry.

$$
\begin{align*}
& N_{x}(x, y)=\frac{\left(x-x_{C}\right) N_{\rho}}{\sqrt{\left(x-x_{C}\right)^{2}+\left(y-y_{C}\right)^{2}}} \\
& N_{y}(x, y)=\frac{\left(y-y_{C}\right) N_{\rho}}{\sqrt{\left(x-x_{C}\right)^{2}+\left(y-y_{C}\right)^{2}}} \tag{7.102}
\end{align*}
$$

In addition, for profiles invariant in z-direction, $\mathrm{N}_{\mathrm{z}}=0$ and $\mathrm{N}_{\mathrm{\rho}}=1$.
The permittivity is given as a piecewise constant function:

$$
\begin{array}{lll}
\varepsilon(\mathrm{x}, \mathrm{y})=\varepsilon_{\text {inside }}, & \text { if } & \left(\mathrm{x}-\mathrm{x}_{\mathrm{C}}\right)^{2}+\left(\mathrm{y}-\mathrm{y}_{\mathrm{C}}\right)^{2}<\mathrm{R}^{2}(\mathrm{z})  \tag{7.103}\\
\varepsilon(\mathrm{x}, \mathrm{y})=\varepsilon_{\text {outside }}, & \text { if } & \left(\mathrm{x}-\mathrm{x}_{\mathrm{C}}\right)^{2}+\left(\mathrm{y}-\mathrm{y}_{\mathrm{C}}\right)^{2} \geq \mathrm{R}^{2}(\mathrm{z})
\end{array}
$$

where $\mathrm{R}(\mathrm{z})$ is the radius of the profile surface for a given z .
Having obtained the values of the normal vector components and permittivity for each x , y enables us to calculate their Fourier transforms, either by Fast Fourier transform (FFT), or analytically.

### 7.6.5. Objects with elliptical cross-section

Similar simplification is possible for systems with elliptical cross-sections that have $\mathrm{N}_{\mathrm{z}}=$ const. for z fixed. Such are the inclusions of vertical cylinders with elliptical cross-section, ellipsoids with one of the axes orientated in z-direction, but also all types of the structures shown schematically in Fig.7.12 that have elliptical or circular cross-sections.

Let us assume that the ellipse axes are parallel to the x and y -axes.. The cross-section curve for $\mathrm{z}=$ const. is given by the equation:

$$
\begin{equation*}
\left(\frac{x_{\mathrm{S}}-\mathrm{x}_{\mathrm{C}}}{a}\right)^{2}+\left(\frac{y_{\mathrm{S}}-\mathrm{y}_{\mathrm{C}}}{\mathrm{~b}}\right)^{2}=\mathrm{R}^{2} \tag{7.104}
\end{equation*}
$$

In order to obtain results similar to eq.(7.102), we introduce an elliptical coordinates, defined as:

$$
\begin{align*}
& \tilde{x}=\frac{x-x_{C}}{a}  \tag{7.105}\\
& \tilde{y}=\frac{y-y_{C}}{b}
\end{align*}
$$

Using these notations, the ellipse becomes a circle, for which the considerations of the previous subsection apply. Thus

$$
\begin{align*}
& N_{x}(x, y)=\frac{\frac{x-x_{C}}{a} N_{\rho}}{\sqrt{\left(\frac{x-x_{C}}{a}\right)^{2}+\left(\frac{y-y_{C}}{b}\right)^{2}}}  \tag{7.106}\\
& N_{y}(x, y)=\frac{\frac{y-y_{C}}{b} N_{\rho}}{\sqrt{\left(\frac{x-x_{C}}{a}\right)^{2}+\left(\frac{y-y_{C}}{b}\right)^{2}}}
\end{align*}
$$

with $N_{\rho}=\sqrt{1-\mathrm{N}_{\mathrm{z}}^{2}}$, which remains constant for each fixed z .
Concerning the permittivity, it is determined in the same way as in eq.(7.103) for circular profile:

$$
\begin{array}{ll}
\varepsilon(\mathrm{x}, \mathrm{y})=\varepsilon_{\text {inside }}, & \text { if }\left(\frac{\mathrm{x}-\mathrm{x}_{\mathrm{C}}}{\mathrm{a}}\right)^{2}+\left(\frac{\mathrm{y}-\mathrm{y}_{\mathrm{C}}}{\mathrm{~b}}\right)^{2}<\mathrm{R}^{2}(\mathrm{z})  \tag{7.107}\\
\varepsilon(\mathrm{x}, \mathrm{y})=\varepsilon_{\text {outside }}, & \text { if } \quad\left(\frac{\mathrm{x}-\mathrm{x}_{\mathrm{C}}}{\mathrm{a}}\right)^{2}+\left(\frac{\mathrm{y}-\mathrm{y}_{\mathrm{C}}}{\mathrm{~b}}\right)^{2} \geq \mathrm{R}^{2}(\mathrm{z})
\end{array}
$$

with $R(z)$ given by eq.(7.104).

## Remark on the prolongation of the normal vector

Special attention has recently been paid to the numerical implementation of the differential method for gratings having 2D periodicity formed by vertical holes or bumps that are invariant in z , and that have arbitrary cross-section in the xOy plane. A detailed study in the case of z-invariant geometry that applies for an eigenvalue/eigenvector technique of integration (FM or RCW method) is given in ref.[7.20], followed by several other works [7.21, 7.22]. It is necessary to note that the technique of prolongation of the normal vector as discussed in [7.20] can be applied also for z -dependent profiles with similar cross section; the difference is the renormalization factor $\sqrt{1-\mathrm{N}_{\mathrm{z}}^{2}}$ for each z .

The authors compare several different formulations of the Fourier Modal method applied to structures with rectangular, circular, or elliptical cross-sections. These formulations include the classical formulation of Moharam and Gaylor that uses only the direct factorization rule, the formulation given by Lifeng Li [7.23] that introduces two different Fourier transforms of the permittivity $\varepsilon$, namely $\lfloor\lceil\varepsilon\rceil\rfloor$ and $\lceil\lfloor\varepsilon\rfloor\rceil$, which are calculated by applying at first the inverse rule along one of the coordinates, and then the direct rule along the other one. This second formulation was made for rectangular and parallelogram crosssections. For circular or elliptical (or other smooth) forms, it introduces a stepwise treatment of the profile, which appears more slowly convergent than the special techniques developed after.


Fig.7.13. Two different prolongations of the normal vector for a circular inclusion. (a) Radial prolongation. (b) Electrostatic continuation of the normal vector for a circular cross-section inclusion inside a square grating cell (after [7.20]).
The third approach to the problem requires a prolongation of the normal vector (NV) to the profile within the entire grating cell. As already stressed, there are several possibilities to make this. A typical example is the radial prolongation, Fig.7.13a, which has been discussed in Sec.7.6.2.1 and 7.6.4 and it includes discontinuities of the normal vector on the cell boundaries, where the permittivity is continuous. Another approach proposed in [7.20] is the electrostatic one, which insures the continuity all over the cell and on its boundaries, except for on single points inside, Fig.7.13b.

Fig.7.14 shows the convergence rates for the transmitted zeroth order of a grating consisting of dielectric cylindrical inclusions with a circular cross-section with refractive index $\mathrm{n}=1.5$, in normal incidence from the substrate. The grating period is $2 \lambda$, the width of the grooves is $\lambda$, and the grating depth is $\lambda /(2 n-1)$. The graph presents the diffraction efficiency in transmission as a function of the truncation order N using the three considered formulations: Moharam's original formulation, Li's formulation, and the formulation using the normal vector ( NV ) field. As usual the Fourier series run from -N to N , which yields $2 \mathrm{~N}+1$ Fourier coefficients for each of the two directions of periodic continuation, or $(2 \mathrm{~N}+1)^{2}$ coefficients in total. As can be expected, both the original approach and the formulation by Li have worse convergence than when correctly taking into account the factorization rules for the tangential electric field and normal displacement components to the profile, where the permittivity is discontinuous [7.20]. It is necessary to stress that the difference in the
convergence rates is even more pronounced for metallic inclusions, having much larger optical contrast.


Fig.7.14. Convergence rates with respect to the truncation of the Fourier series for four different approaches used to model the diffraction by a cylindrical inclusion with circular cross-section.

Recently, Weiss et al. [7.24] proposed another alternative approach to treating smooth inclusions, by changing the coordinate system, so that its planes are parallel to the profile and to the grating sell walls (see Fig.7.15). If the transformed system is orthonormal, its coordinate lines are automatically tangential or perpendicular to the physical walls. If not, the Maxwell equations have to be rewritten in covariant vector form using the covariant and contravariant vector components.


Fig.7.15. Coordinate lines and surfaces according to (a) [7.20] and (b) [7.24]

This approach is somehow equivalent to the normal vector prolongation, due to two main reasons:
(i) The NV approach defines in an unambiguous manner the normal vectors on the profile, giving a liberty to continue them all over the cell. The coordinate transformation is also defined on the profile and the outside boundaries, but can be chosen in different ways around the grating cell.
(ii) The change of the coordinate system introduces in the Maxwell equations the metric tensor $\mathbb{G}$ that multiplies the electric displacement and magnetic induction in the right-hand side of eq.(7.6), so that for the electric field we obtain the substitution:

$$
\begin{equation*}
\overrightarrow{\mathrm{D}}=\varepsilon \overrightarrow{\mathrm{E}} \quad \rightarrow \quad \mathbb{G} \varepsilon \overrightarrow{\mathrm{E}} . \tag{7.108}
\end{equation*}
$$

The normal vector approach acts in a similar manner by introducing the matrix $\mathrm{Q}_{\varepsilon}$, given in eq.7.21, which makes the following substitution in the Fourier space, eq.(7.20):

$$
\begin{equation*}
[\overrightarrow{\mathrm{D}}]=[\varepsilon \overrightarrow{\mathrm{E}}] \rightarrow \mathrm{Q}_{\varepsilon}[\overrightarrow{\mathrm{E}}] . \tag{7.109}
\end{equation*}
$$

### 7.6.6. Multiprofile surfaces

A grating with multiple bumps (or inclusions) inside the single cell could be treated by separation the cell into sub-cells, not necessarily rectangular, containing a single inclusion, as shown in Fig.7.16, where a specific cross-section at $\mathrm{z}=$ const. is separated into three regions $\mathrm{A}, \mathrm{B}$, and C . As far as the Fourier components of the normal vector, of the permeability and the permittivity have to be calculated for each value of $z$ (if they depend on $z$ ), the separation into subcells can vary with z.


Fig.7.16. Cross-section at $z=$ const. of a grating having different inclusions. The three different regions to be treated independently are separated by dashed lines.

The case schematized in Fig.7.17a can result from a surface covered by a thin layer of another substance, a layer that cannot be treated using eq.(7.87). The simplest possibility is to have different continuation of the normal vector inside each region. At first, the angle $\varphi=\arctan \left[\left(y-y_{C}\right) /\left(x-x_{C}\right)\right]$ for the point with coordinates $(x, y)$ is calculated, and it is necessary to determine to which region the point belong. If it lies inside the innermost region C, the values of $\overrightarrow{\mathrm{N}}(\mathrm{x}, \mathrm{y})=\overrightarrow{\mathrm{N}}_{1}(\varphi)$, where $\overrightarrow{\mathrm{N}}_{1}(\varphi)$ is determined using one of the procedures discussed above for a single interface that is defined by the inner profile function.

If the point ( $x, y$ ) lies in the outermost region, we take $\vec{N}(x, y)=\vec{N}_{2}(\varphi)$, where $\overrightarrow{\mathrm{N}}_{2}(\varphi)$ corresponds to the second interface. In-between, we have two possibilities. The first choice is to divide the region into two subregions as indicated in Fig.7.17a with the dashed line. In each of them, $\overrightarrow{\mathrm{N}}(\mathrm{x}, \mathrm{y})$ is taken to be equal to its values on the adjacent profile, so that it is continuous everywhere where the permittivity and/or permeability are discontinuous.


Fig.7.17. Structures with interpenetrating cross-section profiles

The second possibility is to introduce a linear interpolation inside region $B$, but it is necessary to know the distances $\rho_{1}$ and $\rho_{2}$ between the central point and the profiles along the ray with $\varphi=\arctan \left[\left(y-y_{C}\right) /\left(x-x_{C}\right)\right]$ fixed. Then:

$$
\begin{align*}
& \overrightarrow{\mathrm{N}}_{\mathrm{B}}(\mathrm{x}, \mathrm{y})=\overrightarrow{\mathrm{N}}_{1}(\varphi)+\left[\overrightarrow{\mathrm{N}}_{2}(\varphi)-\overrightarrow{\mathrm{N}}_{1}(\varphi)\right] \frac{\rho-\rho_{1}}{\rho_{2}-\rho_{1}}  \tag{7.110}\\
& \text { with } \quad \rho^{2}=\left(x-x_{C}\right)^{2}+\left(\mathrm{y}-\mathrm{y}_{\mathrm{C}}\right)^{2}
\end{align*}
$$

Another specific case that appears in the studies of magnetic resonators is presented in Fig.7.17b. In can be treated in the same way is for the case in Fig.7.17a, but it is necessary to introduce a separate region D indicated in the figure and containing the opening, for which $\overrightarrow{\mathrm{N}}=(0,1,0)$, for example.

### 7.7. Integrating schemes

Numerical solution of a system of ordinary differential equations is a mature domain due to the enormous amount of physical and technical applications. Unfortunately, the grating problem represents one of the worst tasks for the theory of ordinary differential equations, because the system to be integrated is a stiff one. To better understand the problem, let us consider the case of a homogeneous layer that introduces no coupling between the diffraction orders. The solution of the diffraction problem contains waves propagating up- and downwards (in z-direction). These are plane waves, propagating or evanescent inside the layer. In lossless medium, their constant of propagation in z can be real, or imaginary, depending on the number of diffraction order under consideration:

$$
\begin{equation*}
\gamma_{\mathrm{m}}= \pm \sqrt{\left(\mathrm{k}_{0} \mathrm{n}\right)^{2}-\alpha_{\mathrm{m}}^{2}} \tag{7.111}
\end{equation*}
$$

with

$$
\begin{equation*}
\alpha_{\mathrm{m}}=\alpha_{0}+\mathrm{mK} . \tag{7.112}
\end{equation*}
$$

The real values of $\gamma$ are bounded by $\mathrm{k}_{0} \mathrm{n}$, but the imaginary parts are not bounded, as their asymptotic values for large $|\mathrm{m}|$ are given by:

$$
\begin{equation*}
\operatorname{Im}\left(\gamma_{\mathrm{m}}\right)= \pm|\mathrm{m}| \mathrm{K} . \tag{7.113}
\end{equation*}
$$

From the point of view of the theory of ordinary differential equations this means that the eigenvalues of the system differ significantly in magnitude, i.e., the differential system is stiff. The greater the difference, the more unstable the solution. On the other hand, the solution of the diffraction problem requires sufficient number of Fourier components of the profile function and electromagnetic field to be correctly represented by the truncated Fourier series, thus the necessity to work with large number of Fourier components, and thus the increasingly greater the stiffness of the differential system, i.e., more instable the solution with respect to the length and number of integration steps. The theory concludes that the so called explicit integration schemes are most instable for such problem, whatever their order, and implicit methods have to be used. The problem with the implicit methods is that they need one matrix inversion and several more matrix operations on each integration step, when compared with the explicit methods, so that the choice is not evident to ensure the most efficient integration scheme.

Let us recall the basic principle of the first-order explicit and implicit schemes. In a first-order approximation, the solution of the differential system:

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{dz}} \mathrm{~F}(\mathrm{z})=\mathrm{M}(\mathrm{z}) \mathrm{F}(\mathrm{z}) . \tag{7.114}
\end{equation*}
$$

between two consecutive points $\mathrm{z}=\mathrm{z}_{\mathrm{j}}$ and $\mathrm{z}=\mathrm{z}_{\mathrm{j}+1}$ can be searched in developing in series:

$$
\begin{equation*}
F\left(z_{j+1}\right)=F\left(z_{j}\right)+\left(z_{j+1}-z_{j}\right) M(z) F(z) . \tag{7.115}
\end{equation*}
$$

If $M(z)$ and $F(z)$ are evaluated in $z=z_{j}$, this leads to the first-order explicit integration (Euler's) scheme:

$$
\begin{equation*}
F\left(z_{j+1}\right)=\left[\mathbb{I}+h M\left(z_{j}\right)\right] F\left(z_{j}\right) . \tag{7.116}
\end{equation*}
$$

where $\mathbb{I}$ is the unit matrix, and $h=\left(z_{j+1}-z_{j}\right)$.
If $\mathrm{M}(\mathrm{z})$ and $\mathrm{F}(\mathrm{z})$ are evaluated in $\mathrm{z}=\mathrm{z}_{\mathrm{j}+1}$, we obtain the first-order implicit (inverted or backward Euler's) scheme:

$$
\begin{equation*}
F\left(z_{j+1}\right)=\left[\mathbb{I}-h M\left(z_{j+1}\right)\right]^{-1} F\left(z_{j}\right) \tag{7.117}
\end{equation*}
$$

The theory says that this scheme is more stable, but it needs one matrix inversion on each step. A combination of the two must provide even better results, because it uses a half of the previous step:

$$
\begin{equation*}
\mathrm{F}\left(\mathrm{z}_{\mathrm{j}+1}\right)=\left[\mathbb{I}-\frac{1}{2} \mathrm{hM}\left(\mathrm{z}_{\mathrm{j}+1}\right)\right]^{-1}\left[\mathbb{I}+\frac{1}{2} \mathrm{hM}\left(\mathrm{z}_{\mathrm{j}}\right)\right] \mathrm{F}\left(\mathrm{z}_{\mathrm{j}}\right) . \tag{7.118}
\end{equation*}
$$

However, we need one additional matrix multiplication. In what follows we use these two single-point first-order methods under the names Expl 1 (single point explicit Euler integration) and Impl 1, eq.(7.118) and compare the convergence with respect to the total number of integration points with several other more sophisticated integration schemes for two different metallic gratings in TM polarization, the most difficult combination when using the differential method.

The advantage of these formulations is that they all are single-step ones, and do not need a storage of the intermediate results on several integration steps. They can be easily programmed and don't need additional memory storage at each step. However, if we refer to one of the most relevant sources [7.25], we see that "this is the generic disease of stiff equations: We are required to follow the variation in the solution on the shortest length scale to maintain the stability of the integration, even though accuracy requirements allow a much larger stepsize." This means that a priori choice of the integration step without adaptive control and change in the step length cannot produce stable and relevant results. Unfortunately, it is quite difficult to use adaptive-step methods, because they require much
longer computation times, as it is necessary to repeat the integration process several times when changing the integration step length. This is why we concentrate our attention to fixedstep algorithms.

Fixed-step multistep explicit methods have been used from decades in the differential method programming. The best results have been obtained when combined with an implicit correction by using a predictor-corrector scheme, as described further on. However, referring again to [7.25], "high order does not always mean high accuracy." It will be more useful, if larger integration step is obtained with high order or multistep methods, which is not obvious, as we observe on several numerical examples.

We have used three simple integration schemes, the single-point implicit or explicit scheme, as well as a 4-point predictor-corrector method (PCM 4). It contains two steps, the first one representing an Adams-Bashforth explicit 4-point scheme [7.26], described by the equation

$$
\begin{align*}
& F\left(z_{j+5}\right)=F\left(z_{j+4}\right) \\
& \quad+h\left[\frac{1901}{720} F^{\prime}\left(z_{j+4}\right)-\frac{1387}{360} F^{\prime}\left(z_{j+3}\right)+\frac{108}{30} F^{\prime}\left(z_{j+2}\right)-\frac{637}{360} F^{\prime}\left(z_{j+1}\right)+\frac{251}{720} F^{\prime}\left(z_{j}\right)\right] . \tag{7.119}
\end{align*}
$$

with

$$
\begin{equation*}
\mathrm{F}^{\prime}\left(\mathrm{z}_{\mathrm{j}}\right)=\mathrm{M}\left(\mathrm{z}_{\mathrm{j}}\right) \mathrm{F}\left(\mathrm{z}_{\mathrm{j}}\right) . \tag{7.120}
\end{equation*}
$$

The corrector step is a 4-point Adams-Moulton integration:

$$
\begin{align*}
& F\left(z_{j+4}\right)=F\left(z_{j+3}\right) \\
& \quad+h\left[\frac{251}{720} F^{\prime}\left(z_{j+4}\right)+\frac{646}{720} F^{\prime}\left(z_{j+3}\right)-\frac{264}{720} F^{\prime}\left(z_{j+2}\right)+\frac{106}{360} F^{\prime}\left(z_{j+1}\right)-\frac{19}{720} F^{\prime}\left(z_{j}\right)\right], \tag{7.121}
\end{align*}
$$

which is an implicit scheme [7.27]. However, contrary to the other explicit schemes (BDF), it does not require inverting a matrix, it just makes one step back as a corrector.

An extension of eq.(7.117) to a multistep algorithm results in multistep implicit method, called also backward differentiation formulae (BDF) [7.28]. Typical 3-point and 5-point formulae take the form:
BDF3: $\quad F\left(z_{j+3}\right)=\left[\mathbb{I}-h M\left(z_{j+3}\right)\right]^{-1}\left[\frac{18}{11} F\left(z_{j+2}\right)-\frac{9}{11} F\left(z_{j+1}\right)+\frac{2}{11} F\left(z_{j}\right)\right]$.
BDF5: $\quad F\left(z_{j+5}\right)$
$=\left[\mathbb{I}-h M\left(z_{j+5}\right)\right]^{-1}\left[\frac{300}{137} F\left(z_{j+4}\right)-\frac{300}{137} F\left(z_{j+3}\right)+\frac{200}{137} F\left(z_{j+2}\right)-\frac{75}{137} F\left(z_{j+1}\right)+\frac{12}{137} F\left(z_{j}\right)\right]$.
It is evident that BDF3 (called further on Impl 3) requires a starting method for the first two points, and BDF5 (Impl 5), for the first 4 points. The same is valid for PCM in eqs.(7.119) (7.121).

The second-order Runge-Kutta method is given in the form:

$$
\begin{align*}
\mathrm{k}_{1} & =\mathrm{hM}\left(\mathrm{z}_{\mathrm{j}}\right) \mathrm{F}\left(\mathrm{z}_{\mathrm{j}}\right) \\
\mathrm{k}_{2} & =\mathrm{hM}\left(\mathrm{z}_{\mathrm{j}+1 / 2}\right)\left[\mathrm{F}\left(\mathrm{z}_{\mathrm{j}}\right)+\frac{1}{2} \mathrm{k}_{1}\right]  \tag{7.124}\\
\mathrm{F}\left(\mathrm{z}_{\mathrm{j}+1}\right) & =\mathrm{F}\left(\mathrm{z}_{\mathrm{j}}\right)+\mathrm{k}_{2}
\end{align*}
$$

which has an error proportional to $h^{3}$. It is also called a midpoint point, because it requires the evaluation of the functions at the middle of the step, i.e., twice the number of the steps of the other tree methods discussed above.

The classical fourth-order Runge-Kutta method also uses midpoint values:

$$
\begin{align*}
\mathrm{k}_{1} & =\mathrm{hM}\left(\mathrm{z}_{\mathrm{j}}\right) \mathrm{F}\left(\mathrm{z}_{\mathrm{j}}\right) \\
\mathrm{k}_{2} & =\mathrm{hM}\left(\mathrm{z}_{\mathrm{j}+1 / 2}\right)\left[\mathrm{F}\left(\mathrm{z}_{\mathrm{j}}\right)+\frac{1}{2} \mathrm{k}_{1}\right] \\
\mathrm{k}_{3} & =\mathrm{hM}\left(\mathrm{z}_{\mathrm{j}+1 / 2}\right)\left[\mathrm{F}\left(\mathrm{z}_{\mathrm{j}}\right)+\frac{1}{2} \mathrm{k}_{2}\right]  \tag{7.125}\\
\mathrm{k}_{4} & =\mathrm{hM}\left(\mathrm{z}_{\mathrm{j}+1 / 2}\right)\left[\mathrm{F}\left(\mathrm{z}_{\mathrm{j}}\right)+\mathrm{k}_{3}\right] \\
\mathrm{F}\left(\mathrm{z}_{\mathrm{j}+1}\right) & =\mathrm{F}\left(\mathrm{z}_{\mathrm{j}}\right)+\frac{1}{6} \mathrm{k}_{1}+\frac{1}{3} \mathrm{k}_{2}+\frac{1}{3} \mathrm{k}_{3}+\frac{1}{6} \mathrm{k}_{4}
\end{align*}
$$

with an error of the order of $h^{5}$. These two methods are higher-order explicit methods that do not need matrix inversions during the integration. However, as already stressed, higher order does not mean larger steps.

The highest possible order in h can be obtained theoretically, by using the eigentechnique:

$$
\begin{equation*}
\mathrm{F}\left(\mathrm{z}_{\mathrm{j}+1}\right)=\mathrm{V}\left(\mathrm{e}^{\gamma \mathrm{h}}\right) \mathrm{V}^{-1} \mathrm{~F}\left(\mathrm{z}_{\mathrm{j}}\right) \tag{7.126}
\end{equation*}
$$

where V is a matrix containing the eigenvectors of M and the exponential term in the round brackets represents a diagonal matrix constructed using the eigenvalues $\gamma$ of M . In practice, this method does not increase the stability, because it remains a single-point explicit method. Moreover, it requires much longer computation time because of the requirement to solve an eigenvalue/eigenvector problem at each integration step.

The first example concerns a typical commercial sinusoidal aluminum grating that has very high efficiency in TM polarization. It supports a single diffracted order in $-1^{\text {st }}$ order Littrow mount and has a modulation depth-to-period ratio of $40 \%$. Fig. 7.20 presents a numeric test of the efficiency calculated for a different number of integration points using several integrations schemes. Due to the polarization and the grating material, it is necessary to separate the integration into several slices ( 5 in this case) in order to avoid numerical loss of precision, the results of the integration in two consecutive slices connected to each other by the use of the S-matrix algorithm. In Fig. 7.20 b we have presented a part of the results, obtained with 20 slices, instead of 5 . The comparison between the two cases show that 5 slices are sufficient, the weaker oscillation for 20 slices are due mainly to the fact that the horizontal scale is less dense, because the step in the total number of integration points is an integer times the number of slices. The truncation parameter $\mathrm{N}=20$, i.e., totally 41 Fourier harmonics of the field are used in the calculations.

As can be concluded, an absolute precision within $1 \%$ is rapidly obtained whatever the method used, with the total number of points of the order of 200 . However, the predictor corrector method is less stable when the number of points is smaller than 300. Implicit methods are more stable, as expected, and result in an error smaller than $0.1 \%$ even for the number of points less than 200. It is interesting to observe than the first-order implicit method is more stable than the higher implicit methods, probably because it contains a middle-point evaluation of the field derivative, as seen in eq.(7.118). It requires a little bit longer computation time than the other two implicit methods, because of the additional matrix multiplication. The explicit method, which is the fastest one, shows slower convergence, as expected, whereas the performance of the higher-order RK methods competes with the implicit methods.


Fig.7.18. Aluminum grating with period $0.5 \mu \mathrm{~m}$ and depth $0.2 \mu \mathrm{~m}$ used in $-1^{\text {st }}$ order Littrow mount at $0.6328 \mu \mathrm{~m}$ wavelength in TM polarization. Convergence with respect to the total number of integration points, truncation to 41 Fourier harmonics and using 5 (a) and 20 (b) slices in the $S$ matrix algorithm. The acronyms for the methods are defined in the text.

Table 7.1 compares the computation times of the different methods for the two investigated cases (Fig.7.18 and the following Fig.7.19). For comparison, (null) indicates the time without any operation due to the integration, and that is necessary for the construction of the M-matrix and the use of the S-matrix propagation algorithm, as described in Appendix 7.A. The fastest method is the single-point explicit method, but as expected it is less precise given the same number of integration points, Fig.7.18. The implicit single-step middle-point method shows stability similar to the 4 -th order Runge-Kutta method, but is slightly more rapid. The predictor-corrector method is less stable and requires longer computation times.

Table 7.1. Computation times of the different methods described in the text for the two cases with groove depth to groove period equal to $40 \%$ and $200 \%$

| Method | $\mathrm{N}=20$, slices 5 <br> int. points 400, <br> modulation $40 \%$ | $\mathrm{N}=50$, slices 35 <br> int.points 1500, <br> modulation $200 \%$ |
| :---: | :---: | :---: |
| Expl. 1 | 1.41 s | 72.3 s |
| Impl.1 | 3.55 s | 187.6 s |
| Impl.3 | 2.67 s | 157.8 s |
| Impl.5 | 2.81 s | 168.9 s |
| PCM 4 | 2.14 s | 120.0 s |
| RK 2 | 2.70 s | 144.9 s |
| RK 4 | 4.70 s | 252.5 s |
| eigentechnique | 7.54 s | 360.0 s |
| (null) | 0.68 s | 38.5 s |

It is necessary to stress out that in reality, the computation times are shorter than listed in the Table, because when the truncation N is smaller (usually 20 is sufficient), the number of slices for the S-matrix algorithm is smaller (due to the smaller number of evanescent orders taken into account); in addition, the total number of integration points used for constructing the Table are chosen to obtain $0.1 \%$ relative error, whereas in most of the cases just $1 \%$ is sufficient. The computation time grows linearly with the total number of integration points, as well as with the number of slices used in the S -matrix algorithm. The time dependence concerning the truncation N in the Fourier series grows as $\mathrm{N}^{3}-\mathrm{N}^{3.5}$, because this parameter determines the size of the matrices.

When the total integration length is multiplied by 5 , the number of integration points required is also multiplied by the same factor, as observed in Fig.7.19. A grating twice as deep as the period, acts almost like a flat mirror in TM polarization, with the efficiency in order -1 hardly exceeding $1 \%$. Due to the large depth, the absorption is increased, so that the reflectivity in order 0 is equal to $56.78 \%$. We compare the convergence in the weak $-1^{\text {st }}$ order, so that even a small absolute error appears as a large relative error that can be easily observed in the figures. The number of Fourier harmonics (truncation parameter $2 \mathrm{~N}+1$ ) also has to be increased by a factor of 2.5 to $101(\mathrm{~N}=50)$. The number of slices in the S-matrix algorithm is increased seventh-fold to 35 .

The first Fig.7.19 compares several explicit integration schemes with the single-point implicit method. The main conclusion to be drawn is that the best scheme remains the implicit method, only the explicit Runge-Kutta fourth-order scheme seems to compete in convergence rate with respect to the total number of integration points, but somehow slower.

The comparison of several implicit methods confirms the general idea that multistep choice does not necessarily improve the stability (Fig.7.20). When compared with Fig.7.19, the implicit methods are characterized by smaller oscillations when the number of steps is increased, but the most rapid convergence is obtained with the simplest procedure, single-step method (let us remark again that we use the middle-point calculations, as in eq.(7.118)). Like all the other implicit methods, it requires a single matrix inversion on each integration step, but needs less memory storage, and avoids several matrix sums and multiplication by different constants, necessary for the multipoint methods.


Fig.7.19. Aluminum grating with period $0.5 \mu \mathrm{~m}$ and depth $1 \mu \mathrm{~m}$ used in $-1^{s t}$ order Littrow mount at $0.6328 \mu \mathrm{~m}$ wavelength in TM polarization. Convergence with respect to the total number of integration points, truncation to 101 Fourier harmonics and using 35 slices in the $S$-matrix algorithm.A.M. 4 - forth-order Adams-Moulton scheme, Implicit 1 - single-point implicit scheme, Explicit 1 - single-point explicit scheme, Expl2.R.K. 2 and 4 - explicit Runge-Kutta method of order 2 and4, respectively.


Fig.7.20. Same as in Fig.7.19 but for three implicit methods of different order.

### 7.8. Staircase approximation

As already discussed, if the surface interface is z-invariant (entirely or piecewisely), the integration of the system of ordinary differential equations along $z$ can be done via eigenvalue/eigenvector technique, eq.(7.126), because the M-matrix containing the coefficients of the differential equations does not depend on $z$. The enormous interest in this approach can be explained by the simple technique of integration, much easier to understand
and apply than the theory of numerical methods of integration of ordinary differential equations.

The idea is sketched in Fig.7.21, where a sinusoidal surface-relief grating is approximated by a 5 -stairs profile. While this approximation (with sufficient number of steps M , depending on the groove depth) works quite well in TE polarization, the TM case presents a convergence rate with respect the truncation number of Fourier components of the field much slower than the ordinary differential method (no staircase approximation), see Fig.7.22. Moreover, the greater the number of vertical slices M , the greater the truncation number required.


Fig.7.21. Schematical approximation of a sinusoidal grating profile, approximated by a 5-step staircase profile.(after [7.29]).


Fig.7.22. Convergence of the minus-first-order efficiency in TM polarization of the FMM (RCW) and the exact modal method (indicated on the figure) for a sinusoidal grating in a staircase presentation with $\tilde{\mathrm{M}}=20$, as compared to the convergence of the differential method for a smooth sinusoidal profile (curve "diff."). Period $d=0.5 \mu \mathrm{~m}$, groove depth $a=0.2 \mu \mathrm{~m}$, aluminum refractive index $n_{A l}=1.3+$ i7.6, illuminated at $40^{\circ}$ incidence with wavelength $\lambda=0.6328 \mu \mathrm{~m}$, (after [7.29)].

A detailed analysis of this problem can be found in [7.14, 7.29], but the basic idea is quite simple. The staircase approximation substitutes the otherwise smooth sinusoidal profile by a profile that has sharp edges. The greater the number of stairs, the greater is the number of edges. It is well-known from general electromagnetism that edges introduce electric field singularities. While in TE polarization the only electric field components are tangential to the profile (in y-direction), thus have no discontinuities and singularities, this is not the case in

TM polarization. This can be observed in Fig.7.23. At the edges of each step, a sharp maximum of the electric field is observed. These maxima are not a numerical artifact, they represent the physical effect of introducing edges to replace a smooth profile. These sharp variations of the field require larger number of Fourier components to be correctly represented. Moreover, the greater the number of slices (stairs), the greater the number of the maxima, thus the greater the truncation number required. Numerical experiment has shown that this phenomenon has nothing to do with the integration (eigenvalue/vector) technique, because the results of the convergence rate and field maps are the same for the staircase approximation when using the RCW technique or the differential method.


Fig.7.23. Spatial field distribution of $\left|E_{x}\right|^{2}$ in the vicinity of several steps inside a groove of a 10step staircase profile, used to approximate the sinusoidal grating under study in TM polarization. The grating parameters are the same as in Fig.7.22, after [7.29].

On the contrary, if the true smooth profile is treated by the differential method by using a numerical integration of the ordinary differential system with the elements of the M-matrix depending on z , there is no such singularities of the electric field (Fig.7.24), so that the convergence with respect to the number of Fourier harmonics is much faster, provide the correct factorization rules are used (Fig.7.22).

Recently, some authors [7.30] have proposed to maintain the eigenvalue/vector technique, but to use the correctly determined Fourier presentation of the profile, i.e., the correct factorization rules, as presented in eqs. (7.39) - (7.44), instead of lamellar-profile factorization, eqs.(7.56) - (7.65), at each step. This is equivalent to using the formulation proposed by the differential method for a smooth profile, i.e. avoiding the field singularities at the edges, but to use the eigenvalue/vector technique of integration by assuming that the modified M-matrix, as given by eqs. (7.39) - (7.44), is z-invariant across each step height. We have already tried this in [7.29] and the conclusion was that using this approach, the number of steps (stairs) has to be relatively larger that by using some better adapted integration technique. And indeed, the eigenvalue/approach to a z-dependent system is equivalent to the rectangular rule with equidistant points of integration, one of the worst choices, as known from the theory of ordinary differential equations. In addition, due to
eigenvalue/vector evaluation on each integration step, its computation times are several times longer than for the other methods (see Table 7.1 in the previous section), known from the theory of ordinary differential equations. This is why the authors of [7.30] need more than 2000 equidistant points of integration for a trapezoidal profile, for which the better adapted numerical integration scheme can suffice with 300 points. Unfortunately, the authors of [7.30] do not consider the differential method as a "reference method" in their work.


Fig.7.24. The same as in Fig.7.23 but calculated using the differential method, after [7.29].

## Appendix 7.A: S-matrix propagation algorithm

Almost all electromagnetic theories work by providing the link between the electromagnetic field amplitude values established on two different interfaces. These values could be calculated in the real or the inverted space, or the projections of the field on some functional basis, etc. Whatever the theory, if the media are linear, the link can be expressed in a matrix form:

$$
\begin{equation*}
\mathrm{A}_{\mathrm{p}}=\mathrm{T}_{\mathrm{p}} \mathrm{~A}_{\mathrm{p}-1} \tag{7.127}
\end{equation*}
$$

Here, A stands for a column vector containing the field amplitudes in the given basis, the first interface has a number $\mathrm{p}-1$, and the second on, p . $\mathrm{T}_{\mathrm{p}}$ is called transmission matrix between the interface ( $\mathrm{p}-1$ ) and p .

Numerical problem arises due to the fact, that the "propagation" between different interfaces contains, in general, both growing and decreasing terms, due to both absorption losses or/and evanescent character of some field components. If a real field term propagates from $\mathrm{p}-1$ to p (the green arrow in Fig.7.A.1), it never grows (unless media with optical gains). Same is valid for the true propagation from interface $p$ to $p-1$. However, eq. (7.127) is asymmetrical, i.e., it contains propagation only from interface $p-1$ to $p$, thus a naturally decreasing field that propagates in the opposite direction (from p to $\mathrm{p}-1$ ), will be expressed in the T-matrix in the form of growing terms (the red arrow in Fig. 7.A.1). If the propagation length is sufficiently large, these artificial growing terms can overweight the other terms, mainly due to the finite numerical length of the computer word.


Fig.7.A.1. Schematical representation of the action of the T-matrix between interfaces $p-1$ and $p$
One approach that overcomes this problem and that has become quite popular during the last 15 years is the so-called S-matrix propagation algorithm, S staying for 'scattering'. The basic idea is quite simple: As far as the problem of growing terms has been identified, let us try to do as Nature, by determining another link between the field amplitudes, by separating them into terms propagating (or decreasing) in direction ( $\mathrm{p}-1 \rightarrow \mathrm{p}$ ) or in direction ( $\mathrm{p} \rightarrow \mathrm{p}-1$ ). Let us denote the first set with superscript + , and the second set by a superscript - . The S-matrix between the two interfaces provides the following link:

$$
\begin{equation*}
\binom{A_{p}^{+}}{A_{p-1}^{-}}=S_{p, p-1}\binom{A_{p-1}^{+}}{A_{p}^{-}} \tag{7.128}
\end{equation*}
$$



Fig.7.A.2. Action of the S-matrix between interface p-1 and interface p .

The physical meaning is that amplitude $\mathrm{A}_{\mathrm{p}-1}^{+}$, which propagates from $\mathrm{p}-1$ to p is defined on $\mathrm{p}-1$ and is not growing in-between $\mathrm{p}-1$ and p . In the same manner, the amplitude $A_{p}^{-}$that represents propagation from $p$ to $p-1$ is defined on the interface $p$ and is not growing in direction of interface $\mathrm{p}-1$. To say in other words, the amplitude $\mathrm{A}_{\mathrm{p}-1}^{+}$is incident on the interface $\mathrm{p}-1$ from the previous interface $\mathrm{p}-2$, the second amplitude $A_{p}^{-}$is incident on p from $\mathrm{p}+1$, while the amplitudes on the left-hand side of eq.(7.128) are the amplitudes that are scattered in direction to the outside interfaces ( $p-2$ and $p+1$ ), thus the name of the scattering matrix S. As observed in Fig.7.A.2., the blocks $\mathrm{S}_{\mathrm{p}}^{--}$and $\mathrm{S}_{\mathrm{p}}^{++}$describe the physically correct transmission between p and $\mathrm{p}-1$ or between $\mathrm{p}-1$ and p , respectively, while the other two blocks, $\mathrm{S}_{\mathrm{p}}^{-+}$and $\mathrm{S}_{\mathrm{p}}^{-+}$describe the reflection on the interface p or $\mathrm{p}-1$, respectively. This interpretation explains why there are no numerical problems due to the growing non-physical interactions when using the S-matrix.

The advantage of this formalism is the absence of artificially growing terms in S . The inconvenience is that electromagnetic theories cannot give a direct expression of the matrix S . However, it is possible to express it by using the T-matrix elements, if it is possible to calculate them correctly. If the 'distance' between interface $\mathrm{p}-1$ and p is quite large (with respect to the growing speed of the growing terms), there is loss of precision in determining the T-matrix. The problem can be solved by introducing additional artificial interfaces between $\mathrm{p}-1$ and p in a such manner that to be able to correctly calculate the T-matrix in each subslice. Once the T-matrix calculated, the S-matrix can be obtained in a closed form. However, the total electromagnetic problem of diffraction (or scattering) requires the knowledge of the entire S-matrix of the system, because the physical problem to be solved needs to express the scattered fields as a function of the fields incident on the system (or generated inside, as is the case for electromagnetic antennas). There exists an iterative algorithm that enables us to establish the total S-matrix without calculating the elementary Smatrix between each consecutive pairs of interfaces, as stated in eq. (7.128). For that sake, we
define another intermediate S-matrix that corresponds to the scattering between some initial interface (numbered as 0 ) and the interface with number $\mathrm{p}, \mathrm{S}_{\mathrm{p}} \equiv \mathrm{S}_{\mathrm{p}, 0}$ :

$$
\begin{equation*}
\binom{A_{p}^{+}}{A_{0}^{-}}=S_{p}\binom{A_{0}^{+}}{A_{p}^{-}} \tag{7.129}
\end{equation*}
$$

The initializing values of $\mathrm{S}_{0}$ for $\mathrm{p}=0$ are just the elements of the unity matrix.
As already said, it is necessary to be able to calculate the T-matrices for each intermediate medium between the interfaces. When advancing from the interface $p$ to $p+1$, we obtain the T-matrix with subscript $\mathrm{p}+1$ :

$$
\begin{equation*}
\binom{A_{p+1}^{+}}{A_{p+1}^{-}}=T_{p+1}\binom{A_{p}^{+}}{A_{p}^{-}} \tag{7.130}
\end{equation*}
$$

That will be expanded in the form:

$$
\binom{A_{p+1}^{+}}{A_{p+1}^{-}}=\left(\begin{array}{cc}
T_{p+1}^{++} & T_{p+1}^{+-}  \tag{7.131}\\
T_{p+1}^{++} & T_{p+1}^{--}
\end{array}\right)\binom{A_{p}^{+}}{A_{p}^{-}}
$$

It is obvious from the previous considerations that the growing terms are potentially present in the block $\mathrm{T}_{\mathrm{p}+1}^{--}$('antipropagation' from p to $\mathrm{p}+1$ ), while the blocks $\mathrm{T}_{\mathrm{p}+1}^{++}$and $\mathrm{T}_{\mathrm{p}+1}^{-+}$can contain decreasing terms ('propagation from p to $\mathrm{p}+1$ ), i.e., it could be numerically instable to invert them.

On the other hand, the 'next' S-matrix will link the amplitudes with index 0 to the amplitudes ( $\mathrm{p}+1$ ):

$$
\begin{equation*}
\binom{\mathrm{A}_{\mathrm{p}+1}^{+}}{\mathrm{A}_{0}^{-}}=\mathrm{S}_{\mathrm{p}+1}\binom{\mathrm{~A}_{0}^{+}}{\mathrm{A}_{\mathrm{p}+1}^{-}} \tag{7.132}
\end{equation*}
$$

Eqs. (7.129)-(7.132) enable us to express the matrix $S_{p+1}$ as a function of $S_{p}$ and $T_{p+1}$. At first, we express $A_{p+1}^{-}$from eq.(7.131) and substitute $A_{p}^{+}$from eq.(7.129):

$$
\begin{equation*}
A_{p+1}^{-}=T_{p+1}^{-+} A_{p}^{+}+T_{p+1}^{--} A_{p}^{-}=T_{p+1}^{-+} S_{p}^{++} A_{0}^{+}+\left(T_{p+1}^{-+} S_{p}^{+-}+T_{p+1}^{--}\right) A_{p}^{-} \tag{7.133}
\end{equation*}
$$

Let us denote as $\mathbb{Z}_{\mathrm{p}+1}=\left(\mathrm{T}_{\mathrm{p}+1}^{-+} \mathrm{S}_{\mathrm{p}}^{+-}+\mathrm{T}_{\mathrm{p}+1}^{--}\right)^{-1}$ in order to eliminate $\mathrm{A}_{\mathrm{p}}^{-}$:

$$
\begin{equation*}
A_{p}^{-}=\mathbb{Z}_{p+1} A_{p+1}^{-}-\mathbb{Z}_{p+1} T_{p+1}^{-+} S_{p}^{++} A_{0}^{+} \tag{7.134}
\end{equation*}
$$

The next step is to expand the first line of eq.(7.148):

$$
\begin{align*}
A_{p+1}^{+} & =T_{p+1}^{++} A_{p}^{+}+T_{p+1}^{+-} A_{p}^{-}=T_{p+1}^{++} S_{p}^{++} A_{0}^{+}+\left(T_{p+1}^{+-}+T_{p+1}^{++} S_{p}^{+-}\right) A_{p}^{-} \\
& =T_{p+1}^{++} S_{p}^{++} A_{0}^{+}+\left(T_{p+1}^{+-}+T_{p+1}^{++} S_{p}^{+-}\right)\left(\mathbb{Z}_{p+1} A_{p+1}^{-}-\mathbb{Z}_{p+1} T_{p+1}^{-+} S_{p}^{++} A_{0}^{+}\right)  \tag{7.135}\\
& =\left[T_{p+1}^{++} S_{p}^{++}-\left(T_{p+1}^{+-}+T_{p+1}^{++} S_{p}^{+-}\right) \mathbb{Z}_{p+1} T_{p+1}^{-+} S_{p}^{++}\right] A_{0}^{+}+\left(T_{p+1}^{+-}+T_{p+1}^{++} S_{p}^{+-}\right) \mathbb{Z}_{p+1} A_{p+1}^{-}
\end{align*}
$$

The comparison with eq.(7.132) gives the first two block-elements of $\mathrm{S}_{\mathrm{p}+1}$ :

$$
\begin{gather*}
\mathrm{S}_{\mathrm{p}+1}^{+-}=\left(\mathrm{T}_{\mathrm{p}+1}^{+-}+\mathrm{T}_{\mathrm{p}+1}^{++} \mathrm{S}_{\mathrm{p}}^{+-}\right) \mathbb{Z}_{\mathrm{p}+1} .  \tag{7.136}\\
\mathrm{S}_{\mathrm{p}+1}^{++}=\mathrm{T}_{\mathrm{p}+1}^{++} \mathrm{S}_{\mathrm{p}}^{++}-\left(\mathrm{T}_{\mathrm{p}+1}^{+-}+\mathrm{T}_{\mathrm{p}+1}^{++} \mathrm{S}_{\mathrm{p}}^{+-}\right) \mathbb{Z}_{\mathrm{p}+1} \mathrm{~T}_{\mathrm{p}+1}^{++} \mathrm{S}_{\mathrm{p}}^{++}  \tag{7.137}\\
=\left(\mathrm{T}_{\mathrm{p}+1}^{++}-\mathrm{S}_{\mathrm{p}+1}^{+-} \mathrm{T}_{\mathrm{p}+1}^{-+}\right) \mathrm{S}_{\mathrm{p}}^{++}
\end{gather*} .
$$

From eq.(7.129)

$$
\begin{align*}
\mathrm{A}_{0}^{-}= & \mathrm{S}_{\mathrm{p}}^{-+} \mathrm{A}_{0}^{+}+\mathrm{S}_{\mathrm{p}}^{--} \mathrm{A}_{\mathrm{p}}^{-}=\mathrm{S}_{\mathrm{p}}^{-+} \mathrm{A}_{0}^{+}+\mathrm{S}_{\mathrm{p}}^{--}\left(\mathbb{Z}_{\mathrm{p}+1} \mathrm{~A}_{\mathrm{p}+1}^{-}-\mathbb{Z}_{\mathrm{p}+1} \mathrm{~T}_{\mathrm{p}+1}^{-+} \mathrm{S}_{\mathrm{p}}^{++} \mathrm{A}_{0}^{+}\right)  \tag{7.138}\\
& =\left(\mathrm{S}_{\mathrm{p}}^{-+}-\mathrm{S}_{\mathrm{p}}^{--} \mathbb{Z}_{\mathrm{p}+1} \mathrm{~T}_{\mathrm{p}+1}^{-+} \mathrm{S}_{\mathrm{p}}^{++}\right) \mathrm{A}_{0}^{+}+\mathrm{S}_{\mathrm{p}}^{--} \mathbb{Z}_{\mathrm{p}+1} \mathrm{~A}_{\mathrm{p}+1}^{-}
\end{align*}
$$

so that

$$
\begin{gather*}
\mathrm{S}_{\mathrm{p}+1}^{--}=\mathrm{S}_{\mathrm{p}}^{--} \mathbb{Z}_{\mathrm{p}+1}  \tag{7.139}\\
\mathrm{~S}_{\mathrm{p}+1}^{-+}=\mathrm{S}_{\mathrm{p}}^{-+}-\mathrm{S}_{\mathrm{p}}^{--} \mathbb{Z}_{\mathrm{p}+1} \mathrm{~T}_{\mathrm{p}+1}^{-+} \mathrm{S}_{\mathrm{p}}^{++}  \tag{7.140}\\
=\mathrm{S}_{\mathrm{p}}^{-+}-\mathrm{S}_{\mathrm{p}+1}^{--} \mathrm{T}_{\mathrm{p}+1}^{-+} \mathrm{S}_{\mathrm{p}}^{++}
\end{gather*}
$$

These relations exist in several possible forms, but this one is quite well adapted to the case without incident waves on interface 0 , because in the iterative algorithm we need to calculate only the half of the blocks, namely the two given by eqs. (7.136) and (7.139).

The only matrix inversion in the iterative algorithm concerns the procedure to obtain the matrix $\mathbb{Z}$. The initial matrix $\mathbb{Z}^{-1}$ contains the potentially large terms from $T_{p+1}^{--}$, so that its inversion creates neither numerical problems to be inverted, nor growing terms to create numerical instabilities.

## Appendix 7.B: Inverted S-matrix propagation algorithm

In Appendix A we have seen how to avoid numerical instabilities due to the artificially growing terms that appear when the propagation of the field amplitudes from one interface to another is made in the wrong direction, a typical property of a half of the field amplitudes used in the transmission matrix approach.

In some cases (for example, the Integral method applied to multilayer grating, but also coordinate transformation method used for a stack containing different profiles), the numerical solution that has been obtained provides a link, having a form inverse to eq.(7.130)

$$
\begin{equation*}
\tilde{T}_{p+1}\binom{A_{p+1}^{+}}{A_{p+1}^{-}}=\binom{A_{p}^{+}}{A_{p}^{-}} . \tag{7.141}
\end{equation*}
$$

Of course, it is easy to obtain the form of eq. (7.130) by simply inverting $T_{p+1}$, but better to avoid this, because some blocks of the matrix contain large terms compared to the others. In particular, the block $\tilde{\mathrm{T}}_{\mathrm{p}+1}^{++}$is responsible for a physical 'antipropagation' from $\mathrm{p}+1$ to p , so that potentially it contains growing terms (as it was that case with $\mathrm{T}_{\mathrm{p}+1}^{--}$) in Appendix A.

We can avoid the direct inversion of $\tilde{T}_{p+1}$ by applying a similar procedure as in Appendix 7.A in order to obtain the S-matrix of the stack. Equation (7.141) is expanded in the form:

$$
\begin{align*}
& A_{p}^{+}=\tilde{T}_{p+1}^{++} A_{p+1}^{+}+\tilde{T}_{p+1}^{+-} A_{p+1}^{-}  \tag{7.142}\\
& A_{p}^{-}=\tilde{T}_{p+1}^{-+} A_{p+1}^{+}+\tilde{T}_{p+1}^{--} A_{p+1}^{-} \tag{7.143}
\end{align*}
$$

On the other hand, from eq.(7.129) we have:

$$
\begin{equation*}
\mathrm{A}_{\mathrm{p}}^{+}=\mathrm{S}_{\mathrm{p}}^{++} \mathrm{A}_{0}^{+}+\mathrm{S}_{\mathrm{p}}^{+-} \mathrm{A}_{\mathrm{p}}^{-} \tag{7.144}
\end{equation*}
$$

so that

$$
\begin{gather*}
\mathrm{S}_{\mathrm{p}}^{++} \mathrm{A}_{0}^{+}+\mathrm{S}_{\mathrm{p}}^{+-} \mathrm{A}_{\mathrm{p}}^{-}=\tilde{\mathrm{T}}_{\mathrm{p}+1}^{++} \mathrm{A}_{\mathrm{p}+1}^{+}+\tilde{\mathrm{T}}_{\mathrm{p}+1}^{+-} \mathrm{A}_{\mathrm{p}+1}^{-}  \tag{7.145}\\
\mathrm{S}_{\mathrm{p}}^{++} \mathrm{A}_{0}^{+}+\mathrm{S}_{\mathrm{p}}^{+-}\left(\tilde{\mathrm{T}}_{\mathrm{p}+1}^{-+} \mathrm{A}_{\mathrm{p}+1}^{+}+\tilde{\mathrm{T}}_{\mathrm{p}+1}^{--} \mathrm{A}_{\mathrm{p}+1}^{-}\right)=\tilde{\mathrm{T}}_{\mathrm{p}+1}^{++} \mathrm{A}_{\mathrm{p}+1}^{+}+\tilde{\mathrm{T}}_{\mathrm{p}+1}^{+-} \mathrm{A}_{\mathrm{p}+1}^{-} \tag{7.146}
\end{gather*}
$$

and

$$
\begin{equation*}
\mathrm{S}_{\mathrm{p}}^{++} \mathrm{A}_{0}^{+}=\left(\tilde{\mathrm{T}}_{\mathrm{p}+1}^{++}-\mathrm{S}_{\mathrm{p}}^{+-} \tilde{\mathrm{T}}_{\mathrm{p}+1}^{-+}\right) \mathrm{A}_{\mathrm{p}+1}^{+}+\left(\tilde{\mathrm{T}}_{\mathrm{p}+1}^{+-}-\mathrm{S}_{\mathrm{p}}^{+-\tilde{T}_{\mathrm{p}+1}^{--}}\right) \mathrm{A}_{\mathrm{p}+1}^{-} \tag{7.147}
\end{equation*}
$$

Now we can identify half of the blocks of $S_{p+1}$ from eq.(7.132):

$$
\begin{gather*}
\mathrm{S}_{\mathrm{p}+1}^{++}=\tilde{\mathbb{Z}}_{\mathrm{p}+1} \mathrm{~S}_{\mathrm{p}}^{++}  \tag{7.148}\\
\mathrm{S}_{\mathrm{p}+1}^{+-}=-\tilde{\mathbb{Z}}_{\mathrm{p}+1}\left(\tilde{\mathrm{~T}}_{\mathrm{p}+1}^{+-}-\mathrm{S}_{\mathrm{p}}^{+-} \tilde{\mathrm{T}}_{\mathrm{p}+1}^{--}\right) \tag{7.149}
\end{gather*}
$$

with $\tilde{\mathbb{Z}}_{\mathrm{p}+1}=\left(\tilde{\mathrm{T}}_{\mathrm{p}+1}^{++}-\mathrm{S}_{\mathrm{p}}^{+-} \tilde{\mathrm{T}}_{\mathrm{p}+1}^{-+}\right)^{-1}$ that contains the numerically dangerous growing terms in $\tilde{\mathrm{T}}_{\mathrm{p}+1}^{++}$, in the same manner that the matrix $\mathbb{Z}_{\mathrm{p}+1}$ in Appendix 7.A 'envelopes' the growing terms in $\mathrm{T}_{\mathrm{p}+1}^{--}$.

The other two block can be obtained by staring with the identity:

$$
\begin{equation*}
\mathrm{A}_{0}^{-}=\mathrm{S}_{\mathrm{p}}^{-+} \mathrm{A}_{0}^{+}+\mathrm{S}_{\mathrm{p}}^{--} \mathrm{A}_{\mathrm{p}}^{-} \tag{7.150}
\end{equation*}
$$

and using eq.(7.143) :

$$
\begin{equation*}
\mathrm{A}_{0}^{-}=\mathrm{S}_{\mathrm{p}}^{-+} \mathrm{A}_{0}^{+}+\mathrm{S}_{\mathrm{p}}^{--}\left(\tilde{T}_{\mathrm{p}+1}^{-+} \mathrm{A}_{\mathrm{p}+1}^{+}+\tilde{\mathrm{T}}_{\mathrm{p}+1}^{--} \mathrm{A}_{\mathrm{p}+1}^{-}\right) \tag{7.151}
\end{equation*}
$$

When taking into account that two blocks of $\mathrm{S}_{\mathrm{p}+1}$ are already known and given in eqs. (7.148) and (7.149), we can eliminate $A_{p+1}^{+}=S_{p+1}^{++} \mathrm{A}_{0}^{+}+\mathrm{S}_{\mathrm{p}+1}^{+-} \mathrm{A}_{\mathrm{p}+1}^{-}$:

$$
\begin{equation*}
A_{0}^{-}=\left(S_{p}^{-+}+S_{p}^{--} \tilde{\mathrm{T}}_{\mathrm{p}+1}^{-+} \mathrm{S}_{\mathrm{p}+1}^{++}\right) \mathrm{A}_{0}^{+}+\mathrm{S}_{\mathrm{p}}^{--}\left(\tilde{\mathrm{T}}_{\mathrm{p}+1}^{--}+\tilde{\mathrm{T}}_{\mathrm{p}+1}^{-+} \mathrm{S}_{\mathrm{p}+1}^{+-}\right) \mathrm{A}_{\mathrm{p}+1}^{-} . \tag{7.152}
\end{equation*}
$$

Thus

$$
\begin{gather*}
S_{p+1}^{-+}=S_{p}^{-+}+S_{p}^{--} \tilde{T}_{p+1}^{-+} S_{p+1}^{++} .  \tag{7.153}\\
S_{p+1}^{--}=S_{p}^{--}\left(\tilde{T}_{p+1}^{--}+\tilde{T}_{p+1}^{-+} S_{p+1}^{+-}\right) . \tag{7.154}
\end{gather*}
$$

The expression are quite similar in form to those obtained in Appendix A. Moreover, they allow avoiding the inversion of $\tilde{T}_{p+1}$.

Finally, there exist a combination of expressions including partial T-matrices, treated separately in Appendix 7.A and 7.B. In some cases the link between the amplitudes on two consecutive interfaces or across a single interface that separates two different media can be expressed in the form:

$$
\begin{equation*}
\tilde{\mathfrak{T}}_{p+1}\binom{A_{p+1}^{+}}{A_{p+1}^{-}}=\mathfrak{T}_{p+1}\binom{A_{p}^{+}}{A_{p}^{-}} . \tag{7.155}
\end{equation*}
$$

Such is the case of the Fourier-modal (RCW) method across each interface, with the partial transmission matrices $\tilde{\mathfrak{T}}_{\mathrm{p}+1}$ containing the eigenvectors of the proper modes inside each media. The same expression is obtained in the coordinate transformation method when using eigenvalue technique of integration. Usually, in both approaches, one obtains the full transmission matrix by inverting $\tilde{\mathfrak{T}}_{\mathrm{p}+1}$ and multiplying the result by $\mathfrak{T}_{\mathrm{p}+1}$. If this creates numerical problems (for thick layers), such direct approach is not applicable. In that case it is better advised to apply twice the S-matrix algorithm, at first in each direct form (Appendix 7.A), and then in the currently discussed inverted form. It is quite easy to understand the logic, by introducing a virtual set of amplitudes in eq.(7.155):

$$
\begin{align*}
& \binom{\tilde{A}_{p+1}^{+}}{\tilde{A}_{p+1}^{-}}=\mathfrak{T}_{p+1}\binom{A_{p}^{+}}{A_{p}^{-}},  \tag{7.156}\\
& \tilde{\mathfrak{T}}_{p+1}\binom{A_{p+1}^{+}}{A_{p+1}^{-}}=\binom{\tilde{A}_{p}^{+}}{\tilde{A}_{p}^{-}} . \tag{7.157}
\end{align*}
$$

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Chapter 8:

# Coordinate Transformation Methods 

## Gérard Granet

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## Chapter 8

# Coordinate Transformation Methods 

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### 8.1 Introduction

The C-method was born in the eighties in Clermont-Ferrand, France, from the need to solve rigorously diffraction problems at corrugated periodic surfaces in the resonance regime [1], [2], [3]. The main difficulty of such problems is the matching of boundaries conditions. It is obvious that any method aimed at solving Maxwell's equation is all the more efficient since it is able to fit the geometry of the problem. For that purpose, Chandezon et al introduced the so called translation coordinate system deduced from the Cartesian coordinate system $x, y, z$ ) by the relations $x=x^{1}, y=x^{2}, z=x^{3}+a\left(x^{1}\right)$ where $a\left(x^{1}\right)$ is a continuously differentiable function describing the surface profile. Hence since the boundary of the physical problem coincides with coordinate surfaces, writing boundary conditions is as simple as it is for classical problems in Cartesian, cylindrical, or spherical coordinates. This is the first ingredient of C-method. The second one is to write Maxwell's equation under the covariant form. This formulation comes from relativity where the use of curvilinear non orthogonal coordinate system is essential and natural. The main feature of this formalism is that Maxwell's equations remain invariant in any coordinate system, the geometry being shifted into the constitutive relations. Chandezon et al derived their 3D formulation from the general 4D relativistic Post's formalism [4] and evidently used tensorial calculus. Although it is with no doubt the most elegant and efficient way to deal with electromagnetic in general curvilinear coordinates it is also probably the reason why the theory appeared difficult to understand to many scientists. The third ingredient of C -method is that it is a modal method. This nice property is linked with the translation coordinate system in which a diffraction problem may be expressed as an eigenvalue eigenvector problem with periodic boundary conditions. The last feature of C-method is the numerical method of solution. The matrix operator is obtained by expanding field components into Floquet-Fourier harmonics and by projecting Maxwell's equations onto periodic exponential functions. The above four features may be resumed by saying that C-method is a curvilinear coordinate modal method by Fourier expansion [5]. Since the original papers, The C-method has gone through many stages of extension and improvement. The original theory was formulated for uncoated perfectly conducting gratings in classical mount. Various authors extended the method to conical diffraction mountings [6], [7]. Granet et al [11], Li et al [12] and Preist et al [13] allowed the various profiles
of a stack of gratings to be different from each other, although keeping the periodicity. Solving the vertical faces case in a simple manner, Plumey et al [14] have shown that the method can be applied to overhanging gratings Preist et al obtained the same results by applying the usual coordinate transformation to oblique coordinates [15]. In the numerical context, Li [16] and Cotter et al [17] improved the numerical stability of the C-method by using the S-matrix propagation algorithm for multilayer gratings. It is seen that C-method has been applied to a large class of surface relief gratings and multilayer coated gratings. The key point of C-method is the joint use of curvilinear coordinates and covariant formulation of Maxwell's equations. All the new developments in the modelling of gratings like Adaptive Spatial Resolution [18],[19],[20], and Matched coordinate [21]derive from this fundamental observation.

### 8.2 C-Method

In Euclidean space with origin O and basis vector $\boldsymbol{e}_{x}, \boldsymbol{e}_{y}, \boldsymbol{e}_{z}$, let us consider an infinite cylindrical surface $(\Sigma)$ whose elements are parallel to the $y$ axis. This surface separates two linear homogeneous and isotropic media denoted (1) and (2). In Cartesian coordinates such a surface can be described by equation $z=a(x)$. Any electromagnetic field interacting with this particular geometry satisfies some boundary conditions. For instance, the tangential components of the electric field vector and the normal component of the displacement field vector are continuous at the surface. The point is that boundary conditions involve quantities that obviously depend on the position at which they are considered on the surface. We are thus led to look for a coordinate system which fits the problem and makes it more readily solvable than it is in a Cartesian framework. The so-called translation coordinate system $\left(x^{1}, x^{2}, x^{3}\right)$ introduced by Chandezon and defined from the Cartesian coordinate system by the direct transformation (curvilinear coordinates to Cartesian coordinates) :

$$
\begin{equation*}
x=x^{1}, \quad y=x^{2}, \quad z=x^{3}+a\left(x^{1}\right) \tag{8.1}
\end{equation*}
$$

or the inverse transformation (Cartesian coordinates to curvilinear coordinates):

$$
\begin{equation*}
x^{1}=x, \quad x^{2}=y, \quad x^{3}=z-a(x) \tag{8.2}
\end{equation*}
$$

is one such system. It makes the surface $(\Sigma)$ coincide with the coordinate surface $x^{3}=0$. A point $M\left(x, y, z=a(x)\right.$ at the surface $(\Sigma)$ is now referenced by the triplet $\left(x^{1}, x^{2}, 0\right)$. The coordinate surface $x^{3}=x_{0}^{3}$ is obtained by translating each point at surface $(\Sigma)$ with vector $x_{0}^{3} e_{z}$, hence the name given by Chandezon to this particular coordinate system: translation coordinate system. The change of coordinates may also be considered as a change of variable. This view


Figure 8.1: Geometry of the problem: Two media are separated by a cylindrical periodic surface, with period $d^{1}$ described by the equation $z=a(x)$
point allows a better understanding of the numerical behaviour of C-method and its connection with Rayleigh expansions. There is actually no difference in the way of deriving the elementary solutions of the scalar Helmholtz equation in Cartesian coordinates or in translation coordinate systems. Both are eigenvectors of an eigenvalue problem with pseudo-periodic boundary conditions. In both cases, the operator eigenvalue problem is transformed into a matrix eigenvalue problem thanks to the Galerkin method with pseudo periodic functions as expansion and test functions. Hence solving the scalar Helmholtz equation in any coordinate system is the very first step when implementing C-method. In the next paragraphs we shall focus on this issue before solving a grating problem.

### 8.2.1 Modal equation in the Cartesian coordinate system

Consider an homogeneous region with relative, possibly complex, permittivity, $\varepsilon$. In the harmonic regime with a time dependence of $\exp (-i \omega t)$, it is possible to construct general solutions to the field equations once we have general solutions to the scalar Helmholtz equation. So as a first task, we are going to investigate elementary solutions to the Helmholtz equation written in the translation coordinate system. Let us start from Cartesian coordinates in which 2D scalar Helmholtz equation is

$$
\begin{equation*}
\left(\partial_{x}^{2}+\partial_{z}^{2}+k^{2}\right) \mathscr{F}=0 \tag{8.3}
\end{equation*}
$$

where $k=\omega \sqrt{\mu_{0} \varepsilon}$ is the wave-number. The coefficients of the Helmholtz equation are independent of $z$ so we seek solutions of the form $\mathscr{F}(x, z)=\exp (i \gamma z) F(x)$ The Helmholtz equation becomes:

$$
\begin{equation*}
\left(\partial_{x}^{2}+k^{2}\right) F(x)=\gamma^{2} F(x) \tag{8.4}
\end{equation*}
$$

Function $F(x)$ is thus an eigenmode of equation 8.4). The requirement that the eigenmodes satisfy the pseudo-periodicity condition $F\left(x+d^{1}\right)=\exp \left(i \alpha_{O} d^{1}\right) F(x)$ is automatically fulfilled by their expansion into Floquet-Fourier series:

$$
\begin{equation*}
F(x)=\sum_{-\infty}^{+\infty} F_{m} \exp \left(i \alpha_{m} x\right) \tag{8.5}
\end{equation*}
$$

$\alpha_{m}=\alpha_{0}+m K_{1}, K_{1}=\frac{2 \pi}{d^{1}}, m \in \mathbb{N}$ and $\alpha_{0}$ is some real parameter. By introducing (8.5) into (8.3) and by projecting onto pseudo-periodic functions $\exp \left(i \alpha_{n} x\right)$, one obtains the matrix equation:

$$
\begin{equation*}
\gamma^{2} \boldsymbol{F}=\left[k^{2} \boldsymbol{I}-\boldsymbol{\alpha}\right] \boldsymbol{F} \tag{8.6}
\end{equation*}
$$

where $\boldsymbol{F}$ is a column vector whose elements are the $F_{m}$ and $\boldsymbol{\alpha}$ is a diagonal matrix w hose elements are the $\alpha_{m}$ and $\boldsymbol{I}$ is the identity matrix. The solution to the above matrix eigenvalue equation is of course trivial since the matrix is diagonal. Let us introduce subscript $q$ to number the eigenvalues and the eigenfunctions. The eigenvalues $\gamma_{q}$ are deduced from their squared number:

$$
\begin{equation*}
\gamma_{q}^{2}=k^{2}-\alpha_{q}^{2} \tag{8.7}
\end{equation*}
$$

and the eigenvectors are determined by $F_{m q}=\delta_{m q}$ where $\delta_{m q}$ is the Kronecker symbol. The square root of $\gamma_{q}^{2}$ is defined as follows:

$$
\gamma_{q}= \begin{cases}\sqrt{\gamma_{q}^{2}} & \text { if } \quad \gamma_{q}^{2} \in \mathbb{R}^{+}  \tag{8.8}\\ \sqrt{-\gamma_{q}^{2}} & \text { if } \quad \gamma_{q}^{2} \in \mathbb{R}^{-} \\ \left(\gamma_{q}^{2}\right)^{1 / 2} & \text { with positive imaginary part if } \quad \gamma_{q}^{2} \in \mathbb{C}\end{cases}
$$

Finally, $\mathscr{F}(x, z)$ can be represented by superposition of eigenmodes

$$
\begin{equation*}
\mathscr{F}(x, z)=\mathscr{F}^{+}(x, z)+\mathscr{F}^{-}(x, z) \tag{8.9}
\end{equation*}
$$

with

$$
\begin{array}{r}
\mathscr{F}^{+}(x, z)=\sum_{q=-\infty}^{q=+\infty} A_{q}^{+} \exp \left(i \gamma_{q} z\right) \sum_{m=-\infty}^{m=+\infty} \delta_{m q} \exp \left(i \alpha_{m} x\right) \\
\mathscr{F}^{-}(x, z)=\sum_{q=-\infty}^{q=+\infty} A_{q}^{-} \exp \left(-i \gamma_{q} z\right) \sum_{m=-\infty}^{m=+\infty} \delta_{m q} \exp \left(i \alpha_{m} x\right) \tag{8.11}
\end{array}
$$

There are two sets of modes, the number of which are equal: those propagating or decaying in the positive direction of $z$ and those propagating or decaying in the opposite direction. We denote these modes by superscript + and - respectively. The $z$ dependence of an eigenmode is determined by function $\exp \left(i \gamma_{p} z\right)$ By increasing $z$ to $z+\Delta z, \exp \left(i \gamma_{q} z\right)$ is multiplied by $\exp \left(i \gamma_{q} \Delta z\right)=\exp \left(i \Re\left(\gamma_{q}\right) \Delta z\right) \times \exp \left(-\mathfrak{I}\left(\gamma_{q}\right) \Delta z\right)$ The real eigenvalues have $\mathfrak{I}\left(\gamma_{q}\right)=0$ and correspond therefore to forward modes if $\mathfrak{R}\left(\gamma_{q}\right)>0$ or backward modes if $\mathfrak{R}\left(\gamma_{q}\right)<0$. The complex eigenvalues modes have a non-zero imaginary part and possibly also a non-zero real part. The associated eigenmodes decay forward if $\mathfrak{I}\left(\gamma_{q}\right)>0$ or backward if $\mathfrak{I}\left(\gamma_{q}\right)>0$. These expansions are known as Rayleigh expansions; they are linear combination of eigenvectors that we call hereafter Rayleigh eigenvectors $R_{q}$ :

$$
\begin{equation*}
R_{q}(x)=\sum_{m=-\infty}^{m=+\infty} \delta_{m q} \exp \left(i \alpha_{m} x\right) \tag{8.12}
\end{equation*}
$$

In Cartesian coordinates, the solutions to the Helmholtz equations may be regarded as the eigenvectors of a matrix equation. The eigenvalues $\gamma_{m}$ are determined by the periodic lateral boundary conditions of the problem and are obtained analytically since the matrix is diagonal. The translation coordinate system preserves the $z$ translation symmetry and also periodic lateral boundary conditions. We may then expect a great formal similitude between solutions obtained in each coordinate system.

### 8.2.2 Modal equation in terms of the new variables

In this section, we derive the master equation of C -method by considering the change of coordinates as a change of variables. For the change of variables $x^{1}=x, x^{2}=y, x^{3}=z-a(x)$ the chain rule for derivatives has the form:

$$
\left\{\begin{array}{l}
\partial_{x}=\partial_{1}-\dot{a} \partial_{3}  \tag{8.13}\\
\partial_{y}=\partial_{2} \\
\partial_{z}=\partial_{3}
\end{array}\right.
$$

Substituting the derivatives (8.13) into (8.3) gives:

$$
\begin{equation*}
\left((1+\dot{a} \dot{a}) \partial_{3}^{2}-\dot{a} \partial_{1} \partial_{3}-\partial_{1} \dot{a} \partial_{3}+\partial_{1}^{2}+k^{2}\right) \mathscr{F}\left(x^{1}, x^{3}\right)=0 \tag{8.14}
\end{equation*}
$$

the solution of which are the same as the solutions of (8.13) expressed in terms of the new variables.

$$
\begin{align*}
\mathscr{F}_{a}^{+}\left(x^{1}, x^{3}\right) & =\mathscr{F}^{+}\left(x=x^{1}, x^{3}=z-a(x)\right) \\
& =\sum_{q=-\infty}^{q=+\infty} A_{q}^{+} \exp \left(i \gamma_{q} x^{3}\right) \sum_{m=-\infty}^{m=+\infty} \delta_{m q} \exp \left(i \gamma_{q} a\left(x^{1}\right)\right) \exp \left(i \alpha_{m} x^{1}\right) \tag{8.15}
\end{align*}
$$

$$
\begin{align*}
\mathscr{F}_{a}^{-}\left(x^{1}, x^{3}\right) & =\mathscr{F}^{-}\left(x=x^{2}, x^{3}=z-a(x)\right) \\
& =\sum_{q=-\infty}^{q=+\infty} A_{q}^{-} \exp \left(-i \gamma_{q} x^{3}\right) \sum_{m=-\infty}^{m=+\infty} \delta_{m q} \exp \left(-i \gamma_{q} a\left(x^{1}\right)\right) \exp \left(i \alpha_{m} x^{1}\right) \tag{8.16}
\end{align*}
$$

The subscript $a$ indicates the profile dependence of function $\mathscr{F}$.
We call function $\exp \left( \pm i \gamma_{q} a\left(x^{1}\right)\right) \exp \left(i \alpha_{q} x^{1}\right)$ the generalized Rayleigh eigenvector of order $q$. It is nothing more than plane wave $\exp \left( \pm i \gamma_{q} z\right) \exp \left(i \alpha_{q} x^{1}\right)$ expressed in terms of the new variables $x^{1}$ and $x^{3}$ and is closely linked with function $a\left(x^{1}\right)$. Let us denote it $R_{a, q}^{ \pm}$:

$$
\begin{equation*}
R_{a, q}^{ \pm}=\exp \left( \pm i \gamma_{q} a\left(x^{1}\right)\right) \exp \left(i \alpha_{q} x^{1}\right)=\sum_{m=-\infty}^{m=+\infty} R_{a m q} \exp \left(i \alpha_{m} x^{1}\right) \tag{8.17}
\end{equation*}
$$

It is assumed so far that $a\left(x^{1}\right)$ is periodic with period $d^{1}$ hence:

$$
\begin{equation*}
\exp \left( \pm i \gamma_{q} a\left(x^{1}\right)\right)=\sum_{p=-\infty}^{p+\infty} L_{p}^{ \pm} \exp \left(\frac{i 2 \pi p x^{1}}{d^{1}}\right) \tag{8.18}
\end{equation*}
$$

with:

$$
\begin{equation*}
L_{p}^{ \pm}=\frac{1}{d^{1}} \int_{0}^{d^{1}} \exp \left( \pm i \gamma_{q} a\left(x^{1}\right)\right) \exp \left(\frac{-i 2 \pi p x^{1}}{d^{1}}\right) d x^{1} \tag{8.19}
\end{equation*}
$$

In physical space, the generalized Rayleigh eigenvectors result from the product of a periodic function with a pseudo-periodic one. Thus, in Fourier space, the spectrum of the $q$ th generalized Rayleigh eigenvector is obtained by translating the spectrum of function $\exp \left( \pm i \gamma_{q} z\right)$ with vector $2 \pi q / d^{1}$, that is:

$$
\begin{equation*}
R_{a m q}^{ \pm}=L_{a m-q}^{ \pm} \tag{8.20}
\end{equation*}
$$

Finally:

$$
\begin{align*}
& \mathscr{F}_{a}^{+}\left(x^{1}, x^{3}\right)=\sum_{q=-\infty}^{q=+\infty} A_{q}^{+} \exp \left(i \gamma_{q} x^{3}\right) \sum_{m=-\infty}^{m=+\infty} L_{a m-q}^{+} \exp \left(i \alpha_{m} x^{1}\right)  \tag{8.21}\\
& \mathscr{F}_{a}^{-}\left(x^{1}, x^{3}\right)=\sum_{q=-\infty}^{q=+\infty} A_{q}^{-} \exp \left(-i \gamma_{q} x^{3}\right) \sum_{m=-\infty}^{m=+\infty} L_{a m-q}^{-} \exp \left(i \alpha_{m} x^{1}\right) \tag{8.22}
\end{align*}
$$

Functions (8.21) and (8.22) give the general solution to (8.14). Indeed each element of this solution is a generalized Rayleigh eigenvector associated to index $q$ such that $\gamma_{q}^{2}+\alpha_{q}^{2}=k^{2}$ and thus satisfies (8.14). The reason for that is obvious. It is obtain from (8.12) in which we have introduced the same change of variable as the one that has allowed us to get (8.14) from 8.3). From a practical view point, one can only manipulate finite size expansions and it does not make sense to speak of $R_{a, q}^{ \pm}\left(x^{1}\right)$. That is why one may wonder if a generalized Rayleigh eigenvector is still a valid a solution of (8.14) when only a finite number of spatial Fourier harmonics is retained to represent it. Let us assume for a while the answer is yes and examine the involvements of such a claim. Introducing an integer $M$, hereafter denoted truncation number, and letting $m$ run from $-M$ to $M$ the truncated generalized Rayleigh eigenvector writes:

$$
\begin{equation*}
R_{a, q}^{ \pm(M)}\left(x^{1}\right)=\sum_{m=-M}^{m=+M} R_{a, m q}^{ \pm} \exp \left(i \alpha_{m} x^{1}\right) \tag{8.23}
\end{equation*}
$$

Substituting $\partial_{3}$ with $i \gamma_{q}$, we have:

$$
\begin{equation*}
-(1+\dot{a} \dot{a}) \gamma_{q} \dot{R}_{a}^{ \pm(M)}-i\left(\dot{a} \partial_{1}+\partial_{1} \dot{a}\right) \gamma_{q} R_{a}^{ \pm,(M)}+\left(\partial_{1}^{2}+k^{2}\right) R_{a}^{ \pm,(M)}=0 \tag{8.24}
\end{equation*}
$$

where $\dot{R}_{a}^{ \pm(M)}$ denotes $\gamma_{q} R_{a}^{ \pm,(M)}$. Replacing $\dot{a}$ by the coefficients of its Fourier series $\dot{a}_{p}$ and denoting $\dot{\boldsymbol{a}}$ the toeplitz matrix whose elements $\dot{a}_{m p}$ are the $\dot{a}_{m-p}$, it is easy to see that the matrix form of relation 8.24 is:

$$
\left[\begin{array}{cc}
k^{2} \boldsymbol{I}-\boldsymbol{\alpha}^{2} & 0  \tag{8.25}\\
0 & \boldsymbol{I}
\end{array}\right]\left[\begin{array}{c}
\boldsymbol{R}_{a q}^{ \pm} \\
\dot{R}_{a q}^{ \pm}
\end{array}\right]=\gamma_{q}\left[\begin{array}{cc}
-\dot{\boldsymbol{a}} \boldsymbol{\alpha}-\boldsymbol{\alpha} \dot{\boldsymbol{a}} & \boldsymbol{I}+\dot{\boldsymbol{a}} \dot{\boldsymbol{a}} \\
\boldsymbol{I} & \mathbf{0}
\end{array}\right]\left[\begin{array}{c}
\boldsymbol{R}_{a q}^{ \pm} \\
\dot{R}_{a q}^{ \pm}
\end{array}\right]
$$

where $\boldsymbol{R}_{a q}^{ \pm}$and $\dot{R}_{a q}^{ \pm}$are column vectors formed by the $2 M+1$ Fourier coefficients of $R_{a q}^{ \pm(M)}$ $\dot{R}_{a q}^{ \pm(M)}$ respectively. (8.25) shows that $\gamma_{q}$ and $\left[\begin{array}{c}\boldsymbol{R}_{a q}^{ \pm} \\ \dot{R}_{a q}^{ \pm}\end{array}\right]$are an eigenvalue and an eigenvector of the generalized matrix eigenequation $\boldsymbol{A} \psi=\rho \boldsymbol{B} \psi$. Since $R_{a, q}^{ \pm}$is an exact eigenvector of ( 8.14 ) its truncated part can only approximate the sotution of (8.14) and consequently, mathematically speaking $\gamma_{q}$ cannot be an eigenvalue of (8.25). It follows that our claim was false. Nevertheless, elementary pseudo-periodic solutions to (8.14) do exist and we will derive them in the next paragraph.

### 8.2.3 Fourier expansion of elementary waves in the translation coordinate system

In this paragraph, we derive the generalized eigenvalue eigenvector matrix equation starting from (8.14) the only assumption being the pseudo periodicity of the field and we discuss the obtained solutions. First, the propagation equation is rewritten as a pair of first-order equations:

$$
\left[\begin{array}{cc}
k^{2}+\partial_{1}^{2} & 0  \tag{8.26}\\
0 & 1
\end{array}\right]\left[\begin{array}{c}
\mathscr{F} \\
\partial_{3} \mathscr{F}
\end{array}\right]=-\partial_{3}\left[\begin{array}{cc}
\dot{a} \partial_{1}+\partial_{1} \dot{a} & 1+\dot{a} \dot{a} \\
-1 & 0
\end{array}\right]\left[\begin{array}{c}
\mathscr{F} \\
\partial_{3} \mathscr{F}
\end{array}\right]
$$

The coefficients of this equation do not depend on $x^{3}$ which allows to write the $x^{3}$ dependence as $\exp \left(i \rho x^{3}\right)$. The parameter $\rho$ depends on the boundary conditions that $F\left(x^{1}, x^{3}\right)$ has to satisfy along $x^{1}$ direction. For gratings, periodic with period $d^{1}$ along $x^{1}, \mathscr{F}\left(x^{1}+d^{1}, x^{3}\right)=$ $\exp \left(i \alpha_{0} d^{1}\right) \mathscr{F}\left(x^{1}, x^{3}\right)$ where $\alpha_{0}$ is some real parameter. $\partial_{3} \mathscr{F}$ verifies of course the same property. The above requirements on the solution are all fulfilled by expanding function $\mathscr{F}$ and $\partial_{3} \mathscr{F}$ under the form:

$$
\begin{align*}
& \mathscr{F}\left(x^{1}, x^{3}\right)=\exp \left(i \rho x^{3}\right) F_{a}\left(x^{1}\right)=\exp \left(i \rho x^{3}\right) \sum_{m=-M}^{m=M} F_{a, m} \exp \left(i \alpha_{m} x^{1}\right)  \tag{8.27}\\
& \partial_{3} \mathscr{F}\left(x^{1}, x^{3}\right)=\exp \left(i \rho x^{3}\right) \dot{F}_{a}\left(x^{1}\right)=\exp \left(i \rho x^{3}\right) \sum_{m=-M}^{m=M} \dot{F}_{a, m} \exp \left(i \alpha_{m} x^{1}\right) \tag{8.28}
\end{align*}
$$

Introducing the above expansions into (8.26) and projecting the latter onto $\exp \left(\frac{i 2 \pi n x^{1}}{d^{1}}\right)$ basis, we get the sought algebraic matrix eigenvalue equation from which eigenvalues $\rho_{q}$ and eigenvectors $\boldsymbol{F}_{a, q}$ are readily obtained thanks to standard computer libraries:

$$
\left[\begin{array}{cc}
k^{2} \boldsymbol{I}-\boldsymbol{\alpha}^{2} & 0  \tag{8.29}\\
0 & \boldsymbol{I}
\end{array}\right]\left[\begin{array}{c}
\boldsymbol{F}_{a, q} \\
\dot{\boldsymbol{F}}_{a, q}
\end{array}\right]=\rho_{a, q}\left[\begin{array}{cc}
-\dot{\boldsymbol{a}} \boldsymbol{\alpha}-\boldsymbol{\alpha} \dot{\boldsymbol{a}} & \boldsymbol{I}+\dot{\boldsymbol{a}} \dot{\boldsymbol{a}} \\
\boldsymbol{I} & \mathbf{0}
\end{array}\right]\left[\begin{array}{c}
\boldsymbol{F}_{a, q} \\
\dot{\boldsymbol{F}}_{a, q}
\end{array}\right]
$$

As in the Cartesian coordinate system, it is observed numerically that there are two sets of modes, the number of which are equal: those propagating or decaying in the positive $x^{3}$ direction and those propagating or decaying in the opposite direction. Furthermore, it has been shown numerically and analytically [8], that, as the truncation number increases, the computed real eigenvalues converge to the real Rayleigh eigenvalues $\pm \gamma_{q}$.

$$
\begin{equation*}
\lim _{M \rightarrow \infty} \pm \rho_{a, q}^{M}= \pm \gamma_{q}^{R} \tag{8.30}
\end{equation*}
$$

In the above relation, we have added an extra subscript $M$ to indicate the truncation dependence. Indeed, the truncation order $M$ has to be chosen large enough so that the computed real eigenvectors coincide with a great accuracy with their Rayleigh counterpart. In that case, provided that the eigenvalues are not degenerated, up to a multiplicative constant coefficient, the associated computed eigenvectors tend to the corresponding plane waves expressed in terms of the new variables $\left(x^{1}, x^{2}, x^{3}\right)$.

$$
\begin{equation*}
\lim _{M \rightarrow \infty} F_{a, q}^{ \pm(M)}=R_{a, q}^{ \pm} \tag{8.31}
\end{equation*}
$$

Thus in the translation coordinate system defined by $x^{1}=x, x^{2}=y, x^{3}=z-a(x)$ as in the Cartesian coordinate system $O x y z$, linear combinations of elementary solutions to the Helmholtz equation allow us to express electromagnetic field while giving it a physical meaning in terms of forward and backward waves.We write numerically the solution to the Helmholtz equation as:

$$
\begin{align*}
& \mathscr{F}_{a}^{+}\left(x^{1}, x^{3}\right)=\sum_{q \in U^{+}} A_{q}^{+} \exp \left(i \rho_{a, q}^{+} x^{3}\right) R_{a, q}^{+(M)}\left(x^{1}\right)+\sum_{q \in V^{+}} A_{q}^{+} \exp \left(i \rho_{a, q}^{+}\left(x^{3}\right) F_{a, q}^{+}\left(x^{1}\right)\right.  \tag{8.32}\\
& \mathscr{F}_{a}^{-}\left(x^{1}, x^{3}\right)=\sum_{q \in U^{+}} A_{q}^{-} \exp \left(i \rho_{a, q}^{-}\left(x^{3}\right) R_{a, q}^{-(M)}\left(x^{1}\right)+\sum_{q \in V^{+}} A_{q}^{-} \exp \left(i \rho_{a, q}^{-}\left(x^{3}\right) F_{a, q}^{-}\left(x^{1}\right)\right.\right. \tag{8.33}
\end{align*}
$$

with:

$$
\begin{equation*}
F_{a, q}^{ \pm}\left(x^{1}\right)=\sum_{m=-M}^{m=+M} F_{a, m q}^{ \pm} \exp \left(i \alpha_{m} x^{1}\right) \tag{8.34}
\end{equation*}
$$

$U^{ \pm}, V^{ \pm}$denote the sets of indices for the propagating and decaying orders in the positive and negative direction respectively.

$$
\begin{gather*}
U^{+}=\left\{q / \mathfrak{R}\left(\rho_{a, q}\right)>0 \text { and } \mathfrak{I}\left(\rho_{a, q}\right)=0\right\} \quad U^{-}=\left\{q / \mathfrak{R}\left(\rho_{a, q}\right)<0 \text { and } \mathfrak{I}\left(\rho_{a, q}\right)=0\right\}  \tag{8.35}\\
V^{+}=\left\{q / \mathfrak{I}\left(\rho_{a, q}\right)>0\right\} \tag{8.36}
\end{gather*}
$$

### 8.3 Application to a grating problem

Let's come back to the one-dimensional grating problem. Consider the electromagnetic problem in which two homogeneous non magnetic media are separated by a cylindrical periodic surface with period $d^{1}$ which is invariant along the $y$ axis in the Cartesian coordinate system $O x y z$. Such a surface, described by equation $z=a(x)$ is illuminated from above by a unit amplitude linear polarized monochromatic plane wave with vacuum wavelength $\lambda_{0}$, angular frequency $\omega$ and vacuum wave number $k_{0}=2 \pi / \lambda_{0}$. The wave vector is inclined at $\theta$ to the $O z$ axis. Medium (1)


Figure 8.2: Geometry of the diffraction problem. Sketch of the coefficients for scattering matrix
and medium (2) have relative permittivity $\varepsilon_{1}$ and $\varepsilon_{2}$ respectively. Time dependence is expressed by the factor $\exp (-i \omega t)$ Such a problem is reduced to the study of the two fundamental cases of polarisation and the unknown function $\mathscr{F}(x, z)$ is the $y$ component of the electric or the magnetic field for TE and TM polarization respectively. We solved half the problem since we already determined the general solution to the scalar Helmholtz equation as a linear combination of elementary waves the coefficients of which remain to calculate. The situation is very common in electromagnetic theory: the fields on both side of the grating are expanded in terms of the modes in the respective regions with unknown coefficients. A method of solution known as mode-matching method was developed in the context of guided waves in the micro-wave range. The grating may be considered as a generalized multi-port whose inputs are excited by waves that propagate or decay towards it giving rise to a response at the outputs that consists of the waves that propagate or decay away from it [25],[24]. The mode coupling is caused by the modulation of the interface and by the different constitutive parameters in either side of it. The so-called scattering matrix $S_{a}$ defined as

$$
\left[\begin{array}{l}
\boldsymbol{A}^{(1)+}  \tag{8.37}\\
\boldsymbol{A}^{(2)-}
\end{array}\right]=\boldsymbol{S}_{a}\left[\begin{array}{l}
\boldsymbol{A}^{(1)-} \\
\boldsymbol{A}^{(2)+}
\end{array}\right]
$$

provides a linear relation between the output and input coefficients. In a grating problem, the vector formed by the amplitudes of the incoming waves has only one non null component: that corresponding to the incident wave which was assumed enforced to one. The subscript $a$ indicates that the $S$ matrix depends on the profile function $a(x)$. We call $S_{a}$ matrix an interface scattering matrix. The $\boldsymbol{S}_{a}$ matrix is derived from boundary conditions at the surface $x^{3}=x_{0}^{3}$. The change of variable makes it easy to write them. We have solved the scalar Helmholtz equation, the scalar field being a field component tangent to the surface; indeed $F$ coincides with $H_{y}$ and $E_{y}$ in TM polarisation and TE polarisation respectively. For simplicity, let us consider TE polarisation where the non null components of the electromagnetic field are $E_{y}, H_{x}, H_{z}$. Boundary conditions require matching the tangential components of the magnetic and electric field. We have already derived one of them, $E_{y}$, we have to derive the tangential component $H_{t}$ of the electric field given by :

$$
\begin{equation*}
H_{t}=\boldsymbol{H} . \boldsymbol{t} \tag{8.38}
\end{equation*}
$$

where $t$ is the unit vector at point $P$ which is tangential to the grating profile function. It is defined in terms of the $\boldsymbol{e}_{x}$ and $\boldsymbol{e}_{z}$ Cartesian unit vectors by:

$$
\begin{equation*}
t=\frac{1}{\sqrt{1+\dot{a}^{2}}}\left(e_{x}+\dot{a} e_{z}\right) \tag{8.39}
\end{equation*}
$$

The square root in the denominator represents a normalization factor that can be omitted since at a given point, it is identical on both sides of the boundary surface. let us introduce $G$ such
that:

$$
\begin{equation*}
\mathscr{G}=i Z H \cdot t \tag{8.40}
\end{equation*}
$$

Where $Z=\sqrt{\mu_{0} / \varepsilon}$ is the wave impedance. From Maxwells equation we have $i \omega \mu_{0} H_{x}=-\partial_{z} E_{y}$ and $i \omega \mu_{0} H_{z}=\partial_{x} E_{y}$ thus :

$$
\begin{equation*}
\mathscr{G}(x, z)=-\frac{1}{k}\left(\partial_{z} \mathscr{F}(x, z)-\dot{a} \partial_{x} \mathscr{F}(x, z)\right) \tag{8.41}
\end{equation*}
$$

substituting $\partial_{3}$ for $\partial_{z}$ and $\partial_{1}-\dot{a} \partial_{3}$ for $\partial_{x}$ we get:

$$
\begin{equation*}
\mathscr{G}\left(x^{1}, x^{3}\right)=-\frac{1}{k}\left((1+\dot{a} \dot{a}) \partial_{3}-\dot{a} \partial_{1}\right) \mathscr{F}\left(x^{1}, x^{3}\right) \tag{8.42}
\end{equation*}
$$

Similarly to $\mathscr{F}, \mathscr{G}$ depends on $x^{3}$ as $\exp \left(i \rho x^{3}\right)$ and we may write:

$$
\begin{equation*}
\mathscr{G}\left(x^{1}, x^{3}\right)=\exp \left(i \rho x^{3}\right) G\left(x^{1}\right) \tag{8.43}
\end{equation*}
$$

We are now familiar with the operational rules that allow to associate in Fourier space a matrix with an operator. We have:

$$
\begin{equation*}
1+\dot{a} \dot{a} \rightarrow \boldsymbol{I}+\dot{\boldsymbol{a}} \dot{\boldsymbol{a}}, \quad \dot{a} \partial_{1} \rightarrow i \dot{\boldsymbol{a}} \boldsymbol{\alpha} \tag{8.44}
\end{equation*}
$$

From which we deduce:

$$
\begin{equation*}
i k \boldsymbol{G}_{a}^{ \pm}=(I+\dot{\boldsymbol{a}} \dot{\boldsymbol{a}}) \boldsymbol{F}_{a}^{ \pm} \boldsymbol{\rho}_{a}-\dot{\boldsymbol{a}} \boldsymbol{\alpha} \boldsymbol{F}_{a}^{ \pm} \tag{8.45}
\end{equation*}
$$

where $\rho$ is a diagonal matrix whose elements are the eigenvalues $\rho_{a, q}$. Writing the continuity of $\mathscr{F}^{(1)}$ and $\mathscr{F}^{(2)}$ and $\mathscr{G}^{(1)} / Z^{(1)}$ and $\mathscr{G}^{(2)} / Z^{(2)}$ at $x^{3}=x_{0}^{3}$ is straightforward and leads to the following expression of the $\boldsymbol{S}_{a}$ matrix:

$$
\boldsymbol{S}_{a}=\left[\begin{array}{ll}
\boldsymbol{F}_{a}^{(1)+} & -\boldsymbol{F}_{a}^{(2)-}  \tag{8.46}\\
\boldsymbol{G}_{a}^{(1)+} & -\boldsymbol{G}_{a}^{(2)-}
\end{array}\right]^{-1}\left[\begin{array}{ll}
-\boldsymbol{F}_{a}^{(1)-} & \boldsymbol{F}_{a}^{(2)+} \\
-\boldsymbol{G}_{a}^{(1)-} & \boldsymbol{G}_{a}^{(2)+}
\end{array}\right]
$$

The knowledge of $\boldsymbol{S}_{a}$ matrix allows to calculate the constant coefficients of outgoing waves. Since the spectrum of the solutions of the transformed Helmholtz equation include the generalized Rayleigh eigenvectors associated to real Rayleigh eigenvalues the efficiencies may be calculated in the very same way as in the Cartesian coordinate system.

$$
\begin{equation*}
R_{q}=\left|A_{q}^{(1)+}\right|^{2} \frac{\gamma_{q}^{(1)}}{\gamma_{0}^{(1)}} \quad T_{p}=\left|A_{p}^{(2)-}\right|^{2} \frac{\gamma_{p}^{(2)}}{\gamma_{0}^{(1)}} \tag{8.47}
\end{equation*}
$$

with:

$$
\begin{equation*}
\gamma_{q}^{(1)}=\sqrt{k_{0}^{2} \varepsilon_{1}-\left(k_{0} \sqrt{\varepsilon_{1}} \sin \theta+q \frac{2 \pi}{d^{1}}\right)^{2}} \gamma_{p}^{(2)}=\sqrt{k_{0}^{2} \varepsilon_{2}-\left(k_{0} \sqrt{\varepsilon}_{1} \sin \theta+p \frac{2 \pi}{d^{1}}\right)^{2}} \tag{8.48}
\end{equation*}
$$

The values of integers $p$ and $q$ are such that $\gamma_{q}^{(1)}$ and $\gamma_{p}^{(2)}$ are real.

### 8.3.1 Implementation of $C$-Method

The main interest to consider C-Method through a simple change of variable is to make us understand its numerical link with Rayleigh expansions and to calculate efficiencies as in the Cartesian Coordinate system. Within the framework of translation coordinate systems, starting from Maxwell's equations written under the covariant form, we have shown that all tangential components of the field at surface $S$ could be generated from the longitudinal covariant components along the axis of invariance. Moreover,these components are solutions of the scalar Helmholtz equation. Therefore, one clearly understands that finding the elementary solutions of the scalar Helmholtz equation is the kernel of C-method. To summarize we may enunciate the different steps for solving a grating problem with C-method:

- Define a translation coordinate system.
- Find the elementary waves of the Helmholtz equation. For that purpose use the Galerkin method with $\exp \left(i \alpha_{n} x^{1}\right)$ as expansion and test functions. Substitute the generalized Rayleigh eigenvectors for the computed eigenvectors associated to real eigenvalues. Sort the elementary waves into forward and backward waves.
- Write boundary conditions at surface $\Sigma$ and calculate efficiencies as in the Cartesian coordinate system.


### 8.4 Various formulations of C-method

So far, the Helmholtz equation in the translation coordinate system was derived by using the chain rule for derivatives in the Helmholtz equation written in the translation coordinate system. In this section,we start from the covariant Maxwell's equations and we show that they lead to several operators one of them being the propagation equation. In a homogeneous isotropic medium with permittivity $\varepsilon$ and permeability $\mu$, with a time dependence $\exp (-i \omega t)$,the symmetrized Maxwell equations write:

$$
\begin{align*}
& \xi^{\alpha \beta \gamma} \partial_{\beta} \mathscr{F}_{\gamma}=k \sqrt{g} g^{\alpha \beta} \mathscr{G}_{\alpha}  \tag{8.49}\\
& \xi^{\alpha \beta \gamma} \partial_{\beta} \mathscr{G}_{\gamma}=k \sqrt{g} g^{\alpha \beta} \mathscr{F}_{\alpha}
\end{align*}
$$

where $k=\omega \sqrt{\mu \varepsilon}$, the $\mathscr{F}_{\gamma}$ and the $\mathscr{G}_{\beta}$ are the complex amplitudes of the electric field and of a renormalized magnetic field respectively. We restrict our analysis to 1D problems in which both the geometry and the solution are independent of $y$. Practically this means that $\partial_{2}$ is null as well as $g^{12}, g^{21}, g^{32}$ and $夕^{23}$. 1 follows that (8.49) decouple into two fully identical systems where the non null components are $\mathscr{F}_{2}, \mathscr{G}_{1}, \mathscr{G}_{3}$, and $\mathscr{G}_{2}, \mathscr{F}_{1}, \mathscr{F}_{3}$ respectively.The first set of three components corresponds to $T E$ polarisation, the second one to $T M$ polarisation. Both polarisations obey the same first order differential equations system written hereafter for $T E$ polarisation:

$$
\begin{array}{r}
-\partial_{3} \mathscr{F}_{2}=k\left(\sqrt{g} g^{11} \mathscr{G}_{1}+\sqrt{g} g^{13} \mathscr{G}_{3}\right) \\
\partial_{1} \mathscr{F}_{2}=k\left(\sqrt{g} g^{31} \mathscr{G}_{1}+\sqrt{g} g^{33} \mathscr{G}_{3}\right) \\
\partial_{3} \mathscr{G}_{1}-\partial_{1} \mathscr{G}_{3}=k \sqrt{g} g^{22} F_{2} \tag{8.50c}
\end{array}
$$

For $T M$ polarisation, it is enough to permute $\mathscr{F}$ and $\mathscr{G}$. Among the three components of each system, two play a particular role. Let us assume that $x^{3}=x_{0}^{3}$ separates two isotropic homogeneous media. Then, in $T E$ polarisation $\mathscr{F}_{2}$ and $\mathscr{G}_{1} / \sqrt{\frac{\mu}{\varepsilon}}$ have to be continuous at surface $x^{3}=x_{0}^{3}$. The same conclusions holds for $\mathscr{G}_{2} / \sqrt{\frac{\mu}{\varepsilon}}$ and $\mathscr{F}_{1}$ for $T M$ polarisation. So, we have to solve (8.50) for the components labelled by two and by one. C-method is a Fourier based method which means that constitutive relations have to be written in Fourier space. In other words a matrix is to be associated to each element $\sqrt{g} g^{\alpha \beta}$ of the constitutive tensors. The way for doing so should follow the so-called "Fourier factorization" rules derived by Li [22], [23]. let us denote by $\left(\sqrt{g} g^{\alpha \beta}\right)$ the matrix associated to coefficient $\sqrt{g} g^{\alpha \beta}$. According to Li's rules, the $\left(\sqrt{g} g^{\alpha \beta}\right)$ write:

$$
\begin{align*}
& \left(\sqrt{g} g^{11}\right)=\left[\frac{1}{\sqrt{g} g^{11}}\right]^{-1} \\
& \left(\sqrt{g} g^{13}\right)=\left[\frac{1}{\sqrt{g} g^{11}}\right]^{-1}\left[\frac{g^{13}}{g^{11}}\right] \\
& \left(\sqrt{g} g^{31}\right)=\left[\frac{g^{31}}{g^{11}}\right]\left[\frac{1}{\sqrt{g} g^{11}}\right]^{-1}  \tag{8.51}\\
& \left(\sqrt{g} g^{33}\right)=\left[\frac{1}{\sqrt{g} g^{11}}\right]+\left[\frac{g^{31}}{g^{11}}\right]\left[\frac{1}{\sqrt{g} g^{11}}\right]^{-1}\left[\frac{g^{13}}{g^{11}}\right] \\
& \left(\sqrt{g} g^{22}\right)=\left[\sqrt{g} g^{22}\right]
\end{align*}
$$

The notation $[f]$ designates the toeplitz matrix whose elements $f_{m p}$ are the $f_{m-p}$ elements of the Fourier series of function $f\left(x^{1}\right)$. For the translation coordinate $\left(x^{1}, x^{2}, x^{2}\right)$ such that $x=x^{1}$, $y=x^{2}, z=x^{3}+a\left(x^{1}\right)$, we have:

$$
\begin{align*}
\left(\sqrt{g} g^{11}\right) & \rightarrow \boldsymbol{I}  \tag{8.52}\\
\left(\sqrt{g} g^{13}\right) & \rightarrow-\dot{\boldsymbol{a}} \\
\left(\sqrt{g} g^{31}\right) & \rightarrow-\dot{\boldsymbol{a}} \\
\left(\sqrt{g} g^{33}\right) & \rightarrow[\boldsymbol{I}+\dot{\boldsymbol{a}} \dot{\boldsymbol{a}}] \\
\left(\sqrt{g} g^{22}\right) & \rightarrow \boldsymbol{I}
\end{align*}
$$

In Fourier space, the derivative operator $\partial_{1}$ is associated to the diagonal matrix $i \boldsymbol{\alpha}$ the elements of which are the $i \alpha_{m}$ such that:

$$
\begin{equation*}
\alpha_{m}=\alpha_{0}+m \frac{2 \pi}{d^{1}} \tag{8.53}
\end{equation*}
$$

Setting

$$
\begin{align*}
& \mathscr{F}_{2}\left(x^{1}, x^{3}\right)=\sum_{m=-M}^{m=+M} F_{2 m}\left(x^{3}\right) \exp \left(i \alpha_{m} x^{1}\right)  \tag{8.54a}\\
& \mathscr{G}_{1}\left(x^{1}, x^{3}\right)=\sum_{m=-M}^{m=+M} G_{1 m}\left(x^{3}\right) \exp \left(i \alpha_{m} x^{1}\right)  \tag{8.54b}\\
& \mathscr{G}_{3}\left(x^{1}, x^{3}\right)=\sum_{m=-M}^{m=+M} G_{3 m}\left(x^{3}\right) \exp \left(i \alpha_{m} x^{1}\right) \tag{8.54c}
\end{align*}
$$

We are now able to write 8.50 ) in Fourier space:

$$
\begin{align*}
-\partial_{3} \boldsymbol{F}_{2}= & k\left(\left(\sqrt{g} g^{11}\right) \boldsymbol{G}_{1}+\left(\sqrt{g} g^{13}\right) \boldsymbol{G}_{3}\right)  \tag{8.55a}\\
i \boldsymbol{\alpha} \boldsymbol{F}_{2}= & k\left(\left(\sqrt{g} g^{31}\right) \boldsymbol{G}_{1}+\left(\sqrt{g} g^{33}\right) \boldsymbol{G}_{3}\right)  \tag{8.55b}\\
& \partial_{3} \boldsymbol{G}_{1}-i \boldsymbol{\alpha} \boldsymbol{G}_{3}=k\left(\sqrt{g} g^{22}\right) \boldsymbol{F}_{2} \tag{8.55c}
\end{align*}
$$

where $\boldsymbol{F}_{2}, \boldsymbol{G}_{1} \boldsymbol{G}_{3}$ are column vectors of size $2 M+1$ whose components are the $F_{2 m}\left(x^{3}\right)$, $G_{1 m}\left(x^{3}\right), G_{3 m}\left(x^{3}\right)$ respectively:

$$
\begin{array}{r}
\boldsymbol{F}_{2}\left(x^{3}\right)=\left[F_{2,-M}\left(x^{3}\right), F_{2,-M+1}\left(x^{3}\right), \cdots, F_{2,0}\left(x^{3}\right), \cdots, F_{2, M-1}\left(x^{3}\right), F_{2, M}\left(x^{3}\right)\right]^{T} \\
\boldsymbol{G}_{1}\left(x^{3}\right)=\left[G_{1,-M}\left(x^{3}\right), G_{1,-M+1}\left(x^{3}\right), \cdots, G_{1,0}\left(x^{3}\right), \cdots, G_{1, M-1}\left(x^{3}\right), G_{1, M}\left(x^{3}\right)\right]^{T} \\
\boldsymbol{G}_{3}\left(x^{3}\right)=\left[G_{3,-M}\left(x^{3}\right), G_{3,-M+1}\left(x^{3}\right), \cdots G_{3,0}\left(x^{3}\right), \cdots, G_{3, M-1}\left(x^{3}\right), G_{3, M}\left(x^{3}\right)\right]^{T} \tag{8.56c}
\end{array}
$$

where the exponent $T$ is for the transposition.

### 8.4.1 Propagation equation in curvilinear coordinates

From Eqs 8.50 a and (8.50b), $\boldsymbol{G}_{1}$ and $\boldsymbol{G}_{3}$ may be expressed in terms of $\boldsymbol{F}_{2}$

$$
\begin{gather*}
k \boldsymbol{G}_{1}=-\left(\sqrt{g} g^{33}\right) \partial_{3} \boldsymbol{F}_{2}+\left(\sqrt{g} g^{13}\right) i \boldsymbol{\alpha} \boldsymbol{F}_{2}  \tag{8.57}\\
k \boldsymbol{G}_{3}\left(x^{3}\right)=\left(\sqrt{g} g^{31}\right) \partial_{3} \boldsymbol{F}_{2}\left(x^{3}\right)+\left(\sqrt{g} g^{11}\right) i \boldsymbol{\alpha} \boldsymbol{F}_{2}\left(x^{3}\right) \tag{8.58}
\end{gather*}
$$

8.50c), in which we substitute $G_{1}$ and $G_{3}$ with expressions (8.57) and 8.58), gives the propagation equation:

$$
\begin{equation*}
\left(-\boldsymbol{\alpha}\left(\sqrt{g} g^{11}\right) \boldsymbol{\alpha}+\partial_{3}\left(\sqrt{g} g^{33}\right) \partial_{3}+i \boldsymbol{\alpha}\left(\sqrt{g} g^{13}\right) \partial_{3}+\partial_{3}\left(\sqrt{g} g^{31}\right) i \boldsymbol{\alpha}+k^{2}\left(\sqrt{g} g^{22}\right)\right) \boldsymbol{F}_{2}\left(x^{3}\right)=0 \tag{8.59}
\end{equation*}
$$

which is rewritten as a pair of first-order differential equation as:

$$
-i \partial_{3} \boldsymbol{A}\left[\begin{array}{c}
\boldsymbol{F}_{2}\left(x^{3}\right)  \tag{8.60}\\
-i \partial_{3} \boldsymbol{F}_{2}\left(x^{3}\right)
\end{array}\right]=\boldsymbol{B}\left[\begin{array}{c}
\boldsymbol{F}_{2}\left(x^{3}\right) \\
-i \partial_{3} \boldsymbol{F}_{2}\left(x^{3}\right)
\end{array}\right]
$$

with:

$$
\begin{gather*}
\boldsymbol{A}=\left[\begin{array}{cc}
\alpha\left(\sqrt{g} g^{13}\right)+\left(\sqrt{g} g^{13}\right) \alpha & \left(\sqrt{g} g^{33}\right) \\
\boldsymbol{I} & 0
\end{array}\right]  \tag{8.61}\\
\boldsymbol{B}=\left[\begin{array}{cc}
-\alpha\left(\sqrt{g} g^{11}\right) \alpha+k^{2}\left(\sqrt{g} g^{22}\right) & 0 \\
0 & I
\end{array}\right] \tag{8.62}
\end{gather*}
$$

Since the coefficients of matrices $\boldsymbol{A}$ and $\boldsymbol{B}$ are independent of variable $x^{3}$, we may seek vectors $\boldsymbol{F}_{2}\left(x^{3}\right), \boldsymbol{G}_{1}\left(x^{3}\right), \boldsymbol{G}_{3}\left(x^{3}\right)$ under the form:

$$
\begin{align*}
\boldsymbol{F}_{2}\left(x^{3}\right) & =\boldsymbol{F}_{2} \exp \left(i \rho x^{3}\right)  \tag{8.63a}\\
-i \partial_{3} \boldsymbol{F}_{2}\left(x^{3}\right) & =\dot{\boldsymbol{F}}_{2} \exp \left(i \rho x^{3}\right)  \tag{8.63b}\\
\boldsymbol{G}_{1}\left(x^{3}\right) & =\boldsymbol{G}_{1} \exp \left(i \rho x^{3}\right)  \tag{8.63c}\\
\boldsymbol{G}_{3}\left(x^{3}\right) & =\boldsymbol{G}_{3} \exp \left(i \rho x^{3}\right) \tag{8.63d}
\end{align*}
$$

This last step transforms (8.60) into a generalized eigenvalue eigenvector matrix equation:

$$
\boldsymbol{A} \rho\left[\begin{array}{l}
\boldsymbol{F}_{2}  \tag{8.64}\\
\dot{\boldsymbol{F}}_{2}
\end{array}\right]=\boldsymbol{B}\left[\begin{array}{l}
\boldsymbol{F}_{2} \\
\dot{\boldsymbol{F}}_{2}
\end{array}\right]
$$

It is then easy to check that (8.64) is the same as (8.29). After $\boldsymbol{F}_{2}$ is determined, it remains to deduce $G_{1}$ from 8.57).

### 8.4.2 "Classical" C-method operator

We call "classical" operator the operator derived by Chandezon in his early work. From 8.50b and taking into account $(8.63)$ we find an expression for $G_{3}$ as follows:

$$
\begin{equation*}
\boldsymbol{G}_{3}=\frac{1}{k}\left(\sqrt{g} g^{33}\right)^{-1} i \alpha \boldsymbol{F}_{2}-\left(\sqrt{g} g^{33}\right)^{-1}\left(\sqrt{g} g^{31}\right) \boldsymbol{G}_{1} \tag{8.65}
\end{equation*}
$$

Substituting $G_{3}$ in 8.50a and 8.50 c with the above expression yields:

$$
\begin{gather*}
{\left[\begin{array}{cc}
-\left(\sqrt{g} g^{13}\right)\left(\sqrt{g} g^{33}\right)^{-1} \boldsymbol{\alpha} & i k\left(\left(\sqrt{g} g^{11}\right)-\left(\sqrt{g} g^{13}\right)\left(\sqrt{g} g^{33}\right)^{-1}\left(\sqrt{g} g^{31}\right)\right) \\
-i k\left(\left(\sqrt{g} g^{22}\right)-\frac{1}{k^{2}} \boldsymbol{\alpha}\left(\sqrt{g} g^{33}\right)^{-1} \boldsymbol{\alpha}\right) & -\boldsymbol{\alpha}\left(\sqrt{g} g^{33}\right)^{-1}\left(\sqrt{g} g^{31}\right)
\end{array}\right]\left[\begin{array}{l}
\boldsymbol{F}_{2} \\
\boldsymbol{G}_{1}
\end{array}\right]} \\
=\rho\left[\begin{array}{c}
\boldsymbol{F}_{2} \\
\boldsymbol{G}_{1}
\end{array}\right] \tag{8.66}
\end{gather*}
$$

### 8.5 Multilayer grating

The extension of C-method to multilayer gratings is straightforward provided the interfaces which separate the layers share the same periodicity. It is just a generalization of the theory of planar stratified media. As a canonical case, let us consider a layer made of isotropic homogeneous media limited on the top by surface $z=a_{j}\left(x^{1}\right)$ and on the bottom by surface $z=a_{j+1}\left(x^{1}\right)=a_{j}\left(x^{1}\right)-t_{j}\left(x^{1}\right)$. When $t_{j}\left(x^{1}\right)$ is constant the two surfaces are parallel to each other. In homogeneous media the field is a superposition of forward and backward waves. The only places where coupling occurs are the interfaces. Thus, we have to describe two different phenomena: on the one hand scattering at the interfaces and on the other hand propagation or attenuation in the layer. To summarize we assimilate an interface to a 4 N -port local network ( $N=4 M+1, M$ being the truncation number) and a layer to a multi-channel pipe connecting the $2 N$-ports of its input network and output network [24]. We have already defined interface scattering matrices which are local matrices in the sense they depend on the profile. In other words, in the context of C-method they depend on the coordinate system. Thus, for a layer bounded by two non parallel surfaces, we have to solve two eigenvalue problems for each surface which allows to calculate interface matrices $\boldsymbol{S}_{a_{j}}$ and $\boldsymbol{S}_{a_{j+1}}$. It remains to define and to calculate layer scattering matrices. Although two cases have to be considered according to whether the layer separates two identical surfaces or not, the line of reasoning is the same. As already mentioned, we have two coordinate systems such that $z=x_{j}^{3}+a_{j}\left(x^{1}\right)$ and $z=x_{j+1}^{3}+a_{j+1}\left(x^{1}\right)$. They are linked by the following relation:

$$
\begin{equation*}
x_{j}^{3}=x_{j+1}^{3}+a_{j+1}\left(x^{1}\right)-a_{j}\left(x^{1}\right)=x_{j+1}^{3}-t_{j}\left(x^{1}\right) \tag{8.67}
\end{equation*}
$$



Figure 8.3: Schematic representation of diffraction at two surfaces separated by a layer

In medium $j$, located in between surfaces $z=a_{j}\left(x^{1}\right)$ and $z=a_{j+1}\left(x^{1}\right)$, we may express the linear combination of forward and backard waves with coordinate $x_{j}^{3}=0$ as local origin (that is $z=a_{j}\left(x^{1}\right)$ and write:

$$
\begin{equation*}
\left.\mathscr{F}_{a_{j}}^{(j)}\left(x_{j}^{3}, x^{1}\right)=\sum_{q} A_{j, q}^{(j)+} \exp \left(i \rho_{a_{j}, q}^{(j)+} x_{j}^{3}\right) F_{a_{j}, q}^{(j)+}\left(x^{1}\right)+\sum_{q} A_{j, q}^{(j)-} \exp \left(i \rho_{a_{j}, q}^{(j)-} x_{j}^{3}\right)\right) F_{a_{j}, q}^{(j)-}\left(x^{1}\right) \tag{8.68}
\end{equation*}
$$

In the same medium $j$ we may also choose $x_{j+1}^{3}=0$ as local origin (that is $z=a_{j+1}\left(x^{1}\right)$ ) which gives:

$$
\begin{equation*}
\left.\mathscr{F}_{a_{j+1}}^{(j)}\left(x_{j+1}^{3}, x^{1}\right)\right)=\sum_{q} A_{j+1, q}^{(j)+} \exp \left(i \rho_{a_{j+1}, q}^{(j)+} x_{j+1}^{3}\right) F_{a_{j+1}, q}^{(j)+}\left(x^{1}\right)+\sum_{q} A_{j+1, q}^{(j)-} \exp \left(i \rho_{a_{j+1}, q}^{(j)-}\left(x_{j+1}^{3}\right) F_{a_{j+1}, q}^{(j)-}\left(x^{1}\right)\right. \tag{8.69}
\end{equation*}
$$

The layer is considered as a $4 N$ - ports which connects input waves $\mathscr{F}_{a_{j}}^{(j)-}$ and $\mathscr{F}_{a_{j+1}}^{(j)+}$ to output waves $\mathscr{F}_{a_{j}}^{(j)+}$ and $\mathscr{F}_{a_{j+1}}^{(j)-}$, hence the definition of the layer $S$ matrix:

$$
\left[\begin{array}{c}
A_{j, q}^{(j)+}  \tag{8.70}\\
A_{j+1, q}^{(j)-}
\end{array}\right]=\boldsymbol{S}_{j, j+1}\left[\begin{array}{c}
A_{j, q}^{(j)-} \\
A_{j+1, q}^{(j)+}
\end{array}\right]
$$

At the input of the layer, that is at $x_{j}^{3}=0$, the outgoing waves correspond to the incoming wave of the output plane:

$$
\begin{equation*}
\mathscr{F}_{a_{j}}^{(j)+}\left(x_{j}^{3}=0\right)=\mathscr{F}_{a_{j+1}}^{(j)+}\left(x_{j+1}^{3}=t_{j}\left(x^{1}\right)\right) \tag{8.71}
\end{equation*}
$$

Similarly, at the output of the layer, that is $x_{j+1}^{3}=0$, the outgoing waves correspond to the incoming waves of the input plane:

$$
\begin{equation*}
\mathscr{F}_{a_{j+1}}^{(j)-}\left(x_{j+1}^{3}=0\right)=\mathscr{F}_{a_{j}}^{(j)-}\left(x_{j}^{3}=-t_{j}\left(x^{1}\right)\right. \tag{8.72}
\end{equation*}
$$

At this stage, we infer that layer $S$ matrix looks like:

$$
\boldsymbol{S}_{j, j+1}=\left[\begin{array}{cc}
0 & \boldsymbol{P}^{(j)+}  \tag{8.73}\\
\boldsymbol{P}^{(j)-} & 0
\end{array}\right]
$$

The sought sub-matrices $\boldsymbol{P}^{(j)+}, \boldsymbol{P}^{(j)-}$ depend on whether the layer faces are parallel or not.

### 8.5.1 Layer with non parallel faces

Consider equation (8.75) and write it in terms of the eigenvectors of both coordinate systems:

$$
\begin{equation*}
\sum_{m} \sum_{q} A_{j, q}^{(j)+} F_{a_{j}, m q}^{(j)+} \exp \left(i \alpha_{m} x^{1}\right)=\sum_{m} \sum_{q} A_{j+1, q}^{(j)+} \exp \left(i \rho_{a_{j}+1, q}^{+} t_{j}\left(x^{1}\right)\right) F_{a_{j+1, q}}^{(j)+} \exp \left(i \alpha_{m} x^{1}\right) \tag{8.74}
\end{equation*}
$$

The left hand side purely consists of a linear combination of eigenvectors expanded onto the $\exp \left(i \alpha_{m} x^{1}\right)$ basis whereas the right hand side consists of a linear combination of eigenvectors each of which being multiplied by a periodic functions of the $x^{1}$ variable. In order to get a matrix relation between the $A_{j, q}^{(j)+}$ and the $A_{j+1, q}^{(j)+}$, we project (8.74) onto $\exp \left(i \alpha_{m} x^{1}\right)$. We get:

$$
\begin{equation*}
\sum_{m} \sum_{q} A_{j, q}^{(j)+} F_{a_{j}, m q}^{(j)+}=\sum_{m} \sum_{q} A_{j+1, q}^{(j)+} \tilde{F}_{a_{j+1, q}}^{(j)+} \exp \left(i \alpha_{m} x^{1}\right) \tag{8.75}
\end{equation*}
$$

with:

$$
\begin{equation*}
\tilde{F}_{a_{j+1} m q}^{(j)+}=\frac{1}{d^{1}} \int_{0}^{d^{1}}\left(\sum_{l} F_{a_{j+1}, l q}^{(j)+} \exp \left(i \alpha_{l} x^{1}\right)\right) \exp \left(i \rho_{a_{j}+1, q^{\prime}}^{(j)+}\left(x^{1}\right)\right) \exp \left(-i \alpha_{m} x^{1}\right) d x^{1} \tag{8.76}
\end{equation*}
$$

Then, the $\boldsymbol{P}^{(j)+}$ matrix is readily obtained as

$$
\begin{equation*}
\boldsymbol{P}^{(j)+}=\left(F_{a_{j}}^{(j)+}\right)^{-1} \tilde{\boldsymbol{F}}_{a_{j+1}}^{(j)+} \tag{8.77}
\end{equation*}
$$

where $\boldsymbol{F}_{a_{j}}^{(j)+}$ (respectively $\tilde{\boldsymbol{F}}_{a_{j+1}}^{(j)+}$ ) is the matrix formed by juxtaposition of vectors $F_{a_{j}, q}^{(j)+}$ (respectively $\tilde{F}_{a_{j+1} q}^{(j)+}$. Similarly we have:

$$
\begin{equation*}
\sum_{m} \sum_{q} A_{j+1, q}^{(j)-} F_{a_{j+1}, m q}^{(j)-} \exp \left(i \alpha_{m} x^{1}\right)=\sum_{m} \sum_{q} A_{j, q}^{(j)-} \exp \left(i \rho_{a_{j, q}}^{(j)-} t_{j}\left(x^{1}\right)\right) F_{a_{j, q}}^{(j)-} \exp \left(i \alpha_{m} x^{1}\right) \tag{8.78}
\end{equation*}
$$

and

$$
\begin{equation*}
\sum_{m} \sum_{q} A_{j+1, q}^{-(j)} F_{a_{j+1}, m q}^{(j)-}=\sum_{m} \sum_{q} A_{j, q}^{(j)-} \tilde{F}_{a_{j, q}}^{(j)-} \tag{8.79}
\end{equation*}
$$

with:

$$
\begin{equation*}
\tilde{F}_{a_{j m q}}^{(j)-}=\frac{1}{d^{1}} \int_{0}^{d^{1}}\left(\sum_{l} F_{a_{j}, l q}^{(j)-} \exp \left(i \alpha_{l} x^{1}\right)\right) \exp \left(-i \rho_{a_{j}+1, q}^{(j)-} t_{j}\left(x^{1}\right)\right) \exp \left(-i \alpha_{m} x^{1}\right) d x^{1} \tag{8.80}
\end{equation*}
$$

from which we derive:

$$
\begin{equation*}
\boldsymbol{P}^{(j)-}=\left(F_{a_{j+1}}^{(j)-}\right)^{-1} \tilde{\boldsymbol{F}}_{a_{j}}^{(j)-} \tag{8.81}
\end{equation*}
$$

### 8.5.2 Layer with parallel faces

In that case, the two coordinate systems are identical and $t_{j}\left(x^{1}\right)$ is a constant. Equations (8.72) and (8.75) reduce to:

$$
\begin{array}{r}
\sum_{m} \sum_{q} A_{j, q}^{(j)+} F_{a_{j}, m q}^{(j)+} \exp \left(i \alpha_{m} x^{1}\right)=\sum_{m} \sum_{q} A_{j+1, q}^{+(j)} \exp \left(i \rho_{a_{j}, q}^{(j)+} t_{j}\right) F_{a_{j}, m q}^{(j)+} \exp \left(i \alpha_{m} x^{1}\right) \\
\sum_{m} \sum_{q} A_{j+1, q}^{(j)-} F_{a_{j}, q}^{(j)+}=\sum_{m} \sum_{q} A_{j, q}^{(j)-} \exp \left(-i \rho_{a_{j}, q}^{(j)-} t_{j}\right) F_{a_{j}, m q}^{(j)-} \exp \left(i \alpha_{m} x^{1}\right) \tag{8.83}
\end{array}
$$

from which we easily deduce:

$$
\begin{gather*}
A_{j, q}^{(j)+}=A_{j+1, q}^{(j)+} \exp \left(i \rho_{a_{j}, q}^{(j)+} t_{j}\right) \text { or } \boldsymbol{P}^{(j)+}=\operatorname{diag}\left(\exp \left(i \rho_{a_{j}, q}^{(j)+} t_{j}\right)\right)  \tag{8.84}\\
A_{j+1, q}^{(j)-}=A_{j, q}^{(j)-} \exp \left(-i \rho_{a_{j}, q}^{(j)-} t_{j}\right) \text { or } \boldsymbol{P}^{(j)-}=\operatorname{diag}\left(\exp \left(-i \rho_{a_{j}, q}^{(j)-} t_{j}\right)\right) \tag{8.85}
\end{gather*}
$$

It should be noted that when $\rho_{a_{j}, q}^{(j)+}$ (respectively $\left.\rho_{a_{j}, q}^{(j)-}\right)$ is complex valued, its imaginary part is negative (respectively positive). Since $t_{j}$ is positive, exponential functions $\exp \left( \pm i \rho_{a_{j}, q}^{(j) \pm} t_{j}\right)$ associated to complex eigenvalues always decay when the layer thickness increases.

### 8.5.3 Combination of $S$ matrices

The final step for analysing reflection and transmission by a layer is to combine the two interfaces S matrix and the layer S matrix. The tool for doing this is the Redheffer star product which gives the composition rules of two cascaded $S$ matrices [26]. Consider two S matrices and partition them into four blocks:

$$
S_{1}=\left[\begin{array}{cc}
S_{1}^{11} & S_{1}^{12}  \tag{8.86}\\
S_{1}^{21} & S_{1}^{22}
\end{array}\right] \boldsymbol{S}_{2}=\left[\begin{array}{cc}
S_{2}^{11} & S_{2}^{12} \\
S_{2}^{21} & S_{2}^{22}
\end{array}\right]
$$

The star product $*$ is defined by:

$$
\begin{array}{r}
S=S_{1} * S_{2} \\
S^{11}=S_{1}^{11}+S_{1}^{12}\left(\boldsymbol{I}-S_{2}^{11} S_{1}^{22}\right)^{-1} S_{2}^{11} \times S_{1}^{21} \\
S^{12}=S_{1}^{12} \times\left(\boldsymbol{I}-S_{2}^{11} S_{1}^{22}\right)^{-1} \times S_{2}^{12} \\
S^{21}=S_{2}^{21} \times\left(\boldsymbol{I}-S_{1}^{22} S_{2}^{11}\right)^{-1} \times S_{1}^{21} \\
S^{22}=S_{2}^{22}+S_{2}^{21} \times\left(I-S_{1}^{22} S_{2}^{11}\right)^{-1} \times S_{1}^{22} \times \boldsymbol{S}_{2}^{12} \tag{8.91}
\end{array}
$$

where $I$ is the identity matrix. The combined $S$ matrix of the top and bottom interfaces and of the layer is given by:

$$
\begin{equation*}
\boldsymbol{S}=\boldsymbol{S}_{a_{j}} * \boldsymbol{S}_{j, j+1} * \boldsymbol{S}_{a_{j+1}}=\left(\boldsymbol{S}_{a_{j}} * \boldsymbol{S}_{j, j+1}\right) * \boldsymbol{S}_{a_{j+1}}=\boldsymbol{S}_{a_{j}} *\left(\boldsymbol{S}_{j, j+1} * \boldsymbol{S}_{a_{j+1}}\right) \tag{8.92}
\end{equation*}
$$

and finally, it turns out that

$$
\begin{array}{r}
\boldsymbol{S}^{11}=\boldsymbol{S}_{a_{j}}^{11}+\boldsymbol{S}_{a_{j}}^{12} \boldsymbol{P}(j)^{+} \boldsymbol{U}_{2} \boldsymbol{P}^{(j)-} \boldsymbol{S}_{a_{j+1}}^{11} \\
\boldsymbol{S}^{12}=\boldsymbol{S}_{a_{j}}^{12} \boldsymbol{P}^{(j)+} \boldsymbol{U}_{2} \boldsymbol{S}_{a_{j+1}}^{12} \\
\boldsymbol{S}^{21}=\boldsymbol{S}_{a_{j+1}}^{21} \boldsymbol{U}_{1} \boldsymbol{P}^{(j)-} \boldsymbol{S}_{a_{j}}^{21} \\
\boldsymbol{S}^{22}=\boldsymbol{S}_{a_{j+1}}^{22}+\boldsymbol{S}_{a_{j+1}}^{21} \boldsymbol{P}^{(j)-} \boldsymbol{U}_{1} \boldsymbol{P}^{(j)+} \boldsymbol{S}_{a_{j}}^{22} \tag{8.96}
\end{array}
$$

where

$$
\boldsymbol{U}_{1}=\left(\boldsymbol{I}-\boldsymbol{S}_{a_{j}}^{22} \boldsymbol{P}^{(j)+} \boldsymbol{S}_{a_{j+1}}^{11} \boldsymbol{P}^{(j)-}\right)^{-1} \boldsymbol{U}_{2}=\left(\boldsymbol{I}-\boldsymbol{S}_{a_{j+1}}^{11} \boldsymbol{P}^{(j)-} \boldsymbol{S}_{a_{j}}^{22} \boldsymbol{P}^{(j)+}\right)^{-1}
$$

### 8.6 Extensions of C Method

The key idea of C-method as applied to diffraction by surface-relief gratings is to map the surface of the grating to a plane. Until now, we have only described profiles under the form $z=a(x)$. However, in Cartesian coordinates $(x, y, z)$, a cylindrical surface whose generating line is parallel to the $O y$ axis may be described by the parametric equations:

$$
\begin{equation*}
x=f\left(x^{1}\right) z=g\left(x^{1}\right) \tag{8.97}
\end{equation*}
$$

where $f$ and $g$ are two continuous functions. Now consider the following relations:

$$
\begin{equation*}
x=f\left(x^{1}\right)+c_{1} x^{3} \quad y=x^{2} z=g\left(x^{1}\right)+x^{3} \tag{8.98}
\end{equation*}
$$

where $c_{1}$ is a real constant. They define an additive change of coordinates whose metric tensor is given by:

$$
g_{i j}=\left[\begin{array}{ccc}
\partial_{1}^{2} f+\partial_{1}^{2} g & 0 & c_{1} \partial_{1} f+\partial_{1} g  \tag{8.99}\\
0 & 1 & 0 \\
c_{1} \partial_{1} f+\partial_{1} g & 0 & 1+\partial_{1}^{2} g
\end{array}\right]
$$

Actually, the above matrix corresponds to a change of coordinates provided the Jacobian determinant $J$ of the transformation does not go to zero.

$$
J=\left|\begin{array}{cc}
\partial_{1} x & \partial_{3} x  \tag{8.100}\\
\partial_{1} z & \partial_{3} z
\end{array}\right|=\left|\begin{array}{cc}
\partial_{1} f & c_{1} \\
\partial_{1} g & 1
\end{array}\right|=\partial_{1} f-c_{1} \partial_{1} g
$$

More over, the metric tensor is independent of coordinate $x^{3}$ which means there exists a translation symmetry along $x^{3}$ axis. Hence Equations (8.98) define in a general way translation coordinate systems which allow to solve new classes of problems.

### 8.6.1 Oblique transformations

In Cartesian coordinates, usual coordinates lines of a plane are two straight lines orthogonal to each other. One can also imagine having straight lines which make an angle different from $\pi / 2$. Consider the straight line $\Delta$ given by $z=\tan (\phi) x$ and let us call $\phi$ the obliquity angle. The following sets of relations define a coordinate system $\left(x^{1}, x^{3}\right)$ in which lines parallel to $\Delta$ are coordinate lines $x^{1}=$ constant and lines $x^{3}=$ constant remain parallel to $O x$

$$
\begin{align*}
& x=x^{1}+\frac{1}{\tan \phi} x^{3}  \tag{8.101}\\
& z=x^{3}
\end{align*}
$$

Such oblique transformation allow to model an extended class of surface shapes which would otherwise be numerically inefficient (very blazed gratings) or even impossible like overhanging gratings. As an illustrative example, consider in the coordinate system $\left(x^{1}, x^{3}\right)$ the symmetric triangular function.

$$
t\left(x^{1}\right)= \begin{cases}2 x^{1} & 0<x^{1}<.5  \tag{8.102}\\ 2\left(1-x^{1}\right) & .5<x^{1}<1 \\ 0 & \text { elsewhere }\end{cases}
$$



Figure 8.4: Coordinate system in which coordinate lines are parallel to $\Delta$ and to $O x$ axis

Using an oblique transformation one gets:

$$
\begin{align*}
& x=x^{1}+\frac{1}{\tan \phi}\left(x^{3}+t\left(x^{1}\right)\right)  \tag{8.103}\\
& z=x^{3}+t\left(x^{1}\right)
\end{align*}
$$



Figure 8.5: Echelette grating in three different oblique coordinate systems
Figure (8.5) shows three typical grating surfaces obtained with (8.103) and with $\phi=90$, 63.4349 and 40 respectively, the latter demonstrating the extreme overhanging forms possible for small $\phi$ without the double value problem implicit with Cartesian coordinates.

### 8.6.2 Stretched coordinates

The essence of C-method is to choose a coordinate system that facilitate the solution of a given problem. Oblique transformations are a typical example of the usefulness of this technique. Indeed they provide an easy and elegant way to handle gratings with one vertical facet and also overhanging gratings. Similarly, we have believed for a long time that sharp edges were an intrinsic limitation of the C method. Actually, it turns out that transformations which stretch coordinates around the edges overcome the problem. With C-Method, the solution of Maxwell's equations is reduced to the solution of an algebraic eigenvalue problem in discrete Fourier
space. The derivation of the matrix operator involves two steps: (1) The electromagnetic field is expanded into Floquet-Fourier series, and (2) the derivative of the grating profile function is expanded into Fourier series. When the latter function is discontinuous, the Fourier method is known to converge slowly.This weakness remains even when the correct Fourier factorization of products of discontinuous periodic functions, as given by Li , is applied. The reason for slow convergence is that the spatial resolution of the Fourier expansion remains uniform within a grating period whatever the grating profile function may be. On the contrary, stretched coordinates allow a mapping of space that increases spatial resolution around the discontinuities of the derivative of the profile function. For this reason the technique is known as adaptive spatial resolution.

### 8.6.3 Parametric C-method

Whether for mandatory reasons as is the case for overhanging gratings or simply to improve convergence speed, the most general representation of a one dimensional profile happens to be a parametric one. Adding an additional degree of freedom with an obliquity angle, a class of translation coordinate systems has the form given by (8.98). Due to the translational symmetry along vector $\boldsymbol{e}_{3}=c_{1} \boldsymbol{e}_{x}+\boldsymbol{e}_{z}$, a numerical solution in terms of eigenvectors and eigenvalues is possible. Equations (8.98) describe a coordinate system where coordinates lines $x^{3}=$ constant coincide with functions which are periodic with period $d^{1}$ along direction $O x$. Compared to the non-oblique coordinate system, the period $d^{1}$ and the direction of periodicity remain unaffected by the introduction of parameter $c_{1}$. Thus, assuming an incident plane wave vector $\boldsymbol{k}$ such that $k . e_{x}=\alpha_{0}$ the $x^{1}$ dependence is of the form $\exp \left(i \alpha_{m} x^{1}\right)$ with $\alpha_{m}=\alpha_{0}+m \frac{2 \pi}{d^{1}}$. So we have all the ingredients to determine the matrix from which eigenvectors and eigenvalues will be sought. In Fourier space, the matrices associated to the elements of the metric tensor are:

$$
\begin{align*}
\left(\sqrt{g} g^{11}\right) & =\left(c_{1}^{2}+1\right)\left[\dot{\boldsymbol{f}}-c_{1} \dot{\boldsymbol{g}}\right]^{-1} \\
\left(\sqrt{g} g^{13}\right) & =\left[\dot{\boldsymbol{f}}-c_{1} \dot{\boldsymbol{g}}\right]^{-1}\left[c_{1} \dot{\boldsymbol{f}}+\dot{\boldsymbol{g}}\right] \\
\left(\sqrt{g} g^{31}\right) & =\left[c_{1} \dot{\boldsymbol{f}}+\dot{\boldsymbol{g}}\right]\left[\dot{\boldsymbol{f}}-c_{1} \dot{\boldsymbol{g}}\right]^{-1}  \tag{8.104}\\
\left(\sqrt{g} g^{33}\right) & =\left[\dot{\boldsymbol{f}}-c_{1} \dot{\boldsymbol{g}}\right]+\left(c_{1}^{2}+1\right)\left[c_{1} \dot{\boldsymbol{f}}+\dot{\boldsymbol{g}}\right]\left[\dot{\boldsymbol{f}}-c_{1} \dot{\boldsymbol{g}}\right]^{-1}\left[c_{1} \dot{\boldsymbol{f}}+\dot{\boldsymbol{g}}\right] \\
\left(\sqrt{g} g^{22}\right) & =\dot{\boldsymbol{f}}-c_{1} \dot{\boldsymbol{g}}
\end{align*}
$$

where $\dot{f}$ and $\dot{g}$ designates the toeplitz matrices formed by the elements of the Fourier series of $\partial_{1} f$ and $\partial_{1} g$ respectively.

### 8.6.4 Plane waves and parametric C-method

More over since the physics remains the same compared to non-oblique translation coordinate systems, eigenvectors separate into forward and backward waves as was already the case:

$$
\begin{equation*}
\mathscr{F}\left(x^{1}, x^{3}\right)=\sum_{q} A_{q}^{ \pm} F_{q}^{ \pm}\left(x^{1}\right) \exp \left(i \rho_{q}^{ \pm} x^{3}\right) \tag{8.105}
\end{equation*}
$$

As in the classical translation coordinate system, we substitute the computed propagative forward and backward eigenvectors with the corresponding transformed plane waves. Consider
plane waves

$$
\exp \left(i \alpha_{n} x\right) \exp \left( \pm i \gamma_{n} z\right)
$$

such that $\pm \gamma_{n} \in R$. Taking into account (8.98), their expression in oblique coordinates is:

$$
\begin{equation*}
\exp \left(i\left(\alpha_{n} c_{1} \pm \gamma_{n}\right) x^{3}\right) \exp \left(i\left(\alpha_{n} f\left(x^{1}\right) \pm \gamma_{n} g\left(x^{1}\right)\right)\right. \tag{8.106}
\end{equation*}
$$

hence the following correspondences between propagative waves in Cartesian coordinates and their computed counterparts in oblique coordinates:

$$
\begin{equation*}
\rho_{(p), n}^{ \pm,(M)} \longleftrightarrow\left( \pm \gamma_{n}+\alpha_{n} c_{1}\right) ; \quad F_{(p), n}^{ \pm}\left(x^{1}\right) \longleftrightarrow \exp \left(i \alpha_{n} f\left(x^{1}\right) \pm i \gamma_{n}\left(g\left(x^{1}\right)\right)\right. \tag{8.107}
\end{equation*}
$$

We have added an extra subscript $(p)$ and a superscript $(M)$ to indicate that we only care about the above correspondence for propagative waves and that $\rho$ depends on the truncation number.

### 8.6.5 Illustrative example

Consider a right angled triangular profile whose base is aligned on $O x$. Other parameters are period $d^{1}$ and height $h$. It is illuminated by a plane wave inclined at $\theta$ to the $O z$ axis. In the context of C-method we ask ourselves which coordinate system choosing for modelling diffraction by such a grating. Here the main difficulty comes from the vertical facet located at $x=d^{1}$. The operator associated with C-method involves the derivative of the profile function. With a description of the profile by a function of the kind $z=a(x)$, the derivative is constant and everything happens as if the vertical did not exist. Should the vertical be replaced by a very sloping facet, then a highly located and large discontinuity in the derivative would appear. None of the situation is satisfactory. An easy way to overcome the problem consists in introducing an oblique coordinate system in which the vertical is transformed into a straight line with a "reasonable" slope. Actually, doing so amounts to parametrizing the profile in the Cartesian coordinate system.

### 8.6.5.1 Obliquity angle and parametrization of the profile

Since one of the facets of the grating is vertical, an inclined coordinate system is needed. On the one hand, the parameter $c_{1}$ is linked to the obliquity angle $\phi$ by $c_{1}=1 / \tan \phi$ and on the other hand, according to (8.100) it should satisfy the constraint $1-c_{1} \partial_{x} a>0$. Hence, in principle $\phi$ may be any angle such that $\tan \phi<h / d$ Let $t_{1}$ be $\tan (\phi)$. On the first facet we have:

$$
\begin{equation*}
x=x^{1}+\frac{1}{t_{1}} y, \quad y=\frac{h}{d} x \tag{8.108}
\end{equation*}
$$

and on the second one

$$
\begin{equation*}
x=d, \quad d=x^{1}+\frac{1}{t_{1}} y \tag{8.109}
\end{equation*}
$$

Thus the parametrization of the profile is:

$$
\begin{array}{ll}
f\left(x^{1}\right)=\frac{d}{d-t_{1} h} x^{1} & g\left(x^{1}\right)=\frac{h}{d-t_{1} h} x^{1} \tag{8.110}
\end{array} \quad \text { if } x^{1} \leqslant x_{0}^{1} .
$$

with :

$$
\begin{equation*}
x_{0}^{1}=d\left(1-\frac{1}{t_{1}} \frac{h}{d}\right) \tag{8.111}
\end{equation*}
$$

Now that we have parametrized the profile, it remains to define a translation direction. The direction which served to parametrize the profile is a natural choice although not mandatory. Once more, the only constraint is that the tangent of the chosen obliquity angle is smaller than $h / d$.

### 8.6.5.2 Stretched coordinates and parametrization of the profile

At $x^{1}=0$ and $x^{1}=x_{0}^{1}$, the parametric functions $f\left(x^{1}\right)$ and $g\left(x^{1}\right)$ have jumps which can be reduced if one introduces an additional change of coordinates aimed at increasing spatial resolution around these points. Let $x^{1}$ be a function of a new variable $u$ : $x^{1}=s(u)$. The chain rule for derivative gives :

$$
\begin{equation*}
\partial_{u} x=\partial_{1} f\left(x^{1}(s(u)) \partial_{u} s, \quad \partial_{u} y=\partial_{1} g\left(x^{1}(s(u)) \partial_{u} s\right.\right. \tag{8.112}
\end{equation*}
$$

Compared to the initial parametrization, spatial resolution is modulated by the multiplicative factor $\partial_{u} s$. The smaller the latter, the higher the spatial resolution. A possible stretching function is as follows :

$$
s(u)= \begin{cases}u-\frac{\eta x_{0}}{2 \pi} \sin \left(\frac{2 \pi u}{x_{0}}\right) & \text { if } 0 \leqslant u<x_{0}^{1}  \tag{8.113}\\ \left(u-x_{0}\right)-\frac{\eta\left(d^{1}-x_{0}^{1}\right)}{2 \pi} \sin \left(\frac{2 \pi\left(u-x_{0}\right)}{d^{1}-x_{0}^{1}}\right) & \text { if } x_{0}^{1} \leqslant u<d^{1}\end{cases}
$$

The parameter $\eta$ between zero and one controls the density of coordinate lines around the transition points. It allows to stretch space thinner where discontinuities of coefficients in Maxwell's equations occur. The larger $\eta$, the smaller $\partial_{u} s$ and thus the higher the spatial resolution. In principle the parameter $\eta$ does not have to reach one because, in that case, the Jacobian would be zero. Figure (8.6) shows four possible parametrization of the considered right angle triangle. Case (a) corresponds to the usual representation $z=a(x)$.


Figure 8.6: Various parametric representations of a right angle triangular profile

Finally, figure 8.7) shows the speed of convergence of the specular reflected order for a perfectly conducting right angle triangular profile for two different parametrizations. It has to be emphasized that modelling this kind of profile is out of reach for the "classical" C-method since it has a vertical facet.


Figure 8.7: Comparison of speed of convergence for two parametric representation of a right angle triangular profile. Full line: $x_{0}^{1}=.5, \eta=0$, dashed line: $x_{0}^{1}=.4, \eta=.9$. Other parameters are: $\theta=25^{\circ}, \lambda=1, h=d^{1}=1$

## Appendix 8.A: Curvilinear Coordinates

In Cartesian coordinates we deal with three mutually perpendicular families of planes: $\mathrm{x}=$ constant, $y=$ constant, $z=$ constant. Imagine that we superimpose on this system three other families of surfaces. We may reference any variable point $M$ by the intersection of three planes in Cartesian coordinates, ie by the triplet $(x, y, z)$ or as the intersection of the three surfaces that form our new, curvilinear coordinates. Describing the curvilinear coordinates surfaces by $x^{1}=$ constant,$x^{2}=$ constant,$x^{3}=$ constant we may identify our point by the triplet $x, y, z$ as well as by $x^{1}, x^{2}, x^{3}$. This means that in principle we may define a curvilinear coordinate system from the Cartesian system $(x, y, z)$ by:

$$
\begin{equation*}
x=x^{1^{\prime}}=x^{1^{\prime}}\left(x^{1}, x^{2}, x^{3}\right), y=x^{2^{\prime}}=x^{2^{\prime}}\left(x^{1}, x^{2}, x^{3}\right), z=\bar{x}^{3}=x^{3^{\prime}}\left(x^{1}, x^{2}, x^{3}\right) \tag{8.114}
\end{equation*}
$$

or by the inverse relations

$$
\begin{equation*}
x^{1}=x^{1}\left(x^{1^{\prime}}, x^{2^{\prime}}, x^{3^{\prime}}\right), x^{2}=x^{2}\left(x^{1^{\prime}}, x^{2^{\prime}}, x^{3^{\prime}}\right), x^{3}=x^{3}\left(x^{1^{\prime}}, x^{2^{\prime}}, x^{3^{\prime}}\right) \tag{8.115}
\end{equation*}
$$

$x^{1^{\prime}}, x^{2^{\prime}}, x^{3^{\prime}}$ respectively $x^{1}, x^{2}, x^{3}$ are regarded as independent and continuously differentiable functions of $x^{1}, x^{2}$ and $x^{3}$ respectively $x^{1^{\prime}}, x^{2^{\prime}}, x^{3^{\prime}}$. let $M$ denote a variable point referenced by the rectangular coordinates $(x, y, z)$. At M the so-called natural referential $\left(M, e_{1}, \boldsymbol{e}_{2}, \boldsymbol{e}_{3}\right)$ is defined by the the following basis vectors:

$$
\begin{equation*}
e_{\alpha}=\sum_{\beta^{\prime}=1}^{\beta^{\prime}=3} \frac{\partial x^{\beta^{\prime}}}{\partial x^{\alpha}} e_{\beta^{\prime}} \tag{8.116}
\end{equation*}
$$

with $e_{1^{\prime}}=e_{x}, e_{2^{\prime}}=e_{y}, e_{3^{\prime}}=e_{z}, e_{x}, e_{y}$ and $e_{z}$ being the unit vectors of an orthogonal Cartesian referential. In a similar way we may write

$$
\begin{equation*}
e_{\alpha^{\prime}}=\sum_{\beta=1}^{\beta=3} \frac{\partial x^{\beta}}{\partial x^{\alpha^{\prime}}} e_{\beta} \tag{8.117}
\end{equation*}
$$

Moreover introducing $\Lambda_{\alpha}^{\beta^{\prime}}=\frac{\partial x^{\beta^{\prime}}}{\partial x^{\alpha}}$ and Einsteins' summation convention $\operatorname{Eq}(8.116)$ and $\operatorname{Eq}(8.117)$ write:

$$
\begin{equation*}
e_{\alpha}=\Lambda_{\alpha}^{\beta^{\prime}} e_{\beta^{\prime}} \quad e_{\alpha}=\Lambda_{\alpha}^{\beta} e_{\beta^{\prime}} \tag{8.118}
\end{equation*}
$$

vectors $e_{\alpha}$ are tangent vectors along coordinate curve $x^{\alpha}$. The matrix formed by the coefficient $\Lambda_{\alpha}^{\beta^{\prime}}$ is the Jacobian matrix $J$ of the change of coordinates. Since functions $x^{1^{\prime}}, x^{2^{\prime}}, x^{3^{\prime}}$ are independent $\boldsymbol{J}$ is inverible and its inverse is formed by the coefficients $\Lambda_{\alpha^{\prime}}^{\beta}$

$$
\boldsymbol{J}=\left[\begin{array}{ccc}
\Lambda_{1}^{1^{\prime}} & \Lambda_{1}^{2^{\prime}} & \Lambda_{1}^{3^{\prime}}  \tag{8.119}\\
\Lambda_{2}^{1^{\prime}} & \Lambda_{2}^{2^{\prime}} & \Lambda_{2}^{3^{\prime}} \\
\Lambda_{3}^{1^{\prime}} & \Lambda_{3}^{2^{\prime}} & \Lambda_{3}^{3^{\prime}}
\end{array}\right] \quad \boldsymbol{J}^{-1}=\left[\begin{array}{ccc}
\Lambda_{1^{\prime}}^{1} & \Lambda_{1^{\prime}}^{2} & \Lambda_{1^{\prime}}^{3} \\
\Lambda_{2^{\prime}}^{1} & \Lambda_{2^{\prime}}^{2} & \Lambda_{2^{\prime}}^{3} \\
\Lambda_{3^{\prime}}^{1} & \Lambda_{3^{\prime}}^{2} & \Lambda_{3^{\prime}}^{3}
\end{array}\right]
$$

One can also define basis vectors $e^{\alpha}$ that are normal to coordinate surfaces $x^{\alpha}=$ constant by

$$
\begin{equation*}
e^{\alpha}=\frac{\partial x^{\alpha}}{\partial x^{\alpha^{\prime}}} e^{\alpha^{\prime}}, \quad \text { with } e^{\alpha^{\prime}}=e_{\alpha^{\prime}} \tag{8.120}
\end{equation*}
$$

The vectors $e_{\alpha}$ and $e^{\beta}$ form a set of reciprocal basis with $e_{\alpha} \cdot e^{\beta}=\delta_{\alpha}^{\beta}$, where $\delta_{\alpha}^{\beta}$ is the Kronecker delta. The representation of any vector $\boldsymbol{A}$ in one of these bases is:

$$
\begin{equation*}
\boldsymbol{A}=A^{\alpha} e_{\alpha}=A_{\alpha} e^{\alpha} \tag{8.121}
\end{equation*}
$$



Figure 8.8: Curvilinear Coordinates:covariant and contravariant components of a vector in a plane

The $A^{\alpha}$ and the $A_{\alpha}$ are the contravariant components and the covariant components of vector $\boldsymbol{A}$ respectively. The nullity of a component of vector $\boldsymbol{A}$ may be geometrically interpreted as follows:
$A_{\alpha}=0: \boldsymbol{A}$ is orthogonal to the tangent at point $M$ to the coordinate line $x^{\alpha}$
$A^{\alpha}=0: \boldsymbol{A}$ belongs to the tangential plane at point $M$ to coordinate surface $x^{\alpha}$
In normalized orthogonal Cartesian coordinates differentiate contravariant and covariant components of a vector is generally not necessary. The Jacobian matrix allows to express the Cartesian components $A^{\alpha^{\prime}}=A_{\alpha^{\prime}}$ of vector $\boldsymbol{A}$ in terms of its local contravariant components $A^{\alpha}=\boldsymbol{A} \cdot \boldsymbol{e}^{\alpha}$ or covariants components $A_{\alpha}=\boldsymbol{A} \cdot \boldsymbol{e}_{\alpha}$ :

$$
\begin{gather*}
A^{\alpha^{\prime}}=A_{\alpha^{\prime}}=\Lambda_{\alpha}^{\alpha^{\prime}} A^{\alpha} \quad \text { or } \quad A^{\alpha}=\Lambda_{\alpha^{\prime}}^{\alpha} A^{\alpha^{\prime}}  \tag{8.122}\\
A_{\alpha^{\prime}}=\Lambda_{\alpha^{\prime}}^{\alpha} A_{\alpha} \quad \text { or } \quad A_{\alpha}=\Lambda_{\alpha}^{\alpha^{\prime}} A_{\alpha^{\prime}} \tag{8.123}
\end{gather*}
$$

The quantities

$$
\begin{equation*}
g_{\alpha \beta}=e_{\alpha} \cdot e_{\beta}=\Lambda_{\alpha}^{\alpha^{\prime}} \Lambda_{\beta}^{\beta^{\prime}} g_{\alpha^{\prime} \beta^{\prime}} \tag{8.124}
\end{equation*}
$$

define the metric of the coordinate system. In matrix form we have:

$$
\begin{equation*}
\left[g_{\alpha \beta}\right]=\boldsymbol{J}^{t} \boldsymbol{J} \tag{8.125}
\end{equation*}
$$

and

$$
\begin{equation*}
g=\operatorname{det}\left(\left[g_{\alpha \beta}\right]\right)=\operatorname{det}(\boldsymbol{J})^{2}=\operatorname{det}\left(\Lambda_{\beta}^{\alpha^{\prime}}\right) \tag{8.126}
\end{equation*}
$$

The $g_{\alpha \beta}$ establish a connexion between the $A^{\alpha}$ and the $A_{\beta}$

$$
\begin{equation*}
A_{\alpha}=e_{\alpha} \cdot\left(A^{\beta} \cdot e_{\beta}\right)=g_{\alpha \beta} A^{\beta} \quad \text { or } \quad A^{\beta}=g^{\beta \alpha} A_{\alpha} \tag{8.127}
\end{equation*}
$$

## Appendix 8.B: Transformation of Maxwell's equations

We have seen that the natural referentiel gives the tools to easily maipulate tangential and normal components of a vector field.Therefore, writting boundary conditions at a surface should be
straightforward. We need now to express Maxwell's equation in the new coordinate system. For that purpose, we may follow a tensorial approach or stay at an elementary level and make a simple change of coordinates and components in the usual Maxwell's equations. We present briefly both points of view. A time dependence of the form $\exp (-i \omega t)$ is assumed.

## Vectorial approch

Let us start from one of the Maxwell's curl equation written in the Cartesian coordinate system and in an homogeneous medium with permittivity $\varepsilon$ and permeability $\mu$

$$
\begin{equation*}
\xi^{\alpha^{\prime} \beta^{\prime} \gamma^{\prime}} \partial_{\beta^{\prime}} H_{\gamma^{\prime}}=-i \omega \varepsilon E_{\alpha^{\prime}} \tag{8.128}
\end{equation*}
$$

where $\xi^{\alpha^{\prime} \beta^{\prime} \gamma^{\prime}}$ stands for the Levi-Civita indicator :

$$
\xi^{\alpha^{\prime} \beta^{\prime} \gamma^{\prime}}=\left\{\begin{array}{l}
1 \text { for } \alpha^{\prime} \beta^{\prime} \gamma^{\prime}=123,231,312  \tag{8.129}\\
-1 \text { for } \alpha^{\prime} \beta^{\prime} \gamma^{\prime}=321,213,132 \\
0 \text { otherwise }
\end{array}\right.
$$

Then let us change the coordinates: $\partial_{\beta^{\prime}}=\Lambda_{\beta^{\prime}}^{\alpha} \partial_{\alpha}$ and the components $H_{\gamma^{\prime}}=\Lambda_{\gamma^{\prime}}^{\beta} H_{\beta}$

$$
\begin{equation*}
\xi^{\alpha^{\prime} \beta^{\prime} \gamma^{\prime}} \Lambda_{\beta^{\prime}}^{\alpha} \partial_{\alpha}\left(\Lambda_{\gamma^{\prime}}^{\beta} H_{\beta}\right)=-i \omega \varepsilon E_{\alpha^{\prime}} \tag{8.130}
\end{equation*}
$$

The left hand side of the above equation is equal to:

$$
\begin{equation*}
\xi^{\alpha^{\prime} \beta^{\prime} \gamma^{\prime}} \Lambda_{\beta^{\prime}}^{\alpha}\left(\Lambda_{\gamma^{\prime}}^{\beta} \partial_{\alpha} H_{\beta}\right)+\xi^{\alpha^{\prime} \beta^{\prime} \gamma^{\prime}} \Lambda_{\beta^{\prime}}^{\alpha}\left(\partial_{\alpha} \Lambda_{\gamma^{\prime}}^{\beta}\right) H_{\beta} \tag{8.131}
\end{equation*}
$$

on the one hand we have

$$
\begin{equation*}
\Lambda_{\beta^{\prime}}^{\alpha} \partial_{\alpha} \Lambda_{\gamma^{\prime}}^{\beta}=\partial_{\beta^{\prime}} \Lambda_{\gamma^{\prime}}^{\beta} \tag{8.132}
\end{equation*}
$$

and the other hand this term is symmetrical with respect to $\beta^{\prime}$ and $\gamma^{\prime}$. Thus by applying the operator $\xi^{\alpha^{\prime} \beta^{\prime} \gamma^{\prime}}$ which is antisymmetric with respect to $\beta^{\prime}$ and $\gamma^{\prime}$ we obtain 0 . Thus, the Maxwell curl equation reduces to:

$$
\begin{equation*}
\xi^{\alpha^{\prime} \beta^{\prime} \gamma^{\prime}} \Lambda_{\beta^{\prime}}^{\alpha} \Lambda_{\gamma^{\prime}}^{\beta} \partial_{\alpha} H_{\beta}=-i \omega \varepsilon E_{\alpha^{\prime}} \tag{8.133}
\end{equation*}
$$

let us multiply both sides by $\Lambda_{\alpha^{\prime}}^{\gamma}$ and make summation on dummy index $\alpha^{\prime}$. We obtain:

$$
\begin{equation*}
\xi^{\gamma \alpha \beta} \operatorname{det}\left(\Lambda_{\beta^{\prime}}^{\alpha}\right) \partial_{\alpha} H_{\beta}=-i \omega \varepsilon \Lambda_{\alpha^{\prime}}^{\gamma} E^{\alpha^{\prime}}=-i \omega \varepsilon E^{\gamma}=-i \omega \varepsilon g^{\gamma \beta} E_{\beta} \tag{8.134}
\end{equation*}
$$

Finally we get:

$$
\begin{equation*}
\xi^{\gamma \alpha \beta} \partial_{\alpha} H_{\beta}=-i \omega \varepsilon \sqrt{g} g^{\gamma \beta} E_{\beta} \tag{8.135}
\end{equation*}
$$

and

$$
\begin{equation*}
\xi^{\gamma \alpha \beta} \partial_{\alpha} E_{\beta}=i \omega \mu \sqrt{g} g^{\gamma \beta} H_{\beta} \tag{8.136}
\end{equation*}
$$

setting

$$
\begin{equation*}
\mathscr{F}_{\alpha}=E_{\alpha} \quad \mathscr{G}_{\alpha}=i Z H_{\alpha} \quad \text { with } Z=\sqrt{\frac{\mu}{\varepsilon}} \tag{8.137}
\end{equation*}
$$

we then obtain a set of equations relating the complex amplitudes of the field components where the $\mathscr{F}_{\alpha}$ and the $\mathscr{G}_{\alpha}$ play a fully symmetric role:

$$
\begin{align*}
& \xi^{\alpha \beta \gamma} \partial_{\beta} \mathscr{F}_{\gamma}=k \sqrt{g} g^{\alpha \beta} \mathscr{G}_{\alpha}  \tag{8.138}\\
& \xi^{\alpha \beta \gamma} \partial_{\beta} \mathscr{G}_{\gamma}=k \sqrt{g} g^{\alpha \beta} \mathscr{F}_{\alpha}
\end{align*}
$$

where $k=\omega \sqrt{\mu \varepsilon}$

## Appendix 8.C: Summary of tensorial approach

In curvilinear coodinates systems, the Maxwell's equations are based on the tensorial formalism deduced from relativity. If we consider only materials which are stationary with respect to the coordinate system, then the four-dimensional formalism developed by Post can be simplified. Maxwell equations are written :

$$
\begin{align*}
& \xi^{\alpha \beta \gamma} \partial_{\beta} E_{\gamma}=-\partial_{t} B^{\alpha} \\
& \xi^{\alpha \beta \gamma} \partial_{\beta} H_{\gamma}=\partial_{t} D^{\alpha}+J^{\alpha}  \tag{8.139}\\
& \partial_{\alpha} D^{\alpha}=\rho \\
& \partial_{\alpha} B^{\alpha}=0
\end{align*}
$$

Post's formalism preserves the affine nature of Maxwell' equations: their expression is independent of the coordinate system. The geometry only appears in the constitutive equations along with the material's properties

$$
\begin{equation*}
D^{\alpha}=\varepsilon^{\alpha \beta} E_{\beta} \quad B^{\alpha}=\mu^{\alpha \beta} H_{\beta} \tag{8.140}
\end{equation*}
$$

In a perfectly linear, isotropic media with permittivity $\varepsilon$ and $\mu$, these relation ships become:

$$
\begin{equation*}
\varepsilon^{\alpha \beta}=\varepsilon \sqrt{g} g^{\alpha \beta} \quad \varepsilon^{\alpha \beta}=\varepsilon \sqrt{g} g^{\alpha \beta} \tag{8.141}
\end{equation*}
$$

where $g^{\alpha \beta}$ are the contravariant components of the metric tensor.

$$
\begin{equation*}
\left(g^{\alpha \beta}\right)=\left(g_{\alpha \beta}\right)^{-1} \quad g=(\operatorname{det})\left(g_{\alpha \beta}\right) \tag{8.142}
\end{equation*}
$$

In an arbitrary curvilinear coordinates systeme $x^{\alpha}$, if the surface separating two materials, denoted (1) and (2), coincides with a surface of coordinates $x^{3}=$ constant, for example, then the conditions of continuity are expressed quite simply

$$
\begin{gather*}
\text { tangential component continuity: }\left\{\begin{array}{l}
a_{1}(1)=a_{1}(2) \\
a_{2}(1)=a_{2}(2)
\end{array}\right.  \tag{8.143}\\
\text { normal component continuity: } a^{3}(1)=a^{3}(2) \tag{8.144}
\end{gather*}
$$

Assuming a time dependence of the form $\exp (-i \omega t)$, in a source free region if we substitute the constitutive equations for the material8.141 into Maxwell equations in the covariant form 8.139 , setting

$$
\begin{equation*}
\mathscr{F}_{\alpha}=E_{\alpha} \quad \mathscr{G}_{\alpha}=-i Z H_{\alpha} \quad \text { with } Z=\sqrt{\frac{\mu}{\varepsilon}} \tag{8.145}
\end{equation*}
$$

we then obtain a set of equations relating the complex amplitudes of the field components where the $\mathscr{F}_{\alpha}$ and the $\mathscr{G}_{\alpha}$ play a fully symmetric role:

$$
\begin{align*}
& \xi^{\alpha \beta \gamma} \partial_{\beta} \mathscr{F}_{\gamma}=k \sqrt{g} g^{\alpha \beta} \mathscr{G}_{\alpha}  \tag{8.146}\\
& \xi^{\alpha \beta \gamma} \partial_{\beta} \mathscr{G}_{\gamma}=k \sqrt{g} g^{\alpha \beta} \mathscr{F}_{\alpha}
\end{align*}
$$

where $k=\omega \sqrt{\mu \varepsilon}$

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Chapter 9:
Finite Difference Time Domain Method for Grating Structures

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and

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## Chapter 9

# Finite Difference Time Domain Method For Grating Structures 

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The Finite Difference Time Domain method (FDTD), based on the Yee's scheme, is one of the most commonly used time methods for the modeling of electromagnetic waves propagation and diffraction. It was first introduced by Yee in 1966 [1] in the context of differential equations resolution and the first articles recommending its futur applications are published from 1975 [2, 3, 4]. Due to the simplicity of its implementation and the rapid growth of computing capacity, the FDTD is gaining users in all areas of electromagnetism applications. It allows a real-time monitoring of the electromagnetic wave evolution in any kind of environment (dielectric, metal, plasma...). Its theoretical formulation is very easy since it requires no matrix inversion and could take into account the more complex geometric shapes of objects in the studied system. In addition, using this time domain method, a wide spectral range characterization can be obtained from one temporal calculation via a simple Fourrier transform.

In this chapter, we present a brief review on the fundamentals of the FDTD method. We show how to adapt it to the calculation of the photonic band gap structures in the case of 2D periodic (invariant in the third direction) structures. The both in-plane, for the TE and TM polarizations, and off-plane propagations are considered. The last part of this chapter is devoted to FDTD general formulation, based on the Split Field Method technique, for the modeling of bi-periodic gratings that are finished according to the third direction.

### 9.1 Fundamentals of the FDTD method

### 9.1.1 The Yee's algorithm

The FDTD method is based on the numerical resolution of the Maxwell's equations using a centered finite difference schema to approximate the partial derivatives both in time and space.

Let us start from these equations expressed in their differential formulation:

$$
\begin{align*}
& \nabla \times \vec{E}=-\frac{\partial \vec{B}}{\partial t}  \tag{9.1}\\
& \nabla \times \vec{H}=\frac{\partial \vec{D}}{\partial t} \tag{9.2}
\end{align*}
$$

The electromagnetic properties of the medium are described through the so-called constitutive relationships:

$$
\begin{align*}
\vec{D} & =\varepsilon \vec{E}  \tag{9.3}\\
\vec{B} & =\mu \vec{H} \tag{9.4}
\end{align*}
$$

$\varepsilon$ and $\mu$ are respectively the dielectric permittivity and magnetic permeability of the medium.
In a Cartesian coordinate system $(O, x, y, z)$, the Maxwell's equations in the time domain are written as:

$$
\begin{align*}
\frac{\partial H_{x}}{\partial t} & =\frac{1}{\mu}\left[\frac{\partial E_{y}}{\partial z}-\frac{\partial E_{z}}{\partial y}\right]  \tag{9.5.a}\\
\frac{\partial H_{y}}{\partial t} & =\frac{1}{\mu}\left[\frac{\partial E_{z}}{\partial x}-\frac{\partial E_{x}}{\partial z}\right]  \tag{9.5.b}\\
\frac{\partial H_{z}}{\partial t} & =\frac{1}{\mu}\left[\frac{\partial E_{x}}{\partial y}-\frac{\partial E_{y}}{\partial x}\right]  \tag{9.5.c}\\
\frac{\partial E_{x}}{\partial t} & =\frac{1}{\varepsilon}\left[\frac{\partial H_{z}}{\partial y}-\frac{\partial H_{y}}{\partial z}\right]  \tag{9.5.d}\\
\frac{\partial E_{y}}{\partial t} & =\frac{1}{\varepsilon}\left[\frac{\partial H_{x}}{\partial z}-\frac{\partial H_{z}}{\partial x}\right]  \tag{9.5.e}\\
\frac{\partial E_{z}}{\partial t} & =\frac{1}{\varepsilon}\left[\frac{\partial H_{y}}{\partial x}-\frac{\partial H_{x}}{\partial y}\right] \tag{9.5.f}
\end{align*}
$$

The numerical treatment of the partial differential equations 9.5 requires a space and time discretization. The calculation volume, shown in figure 9.1 is a rectangular parallelepiped divided into $\left(N_{x} \times N_{y} \times N_{z}\right)$ cells, each one with elementary volume $(\Delta x \times \Delta y \times \Delta z)$ where $\Delta x, \Delta y$ and $\Delta z$ are the spatial discretization steps according to the $O x, O y$ and $O z$ directions respectively.

Each well defined node of the grid is associated with a triplet of integers $(i, j, k)$ so that the coordinates $\left(x_{i}, y_{j}, z_{k}\right)$ of the node satisfy:

$$
\begin{aligned}
x_{i} & =i \cdot \Delta x \\
y_{j} & =j \cdot \Delta y \\
z_{k} & =k \cdot \Delta z
\end{aligned}
$$

The computational time is also discretized with a $\Delta t$ time step. Each computing time $t$ is associated with the integer $n$ defining the number of temporal sampling:

$$
t=n \cdot \Delta t
$$



Figure 9.1: An exemple of the FDTD calculation volume.

Temporal and spatial derivatives of the field components ( $E_{x}, E_{y}, E_{z}, H_{x}, H_{y}, H_{z}$ ) are approximated from their Taylor development to the first order. Thus, if $U$ is one of these components, we will adopt the following notation:

$$
\begin{equation*}
U\left(x_{i}, y_{j}, z_{k}, t\right)=U_{i, j, k}^{n} \tag{9.6}
\end{equation*}
$$

The temporal derivative of the $U$ component at $t$ time and $\left(x_{i}, y_{j}, z_{k}\right)$ node is approximated with finite centred difference as follows:

$$
\begin{equation*}
\left[\frac{\partial U}{\partial t}\right]_{i, j, k}=\frac{U_{i, j, k}^{n+\frac{1}{2}}-U_{i, j, k}^{n-\frac{1}{2}}}{\Delta t}+0\left([\Delta t]^{2}\right) \tag{9.7}
\end{equation*}
$$

The spatial derivatives of the $U$ component are approximated in the same manner:

$$
\begin{align*}
& {\left[\frac{\partial U}{\partial x}\right]_{j, k, n}=\frac{U_{i+\frac{1}{2}, j, k}^{n}-U_{i-\frac{1}{2}, j, k}^{n}}{\Delta x}+0\left([\Delta x]^{2}\right)}  \tag{9.8.a}\\
& {\left[\frac{\partial U}{\partial y}\right]_{i, k, n}=\frac{U_{i, j+\frac{1}{2}, k}^{n}-U_{i, j-\frac{1}{2}, k}^{n}}{\Delta y}+0\left([\Delta y]^{2}\right)}  \tag{9.8.b}\\
& {\left[\frac{\partial U}{\partial z}\right]_{i, j, n}=\frac{U_{i, j, k+\frac{1}{2}}^{n}-U_{i, j, k-\frac{1}{2}}^{n}}{\Delta z}+0\left([\Delta z]^{2}\right)} \tag{9.8.c}
\end{align*}
$$

As explicitly mentioned in equations 9.8 , the use of centered difference scheme allows a precision of the second order even if a first order Taylor development is considered. This greatly enhances the numerical convergence of the FDTD algorithm.


Figure 9.2: Spatial discretization: Yee's cell.


Figure 9.3: Temporal discretization into the Yee's scheme.

## Yee's algorithm

The algorithm proposed by Kane Yee in 1966 [1] uses in a clever way this discretization for solving the system of equations (9.5). In the Yee's scheme, the electromagnetic field components are located at different points in a unit cell (Figure9.2). The electric field components are calculated along the edges of the cell while the perpendicular magnetic field components are calculated at the centers of the cell faces. Thus, each electric field component is surrounded by four magnetic field components and similarly for each magnetic field component.

The temporal increment into the Yee's scheme is done through a "leapfrog" discretization schema. The field components $\vec{H}$ (or $\vec{E}$ ) are calculated at times odd multiples of the half time step $\frac{\Delta t}{2}$, while the field components $\vec{E}$ (respectively $\vec{H}$ ) are updated at the times even multiples of $\frac{\Delta t}{2}$ as shown in figure 9.3 . Such a discretization allows evaluating the time derivatives by keeping a centered difference schema as for spatial derivatives.

Consequently, replacing the partial derivatives in equations (9.5) by central difference (9.7. 9.8 , according to the Yee's scheme leads to the updated equations of electromagnetic components in the FDTD algorithm:

$$
\begin{align*}
& \begin{aligned}
& H_{x_{\left(i, j+\frac{1}{2}, k+\frac{1}{2}\right)}^{n}}^{n+\frac{1}{2}}=H_{x_{\left(i, j+\frac{1}{2}, k+\frac{1}{2}\right)}}^{n-\frac{1}{2}}-\frac{\Delta t}{\mu_{0} \Delta}\left\{\left[E_{z_{\left(i, j+1, k+\frac{1}{2}\right)}}^{n}-E_{z_{\left(i, j, k+\frac{1}{2}\right)}}^{n}\right]+\right. \\
& {\left.\left[E_{y_{\left(i, j+\frac{1}{2}, k\right)}^{n}}^{n}-E_{y_{\left(i, j+\frac{1}{2}, k+1\right)}}^{n}\right]\right\} } \\
& H_{y_{\left(i+\frac{1}{2}, j, k+\frac{1}{2}\right)}^{n+\frac{1}{2}}}^{n}=H_{y_{\left(i+\frac{1}{2}, j, k+\frac{1}{2}\right)}^{n-\frac{1}{2}}}^{n}-\frac{\Delta t}{\mu_{0} \Delta}\left\{\left[E_{x_{\left(i+\frac{1}{2}, j, k+1\right)}}^{n}-E_{x_{\left(i+\frac{1}{2}, j, k\right)}}^{n}\right]+\right. \\
& {\left.\left[E_{z_{\left(i, j, k+\frac{1}{2}\right)}^{n}}^{n}-E_{z_{\left(i+1, j, k+\frac{1}{2}\right)}}^{n}\right]\right\} }
\end{aligned}
\end{align*}
$$

$$
H_{z_{\left(i+\frac{1}{2}, j+\frac{1}{2}, k\right)}^{n+\frac{1}{2}}}^{n+} H_{z_{\left(i+\frac{1}{2}, j+\frac{1}{2}, k\right)}^{n-\frac{1}{2}}}^{n}-\frac{\Delta t}{\mu_{0} \Delta}\left\{\left[E_{y_{\left(i+1, j+\frac{1}{2}, k\right)}^{n}}^{n}-E_{y_{\left(i, j+\frac{1}{2}, k\right)}^{n}}^{n}\right]+\right.
$$

$$
\begin{equation*}
\left.\left[E_{x_{\left(i+\frac{1}{2}, j, k\right)}^{n}}^{n}-E_{x_{\left(i+\frac{1}{2}, j+1, k\right)}^{n}}^{n}\right]\right\} \tag{9.9.c}
\end{equation*}
$$

$$
E_{x_{\left(i+\frac{1}{2}, j, k\right)}}^{n+1}=E_{x_{\left(i+\frac{1}{2}, j, k\right)}}^{n}+\frac{\Delta t}{\varepsilon \Delta}\left\{\left[H_{z_{\left(i+\frac{1}{2}, j+\frac{1}{2}, k\right)}^{n}}^{n}-H_{z_{\left(i+\frac{1}{2}, j-\frac{1}{2}, k\right)}^{n}}^{n}\right]+\right.
$$

$$
\begin{equation*}
\left.\left[H_{\left(i+\frac{1}{2}, j, k-\frac{1}{2}\right)}^{n}-H_{y}^{n}{ }_{\left(i+\frac{1}{2}, j, k+\frac{1}{2}\right)}\right]\right\} \tag{9.9.d}
\end{equation*}
$$

$$
E_{y_{\left(i, j+\frac{1}{2}, k\right)}^{n+1}}^{n+1}=E_{y_{\left(i, j+\frac{1}{2}, k\right)}^{n}}^{n}+\frac{\Delta t}{\varepsilon \Delta}\left\{\left[H_{x_{\left(i, j+\frac{1}{2}, k+\frac{1}{2}\right)}^{n}}-H_{x_{\left(i, j+\frac{1}{2}, k-\frac{1}{2}\right)}^{n}}\right]+\right.
$$

$$
\begin{equation*}
\left.\left[H_{z_{\left(i-\frac{1}{2}, j+\frac{1}{2}, k\right)}^{n}}-H_{z_{\left(i+\frac{1}{2}, j+\frac{1}{2}, k\right)}^{n}}^{n}\right]\right\} \tag{9.9.e}
\end{equation*}
$$

$$
E_{z_{\left(i, j, k+\frac{1}{2}\right)}}^{n+1}=E_{z_{\left(i, j, k+\frac{1}{2}\right)}^{n}}^{n}+\frac{\Delta t}{\varepsilon \Delta}\left\{\left[H_{\left(i+\frac{1}{2}, j, k+\frac{1}{2}\right)}^{n}-H_{y_{\left(i-\frac{1}{2}, j, k+\frac{1}{2}\right)}^{n}}\right]+\right.
$$

$$
\begin{equation*}
\left.\left[H_{x_{\left(i, j-\frac{1}{2}, k+\frac{1}{2}\right)}^{n}}^{n}-H_{x_{\left(i, j+\frac{1}{2}, k+\frac{1}{2}\right)}^{n}}^{n}\right]\right\} \tag{9.9.f}
\end{equation*}
$$

Let us note that this last equation system can be simplified significantly in case of 2D structures (see section 2 of this chapter).

For the modeling of structures with a symmetry of revolution, a basis change from Cartesian to cylindrical coordinates is strongly recommended to accurately describe the fine details of the samples and to make more flexible the FDTD calculation codes. In these so-called BORFDTD (Body of Revolution FDTD) codes, the symmetry of revolution is exploited to express
the azimuthal dependence $(\phi)$ of the electromagnetic fields as Fourrier series. BOR-FDTD algorithm can, in this case, compute solutions for all Fourier modes through one simulation per mode. This code is commonly called 2.5D since the azimuthal field variation is analytically accounted for. Thus, there is no griding in the $\phi$-direction. This implies that the BOR-FDTD algorithm is two-dimensional in terms of computer ressource usage even 3D structures are modeled.

### 9.1.2 Spatiotemporal criteria of convergence

As all explicit schemes, Yee's algorithm is subjected to a stability condition setting the time step from the space discretization. Arbitrary values of spatiotemporal discretization can lead to infinite solutions of the electromagnetic field. Stability problems in explicit numerical methods have been analyzed in detail by Courant, Friedrichs and Levy [5] and Von Neumann, from a mathematically rigorous approach. This analysis shows that the explicit schemes are stable under a condition called CFL (for Current, Friedrich and Levy) and applied to the FDTD method in the case of a regular mesh [6]:

$$
\begin{equation*}
\Delta t \leq\left[v_{\max } \cdot \sqrt{\frac{1}{\Delta x^{2}}+\frac{1}{\Delta y^{2}}+\frac{1}{\Delta z^{2}}}\right]^{-1} \tag{9.10}
\end{equation*}
$$

where $v_{\max }$ is the maximum velocity of light propagation in the studied system, generally the velocity of light in vacuum.

In case of uniforme mesh $(\Delta x=\Delta y=\Delta z=\triangle)$, the CFL criterion becomes:

$$
\begin{array}{ll}
\Delta t \leq \frac{1}{v_{\max }} \cdot \frac{\Delta}{\sqrt{3}} & \text { in } 3 D \\
\Delta t \leq \frac{1}{v_{\max }} \cdot \frac{\Delta}{\sqrt{2}} & \text { in } 2 D \tag{9.12}
\end{array}
$$

However, it is possible to overcome the restrictive assumption of regular mesh that achieves the above result with the following generalized criterion:

$$
\begin{equation*}
\Delta t \leq\left[v_{\max } \cdot{\sqrt{\frac{1}{\Delta x_{\text {min }}^{2}}+\frac{1}{\Delta y_{\text {min }}^{2}}+\frac{1}{\Delta z_{\text {min }}^{2}}}}^{-1}\right. \tag{9.13}
\end{equation*}
$$

where $\Delta x_{\text {min }}, \Delta y_{\text {min }}$ et $\Delta z_{\text {min }}$ are the smallest step in the three directions $x, y$ and $z$ respectively.
In addition to the numerical instability problem, the transition from continuous forms of Maxwell's equations to the discrete numerical approximations can cause a parasitic effect called "numerical dispersion". This is due to the fact that numerical signals are propagated over time in the FDTD grid, with a phase velocity less than the actual velocity. This dispersion varies with frequency, propagation direction in the grid and the spatial discretization [6]. Numerical dispersion errors increase with the signal frequency and size of the computational domain, thus
making the simulation results less reliable. They may appear in various forms: phase error, signal distortion, loss of amplitude, pulse broadening ...

The solution to this problem requires a very fine mesh in the FDTD grid, so that the maximum discretization is of the order $\lambda_{\text {min }} / 20$ [6], $\lambda_{\text {min }}$ being the minimum wavelength of propagating waves in the FDTD grid.

### 9.1.3 Absorbing boundary conditions - Perfectly Matched Layers

Such conditions allow us to describe open systems where emitted or reflected waves propagate to infinity. Indeed, the limited memory space of computers requires users to truncate their FDTD computational domain. At the limits of this truncated domain, components of the electromagnetic field can not be calculated by the discretized equations (9.9). Therefore special treatment at the borders is needed to avoid the incident electromagnetic wave on these "edges" does reflect back and contaminate the actual physical signal. One of the most widely used technique is that proposed by Berenger [7] called Perfectly Matched Layer (PML). This technique consists of adding around the studied domain not necessarily physical layer causing no reflection and almost totally absorbing all the propagating electromagnetic field. Its use is based on the condition of impedance matching of two waves at the interface between two media with the same index but which one is absorbing (with nonzero electrical conductivity $\sigma$ and magnetic equivalent conductivity $\sigma^{\star}$ as shown in figure 9.4).


Figure 9.4: Impedance matching principle.
This condition is expressed as:

$$
\begin{equation*}
\frac{\sigma}{\varepsilon}=\frac{\sigma^{\star}}{\mu_{0}} \tag{9.14}
\end{equation*}
$$

Thus, a magnetic conductivity is needed to fulfill this impedance matching condition. In addition, absorption is needed only for components of the fields that propagates perpendicularly to the interface (the FDTD window border) and not in the parallel direction. Bérenger solved this problem by proposing an artificially biaxial absorbing medium. The absorption is not zero in the direction normal to the interface between the two media and is zero along the axis parallel to the interface. In the PML medium, the incident plane wave is split into two fictitious waves (see figure 9.5):

1) A wave propagating at normal incidence and satisfying the equation 9.14 . This wave is attenuated and absorbed by the PML medium and undergoes only very low reflectivity to the incident medium.


Figure 9.5: Schematic of the PML principle.
2) A second grazing incidence wave that shows no absorption in the PML medium. This wave, propagating parallel to the interface between two media undergoes no reflection and sees a medium identical to that of the main grid window.

Abrupt changes in conductivities at this interface degrade the performances of absorption. This effect is, however, reduced by imposing a progressive variation of the absorption according to a polynomial law given by [7]:

$$
\begin{equation*}
\sigma=\sigma_{\max }\left(\frac{x_{p m l}}{e}\right)^{m} \tag{9.15}
\end{equation*}
$$

where $\sigma_{\max }$ is the maximum value of the conductivity, $x_{p m l}$ represents the depth in the PML region measured from the interface, $e$ denotes the thickness of the PML layer and $m$ is the polynomial order generally fixed to 2 .

Let us note that in the case of gratings such conditions are not necessary according to the periodicity directions. The absorbing boundaries conditions are hereby replaced by FloquetBloch periodic conditions in order to describe periodic structures (see section 2 of this chapter). Nevertheless, for a 2D periodic structure, PML are needed in the third direction where the structure is usually finite.

### 9.1.4 Dispersive media

The dispersive media, such as metals in the optical range, are characterized by a complex permittivity frequency dependent $\varepsilon(\omega)=\varepsilon^{\prime}(\omega)+i \varepsilon^{\prime \prime}(\omega)$. As the FDTD method is temporal, in such environments the direct implementation of the above equations, in which appear explicitly permittivity and hence the frequency, is impossible. The solution for this problem is to calculate the displacement vector $\vec{D}$ components in the classical Yee's scheme and then back to electrical field components using the constitutive equation of the medium established in the frequency domain $\vec{D}(\omega)=\varepsilon(\omega) \vec{E}(\omega)$. The temporal nature of the FDTD needs a temporal constitutive equation written as a convolution product $\vec{D}(t)=\varepsilon(t) \otimes \vec{E}(t)$. It is a non local relationship whose resolution requires the knowledge of the electric field at all previous times. Numerically, this leads to a storage of a very large amount of data and therefore requires to have a very large memory space. This issue can be bypassed using analytical models describing the dielectric function $\varepsilon(\omega)$ of these metals. The choice of adapted analytical model depends on the type of metal as well as the spectral range of study.

### 9.1.4.1 Drude Model

The Drude model of free electrons [8, 9] for the dielectric function which, although based on a purely classical approach, can well account for intraband transitions. In this model, firstly proposed in 1908 by P. Drude, a gas of free electrons moving in a immobile metal ions lattice. Thus, the electron-electron interactions and electron-ions are not taken into account and the movement of all the electron cloud is thus the average of the movements of individual electrons. The relative permittivity given by this model is:

$$
\begin{equation*}
\varepsilon_{D}=\varepsilon_{\infty}-\frac{\omega_{D}^{2}}{\omega^{2}+i \omega \gamma_{D}} \tag{9.16}
\end{equation*}
$$

where $\omega_{D}$ is the "plasma frequency" of the metal and $\varepsilon_{\infty}$ its relative permittivity at infinite frequencies. $\gamma_{D}$ represents a damping term that is inversely proportional to the relaxation time.

## FDTD implementation of the Drude model

The principle consists in replacing the electric field vector $\vec{E}$ by $\vec{D} / \varepsilon$ in Maxwell's equations in order to eliminate $\varepsilon$ term. In dispersive media, equations 9.9.d 9.9.e et 9.9.f) are replaced by:

$$
\begin{align*}
& D_{x_{\left(i+\frac{1}{2}, j, k\right)}}^{n+1}=D_{x_{\left(i+\frac{1}{2}, j, k\right)}}^{n}+\frac{\Delta t}{\Delta}\left\{\left[H_{z_{\left(i+\frac{1}{2}, j+\frac{1}{2}, k\right)}^{n}}^{n}-H_{z_{\left(i+\frac{1}{2}, j-\frac{1}{2}, k\right)}^{n}}^{n}\right]+\right. \\
& \left.\left[H_{\left(i+\frac{1}{2}, j, k-\frac{1}{2}\right)}^{n}-H_{y_{\left(i+\frac{1}{2}, j, k+\frac{1}{2}\right)}^{n}}\right]\right\}  \tag{9.17}\\
& D_{y_{\left(i, j+\frac{1}{2}, k\right)}^{n+1}}^{n+1}=E_{y_{\left(i, j+\frac{1}{2}, k\right)}^{n}}^{n}+\frac{\Delta t}{\Delta}\left\{\left[H_{x_{\left(i, j+\frac{1}{2}, k+\frac{1}{2}\right)}^{n}}-H_{x_{\left(i, j+\frac{1}{2}, k-\frac{1}{2}\right)}^{n}}^{n}\right]+\right. \\
& \left.\left[H_{z_{\left(i-\frac{1}{2}, j+\frac{1}{2}, k\right)}^{n}}-H_{z_{\left(i+\frac{1}{2}, j+\frac{1}{2}, k\right)}^{n}}^{n}\right]\right\}  \tag{9.18}\\
& D_{z_{\left(i, j, k+\frac{1}{2}\right)}^{n+1}}^{n+}=D_{z_{\left(i, j, k+\frac{1}{2}\right)}^{n}}^{n}+\frac{\Delta t}{\Delta}\left\{\left[H_{y_{\left(i+\frac{1}{2}, j, k+\frac{1}{2}\right)}^{n}}-H_{y_{\left(i-\frac{1}{2}, j, k+\frac{1}{2}\right)}^{n}}\right]+\right. \\
& \left.\left[H_{x_{\left(i, j-\frac{1}{2}, k+\frac{1}{2}\right)}^{n}}^{n}-H_{x_{\left(i, j+\frac{1}{2}, k+\frac{1}{2}\right)}^{n}}^{n}\right]\right\} \tag{9.19}
\end{align*}
$$

Once the components of the displacement vector $\vec{D}$ are updated from the previous equations, we proceed to the determination of the $\vec{E}$ components using the relation $\vec{D}=\varepsilon(\omega) \vec{E}$. Replacing $\varepsilon(\omega)$ by its expression given by the Drude model, we get to:

$$
\begin{equation*}
\left(\omega^{2}+i \omega \gamma_{D}\right) \vec{D}=\varepsilon_{0} \varepsilon_{\infty}\left(\omega^{2}+i \omega \gamma_{D}\right) \vec{E}-\varepsilon_{0} \omega_{D}^{2} \vec{E} \tag{9.20}
\end{equation*}
$$

Assuming time dependance of electromagnetic field in $e^{-i \omega t}$, a simple Fourier transform $(\omega \rightarrow$ $t$ ) of this last equation leads to:

$$
\frac{d^{2} \vec{D}}{d t^{2}}+\gamma_{D} \frac{d \vec{D}}{d t}=\varepsilon_{0}\left(\varepsilon_{\infty} \frac{d^{2} \vec{E}}{d t^{2}}+\varepsilon_{\infty} \gamma_{D} \frac{d \vec{E}}{d t}+\omega_{D}^{2} \vec{E}\right)
$$

The partial derivatives of this equations are then replaced by their expressions through the centered finite difference schema. The electric field updated equation in the dispersive media is then obtained:

$$
\begin{equation*}
\xi \vec{E}^{n+1}=-\chi \vec{E}^{n-1}+4 \varepsilon_{\infty} \varepsilon_{0} \vec{E}^{n}+\vec{D}^{n+1}\left[\gamma_{D} \Delta t+2\right]-4 \vec{D}^{n}+\left[-\gamma_{D} \Delta t+2\right] \vec{D}^{n-1} \tag{9.21}
\end{equation*}
$$

with $\xi=\varepsilon_{0}\left[\omega_{D}^{2} \Delta t^{2}+\varepsilon_{\infty} \gamma_{D} \Delta t+2 \varepsilon_{\infty}\right]$ and $\chi=\varepsilon_{0}\left[\omega_{D}^{2} \Delta t^{2}-\gamma_{D} \varepsilon_{\infty} \Delta t+2 \varepsilon_{\infty}\right]$. Due to the dispersion, an additional step of calculation is necessary. It consists of determining the displacement field components for all nodes representing the dispersive media. In addition and as can be seen in equation (9.21), we need to store the $\vec{E}$ and $\vec{D}$ components on two time steps, which has the effect of increasing the memory space to be allocated and the computation time.

### 9.1.4.2 Drude-Lorentz Model

In addition to the conduction electrons, the Drude-Lorentz model takes into account the bound electrons. The interband transition of electrons from filled bands to the conduction band can significantly influence the optical response. In alkali metals, these transitions occur at high frequencies and provide only small corrections to the dielectric function in the optical domain. These metals are well described by the Drude model. On the other side, in noble metals a correction must be made to the dielectric function. It is due to transitions between the bands d and the conduction band s-p. The contribution of bound electrons to the dielectric function can be described by the Lorentz model. To the above Drude dielectric function, a Lorentzian term is added:

$$
\varepsilon_{D L}(\omega)=\varepsilon_{D}(\omega)+\varepsilon_{L}(\omega)
$$

Estimating $\varepsilon_{L}(\omega)$, the bound electrons are described by forced and damped harmonic oscillators. Vial et al. [10] suggested a single oscillator leading to a single Lorentzian additional term to well describe the permittivity of gold in the optical range compared with the classical Drude model. In this case, the relative dielectric function is:

$$
\begin{equation*}
\varepsilon_{D L}(\omega)=\varepsilon_{\infty}-\frac{\omega_{p}^{2}}{\omega^{2}+i \omega \gamma}-\frac{\Delta \varepsilon \cdot \Omega_{L}^{2}}{\left(\omega^{2}-\Omega_{L}^{2}\right)+i \Gamma_{L} \omega} \tag{9.22}
\end{equation*}
$$

where $\Gamma_{L}$ et $\Omega_{L}$ stand for the spectral width and the strength of the Lorentz oscillator respectively. $\Delta \varepsilon$ is a weighting factor.

The FDTD implementation of this model can be done with the Auxilliary Differential Equations (ADE) method previously described above in the case of the Drude model or the so-called Recursive Convolution (RC) method [10]. Because of the additional Lorentzian term, its use requires the introduction of additional intermediate electromagnetic components in the algorithm. Thus, a larger memory space is required compared to the case of the Drude model. In general, many involving multiple oscillators Lorentz terms are needed to accurately model the permittivity of noble metals in the optical range.

### 9.1.4.3 Drude critical points model

The optical properties of some metals, particularly gold, are more difficult to be analytically described in the visible/near-UV region. This comes from much more important role, in the case of gold, played by interband transitions in this region. Some attempts to add Lorentz oscillators to the classical Drude term to account for these transitions rapidly face limitations [11]. In fact, besides the huge simulation time, increasing the number of parameters (mainly non-physical and not well defined) would not provide more insight than quality fit (itself nonphysical) with a polynomial high degree or a simple numerical interpolation of the experimental data.

In order to achieve a reasonable representation of the dielectric function, Etchegoin et al. [12] took inspiration from the parametric critical points model developed for semiconductors [13]. This model is very suitable for the description of optical properties of metals (such as gold) for which the band structure is quite complex. In this approach, the frequency dependence of the optical properties of gold in the visible/near-UV may be well described by an analytical formula with three main contributions that can be expressed as follows:

$$
\begin{equation*}
\varepsilon_{D 2 C P}(\omega)=\varepsilon_{\infty}-\frac{\omega_{D}^{2}}{\omega^{2}+i \omega \gamma_{D}}+\sum_{p=1}^{p=2} G_{p}(\omega) \tag{9.23}
\end{equation*}
$$

with

$$
\begin{equation*}
G_{p}(\omega)=A_{p} \Omega_{p}\left(\frac{e^{i \phi_{p}}}{\Omega_{p}-\omega-i \Gamma_{p}}+\frac{e^{-i \phi_{p}}}{\Omega_{p}+\omega+i \Gamma_{p}}\right) \tag{9.24}
\end{equation*}
$$

The two first terms of equation (9.23) represents the standard contribution of the classical Drude Model. The sum in equation (9.23) is the contribution of the inter-band transitions with the amplitude $A_{p}$, gap energy $\Omega_{p}$, phase $\phi_{p}$ and broadening $\Gamma_{p}$.

In a comparative study of this Drude critical points (CP) model with the so-called L4 model which consists of four Lorentzian terms [14], Vial et al. [15] show the possibility to increase the accuracy of gold and silver permittivity description by using the CP model with fewer parameters to determine and less memory use within the FDTD method.

## Implementation of the CP model in FDTD using ADE technique

As in the previous case of the Drude model, the technique is to calculate the displacement vector components by the FDTD equations $9.17 \mid 9.18$ and 9.19 and determine electrical components using the following relationship:

$$
\begin{equation*}
\vec{D}=\varepsilon_{0} \varepsilon_{D C P} \vec{E} \tag{9.25}
\end{equation*}
$$

In the case of the CP model, $\vec{D}$ can be written as the sum of the electric displacement vectors corresponding to each of the contributions in the dielectric function expression:

$$
\begin{equation*}
\vec{D}=\vec{D}_{D}+\sum_{p=1}^{2} \vec{D}_{C p} \tag{9.26}
\end{equation*}
$$

with

$$
\begin{align*}
\vec{D}_{D} & =\varepsilon_{0}\left[\varepsilon_{\infty}-\frac{\omega_{P}^{2}}{\omega^{2}+\mathrm{i} \gamma \omega}\right] \vec{E}  \tag{9.27.a}\\
\vec{D}_{C p} & =\varepsilon_{0}\left[A_{p} \Omega_{p}\left(\frac{\mathrm{e}^{\mathrm{i} \phi_{p}}}{\Omega_{p}-\omega-\mathrm{i} \Gamma_{p}}+\frac{\mathrm{e}^{-\mathrm{i} \phi_{p}}}{\Omega_{p}+\omega+\mathrm{i} \Gamma_{p}}\right)\right] \vec{E} \tag{9.27.b}
\end{align*}
$$

As before the temporal evolution of the fields in $e^{-i \omega t}$ is considered. By inverse Fourier transform, we obtain:

$$
\begin{align*}
\left(\frac{\partial^{2}}{\partial t^{2}}+\gamma \frac{\partial}{\partial t}\right) \overrightarrow{D_{D}} & =\varepsilon_{0} \varepsilon_{\infty}\left(\frac{\partial^{2}}{\partial t^{2}}+\gamma \frac{\partial}{\partial t}+\frac{\omega_{P}^{2}}{\varepsilon_{\infty}}\right) \vec{E}  \tag{9.28.a}\\
\left(\Omega_{p}^{2}+\Gamma_{p}^{2}+\frac{\partial^{2}}{\partial t^{2}}+2 \Gamma_{p} \frac{\partial}{\partial t}\right) \overrightarrow{D_{C p}} & =2 \varepsilon_{0} A_{p} \Omega_{p}\left(\sqrt{\Gamma_{p}^{2}+\Omega_{p}^{2}} \sin \left(\theta_{p}-\phi_{p}\right)-\sin \phi_{p} \frac{\partial}{\partial t}\right) \vec{E} \tag{9.28.b}
\end{align*}
$$

where: $\theta_{p}=\arctan \left(\frac{\Omega_{p}}{\Gamma_{p}}\right)$
By centered difference discretization of the equation system (9.28) and taking into account the split equation of the displacement vector (9.26), we reach the updated equations system for the electric field vector at each point $(i, j, k)$ of the calculation window:

$$
\begin{align*}
\vec{E}^{n+1}= & \frac{1}{\frac{\chi_{D}}{\alpha_{D}}+\sum_{p=1}^{p=2}\left(\frac{\chi_{p}}{\alpha_{p}}\right)}\left[\vec{D}^{n+1}+\frac{\beta_{D}}{\alpha_{D}} \vec{D}_{D}^{n-1}+\frac{4}{\alpha_{D}} \vec{D}_{D}^{n}-\frac{\delta_{D}}{\alpha_{D}} \vec{E}^{n-1}-\frac{4 \varepsilon_{0} \varepsilon_{\infty}}{\alpha_{D}} \vec{E}^{n}\right. \\
& \left.+\sum_{p=1}^{p=2}\left(\frac{\beta_{p}}{\alpha_{p}} \vec{D}_{C p}^{n-1}-\frac{4}{\alpha_{p}} \vec{D}_{C p}^{n}\right)+\sum_{p=1}^{p=2}\left(\frac{\delta_{p}}{\alpha_{p}}\right) \vec{E}^{n-1}\right]  \tag{9.29.a}\\
\vec{D}_{D}^{n+1}= & \frac{1}{\alpha_{D}}\left[-\beta_{D} \vec{D}_{D}^{n-1}-4 \vec{D}_{D}^{n}+\chi_{D} \vec{E}^{n+1}+\delta_{D} \vec{E}^{n-1}+4 \varepsilon_{0} \varepsilon_{\infty} \vec{E}^{n}\right]  \tag{9.29.b}\\
\vec{D}_{C p}^{n+1}= & \frac{1}{\alpha_{p}}\left[-\beta_{p} \vec{D}_{C p}^{n-1}+4 \vec{D}_{C p}^{n}+\chi_{p} \vec{E}^{n+1}+\delta_{p} \vec{E}^{n-1}\right] \tag{9.29.c}
\end{align*}
$$

with:

$$
\begin{align*}
\alpha_{D} & =-2-\gamma \Delta t  \tag{9.30a}\\
\beta_{D} & =-2+\gamma \Delta t  \tag{9.30b}\\
\chi_{D} & =\varepsilon_{0} \varepsilon_{\infty}\left[-2-\gamma \Delta t-\left(\omega_{p} \Delta t\right)^{2} / \varepsilon_{\infty}\right]  \tag{9.30c}\\
\delta_{D} & =\varepsilon_{0} \varepsilon_{\infty}\left[-2+\gamma \Delta t-\left(\omega_{p} \Delta t\right)^{2} / \varepsilon_{\infty}\right]  \tag{9.30d}\\
\alpha_{p} & =\left[\Omega_{p}^{2}+\Gamma_{p}^{2}\right] \Delta t^{2}+2 \Gamma_{p} \Delta t+2  \tag{9.30e}\\
\beta_{p} & =\left[\Omega_{p}^{2}+\Gamma_{p}^{2}\right] \Delta t^{2}-2 \Gamma_{p} \Delta t+2  \tag{9.30f}\\
\chi_{p} & =2 A_{p} \Omega_{p} \varepsilon_{0}\left[\Delta t^{2} \sqrt{\Omega_{p}^{2}+\Gamma_{p}^{2}} \sin \left(\theta_{p}-\phi_{p}\right)-\Delta t \sin \phi_{p}\right]  \tag{9.30~g}\\
\delta_{p} & =2 A_{p} \Omega_{p} \varepsilon_{0}\left[\Delta t^{2} \sqrt{\Omega_{p}^{2}+\Gamma_{p}^{2}} \sin \left(\theta_{p}-\phi_{p}\right)+\Delta t \sin \phi_{p}\right] \tag{9.30h}
\end{align*}
$$

Let us mention that the displacement vector split into three contributions avoids doing appear derivatives of order higher than 2 in the equations system (9.28). As seen on the equations
system (9.29), taking into consideration the two critical points in the FDTD algorithm does not need to store $\vec{E}$ and $\vec{D}$ components over more than two time steps. However, against the Drude model implementation, additional calculation stages appear in order to determine the two parts of the displacement vector corresponding to the two critical contributions.

### 9.2 Band gap calculation for 2D periodic structures

In this section, we describe how to adapt the FDTD calculation for photonic bandgap structures (PBG) of periodic arrays. The biperiodic structures case is there considered. These 2D structures are photonic crystals $(\mathrm{PhC})$ whose permittivity is periodic in two dimensions ( $x$ and $y$ for example) and remains invariant according to the third one (z). They mainly include three main families that are square, triangular and hexagonal lattices. For this type of structures, we can distinguish two kinds of propagation, in the plane (in-plane, $k_{z}=0$ ) and out of plane (off-plane, nonzero $k_{z}$ ). The system of equations (9.5) becomes easier depending on the type of propagation. To illustrate this, let us assume in what follows that the PhC is periodic along the $x$ and $y$ directions and infinite along $z$ direction.

### 9.2.1 In-plane propagation: TE and TM polarizations

In that case the propagation is done in the plane and the field variation vanishes along the third direction. The system of equations (9.5) is simplified and divided into two independent subsystems giving rise to two polarizations: transverse electric (TE) and transverse magnetic (TM):
$T E$ Polarization

$$
\begin{align*}
\frac{\partial H_{z}}{\partial t} & =\frac{1}{\mu}\left(\frac{\partial E_{x}}{\partial y}-\frac{\partial E_{y}}{\partial x}\right)  \tag{9.31a}\\
\frac{\partial E_{x}}{\partial t} & =\frac{1}{\varepsilon} \frac{\partial H_{z}}{\partial y}  \tag{9.31b}\\
\frac{\partial E_{y}}{\partial t} & =-\frac{1}{\varepsilon} \frac{\partial H_{z}}{\partial x} \tag{9.31c}
\end{align*}
$$

$T M$ Polarization

$$
\begin{align*}
\frac{\partial H_{x}}{\partial t} & =-\frac{1}{\mu} \frac{\partial E_{z}}{\partial y}  \tag{9.32a}\\
\frac{\partial H_{y}}{\partial t} & =\frac{1}{\mu} \frac{\partial E_{z}}{\partial x}  \tag{9.32b}\\
\frac{\partial E_{z}}{\partial t} & =\frac{1}{\varepsilon}\left(\frac{\partial H_{y}}{\partial x}-\frac{\partial H_{x}}{\partial y}\right) \tag{9.32c}
\end{align*}
$$

In case of $T E$ polarization, the electrical components are transverse. They are in the plane of periodicity of the PhC . On the other hand, for the $T M$ polarization, the electric field is perpendicular to the directions of periodicity and the magnetic components are transverse.

Let us note that the two polarizations can be studied by the same system of equations (9.5) without separating it into two sub-systems. But to simplify the calculation codes and gain memory space, it is recommended to study these two polarizations separately.

### 9.2.2 Off-plane propagation

Off-plane propagation is characterized by a nonzero propagation constant $k_{z}$ according to $z$ direction. Diagram dispersion is generally determined for a fixed value of $k_{z}$. Thus, the zderivatives in Maxwell equations become analytical while the electric and magnetic field vectors can be written as follows:

$$
\begin{align*}
\vec{E}(x, y, z, t) & =\overrightarrow{E_{0}}(x, y, t) \exp \left(i k_{z} z\right)  \tag{9.33a}\\
\vec{H}(x, y, z, t) & =\overrightarrow{H_{0}}(x, y, t) \exp \left(i k_{z} z\right) \tag{9.33b}
\end{align*}
$$

The Maxwell's system of equations (9.5) becomes:

$$
\begin{align*}
\frac{\partial H_{x}}{\partial t} & =\frac{1}{\mu}\left(i k_{z} E_{y}-\frac{\partial E_{z}}{\partial y}\right)  \tag{9.34a}\\
\frac{\partial H_{y}}{\partial t} & =\frac{1}{\mu}\left(\frac{\partial E_{z}}{\partial x}-i k_{z} E_{x}\right)  \tag{9.34b}\\
\frac{\partial H_{z}}{\partial t} & =\frac{1}{\mu}\left(\frac{\partial E_{x}}{\partial y}-\frac{\partial E_{y}}{\partial x}\right)  \tag{9.34c}\\
\frac{\partial E_{x}}{\partial t} & =\frac{1}{\varepsilon}\left(\frac{\partial H_{z}}{\partial y}-i k_{z} H_{y}\right)  \tag{9.34d}\\
\frac{\partial E_{y}}{\partial t} & =\frac{1}{\varepsilon}\left(i k_{z} H_{x}-\frac{\partial H_{z}}{\partial x}\right)  \tag{9.34e}\\
\frac{\partial E_{z}}{\partial t} & =\frac{1}{\varepsilon}\left(\frac{\partial H_{y}}{\partial x}-\frac{\partial H_{x}}{\partial y}\right) \tag{9.34f}
\end{align*}
$$

In this case, it is no longer possible to separate the system into two subsystems as before. The $T E$ and $T M$ cases are therefore mixed together and can not be treated separately. However, we can note that the calculation code is simplified since the $z$ derivatives are analytically evaluated so there is no discretization along the $z$ direction. A 2 D algorithm is still needed.

### 9.2.3 Periodic boundary conditions

As the CPU time and space memory is limited, the FDTD calculation window must also be finite. Because of symmetry, only one unit cell is considered. To reproduce the crystal at the truncated domain boundaries, the Floquet-Bloch periodicity conditions [9] are applied to the electric and magnetic components. Despite the fact that these periodicity conditions are general and can be applied to any periodic structure, their expressions depend on the Bravais lattice. Consequently, we will consider the two most used Bravais lattices i.e. the rectangular and the triangular ones.

## Rectangular cell

Let us consider a PhC made of cylinders (refractive index $n_{1}$ ) immersed in a medium of refractive index $n_{2}$. $a$ and $b$ are the lattice constants in the $x$ and $y$ directions respectively (see figure 9.6). The FDTD window calculation is shown in figure 9.6-b.


Figure 9.6: Rectangular structure and FDTD window calculation

The Floquet-Bloch conditions are applied to the electric and magnetic components as follows:

$$
\begin{align*}
& \vec{E}(x=0, y, t)=\vec{E}(x=a, y, t) \exp \left(-i k_{x} \cdot a\right)  \tag{9.35a}\\
& \vec{E}(x, y=0, t)=\vec{E}(x, y=b, t) \exp \left(-i k_{y} \cdot b\right)  \tag{9.35b}\\
& \vec{H}(x=a, y, t)=\vec{H}(x=0, y, t) \exp \left(i k_{x} \cdot a\right)  \tag{9.35c}\\
& \vec{H}(x, y=b, t)=\vec{H}(x, y=0, t) \exp \left(i k_{y} \cdot b\right) \tag{9.35d}
\end{align*}
$$

## Triangular cell

Similarly to the rectangular cell, the calculation FDTD window is limited to a single unit cell. To model the triangular photonic structure (see figure 9.7-a), three choices of the FDTD window are possible. The first one is to take a non-orthogonal unit cell (cell 1 in figure 9.7-a) and implement the periodic boundary conditions in a Non orthogonal-FDTD algorithm [16, 17] for which the classical FDTD developed in an orthogonal coordinate system is not suitable. To bypass this constraint and remaining in the conventional FDTD with orthogonal coordinates, the second rectangular cell (celle 2 in figure 9.7 a) can be used to derive the periodic conditions. Nevertheless, this cell contains two patterns. This means that the rectangular periodic conditions lead to a less-description of all the possible solutions. Consequently, an aliasing effect will appear in the dispersion diagram.

In order to get gain in computational time and space and prevent this band folding while remaining with the orthogonal FDTD algorithm, a rectangular cell can be defined with only one pattern (cell 3 in figure 9.7 a). Within this FDTD calculation cell (9.7-b), the periodic conditions above are therefore replaced by:
-along the $x$ direction :

$$
\begin{align*}
& \vec{E}(x=0, y, z, t)=\vec{E}(x=a, y, z, t) \exp \left(-i k_{x} \cdot a\right)  \tag{9.36a}\\
& \vec{H}(x=a, y, z, t)=\vec{H}(x=0, y, z, t) \exp \left(i k_{x} \cdot a\right) \tag{9.36b}
\end{align*}
$$



Figure 9.7: Triangular structure and FDTD calculation window


Figure 9.8: Brillouin zone

- along the $y$ direction with $x \geq 0$ and $x \leq a / 2$

$$
\begin{align*}
& \vec{E}(x, y=0, z, t)=\vec{E}\left(x+\frac{a}{2}, y=b, z, t\right) \exp \left(i\left(-k_{y} \cdot b-k_{x} \cdot a / 2\right)\right)  \tag{9.37a}\\
& \vec{H}(x, y=b, z, t)=\vec{H}\left(x+\frac{a}{2}, y=0, z, t\right) \exp \left(i\left(k_{y} \cdot b-k_{x} \cdot a / 2\right)\right) \tag{9.37b}
\end{align*}
$$

- along the $y$ direction with $x>a / 2$ and $x \leq a$

$$
\begin{align*}
& \vec{E}(x, y=0, z, t)=\vec{E}\left(x-\frac{a}{2}, y=b, z, t\right) \exp \left(i\left(-k_{y} \cdot b+k_{x} \cdot a / 2\right)\right)  \tag{9.38a}\\
& \vec{H}(x, y=b, z, t)=\vec{H}\left(x-\frac{a}{2}, y=0, z, t\right) \exp \left(i\left(k_{y} \cdot b+k_{x} \cdot a / 2\right)\right) \tag{9.38b}
\end{align*}
$$

By the way, the dispersion diagram of a triangular or honeycomb Bravais lattices can be calculated without modifying the orthogonal Cartesian Yee schema.

### 9.2.4 Some examples of band gap calculation

To obtain a photonic band diagram, several FDTD calculations are necessary done by varying the $\vec{k}$ wavevector that must scan the irreducible Brillouin zone (figure 9.8). $\Gamma X, X M$ and $M \Gamma$ highest symmetry directions are then discretized.

For this band gap calculation, the N-Order FDTD algorithm is used [18, 19]. This basis of this algorithm is quite simple: a signal is injected to excite all possible frequencies of the structure. This signal is introduced in accordance to the Maxwell-Gauss law $(\operatorname{div}(\vec{E})=0)$ and given as follows:

$$
\begin{equation*}
\vec{E}=\sum_{G}(\vec{v} \wedge(\vec{k}+\vec{G}) \exp (i(\vec{k}+\vec{G}) \cdot \vec{r}) \tag{9.39}
\end{equation*}
$$



Figure 9.9: Normalized electromagnetic energy density at $\Gamma$ point for triangular structure of air holes (of radius $r=0.25 \mathrm{a}$ ) into lithium niobate. TM Polarization.
$\vec{v}$ is a random vector, $\vec{k}$ and $\vec{G}$ denote the wavevector and the reciprocal lattice vector respectively.

After injecting this last initial signal, and for a given $\vec{k}$, the FDTD simulation is run and electromagnetic energy density time-evolution is calculated as a function of the frequency. This later is calculated through:

$$
\begin{equation*}
W=\frac{1}{4}\left(\varepsilon_{0} \varepsilon|E|^{2}+\mu|H|^{2}\right) \tag{9.40}
\end{equation*}
$$

Only eigenmodes of the structure persist and evanescent ones gradually disappear. After a large number of time iterations (typically $10^{5}$ ) a permanent regime is then reached and the electromagnetic energy density spectrum exhibits several peaks corresponding to the eigenfrequencies of the studied structure. An example of eigenfrequencies calculation for a triangular structure in the $\Gamma$ point is shown in figure 9.9. The structure is made of air holes $\left(n_{1}=1\right)$ into a dielectric medium which is lithium niobate $\left(\mathrm{LiNbO}_{3}\right)$ with refractive index $n_{2}=2.1421$. The radius of the holes is $r=0.25$ a which corresponds to a filling factor of $0.2267 \%$. The FDTD grid, one PhC period, contains $60 \times 52$ spatial grids. To satisfy the stability criterion and avoid numerical dispersion, the time step is taken as $\triangle t=a /(120 \cdot c)$.

To get the complete photonic band structure, it is necessary to scan the $k$ values over all the contour of the irreducible Brillouin zone $(\Gamma X M)$. Figure 9.10 shows the photonic band diagram calculated for both $T E$ and $T M$ polarizations for a structure parameters similar to those used above in the case of figure 9.9

We can note the emergence of a photonic bandgap for $\omega a / 2 \pi c$ between 0.32 and 0.35 in the case of the TE polarization (figure 9.10-a). This band does not exist in the case of the TM polarization (figure 9.10-b) so it is called "partial".

Note here that, for a dispersive material, the calculation of the electromagnetic energy density is no more given by equation 9.40 that is only valid for dielectrics (no dispersion). In the case of metallic dispersive material, the electromagnetic energy density is given by (no magnetic dispersion):

$$
\begin{equation*}
W=\frac{1}{4}\left(\frac{\partial\left(\omega \varepsilon_{0} \varepsilon\right)}{\partial \omega}|E|^{2}+\mu|H|^{2}\right) \tag{9.41}
\end{equation*}
$$



Figure 9.10: Photonic band diagram for triangular structure of air holes (of radius $r=0.25$ a) into lithium niobate.


Figure 9.11: In-plane photonic band diagram for annular aperture arrays engraved into silver (TE polarization).

The calculation of the energy density depends then on the dispersion model introduced in the FDTD. Accordingly, an analytic expression of $W$ is obtained through the calculation of the frequency derivative in equation 9.41 . Its numerical value is then performed by the determination of the spectral responses of both the two electric and magnetic fields that are determined by the FDTD code.

Another example of band diagram, corresponding to a metallic structure made of annular aperture arrays (AAA) engraved into silver layer and arranged in a square lattice, is shown in figure 9.11. The AAA structure has been proposed by F. Baida and D. Van Labeke [20] for Enhanced Optical Transmission (EOT) applications. It was showed that transmission through AAA sub-wavelength structure could reach $90 \%$ in the visible range [21]. This EOT is due to the excitation and the propagation of a guided mode inside each aperture. The main transmission peak corresponds to the excitation of the $\mathrm{TE}_{11}$ mode at its cutoff wavelength [19]. This later only depends on the value of the inner and the outer radii. For $r_{i}=50 \mathrm{~nm}$ and $r_{e}=75 \mathrm{~nm}$ and a lattice constant of $a=160 \mathrm{~nm}$ one gets the band diagram of figure 9.11 .

In case of the figure 9.11 , corresponding to the $T E$ polarization, we note the presence of two photonic bandgaps. the first is ranging from zero frequency (infinite wavelength) to the frequency value of $0.1835(c / a)(\lambda=872 \mathrm{~nm})$. The second gap is in the visible range between 492 nm and 630 nm . Note that these bandgaps are "total" since the corresponding eigen


Figure 9.12: Off-plane photonic band diagram for annular aperture arrays made in silver.


Figure 9.13: Dispersion curves at $\Gamma$ point for the coaxial structure made in silver (lattice constante a $=160 \mathrm{~nm}$, inner radius $r_{i}=50 \mathrm{~nm}$ and outer radius $r_{e}=75 \mathrm{~nm}$; silver dispersion is modeled by a Drude model).
frequencies of $T M$ polarization are located above $0.45 \times c / a$.
The figure 9.12 illustrates photonic band diagrams for the same considered AAA structure but in the case of off-plane propagation with two different values of $k_{z}$. There is occurrence of an additional photonic band relative to the in-plane case. This is due to the transverse electromagnetic (TEM) mode excited now at a nonzero frequency (far from the cutoff). For $k_{z}=\pi /(3 a)$, the bandgaps are located in the ranges $] 1873 \mathrm{~nm}, \infty[] 723 \mathrm{~nm},, 1668 \mathrm{~nm}[$ and $] 458 \mathrm{~nm}, 575 \mathrm{~nm}[$. These bandgaps become] $653 \mathrm{~nm}, \infty[] 512 \mathrm{~nm},, 574 \mathrm{~nm}[$ and $] 378 \mathrm{~nm}, 431 \mathrm{~nm}\left[\right.$ when $k_{z}=\pi / a$. According to the theory, this band gap shift is due to the fact that the eigenfrequencies of guided modes increase with $k_{z}$.

Figure 9.13, showing the dispersion curves (at $\Gamma$ point) of the guided modes depending on $k_{z}$, clearly confirms the TEM nature of the additional mode excited in the off-plan case. This mode band starts from zero frequency, and therefore has no cutoff frequency. An EOT based on the excitation of this peculiar mode can be obtained under two conditions: an oblique incidence with TM polarization [22]. The last section of this chapter is devoted to the study of EOT obtained through the excitation of this peculiar mode.

An example of time evolution of the electromagnetic energy density is given on figure 9.14 The considered structure is an array of coaxial waveguides made in perfectly electric conductor (PEC). All the geometrical parameters are given in the caption in addition to the FDTD simulation ones. One notes that the main peak corresponds to the $\mathrm{TE}_{21}$ guided mode


Figure 9.14: Time evolution of the electromagnetic energy density spectrum. The modeled structure is an array of coaxial waveguides made in perfectly electric conductor (PEC) and arranged in square lattice. The inner and outer radii is $r_{i}=100 \mathrm{~nm}$ and $r_{e}=140 \mathrm{~nm}$ respectively. The period of the grating is $a=300 \mathrm{~nm}$ but the obtained results are independent on this value because there is non coupling between tow adjacent waveguides. The FDTD simulations are done with a uniform spatial mesh of $\Delta x=\Delta y=\frac{a}{400}$ and the temporal step was fixed to $\Delta t=\frac{\Delta x}{4 c}$ where $c$ is the light velocity in vacuum.
that has a cutoff wavelength of $\lambda_{T E_{21}}^{c}=\frac{\pi\left(r_{i}+r_{e}\right)}{2}$.

### 9.3 Scattering calculation for 3D biperiodic nanostructures

In this section, we will focus on the FDTD modeling of dielectric and metallic bi-periodic structures. For normal incidence, the FDTD method, based on the classical Yee's scheme, is a powerful tool that can simply model such periodic structures [24, 25, 26]. In fact, in this simple case, the Floquet-Bloch periodic boundary conditions (PBC) can be easily applied without any change because these conditions are independent of the frequency. However, at oblique incidence, applying PBC implicitly involves a frequency term that must be integrated into the FDTD algorithm that operates in the temporel domain. Thus, in order to adapt FDTD to oblique incidence case, Veysoglu [27] introduced the field transformation method applied to $\vec{E}$ and $\vec{H}$ toward new $\vec{P}$ and $\vec{Q}$ fields. By the way, the PBC conditions become similar to the ones of normal incidence case nevertheless the immediate consequence of this transformation is the need to modify the Yee's scheme. Several techniques of implementation are then proposed including that of Split-Field Method (SFM) [28].

In the following, we present the reformulation of the FDTD method, based on this SFM technique to adapt it to the case of any incidence. Maxwell's equations are modified and expressed with $\vec{P}$ and $\vec{Q}$ variables. They are then discretized using SFM technique. To avoid reflections at the edges of the computational window, the equations in the Berenger's PML medium are also modified and expressed in the new domain within the SFM technique. In addition, the dispersion models mentioned above (Drude, Drude-Lorentz and Drude Critical point models) are also described by modifying and adapting them to the SFM technique.


Figure 9.15: Sketch of the biperiodic structure illuminated by plane wave propagating along the $\vec{k}$ vector defined by its Euler angles $\theta$ and $\phi$.

### 9.3.1 Position of the problem: New $\vec{P}-\vec{Q}$ variables

Let us consider a bi-periodic structure, finished along the third direction and illuminated by a plane wave propagating at oblique incidence (see figure 9.15).

According to the notations of figure 9.15, the electric and magnetic fields of the incident plane wave can be written as:

$$
\begin{align*}
& \vec{E}_{i}=\vec{E}_{0 i} e^{i\left[k_{x} \cdot x+k_{y} \cdot y+k_{z} \cdot z+\omega \cdot t\right]}  \tag{9.42.a}\\
& \vec{H}_{i}=\vec{H}_{0 i} e^{i\left[k_{x} \cdot x+k_{y} \cdot y+k_{z} \cdot z+\omega \cdot t\right]} \tag{9.42.b}
\end{align*}
$$

where:

$$
\begin{align*}
& k_{x}=\frac{\omega}{v} \sin \theta \cos \varphi  \tag{9.43}\\
& k_{y}=\frac{\omega}{v} \sin \theta \sin \varphi  \tag{9.44}\\
& k_{z}=\frac{\omega}{v} \cos \theta \tag{9.45}
\end{align*}
$$

For the periodic object, a single pattern (one period) is then considered for the FDTD calculation (see figure 9.6). The periodic conditions are then written so that the fields on one side of the calculation window are expressed versus the fields on the opposite side through the FloquetBloch conditions. For $x$ (lattice constant $a$ ) and $y$ (lattice constant $b$ ) periodic structures, these conditions are expressed as follows:

$$
\begin{align*}
\vec{E}(x, y, z, t) & =\vec{E}(x+a, y, z, t) \cdot e^{-i k_{x} \cdot a}  \tag{9.46.a}\\
\vec{E}(x, y, z, t) & =\vec{E}(x, y+b, z, t) \cdot e^{-i k_{y} \cdot b}  \tag{9.46.b}\\
\vec{H}(x+a, y, z, t) & =\vec{H}(x, y, z, t) \cdot e^{i k_{x} \cdot a}  \tag{9.46.c}\\
\vec{H}(x, y+b, z, t) & =\vec{H}(x, y, z, t) \cdot e^{i k_{y} \cdot b} \tag{9.46.d}
\end{align*}
$$

As the FDTD method operates in the temporal domain and $k_{x}$ and $k_{y}$ components explicitly depend of $\omega$, the direct application of these periodic conditions is prohibited. Consequently, a change of variables is performed so that $\vec{E}$ and $\vec{H}$ components are replaced by two new components $\vec{P}$ and $\vec{Q}$ respectively in order to eliminate the $k_{x}$ and $k_{y}$ dependence in the PBC. These new fields are defined as follows:

$$
\begin{align*}
& \vec{P}=\vec{E} \cdot e^{-i k_{x} x} \cdot e^{-i k_{y} y}  \tag{9.47.a}\\
& \vec{Q}=\vec{H} \cdot e^{-i k_{x} x} \cdot e^{-i k_{y} y} \tag{9.47.b}
\end{align*}
$$

Therefore, the new periodic conditions can be applied similarly to the case of normal incidence through the relations:

$$
\begin{align*}
\vec{P}(x, y, z, t) & =\vec{P}(x+a, y, z, t)  \tag{9.48.a}\\
\vec{Q}(x+a, y, z, t) & =\vec{Q}(x, y, z, t)  \tag{9.48.b}\\
\vec{P}(x, y, z, t) & =\vec{P}(x, y+b, z, t)  \tag{9.48.c}\\
\vec{Q}(x, y+b, z, t) & =\vec{Q}(x, y, z, t) \tag{9.48.d}
\end{align*}
$$

Replacing $\vec{E}$ and $\vec{H}$ by their expressions in terms of $\vec{P}$ and $\vec{Q}$ through equations 9.47 in Maxwell's equations system 9.5 leads to:

$$
\begin{align*}
\frac{\partial Q_{x}}{\partial t} & =\frac{1}{\mu_{0}}\left[\frac{\partial P_{y}}{\partial z}-\frac{\partial P_{z}}{\partial y}-i k_{y} P_{z}\right]  \tag{9.49.a}\\
\frac{\partial Q_{y}}{\partial t} & =\frac{1}{\mu_{0}}\left[\frac{\partial P_{z}}{\partial x}+i k_{x} P_{z}-\frac{\partial P_{x}}{\partial z}\right]  \tag{9.49.b}\\
\frac{\partial Q_{z}}{\partial t} & =\frac{1}{\mu_{0}}\left[\frac{\partial P_{x}}{\partial y}+i k_{y} P_{x}-\frac{\partial P_{y}}{\partial x}-i k_{x} P_{y}\right]  \tag{9.49.c}\\
\frac{\partial P_{x}}{\partial t} & =\frac{1}{\varepsilon}\left[\frac{\partial Q_{z}}{\partial y}+i k_{y} Q_{z}-\frac{\partial Q_{y}}{\partial z}\right]  \tag{9.49.d}\\
\frac{\partial P_{y}}{\partial t} & =\frac{1}{\varepsilon}\left[\frac{\partial Q_{x}}{\partial z}-\frac{\partial Q_{z}}{\partial x}-i k_{x} Q_{z}\right]  \tag{9.49.e}\\
\frac{\partial P_{z}}{\partial t} & =\frac{1}{\varepsilon}\left[\frac{\partial Q_{y}}{\partial x}+i k_{x} Q_{y}-\frac{\partial Q_{x}}{\partial y}-i k_{y} Q_{x}\right] \tag{9.49.f}
\end{align*}
$$

We can notice that for a wave propagating at normal incidence, the system (9.49) above is equivalent to the conventional Maxwell' equations expressed in $\vec{E}-\vec{H}$. Nonetheless, in the oblique case, additional terms appear in the second right members of equations (9.49) and they explicitly depend on $k_{x}$ and $k_{y}$ i.e. on the frequency $\omega$. Even if these terms are equivalent to time derivatives, the direct implementation of the FDTD in this case is impossible. Many implementation techniques have been proposed [29, 30, 31, 28, 32, 33] to overcome this problem. One of them is the Split Field Method [32, 28] which will be described below.

### 9.3.2 Split Field Method

SFM technique is based on the split of $\vec{P}$ and $\vec{Q}$ field components. To illustrate the method, let us take for example the split of the $Q_{x}$ component occurring in equation 9.49.a). By reducing the frequency additional term on the left hand, this equation can be written as:

$$
\begin{equation*}
\frac{\partial Q_{x}}{\partial t}+i \omega \frac{k_{y}}{\mu \omega} P_{z}=\frac{1}{\mu}\left[\frac{\partial P_{y}}{\partial z}-\frac{\partial P_{z}}{\partial y}\right] \tag{9.50a}
\end{equation*}
$$

According to (9.42.a) and 9.47.a), equation 9.50a) becomes:

$$
\begin{equation*}
\frac{\partial}{\partial t}\left[Q_{x}+\frac{k_{y}}{\mu \omega} P_{z}\right]=\frac{1}{\mu}\left[\frac{\partial P_{y}}{\partial z}-\frac{\partial P_{z}}{\partial y}\right] \tag{9.51a}
\end{equation*}
$$

This leads to a new component $Q_{x a}=Q_{x}+\frac{k_{y}}{\mu \omega} P_{z}$ which satisfies Maxwell's equation as for normal incidence. Similarly, the split of all the others components in the $\vec{P}-\vec{Q}$ domain gives:

$$
\begin{align*}
Q_{x a} & =Q_{x}+\frac{k_{y}}{\mu \omega} P_{z}  \tag{9.52.a}\\
Q_{y a} & =Q_{y}-\frac{k_{x}}{\mu \omega} P_{z}  \tag{9.52.b}\\
Q_{z a} & =Q_{z}-\frac{k_{y}}{\mu \omega} P_{x}+\frac{k_{x}}{\mu \omega} P_{y}  \tag{9.52.c}\\
P_{x a} & =P_{x}-\frac{k_{y}}{\varepsilon \omega} Q_{z}  \tag{9.52.d}\\
P_{y a} & =P_{y}+\frac{k_{x}}{\varepsilon \omega} Q_{z}  \tag{9.52.e}\\
P_{z a} & =P_{z}-\frac{k_{x}}{\varepsilon \omega} Q_{y}+\frac{k_{y}}{\varepsilon \omega} Q_{x} \tag{9.52.f}
\end{align*}
$$

The six components thereby obtained satisfy the following equations that can be discretized according to the classical Yee's scheme:

$$
\begin{align*}
\frac{\partial Q_{x a}}{\partial t} & =\frac{1}{\mu}\left[\frac{\partial P_{y}}{\partial z}-\frac{\partial P_{z}}{\partial y}\right]  \tag{9.53.a}\\
\frac{\partial Q_{y a}}{\partial t} & =\frac{1}{\mu}\left[\frac{\partial P_{z}}{\partial x}-\frac{\partial P_{x}}{\partial z}\right]  \tag{9.53.b}\\
\frac{\partial Q_{z a}}{\partial t} & =\frac{1}{\mu}\left[\frac{\partial P_{x}}{\partial y}-\frac{\partial P_{y}}{\partial x}\right]  \tag{9.53.c}\\
\frac{\partial P_{x a}}{\partial t} & =\frac{1}{\varepsilon}\left[\frac{\partial Q_{z}}{\partial y}-\frac{\partial Q_{y}}{\partial z}\right]  \tag{9.53.d}\\
\frac{\partial P_{y a}}{\partial t} & =\frac{1}{\varepsilon}\left[\frac{\partial Q_{x}}{\partial z}-\frac{\partial Q_{z}}{\partial x}\right]  \tag{9.53.e}\\
\frac{\partial P_{z a}}{\partial t} & =\frac{1}{\varepsilon}\left[\frac{\partial Q_{y}}{\partial x}-\frac{\partial Q_{x}}{\partial y}\right] \tag{9.53.f}
\end{align*}
$$

Once the updated components of $\vec{P}_{a}$ and $\vec{Q}_{a}$ completed, the second step of the algorithm is to calculate $\vec{P}$ and $\vec{Q}$ components through the system of equations (9.52) that gives after simple algebra the system below:

$$
\begin{align*}
Q_{z} & =\frac{1}{1-\frac{k_{x}^{2}+k_{y}^{2}}{\varepsilon \mu \omega^{2}}}\left[Q_{z a}+\frac{k_{y}}{\mu \omega} P_{x a}-\frac{k_{x}}{\mu \omega} P_{y a}\right]  \tag{9.54.a}\\
P_{z} & =\frac{1}{1-\frac{k_{x}^{2}+k_{y}^{2}}{\varepsilon \mu \omega^{2}}}\left[P_{z a}+\frac{k_{x}}{\varepsilon \omega} Q_{y a}-\frac{k_{y}}{\varepsilon \omega} Q_{x a}\right]  \tag{9.54.b}\\
Q_{x} & =Q_{x a}-\frac{k_{y}}{\mu \omega} P_{z}  \tag{9.54.c}\\
Q_{y} & =Q_{y a}+\frac{k_{x}}{\mu \omega} P_{z}  \tag{9.54.d}\\
P_{x} & =P_{x a}+\frac{k_{y}}{\varepsilon \omega} Q_{z}  \tag{9.54.e}\\
P_{y} & =P_{y a}-\frac{k_{x}}{\varepsilon \omega} Q_{z} \tag{9.54.f}
\end{align*}
$$

This system (9.54) needs to calculate $\vec{P}$ and $\vec{Q}$ components at the same time iteration as $\vec{P}_{a}$ and $\vec{Q}_{a}$ components. This is in contradiction with the traditional Yee's scheme. Consequently, the new $(\vec{P}, \vec{Q})$ and $\left(\vec{P}_{a}, \vec{Q}_{a}\right)$ fields will be calculated at time $n \Delta t$ and time $\left(n+\frac{1}{2}\right) \Delta t$ in order to reach a stable numerical schema. To this end, each component is calculated twice in one time iteration by introducing other intermediate components in the calculation program (see reference [34]).

## Stability criterion

As the transition to the new $\vec{P}-\vec{Q}$ domain, the stability criterion is also modified. Based on the calculation of Kao [29, 30] and in the case of 3D uniform meshing, this later is expressed as:

$$
\begin{align*}
& \frac{\Delta}{\Delta t} \geq \frac{v_{i}}{v_{i}^{2} \mu \varepsilon-\sin ^{2}(\theta)}\{|\sin (\theta) \cdot \cos (\varphi)|+|\sin (\theta) \cdot \sin (\varphi)|+ \\
& \left.\sqrt{3 v_{i}^{2} \mu \varepsilon-2 \cdot \sin ^{2}(\theta)(1-|\sin (\varphi) \cdot \cos (\varphi)|)}\right\} \tag{9.55}
\end{align*}
$$

where $v_{i}$ is the phase velocity of the incident wave and $\varepsilon$ and $\mu$ are chosen to be the characteristics of the less dense medium in the computational domain.

Let us note here that the time step decreases with the incidence angle $\theta$ and hence the computational time becomes very long for large incidence angles. Nonetheless, the computational time is relatively acceptable up to an incidence angle of $80^{\circ}$.

### 9.3.3 Absorbing boundary conditions: PML

The implementation of absorbing boundary conditions in the oblique case requires to make a change of variables on the fields components in the PML medium similarly to the changes made in the main computational grid [34]. For $x$ and $y$ periodic structure, only PML is needed in the third direction $(O z)$. In this case, the new fields components are defined as follows:

$$
\begin{align*}
P_{v \mu} & =E_{v \mu} \cdot e^{-i k_{x} x} \cdot e^{-i k_{y} y}  \tag{9.56.a}\\
Q_{v \mu} & =H_{v \mu} \cdot e^{-i k_{x} x} \cdot e^{-i k_{y y} y}  \tag{9.56.b}\\
P_{z} & =E_{z} \cdot e^{-i k_{x} x} \cdot e^{-i k_{y} y} \tag{9.56.c}
\end{align*}
$$

where $v$ represents $x$ or $y$ and $\mu$ denotes $x, y$ or $z . E_{v \mu}$ and $Q_{v \mu}$ are the field components in the classical PML shell corresponding to the components of the two fictitious waves resulting from the split of the plane wave inside the PML (see section 1 of this chapter). For details of implementing these PML at oblique incidence, the reader can refer to [34].

### 9.3.4 SFM-FDTD in dispersive media

For oblique incidence, and according to the two systems of equations (9.53) and (9.54), the components that require particular treatment in the dispersive medium are: $P_{x a}, P_{y a}, P_{z a}, Q_{z}, P_{z}$, $P_{x}$ and $P_{y}$. Direct calculation of these components by equations (9.53) and (9.54) involves the permittivity term which is frequency-dependent. In this section, we only show how to take into account the media dispersion in FDTD oblique incidence in the case of the Drude critical points model [35]. The implementation details of the other of dispersion models by SFM-FDTD are given in [36] for Debye model, and in [37] for both Drude and Drude-Lorentz models.

Let us quote that equations (9.53) for the calculation of $P_{x a}, P_{y a}$ and $P_{z a}$ are similar to traditional Maxwell's equations. Accordingly, the calculation of these components in the dispersive medium will not require any further treatment compared to the normal incidence case. Contrarily, equations (9.54) for the $Q_{z}, P_{z}, P_{x}$ and $P_{y}$ need a different way to be processed.
$P_{x a}, P_{y a}$ and $P_{z a}$ implementation: These three components are calculated in a similar way. So let us take as an example only the $P_{x a}$ calculation. By analogy with the normal incidence case (equation 9.25), we introduced a new component $L_{x a}$ (equivalent to the $D_{x}$ component in the classical case) defined as:

$$
\begin{equation*}
L_{x a}=\varepsilon_{0} \cdot \varepsilon_{D C P} \cdot P_{x a} \tag{9.57}
\end{equation*}
$$

Equation (9.53.d) is therefore wrote as:

$$
\begin{equation*}
\frac{\partial L_{x a}}{\partial t}=\left[\frac{\partial Q_{z}}{\partial y}-\frac{\partial Q_{y}}{\partial z}\right] \tag{9.58}
\end{equation*}
$$

The discretization of this last equation allows us to calculate the $L_{x a}$ variable as follows:
$L_{x a_{\left(i+\frac{1}{2}, j, k\right)}^{n+1}}^{n+1}=L_{x a_{\left(i+\frac{1}{2}, j, k\right)}^{n}}^{n}+\frac{\Delta t}{\Delta y}\left[Q_{z_{\left(i+\frac{1}{2}, j+\frac{1}{2}, k\right)}^{n}}-Q_{z_{\left(i+\frac{1}{2}, j-\frac{1}{2}, k\right)}^{n}}^{n}\right]+\frac{\Delta t}{\Delta z}\left[Q_{y_{\left(i+\frac{1}{2}, j, k-\frac{1}{2}\right)}^{n}}-Q_{y_{\left(i+\frac{1}{2}, j, k+\frac{1}{2}\right)}^{n}}\right]$
Analogically to equations (9.26), (9.27.a) and 9.27.b), $L_{x a}$ can be expressed as follows:

$$
\begin{equation*}
L_{x a}=L_{x a_{D}}+\sum_{p=1}^{p=2} L_{x a_{C p}} \tag{9.60}
\end{equation*}
$$

with:

$$
\begin{align*}
L_{x a_{D}} & =\varepsilon_{0}\left[\varepsilon_{\infty}-\frac{\omega_{P}^{2}}{\omega^{2}+\mathrm{i} \gamma \omega}\right] P_{x a}  \tag{9.61.a}\\
L_{x a_{p}} & =\varepsilon_{0}\left[A_{p} \Omega_{p}\left(\frac{\mathrm{e}^{\mathrm{i} \phi_{p}}}{\Omega_{p}-\omega-\mathrm{i} \Gamma_{p}}+\frac{\mathrm{e}^{-\mathrm{i} \phi_{p}}}{\Omega_{p}+\omega+\mathrm{i} \Gamma_{p}}\right)\right] P_{x a} \tag{9.61.b}
\end{align*}
$$

As before, after the inverse Fourier transforms and finite centred differences discretization of different partial derivatives, we reach the updated equations for the component $P_{x a}$ :

$$
\begin{align*}
P_{x a}^{n+1}= & \frac{1}{\frac{\chi_{D}}{\alpha_{D}}+\sum_{p=1}^{p=2}\left(\frac{\chi_{p}}{\alpha_{p}}\right)}\left[L_{x a}^{n+1}+\frac{\beta_{D}}{\alpha_{D}} L_{x a_{D}}^{n-1}+\frac{4}{\alpha_{D}} L_{x a_{D}}^{n}-\frac{\delta_{D}}{\alpha_{D}} P_{x a}^{n-1}-\frac{4 \varepsilon_{0} \varepsilon_{\infty}}{\alpha_{D}} P_{x a}^{n}\right. \\
& \left.+\sum_{p=1}^{p=2}\left(\frac{\beta_{p}}{\alpha_{p}} L_{x a_{C p}}^{n-1}-\frac{4}{\alpha_{p}} L_{x a_{C p}}^{n}\right)+\sum_{p=1}^{p=2}\left(\frac{\delta_{p}}{\alpha_{p}}\right) P_{x a}^{n-1}\right]  \tag{9.62.a}\\
L_{x a_{D}}^{n+1}= & \frac{1}{\alpha_{D}}\left[-\beta_{D} L_{x a_{D}}^{n-1}-4 L_{x a_{D}}^{n}+\chi_{D} P_{x a}^{n+1}+\delta_{D} P_{x a}^{n-1}+4 \varepsilon_{0} \varepsilon_{\infty} P_{x a}^{n}\right]  \tag{9.62.b}\\
L_{x a_{C p}}^{n+1}= & \frac{1}{\alpha_{p}}\left[-\beta_{p} L_{x a_{C p}}^{n-1}+4 L_{x a_{C_{p}}}^{n}+\chi_{p} P_{x a}^{n+1}+\delta_{p} P_{x a}^{n-1}\right] \tag{9.62.c}
\end{align*}
$$

$Q_{z}, P_{z}, P_{x}$ and $P_{y}$ implementation: The calculation of the remaining components $Q_{z}, P_{z}$, $P_{x}$ and $P_{y}$ needs the introduction of other variables involving other equations. We consider as an
example the $P_{z}$ component for which implementation equations are detailed. Equation (9.54.b) involves the following one:

$$
\begin{equation*}
\varepsilon \cdot M_{z}=\frac{k_{x}}{\omega} Q_{y a}-\frac{k_{y}}{\omega} Q_{x a}+\frac{k_{x}^{2}+k_{y}^{2}}{\mu \omega^{2}} P_{z} \tag{9.63}
\end{equation*}
$$

with:

$$
\begin{equation*}
M_{z}=P_{z}-P_{z a} \tag{9.64}
\end{equation*}
$$

By setting:

$$
\begin{equation*}
T_{z}=\frac{k_{x}}{\omega} Q_{y a}-\frac{k_{y}}{\omega} Q_{x a}+\frac{k_{x}^{2}+k_{y}^{2}}{\mu \omega^{2}} P_{z} \tag{9.65}
\end{equation*}
$$

equation (9.63) becomes:

$$
\begin{equation*}
T_{z}=\varepsilon \cdot M_{z}=\varepsilon_{0} \varepsilon_{D C P} M_{z} \tag{9.66}
\end{equation*}
$$

As considered above, the $T_{z}$ component can be expressed as:

$$
\begin{equation*}
T_{z}=T_{z D}+\sum_{p=1}^{p=2} T_{z C_{p}} \tag{9.67}
\end{equation*}
$$

with:

$$
\begin{align*}
T_{z D} & =\varepsilon_{0}\left[\varepsilon_{\infty}-\frac{\omega_{P}^{2}}{\omega^{2}+\mathrm{i} \gamma \omega}\right] M_{z}  \tag{9.68.a}\\
T_{z C_{p}} & =\varepsilon_{0}\left[A_{p} \Omega_{p}\left(\frac{\mathrm{e}^{\mathrm{i} \phi_{p}}}{\Omega_{p}-\omega-\mathrm{i} \Gamma_{p}}+\frac{\mathrm{e}^{-\mathrm{i} \phi_{p}}}{\Omega_{p}+\omega+\mathrm{i} \Gamma_{p}}\right)\right] M_{z} \tag{9.68.b}
\end{align*}
$$

Based on the inverse Fourier transforms of the equations (9.68) above, centered difference approximations for the derivatives and taking into account the equations 9.65), 9.67) and (9.66), we get:

$$
\begin{align*}
M_{z}^{n+1}= & \frac{1}{\frac{\chi_{D}}{\alpha_{D}}+\sum_{p=1}^{p=2}\left(\frac{\chi_{p}}{\alpha_{p}}\right)-\frac{k_{x}^{2}+k_{y}^{2}}{\mu \omega^{2}}}\left[\frac{k_{x}^{2}+k_{y}^{2}}{\mu \omega^{2}} P_{z}^{n+1}+\frac{k_{x}}{\omega} Q_{y a}^{n+1}-\frac{k_{y}}{\omega} Q_{x a}^{n+1}+\frac{\beta_{D}}{\alpha_{D}} T_{z_{D}}^{n-1}+\frac{4}{\alpha_{D}} T_{z_{D}}^{n}\right. \\
& \left.-\frac{\delta_{D}}{\alpha_{D}} M_{z}^{n-1}-\frac{4 \varepsilon_{0} \varepsilon_{\infty}}{\alpha_{D}} M_{z}^{n}+\sum_{p=1}^{p=2}\left(\frac{\beta_{p}}{\alpha_{p}} T_{z C_{p}}^{n-1}-\frac{4}{\alpha_{p}} T_{z_{C p}}^{n}\right)+\sum_{p=1}^{p=2}\left(\frac{\delta_{p}}{\alpha_{p}}\right) M_{z}^{n-1}\right]  \tag{9.69.a}\\
T_{z_{D}}^{n+1}= & \frac{1}{\alpha_{D}}\left[-\beta_{D} T_{z_{D}}^{n-1}-4 T_{z D}^{n}+\chi_{D} M_{z}^{n+1}+\delta_{D} M_{z}^{n-1}+4 \varepsilon_{0} \varepsilon_{\infty} M_{z}^{n}\right]  \tag{9.69.b}\\
T_{z C_{p}}^{n+1=}= & \frac{1}{\alpha_{p}}\left[-\beta_{p} T_{z C_{p}}^{n-1}+4 T_{z_{C p}}^{n}+\chi_{p} M_{z}^{n+1}+\delta_{p} M_{z}^{n-1}\right]  \tag{9.69.c}\\
P_{z}^{n+1}= & M_{z}^{n+1}+P_{z a}^{n+1} \tag{9.69.d}
\end{align*}
$$

The equations to update the $Q_{z}, P_{x}$ and $P_{y}$ components are obtained by the same process.


Figure 9.16: Up: Transmission spectra at normal incidence of an AAA structure made in silver film with different thicknesses values $(H)$. The geometrical parameters of the annular apertures are $r_{e}=75 \mathrm{~nm}, r_{i}=50 \mathrm{~nm}$ and the period is fixed to $a=300 \mathrm{~nm}$. Down: Electric field intensity distributions around the apertures showing the interference patterns that take place inside them along the metal thickness direction. For $F P_{0}$ peak, the $T E_{11}$ guided mode is excited at its cutoff wavelength so that the phase velocity tends to infinity and the effective index falls to zero. In this case, EOT occurs whatever is the value of the thickness because the phase matching condition is automatically fulfilled.

### 9.3.5 3D-SFM-FDTD application: EOT at oblique incidence through AAA structures

Let us recall the origin of the EOT through the AAA structure: as mentioned before, at normal incidence it is only due to the excitation of the $\mathrm{TE}_{11}$ guided mode inside each annular aperture. In this case, the obtained EOT is angle and polarization-independent and its spectral position corresponds to the cutoff wavelength of this guided mode. Consequently, it does not depends either on the metal thickness even if some additional peaks appear in the transmission spectrum when the thickness increases (see figure 9.16).

These peaks (named $\mathrm{FP}_{m}, m \in \mathfrak{R}$ on figure 9.16) are Fabry-Perot harmonics of the $\mathrm{TE}_{11}$ mode that occur at fixed values of the wavelength fulfilling a phase matching condition:

$$
\begin{equation*}
\lambda_{T E_{11}}\left(m \pi-\phi_{r}\right)=2 \pi n_{e f f}^{T E_{11}} H \tag{9.70}
\end{equation*}
$$

where $n_{\text {eff }}^{T E_{11}}$ is the real part of the effective index of the guided mode, $\phi_{r}$ is the phase


Figure 9.17: Schematic of a classical annular aperture array (AAA). $r_{e}$ is the outer radius, $r_{i}$ is the inner one, $a$ is the period and $\theta$ is the angle of incidence.
change induced by the reflection on the two ends of the annular aperture and $H$ is the metallic film thickness. At the cutoff, the effective index of the guided mode becomes very small leading to a phase matching that does not depend on the metal thickness. Nevertheless, a small spectral shift can appear between the cutoff value and the position of the transmission peak due to $\phi_{r} \neq 0$. This shift is clearly shown on all the spectra of figure 9.16 but it seems to be more important in the case of thicker plates (here $H=300 \mathrm{~nm}$ ). In fact, the phase $\phi_{r}$ can be seen as the result of the conversion between the incident plane wave and the guide mode through diffraction phenomenon that must depends on the metal thickness.

Let us now consider the case of oblique incidence (see figure 9.17): as mentioned before, EOT can appear through the excitation of both the $\mathrm{TE}_{11}$ and the TEM modes. In fact only few papers have discussed on this mode [38, 39] while its excitation conditions were recently analytically derived reference [22].

Indeed, this later is only excited with the TM polarization component of the incident beam. FDTD simulations in the case of both PEC (see figure 9.18) and real dispersive metal (figure 9 of reference [37, 40]) are done and demonstrate the occurrence of additional transmission peaks due to the excitation of the TEM guided mode. Nevertheless, others configurations such as the Slanted AAA (SAAA), that was proposed first by S. Nosal and J.J. Greffet [41], also demonstrate a possible excitation of the TEM mode for any incidence angle including normal incidence.

Moreover, as for the $\mathrm{TE}_{11}$ mode, the spectral position of the TEM-transmission peaks is driven by a similar phase matching condition given by equation 9.70 , Nonetheless, the zero harmonic $\left(\mathrm{FP}_{0}\right.$ for $\left.m=0\right)$ is now expelled to infinity and only higher orders correspond to a finite value of the wavelength. In this case, the metal thickness becomes a very important parameter that permits to adapt the transmission peak at a desired value of wavelength. Unluckily, only relatively thick metal plates allow the excitation and the propagation of the TEM mode.

Nevertheless, even if the TEM mode is excited in oblique incidence with conventional AAA (see figure 9.20 ) or at normal incidence through SAAA (figure 9.20b), the transmission


Figure 9.18: Transmission spectra through AAA structure made in perfectly electric conductor and illuminated by a TE (left) and TM (right) linearly polarized plane wave. As depicted on figure $9.15 \theta$ and $\phi$ denote the incidence and azimuthal angle respectively. The geometrical parameter of the AAA structure are: $r_{e}=a / 3, r_{i}=a / 4$ and $H=2 a$ (please see figure 9.17 for notations). Two families of TEM peak are pointed out using two white vertical rectangles. The right one corresponds here to the first Fabry-Perot harmonic and the left one frames the second harmonic. Note that other higher harmonics also occur at smaller wavelength values.


Figure 9.19: Real part $n_{e f f}^{\prime}$ and $\log _{10}$ of the imaginary part n"eff of the effective index associated with the TEMlike mode of an infinite coaxial waveguide as a function of the outer radius $r_{e}$. The inner radius is set to $r_{i}=65 \mathrm{~nm}$ and the working wavelength is $\lambda=1550 \mathrm{~nm}$.
efficiency remains very weak with regard to the $\mathrm{TE}_{11}$ mode. This is essentially due to metal losses. In fact, and as it can be shown in figure 9.19 , the imaginary part of the effective index of the TEM-like guided mode is fairly consistent and can not be negligible.

Fortunately, another solution that is currently used in the radio-frequency domain to increase the impedance adaptation between a coaxial antenna and the vacuum can be envisaged to enhance the transmission coefficient: it consists in stretching out the central metallic part of the coaxial waveguide with respect to the outside electrode. This configuration was implicitly proposed in reference [42] to achieve $90 \%$ light transmission thanks to the excitation of the TEM mode. This kind of structure design and fabrication is readily achievable at radio frequencies. Unfortunately, this becomes more difficult in the visible range but remains possible through manufacturing process having nanometric resolution such as new generation of Focused Ion Beam.


Figure 9.20: Zero-order transmission spectra for three different AAA configurations where outer and inner radii are fixed to $r_{e}=130 \mathrm{~nm}$ and $r_{i}=65 \mathrm{~nm}$ respectively. (a) Conventional structure illuminated at $45^{\circ}$ (metal thickness of $h=495 \mathrm{~nm}$ ). (b) SAAA with tilt angle of $25^{\circ}$ with respect to the vertical direction. The thickness of the metallic film $(h=430 \mathrm{~nm})$ is chosen in order to get a TEM peak transmission at $\lambda=1550 \mathrm{~nm}$. (c) SAAA structure with inner metallic parts that stretch out from the metallic film over a distance $e=80 \mathrm{~nm}$. This allows increasing of the impedance matching between the in- and out-coming plane waves with the TEM guided mode inside the apertures. The metal thickness is also adjusted to in order to get a TEM peak at $\lambda=1550 \mathrm{~nm}$ with a transmission efficiency of $48 \%$.

### 9.4 Conclusion

The FDTD is a powerful tool to model periodic and aperiodic structures. The time evolution of the electromagnetic field is directly evaluated and allows to follow the light propagation inside and around the studied structure. The SFM technique extends the FDTD capabilities to treat the diffraction problem for any incidence angle or any polarization. The integration of dispersion models such as Drude critical point allows accurate simulations that take into account the effective dispersion of noble metals in the considered spectral range especially in the visible domain. Nevertheless, the number of electromagnetic field components grows rapidly and can be larger than 100 in some particular cases (in the PML region with Drude-Lorentz dispersion model for instance). In spite of all these criticisms, the FDTD is actually one of the most used method to model experiments in Nano-Optics as attested by the number of publications in this area.

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Chapter 10:
Exact Modal Methods

Boris Gralak

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## Chapter 10

# Exact Modal Methods 

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### 10.1 Introduction

Exact modal method (EMM) has been proposed to take advantage of geometry of lamellar gratings. These gratings are made of rectangular rods periodically spaced which can be considered locally as periodic multilayered stacks (see figure 10.1). This simple geometry makes it possible to expand the electromagnetic field on the basis of "exact modes", and to obtain an exact representation of the permittivity. In this particular case, EMM can be more efficient than similar methods based on Fourier expansion (coupled-wave method [1] or Fourier modal method [2, 3]) which may lead to poor convergence due to the discontinuous nature of both electromagnetic field and permittivity. This advantage of EMM becomes more important when the permittivity contrast is high, e.g. for metallic lamellar gratings.


Figure 10.1: A lamellar grating made of a single lamellar layer on a substrate. The region corresponding to the lamellar layer, between the planes $x_{3}=0$ and $x_{3}=h$, can be considered as the multilayered stack on the left.

Exact modal method has been introduced in 1981 in order to solve Maxwell's equations in presence of lamellar gratings made of dielectrics [4] and metals [5, 6, 7]. Since these pioneering works, a major contribution to this method is certainly its rigorous extension to conical mountings [8], on which is based an EMM for three-dimensional woodpile structures [9]. Another major development is the introduction of perfectly matched layers in order to model aperiodic systems met in integrated optics [10] (information can be found on the website of CAMFR).

In this chapter, a rigorous formulation of the exact modal method for lamellar structures is presented. In section 10.3, a special attention is paid to the continuation of the electromagnetic
field inside a lamellar layer. In combination with the boundary conditions, this continuation provides a large class of solutions of Maxwell's equations in presence of lamellar gratings. In section 10.4, it is shown that, in each lamellar layer, there is a decoupling of the vector field equations into two independent scalar equations, which correspond to the ones of a multilayered stack (see figure 10.1). Numerical stacking algorithms are presented in section 10.5 and a numerical illustration of the EMM efficiency is proposed in section 10.6. Finally, the techniques used for the calculation of the exact modes and the associated exact eigenvalues are reported in the appendix (section 10.7).

Note that the important extensions to woodpile structures [9] and to lamellar gratings including infinitely conducting metal [11] are not considered in this book chapter. These cases will be however included in the next version.

### 10.2 Notations

Throughout this chapter an orthonormal basis $\left(\boldsymbol{e}_{1}, \boldsymbol{e}_{2}, \boldsymbol{e}_{3}\right)$ is used: every vector $\boldsymbol{x}$ in $\mathbb{R}^{3}$ is described by its three components $x_{1}, x_{2}$ and $x_{3}$. It is shown how to obtain in the presence of a stack of lamellar layers, a large class of solutions $\boldsymbol{E}$ of the Helmholtz equation

$$
\begin{equation*}
\left[\omega^{2}-\varepsilon^{-1} \boldsymbol{\nabla} \times \mu^{-1} \boldsymbol{\nabla} \times\right] \boldsymbol{E}=\mathbf{0}, \tag{10.1}
\end{equation*}
$$

where $\varepsilon$ is the permittivity, $\mu$ is the permeability, $\omega$ is the frequency and $\nabla \times$ is the curl operator. All the media considered in this chapter are isotropic, and thus the permittivity and permeability reduce to scalar functions. The considered structure is independent of the variable $x_{2}$, and $x_{1}$ periodic with spatial period $\boldsymbol{d}=d \boldsymbol{e}_{1}$ :

$$
\begin{equation*}
\varepsilon(\boldsymbol{x}+\boldsymbol{d})=\varepsilon(\boldsymbol{x})=\varepsilon\left(x_{1}, x_{3}\right), \quad \mu(\boldsymbol{x}+\boldsymbol{d})=\mu(\boldsymbol{x})=\mu\left(x_{1}, x_{3}\right), \quad \boldsymbol{x} \in \mathbb{R}^{3} . \tag{10.2}
\end{equation*}
$$

The unit cell associated with this grating is $[0, d]$ and the one-dimensional lattice is $\{n \mathbf{d} \mid n \in \mathrm{Z}\}$. Then, a lamellar grating is a stack in the direction $x_{3}$ of lamellar layers where $\varepsilon$ and $\mu$ are both functions of the single variable $x_{1}$ (figure 10.2). In practice, each lamellar layer is made of infinite parallel rods with rectangular cross section (figure 10.2): the functions $\varepsilon$ and $\mu$ are piecewise constant of the solely variable $x_{1}$.


Figure 10.2: A lamellar grating made of a single lamellar layer on a substrate.

In order to obtain a set of first order differential equations from (10.1) a second field is defined:

$$
\begin{equation*}
\boldsymbol{H}=(\omega \mu)^{-1} \boldsymbol{\nabla} \times \boldsymbol{E} \tag{10.3}
\end{equation*}
$$

Note that this quantity differs from the usual "harmonic $\boldsymbol{H}$ field" by the complex number $i$. Solutions $\boldsymbol{E}, \boldsymbol{H}$ are investigated in the space of fields whose restrictions in every horizontal
plane (normal to $\boldsymbol{e}_{3}$ ) are square integrable:

$$
\begin{equation*}
\int_{\mathbb{R}^{2}}\left|\boldsymbol{F}\left(x_{1}, x_{2}, x_{3}\right)\right|^{2} d x_{1} d x_{2}<\infty, \quad x_{3} \in \mathbb{R} \tag{10.4}
\end{equation*}
$$

where $\boldsymbol{F}=\boldsymbol{E}, \boldsymbol{H}$.
The first consequence of (10.4) is the possibility to perform a decomposition of the problem to take advantage of the spatial invariances of the system: a Fourier decomposition with respect to the variable $x_{2}$,

$$
\begin{equation*}
\boldsymbol{F} \longrightarrow \widehat{\boldsymbol{F}}\left(x_{1}, k_{2}, x_{3}\right)=\frac{1}{2 \pi} \int_{\mathbb{R}} \exp \left[-i k_{2} x_{2}\right] \boldsymbol{F}\left(x_{1}, x_{2}, x_{3}\right) d x_{2}, \tag{10.5}
\end{equation*}
$$

and a Floquet-Bloch decomposition with respect to the variable $x_{1}$,

$$
\begin{equation*}
\widehat{\boldsymbol{F}} \longrightarrow \tilde{\boldsymbol{F}}\left(k_{1}, x_{1}, k_{2}, x_{3}\right) \frac{1}{2 \pi} \sum_{n \in \mathrm{Z}} \exp \left[-i k_{1} n d\right] \widehat{\boldsymbol{F}}\left(x_{1}+p d, k_{2}, x_{3}\right) \tag{10.6}
\end{equation*}
$$

where $k_{1}$ is the Bloch wave vector in the first Brillouin zone $[-\pi / d, \pi / d]$. Thus solutions $\tilde{\boldsymbol{E}}, \tilde{\boldsymbol{H}}$ satisfy

$$
\begin{equation*}
\int_{[-\pi / d, \pi / d]}\left|\tilde{\boldsymbol{F}}\left(k_{1}, x_{1}, k_{2}, x_{3}\right)\right|^{2} d x_{1}<\infty, \quad x_{1}, k_{2}, x_{3} \in \mathbb{R} \tag{10.7}
\end{equation*}
$$

with the partial Bloch boundary condition

$$
\begin{equation*}
\tilde{\boldsymbol{F}}\left(k_{1}, x_{1}+d, k_{2}, x_{3}\right)=\exp \left[i k_{1} d\right] \tilde{\boldsymbol{F}}\left(k_{1}, x_{1}, k_{2}, x_{3}\right), \tag{10.8}
\end{equation*}
$$

where $k_{1}$ is fixed in $[-\pi / d, \pi / d]$.
The second consequence of (10.4) [or (10.7)] is that the restrictions to every horizontal plane of $\boldsymbol{\nabla} \times \boldsymbol{E}$ and $\boldsymbol{\nabla} \times \boldsymbol{H}$ are also locally square integrable [from (10.1,10.3)]. Then, for all $i, j=1,2,3$ and $i \neq j, E_{i}$ and $H_{i}$ are continuous functions of the variable $x_{j}$. In particular, the tangential components $E_{1}, E_{2}, H_{1}$ and $H_{2}$ of $\boldsymbol{E}$ and $\boldsymbol{H}$ are continuous functions of the variable $x_{3}$. It follows that it is possible to solve Maxwell's equations in a stack of layers by the following two steps: the first step consists in solving Maxwell's equations in each layer independently and then the second step consists in connecting each independent solution using the continuity of $E_{1}, E_{2}, H_{1}$ and $H_{2}$.

With the definition (10.3), equation (10.1) is equivalent to the set of first order equations

$$
\begin{equation*}
\boldsymbol{E}=(\omega \boldsymbol{\varepsilon})^{-1} \boldsymbol{\nabla} \times \boldsymbol{H}, \quad \boldsymbol{H}=(\omega \mu)^{-1} \boldsymbol{\nabla} \times \boldsymbol{E} . \tag{10.9}
\end{equation*}
$$

Let the $2 \times 2$ matrix $\sigma$ and the two-components vector $F_{j}$ defined by

$$
\sigma=\omega\left[\begin{array}{cc}
0 & \mu  \tag{10.10}\\
\varepsilon & 0
\end{array}\right], \quad F_{j}=\left[\begin{array}{c}
\tilde{E}_{j} \\
\tilde{H}_{j}
\end{array}\right], \quad j=1,2,3
$$

Then, the first order equations (10.9) can be developed as

$$
\begin{align*}
& F_{1}=\sigma^{-1}\left[\partial_{2} F_{3}-\partial_{3} F_{2}\right], \\
& F_{2}=\sigma^{-1}\left[\partial_{3} F_{1}-\partial_{1} F_{3}\right],  \tag{10.11}\\
& F_{3}=\sigma^{-1}\left[\partial_{1} F_{2}-\partial_{2} F_{1}\right],
\end{align*}
$$

where $\partial_{j}$ is the partial derivative with respect to the variable $x_{j}(j=1,2,3)$. This last set of equations is exactly the same as (10.9) and, with some abuse of notations, it can written in the compact way $\boldsymbol{F}=\sigma^{-1} \boldsymbol{\nabla} \times \boldsymbol{F}$, with $\boldsymbol{F}=\left(F_{1}, F_{2}, F_{3}\right)$.

### 10.3 Continuation of the electromagnetic field

In this section two different formulations are presented to solve the equation $\boldsymbol{F}=\sigma^{-1} \boldsymbol{\nabla} \times \boldsymbol{F}$ in a lamellar layer located between the planes $x_{3}=0$ and $x_{3}=h$. In practice, this solution is expressed as a relationship between $\boldsymbol{F}(0)$ and $\boldsymbol{F}(h)$. Note that the formulations presented in this section remain valid in the general case of cross gratings with two-dimensional periodicity [2].

### 10.3.1 Direct formulation: the transfer matrix

The starting point is the set of equations (10.11). Eliminating the components $F_{3}$, one obtains

$$
\partial_{3} F=i M F, \quad F=\left[\begin{array}{l}
F_{1}  \tag{10.12}\\
F_{2}
\end{array}\right], \quad M=-i\left[\begin{array}{lr}
-\partial_{1} \sigma^{-1} \partial_{2} & \sigma+\partial_{1} \sigma^{-1} \partial_{1} \\
-\sigma-\partial_{2} \sigma^{-1} \partial_{2} & \partial_{2} \sigma^{-1} \partial_{1}
\end{array}\right] .
$$

For a lamellar layer located between the planes $x_{3}=0$ and $x_{3}=h$ (see figure 10.1), the functions $\varepsilon$ and $\mu$ are $x_{3}$-independent for $x_{3}$ in $[0, h]$ : then the matrix $\sigma$ and the operator-valued matrix $M$ are also $x_{3}$-independant for $x_{3}$ in $[0, h]$. Let $L$ be the $x_{3}$-independant operator valued matrix which coincides with $M$ in this single layer:

$$
\begin{equation*}
L=M\left(x_{3}\right), \quad x_{3} \in[0, h] . \tag{10.13}
\end{equation*}
$$

In a first step, it is assumed that (see the end of section 10.4 for a justification) the matrix $L$ has a diagonal form and can be written as

$$
\begin{equation*}
L=V \lambda V^{-1}, \tag{10.14}
\end{equation*}
$$

where the matrix $V$ contains the eigenvectors $V_{ \pm, n}$ of $L$, and $\lambda$ is the diagonal matrix made of the associated eigenvalues $\lambda_{ \pm, n}$ :

$$
\begin{equation*}
L V_{ \pm, n}=\lambda_{ \pm, n} V_{ \pm, n} \tag{10.15}
\end{equation*}
$$

The sets of eigenvectors and eigenvalues are split into two parts according to the sign of the imaginary part of $\lambda_{ \pm, n}: \operatorname{Im} \lambda_{+, n}>0$ and $\operatorname{Im} \lambda_{-, n}<0$. Let $\lambda_{+}$(respectively $\lambda_{-}$) be the diagonal matrices containing the eigenvalues $\lambda_{+, n}>0$ of $L$ (respectively $\lambda_{-, n}>0$ ): then

$$
\lambda=\left[\begin{array}{ll}
\lambda_{+} & 0  \tag{10.16}\\
0 & \lambda_{-}
\end{array}\right], \quad \operatorname{Im} \lambda_{+}>0, \quad \operatorname{Im} \lambda_{-}<0
$$

This last condition on the imaginary part of eigenvalues is always realized if there is some absorption, i.e. $\operatorname{Im} \sigma>0$, or if a small positive imaginary part is added to the frequency $\omega$ (in the later case the limit $\omega=\lim _{\eta \downarrow 0}(\omega+i \eta)$ is considered [12, 13]). The combination of (10.12) and (10.13) leads to the equation

$$
\begin{equation*}
\partial_{3} F\left(x_{3}\right)=i L F\left(x_{3}\right), \quad x_{3} \in[0, h], \tag{10.17}
\end{equation*}
$$

where the dependence on other variables has been omitted. Since $L$ is $x_{3}$-independent, the "formal" solution of this equation is just

$$
\begin{equation*}
F\left(x_{3}\right)=\exp \left[i L x_{3}\right] F(0), \quad x_{3} \in[0, h] . \tag{10.18}
\end{equation*}
$$

This solution is denominated by "formal" since, at this stage, it is still necessary to check if it exists. Using the diagonal form (10.14) of the operator $L$, the expression (10.18) becomes

$$
\begin{equation*}
F\left(x_{3}\right)=V \exp \left[i \lambda x_{3}\right] V^{-1} F(0) . \tag{10.19}
\end{equation*}
$$

Actually, the diagonal matrix $\exp \left[i \lambda x_{3}\right]$ is made of the two parts $\exp \left[i \lambda_{ \pm} x_{3}\right]$ which have different behaviour. From (10.16), the part $\exp \left[i \lambda_{+} x_{3}\right]$ is bounded by $\exp \left[-\operatorname{Im} \lambda_{+} x_{3}\right]<1$. On the contrary, the part $\exp \left[i \lambda_{-} x_{3}\right]$ is not bounded and, in general, the corresponding coefficients are growing towards infinity like exponential functions. Consequently, the transfer matrix $T\left(x_{3}\right)$ defined by

$$
\begin{equation*}
F\left(x_{3}\right)=T\left(x_{3}\right) F(0), \quad T\left(x_{3}\right)=V \exp \left[i \lambda x_{3}\right] V^{-1} \tag{10.20}
\end{equation*}
$$

has "infinite" coefficients and thus expressions like (10.18), (10.19) and (10.20) have to be considered as purely "formal" and have to be handled cautiously. Numerically, the transfer matrix is truncated and, because its coefficients tend to infinity like exponential functions, it presents numerical instabilities which makes it difficult to use it. A numerical solution has been found to solve this problem with the definition of the $S$ - and $R$-algorithms [14] (see section 10.5 for the numerical solution in the present case).

Finally, notice that the transfer matrix is occasionally derived from the matrix $L^{2}$ instead of $L$. Indeed, the assumption (10.14) on the diagonal form of $L$ might be too strong and not rigorously true. According to notations (10.14), we denote by $V$ and $\lambda^{2}$ the matrices containing the eigenvectors and eigenvalues of $L^{2}$ :

$$
\begin{equation*}
L^{2}=V \lambda^{2} V^{-1} \tag{10.21}
\end{equation*}
$$

In that case, it used that equations (10.12) and (10.13) imply

$$
\begin{equation*}
\partial_{3}^{2} F=-L^{2} F, \quad x_{3} \in[0, h] \tag{10.22}
\end{equation*}
$$

The combination of the two last equations leads to

$$
\begin{equation*}
F\left(x_{3}\right)=V \cos \left[\lambda x_{3}\right] V^{-1} F(0)+V \lambda^{-1} \sin \left[\lambda x_{3}\right] V^{-1}\left(\partial_{3} F\right)(0) \tag{10.23}
\end{equation*}
$$

Replacing $\left(\partial_{3} F\right)(0)$ by $i L F(0)$, one obtains for the transfer matrix the following "formal" expression

$$
\begin{equation*}
T\left(x_{3}\right)=V \cos \left[\lambda x_{3}\right] V^{-1}+i V \lambda^{-1} \sin \left[\lambda x_{3}\right] V^{-1} L \tag{10.24}
\end{equation*}
$$

This equation is not "formally" equivalent to the first expression (10.20) derived from $L$. This equivalence requires the assumption (10.14) to become true, so that $L$ can be replaced by $V \lambda V^{-1}$ above (and next the formal identity $\cos \left[\lambda x_{3}\right]+i \sin \left[\lambda x_{3}\right]=\exp \left[i \lambda x_{3}\right]$ has to be used).

### 10.3.2 Rigorous derivation of the continuation procedure

A rigorous formulation is based on the use of the Fourier transform with respect to the variable $x_{3}$ defined by

$$
\begin{equation*}
\mathscr{F}[F]\left(k_{3}\right)=\frac{1}{2 \pi} \int_{\mathbb{R}} \exp \left[-i k_{3} x_{3}\right] F\left(x_{3}\right) d x_{3} \tag{10.25}
\end{equation*}
$$

The function $F$ is then deduced from its Fourier transform $\mathscr{F}[F]$ by

$$
\begin{equation*}
F\left(x_{3}\right)=\int_{\mathbb{R}} \exp \left[i k_{3} x_{3}\right] \mathscr{F}[F]\left(k_{3}\right) d k_{3} \tag{10.26}
\end{equation*}
$$

It is not suitable to perform directly the Fourier transform of the equation (10.12) since the matrix $\sigma$ (and then $M$ ) is not independent of $x_{3}$ in $\mathbb{R}$. However, if equation (10.12) is multiplied by the characteristic function $\Psi$ of the lamellar layer [so $\Psi\left(x_{3}\right)=1$ for $x_{3}$ in $[0, h]$ and vanishes otherwise], then

$$
\begin{equation*}
\Psi\left(x_{3}\right) \partial_{3} F\left(x_{3}\right)=\Psi\left(x_{3}\right) i M F\left(x_{3}\right)=\Psi\left(x_{3}\right) i L F\left(x_{3}\right)=i L \Psi\left(x_{3}\right) F\left(x_{3}\right) . \tag{10.27}
\end{equation*}
$$

After this multiplication, a partial differential equation with the $x_{3}$-independent matrix $L$ is obtained. The Fourier transform (10.25) of $\Psi \partial_{3} F$ is

$$
\begin{align*}
\mathscr{F}\left[\Psi \partial_{3} F\right]\left(k_{3}\right) & =\frac{1}{2 \pi} \int_{\mathbb{R}} \exp \left[-i k_{3} x_{3}\right] \Psi\left(x_{3}\right) \partial_{3} F\left(x_{3}\right) d x_{3} \\
& =\frac{1}{2 \pi} \int_{0}^{h} \exp \left[-i k_{3} x_{3}\right] \partial_{3} F\left(x_{3}\right) d x_{3}  \tag{10.28}\\
& =\frac{1}{2 \pi}\left\{\exp \left[-i k_{3} h\right] F(h)-F(0)\right\}+i k_{3} \mathscr{F}[\Psi F]\left(k_{3}\right),
\end{align*}
$$

where the last line comes from an integration by parts. After this Fourier transform, equation (10.27) becomes

$$
\begin{equation*}
\frac{1}{2 \pi}\left\{\exp \left[-i k_{3} h\right] F(h)-F(0)\right\}+i k_{3} \mathscr{F}[\Psi F]\left(k_{3}\right)=i L \mathscr{F}[\Psi F]\left(k_{3}\right) \tag{10.29}
\end{equation*}
$$

or

$$
\begin{equation*}
\left[k_{3}-L\right] \mathscr{F}[\Psi F]\left(k_{3}\right)=\frac{1}{2 i \pi}\left\{F(0)-\exp \left[-i k_{3} h\right] F(h)\right\} . \tag{10.30}
\end{equation*}
$$

The operator $\left[k_{3}-L\right]$ is always invertible if there is some absorption, i.e. $\operatorname{Im} \sigma>0$, or if the limit $\omega=\lim _{\eta \downarrow 0}(\omega+i \eta)$ is considered (see [12, 13], this is equivalent to the property (10.16) on the eigenvalues $\lambda$ since $k_{3}$ is purely real). Hence it is possible to write

$$
\begin{equation*}
\mathscr{F}[\Psi F]\left(k_{3}\right)=\frac{1}{2 i \pi} \frac{1}{k_{3}-L}\left\{F(0)-\exp \left[-i k_{3} h\right] F(h)\right\}, \tag{10.31}
\end{equation*}
$$

The final step is to apply the inverse Fourier transform (10.26): for $x_{3}$ in $[0, h]$,

$$
\begin{align*}
\Psi\left(x_{3}\right) F\left(x_{3}\right) & =\frac{1}{2 i \pi}\left[\int_{\mathbb{R}} \exp \left[i k_{3} x_{3}\right] \frac{1}{k_{3}-L} d k_{3}\right] F(0) \\
& -\frac{1}{2 i \pi}\left[\int_{\mathbb{R}} \exp \left[i k_{3}\left(x_{3}-h\right)\right] \frac{1}{k_{3}-L} d k_{3}\right] F(h) \tag{10.32}
\end{align*}
$$

Again, it is assumed that the operator $L$ can be written $L=V \lambda V^{-1}$ (10.14). Replacing the matrix $L$ by its diagonal form, the last expression becomes

$$
\begin{align*}
\Psi\left(x_{3}\right) F\left(x_{3}\right) & =\frac{1}{2 i \pi} V\left[\int_{\mathbb{R}} \exp \left[i k_{3} x_{3}\right] \frac{1}{k_{3}-\lambda} d k_{3}\right] V^{-1} F(0) \\
& -\frac{1}{2 i \pi} V\left[\int_{\mathbb{R}} \exp \left[i k_{3}\left(x_{3}-h\right)\right] \frac{1}{k_{3}-\lambda} d k_{3}\right] V^{-1} F(h) . \tag{10.33}
\end{align*}
$$

The integrations above are performed by adding to the real axis of $k_{3}$ a semi-circle with infinite radius (in the complex plane of $k_{3}$ ) on which the integrals vanish. For the first term with $F(0)$,
the complex number $k_{3}$ must have positive imaginary part ( $x_{3}$ is positive), so that the real axis is closed by a semi-circle in the upper half plane (see the red path on figure 10.3). In this case, the solely eigenvalues contained in $\lambda_{+}$generate contributions in the integral. For the second term with $F(h)$, the complex number $k_{3}$ must have negative imaginary part ( $x_{3}-h$ is negative), so that the real axis is closed by a semi-circle in the lower half plane (see the blue path on figure 10.3). Here, the integral is given by the eigenvalues contained in $\lambda_{-}$. Let $P_{ \pm}$be the projectors


Figure 10.3: Integration in the complex plane of $k_{3}$.
upon the spaces corresponding respectively to eigenvalues $\lambda_{ \pm}$:

$$
P_{+} \lambda=\left[\begin{array}{ll}
\lambda_{+} & 0  \tag{10.34}\\
0 & 0
\end{array}\right], \quad P_{-} \lambda=\left[\begin{array}{cc}
0 & 0 \\
0 & \lambda_{-}
\end{array}\right]
$$

Then, after the integration over $k_{3}$, expression (10.33) yields

$$
\begin{equation*}
\Psi\left(x_{3}\right) F\left(x_{3}\right)=V P_{+} \exp \left[i \lambda_{+} x_{3}\right] V^{-1} F(0)+V P_{-} \exp \left[i \lambda_{-}\left(x_{3}-h\right)\right] V^{-1} F(h) \tag{10.35}
\end{equation*}
$$

This expression is always well defined since the integration in the complex plane of $k_{3}$ imposes that all the complex exponential functions decrease:

$$
\begin{equation*}
\left\|\exp \left[i \lambda_{+} x_{3}\right]\right\| \leq 1, \quad\left\|\exp \left[i \lambda_{-}\left(x_{3}-h\right)\right]\right\| \leq 1 \tag{10.36}
\end{equation*}
$$

Considering the rigorous expression (10.35) at $x_{3}=0$ and $x_{3}=h$ and using that $P_{-}+P_{+}$is the identity, one obtains

$$
\begin{align*}
& V P_{-} V^{-1} F(0)=V P_{-} \exp \left[-i \lambda_{-} h\right] V^{-1} F(h)  \tag{10.37}\\
& V P_{+} V^{-1} F(h)=V P_{+} \exp \left[i \lambda_{+} h\right] V^{-1} F(0)
\end{align*}
$$

These two relationships provides a rigorous way to deduce $F(0)$ from $F(h)$ and conversely.
As in the previous section, the continuation procedure is also derived from the diagonal form (10.21) of $L^{2}$. Equation (10.22) is multiplied by $\Psi\left(x_{3}\right)$ to provide an expression similar to (10.27)

$$
\begin{equation*}
\Psi\left(x_{3}\right) \partial_{3} F\left(x_{3}\right)=-L^{2} \Psi\left(x_{3}\right) F\left(x_{3}\right) \tag{10.38}
\end{equation*}
$$

Then, the Fourier transform (10.25) is applied to this equation. Using that

$$
\begin{align*}
\mathscr{F}\left[\Psi \partial_{3}^{2} F\right]\left(k_{3}\right) & =\frac{1}{2 \pi} \int_{\mathbb{R}} \exp \left[-i k_{3} x_{3}\right] \Psi\left(x_{3}\right) \partial_{3}^{2} F\left(x_{3}\right) d x_{3} \\
& =-k_{3}^{2} \mathscr{F}[\Psi F]\left(k_{3}\right)+i k_{3} \frac{1}{2 \pi}\left\{\exp \left[-i k_{3} h\right] F(h)-F(0)\right\}  \tag{10.39}\\
& +\frac{1}{2 \pi}\left\{\exp \left[-i k_{3} h\right]\left(\partial_{3} F\right)(h)-\left(\partial_{3} F\right)(0)\right\},
\end{align*}
$$

and replacing $\left(\partial_{3} F\right)\left(x_{3}\right)$ by $\operatorname{iLF}\left(x_{3}\right)$, equation (10.38) implies

$$
\begin{align*}
\mathscr{F}[\Psi F]\left(k_{3}\right) & =\frac{1}{2 i \pi} \frac{1}{k_{3}^{2}-L^{2}} k_{3}\left\{F(0)-\exp \left[-i k_{3} h\right] F(h)\right\} \\
& =\frac{1}{2 i \pi} \frac{1}{k_{3}^{2}-L^{2}}\left\{L F(0)-\exp \left[-i k_{3} h\right] L F(h)\right\} . \tag{10.40}
\end{align*}
$$

Next, the inverse Fourier transform (10.26) is performed for $x_{3}$ in $[0, h]$ and the diagonal form (10.21) is used:

$$
\begin{align*}
\Psi\left(x_{3}\right) F\left(x_{3}\right) & =\frac{1}{2 i \pi} V \int_{\mathbb{R}} \frac{k_{3}}{k_{3}^{2}-\lambda^{2}}\left\{\exp \left[i k_{3} x_{3}\right] V^{-1} F(0)-\exp \left[i k_{3}\left(x_{3}-h\right)\right] V^{-1} F(h)\right\} d x_{3} \\
& +\frac{1}{2 i \pi} V \int_{\mathbb{R}} \frac{1}{k_{3}^{2}-\lambda^{2}}\left\{\exp \left[i k_{3} x_{3}\right] V^{-1} L F(0)-\exp \left[i k_{3}\left(x_{3}-h\right)\right] V^{-1} L F(h)\right\} d x_{3} . \tag{10.41}
\end{align*}
$$

Again, the integrations above are calculated by adding to the real axis of $k_{3}$ a semi-circle with infinite radius (in the complex plane of $k_{3}$ ) on which the integrals vanish. Without loss of generality, it is considered that the square root of the eigenvalues in $\lambda^{2}$ have non-zero imaginary part: let $\sqrt{\lambda^{2}}$ be the square root of $\lambda^{2}$ with positive imaginary part. For the terms with $F(0)$, the real axis is closed by a semi-circle in the upper half plane, and the eigenvalues with positive imaginary part $\sqrt{\lambda^{2}}$ lead to contributions in the integrals. For the terms with $F(h)$, the real axis is closed by a semi-circle in the lower half plane, and the integrals are given by the eigenvalues with negative imaginary part, i.e. $-\sqrt{\lambda^{2}}$. Calculations of integrals over $k_{3}$ lead to

$$
\begin{align*}
\Psi\left(x_{3}\right) F\left(x_{3}\right) & =V \frac{1}{2} \exp \left[i \sqrt{\lambda^{2}} x_{3}\right] V^{-1} F(0)+V \frac{1}{2} \exp \left[-i \sqrt{\lambda^{2}}\left(x_{3}-h\right)\right] V^{-1} F(h) \\
& +V \frac{1}{2 \sqrt{\lambda^{2}}} \exp \left[i \sqrt{\lambda^{2}} x_{3}\right] V^{-1} L F(0)-V \frac{1}{2 \sqrt{\lambda^{2}}} \exp \left[-i \sqrt{\lambda^{2}}\left(x_{3}-h\right)\right] V^{-1} L F(h) . \tag{10.42}
\end{align*}
$$

This equation is evaluated at $x_{3}=0$ and $x_{3}=h$ :

$$
\begin{align*}
& F(0)=V \exp \left[i \sqrt{\lambda^{2}} h\right] V^{-1} F(h)+V \frac{1}{\sqrt{\lambda^{2}}} V^{-1} L F(0)-V \frac{1}{\sqrt{\lambda^{2}}} \exp \left[i \sqrt{\lambda^{2}} h\right] V^{-1} L F(h), \\
& F(h)=V \exp \left[i \sqrt{\lambda^{2}} h\right] V^{-1} F(0)+V \frac{1}{\sqrt{\lambda^{2}}} \exp \left[i \sqrt{\lambda^{2}} h\right] V^{-1} L F(0)-V \frac{1}{\sqrt{\lambda^{2}}} V^{-1} L F(h) \tag{10.43}
\end{align*}
$$

Thanks to the technique based on the Fourier transform, all the exponential functions in these expressions must be well-defined. Indeed, the imaginary part of $\sqrt{\lambda^{2}}$ is positive and all the exponential functions decrease. Equations (10.43) will be used to construct a stable numerical algorithm to stack several lamellar layers.

### 10.4 Exact eigenmodes and eigenvalues method

The different solutions (10.20), (10.24, (10.37) and (10.43), established in the previous section, are provided from the knowledge of the sets of eigenmodes and eigenvalues of the operator $L$ (or $L^{2}$ ). In this section, it is shown how these eigenmodes $V_{ \pm, n}$ and eigenvalues $\lambda_{ \pm, n}$ of $L$ can be exactly determined in a lamellar layer located between the planes $x=0$ and $x_{3}=h$. The starting point is equation (10.17):

$$
\begin{equation*}
\partial_{3} F=i L F, \quad L=M\left(x_{3}\right), \quad x_{3} \in[0, h] . \tag{10.44}
\end{equation*}
$$

Let $\varepsilon_{1}, \mu_{1}$ and $\sigma_{1}$ be the functions which coincide with respectively $\varepsilon, \mu$ and $\sigma$ in the considered lamellar layer for $x_{3}$ in $[0, h]$. In the considered lamellar layer, they are functions of the solely variable $x_{1}$ (see figure 10.2):

$$
\begin{equation*}
\varepsilon_{1}\left(x_{1}\right)=\varepsilon\left(x_{1}, x_{3}\right), \quad \mu_{1}\left(x_{1}\right)=\mu\left(x_{1}, x_{3}\right), \quad \sigma_{1}\left(x_{1}\right)=\sigma\left(x_{1}, x_{3}\right), \quad x_{3} \in[0, h] . \tag{10.45}
\end{equation*}
$$

According to expression (10.12), the operator $L$ is now

$$
L=-i\left[\begin{array}{lr}
-\partial_{1} \sigma_{1}^{-1} \partial_{2} & \sigma_{1}+\partial_{1} \sigma_{1}^{-1} \partial_{1}  \tag{10.46}\\
-\sigma_{1}-\partial_{2} \sigma_{1}^{-1} \partial_{2} & \partial_{2} \sigma_{1}^{-1} \partial_{1}
\end{array}\right],
$$

and its square is

$$
L^{2}=-\left[\begin{array}{lr}
-\sigma_{1}^{2}-\partial_{1} \sigma_{1}^{-1} \partial_{1} \sigma_{1}-\sigma_{1} \partial_{2} \sigma_{1}^{-1} \partial_{2} & \sigma_{1} \partial_{2} \sigma_{1}^{-1} \partial_{1}-\partial_{1} \sigma_{1}^{-1} \partial_{2} \sigma_{1}  \tag{10.47}\\
\sigma_{1} \partial_{1} \sigma_{1}^{-1} \partial_{2}-\partial_{2} \sigma_{1}^{-1} \partial_{1} \sigma_{1} & -\sigma_{1}^{2}-\partial_{2} \sigma_{1}^{-1} \partial_{2} \sigma_{1}-\sigma_{1} \partial_{1} \sigma_{1}^{-1} \partial_{1}
\end{array}\right] .
$$

Since the matrix $\sigma_{1}$ is $x_{2}$-independent, the equality $\sigma_{1} \partial_{2} \sigma_{1}^{-1}=\partial_{2}=\sigma_{1}^{-1} \partial_{2} \sigma_{1}$ holds, and the expression above becomes

$$
L^{2}=\left[\begin{array}{lr}
\sigma_{1}^{2}+\partial_{2}^{2}+\partial_{1} \sigma_{1}^{-1} \partial_{1} \sigma_{1} & 0  \tag{10.48}\\
\partial_{2} \sigma_{1}^{-1} \partial_{1} \sigma_{1}-\sigma_{1} \partial_{1} \sigma_{1}^{-1} \partial_{2} & \sigma_{1}^{2}+\partial_{2}^{2}+\sigma_{1} \partial_{1} \sigma_{1}^{-1} \partial_{1}
\end{array}\right] .
$$

This expression shows that the components $F_{1}$ can be decoupled from the components $F_{2}$ in the lamellar layer. Indeed, it implies

$$
\begin{equation*}
\partial_{3}^{2} F_{1}=-K F_{1}, \quad K=\sigma_{1}^{2}+\partial_{2}^{2}+\partial_{1} \sigma_{1}^{-1} \partial_{1} \sigma_{1} . \tag{10.49}
\end{equation*}
$$

Moreover, each component of $F_{1}$, i.e. $E_{1}$ and $H_{1}$, can be also decoupled since the operator $K$ is diagonal:

$$
K=\left[\begin{array}{lr}
K_{\varepsilon_{1}} & 0  \tag{10.50}\\
0 & K_{\mu_{1}}
\end{array}\right]
$$

where

$$
\begin{align*}
& K_{\varepsilon_{1}}=\omega^{2} \varepsilon_{1} \mu_{1}+\partial_{2}^{2}+\partial_{1} \varepsilon_{1}^{-1} \partial_{1} \varepsilon_{1} \\
& K_{\mu_{1}}=\omega^{2} \varepsilon_{1} \mu_{1}+\partial_{2}^{2}+\partial_{1} \mu_{1}^{-1} \partial_{1} \mu_{1} \tag{10.51}
\end{align*}
$$

Here, it is important to notice that the two operators $K_{\varepsilon_{1}}$ and $K_{\mu_{1}}$ correspond to the ones of a onedimensional multilayered stack for respectively $p$ - and $s$-polarization. This makes it possible to calculate the exact eigenmodes and eigenvalues of $K_{\mathcal{E}_{1}}$ and $K_{\mu_{1}}$ (see appendix) and thus the ones of $K$. Thus the continuation procedure presented in section 10.3.2 can be applied to
equation (10.49). It provides relationships between the fields $F_{1}(0), F_{1}(h)$ and their derivative with respect to $x_{3}$, i.e. $\left[\partial_{3} F_{1}\right](0)$ and $\left[\partial_{3} F_{1}\right](h)$ (it is recalled that $\partial_{3} F=i L F$ in section 10.3.2).

To complete the derivation of the method, it is necessary to express the component $F_{2}$ of the field from $F_{1}$ and $\partial_{3} F_{1}$. A starting relationship is obtained from (10.44) and (10.47):

$$
\begin{equation*}
\partial_{3} F_{1}=-\partial_{1} \sigma_{1}^{-1} \partial_{2} F_{1}+\left[\sigma_{1}+\partial_{1} \sigma_{1}^{-1} \partial_{1}\right] F_{2} . \tag{10.52}
\end{equation*}
$$

This equation is equivalent to

$$
\begin{equation*}
\left[1+\partial_{1} \sigma_{1}^{-1} \partial_{1} \sigma_{1}^{-1}\right] \sigma_{1} F_{2}=\partial_{3} F_{1}+\partial_{1} \sigma_{1}^{-1} \partial_{2} F_{1} \tag{10.53}
\end{equation*}
$$

Here, it is remarked that the operator $\left[1+\partial_{1} \sigma_{1}^{-1} \partial_{1} \sigma_{1}^{-1}\right]$ is invertible since it equals $\left[K-\partial_{2}^{2}\right] \sigma_{1}^{-2}$ where $\sigma_{1}$ is invertible as well as $\left[K-\partial_{2}^{2}\right]$ (eigenvalues of $K$ have non-zero imaginary part). Moreover, the operator $\left[1+\partial_{1} \sigma_{1}^{-1} \partial_{1} \sigma_{1}^{-1}\right]$ commutes with $\partial_{3}, \partial_{2}$ and $\partial_{1} \sigma_{1}^{-1}$ : hence

$$
\begin{align*}
\sigma_{1} F_{2} & =\frac{1}{1+\partial_{1} \sigma_{1}^{-1} \partial_{1} \sigma_{1}^{-1}} \partial_{3} F_{1}+\frac{1}{1+\partial_{1} \sigma_{1}^{-1} \partial_{1} \sigma_{1}^{-1}} \partial_{1} \sigma_{1}^{-1} \partial_{2} F_{1} \\
& =\frac{1}{1+\partial_{1} \sigma_{1}^{-1} \partial_{1} \sigma_{1}^{-1}} \partial_{3} F_{1}+\partial_{2} \partial_{1} \sigma_{1}^{-1} \frac{1}{1+\partial_{1} \sigma_{1}^{-1} \partial_{1} \sigma_{1}^{-1}} F_{1}  \tag{10.54}\\
& =\sigma_{1}^{2} \frac{1}{\sigma_{1}^{2}+\partial_{1} \sigma_{1}^{-1} \partial_{1} \sigma_{1}} \partial_{3} F_{1}+\partial_{2} \partial_{1} \sigma_{1} \frac{1}{\sigma_{1}^{2}+\partial_{1} \sigma_{1}^{-1} \partial_{1} \sigma_{1}} F_{1} .
\end{align*}
$$

After the Fourier decomposition with respect to the variable $x_{2}$, the operator $\partial_{2}$ becomes $i k_{2}$. Then, the inverse operator above is expressed using the diagonal form $K=V \lambda^{2} V^{-1}$, and the component $F_{2}$ becomes

$$
\begin{equation*}
F_{2}=\sigma_{1} V \frac{1}{\lambda^{2}+k_{2}^{2}} V^{-1} \partial_{3} F_{1}+i k_{2} \sigma_{1}^{-1} \partial_{1} \sigma_{1} V \frac{1}{\lambda^{2}+k_{2}^{2}} V^{-1} F_{1} . \tag{10.55}
\end{equation*}
$$

Finally, it is stressed that the coefficients of the operators $\sigma_{1} V$ and $\sigma_{1}^{-1} \partial_{1} \sigma_{1} V$ above can be calculated exactly from the knowledge of the exact eigenmodes. Indeed, the technique presented in appendix shows that the determination of the eigenmodes $\phi_{n}$ lies on the calculations of functions $\sigma_{1} \phi_{n}$ and $\sigma_{1}^{-1} \partial_{1} \sigma_{1} \phi_{n}$. The final expression of $F_{2}$ in terms of $F_{1}$ and $\partial_{3} F_{1}$ is then

$$
\begin{equation*}
F_{2}=U \frac{1}{\lambda^{2}+k_{2}^{2}} V^{-1} \partial_{3} F_{1}+i k_{2} W \frac{1}{\lambda^{2}+k_{2}^{2}} V^{-1} F_{1}, \quad U=\sigma_{1} V, \quad W=\sigma_{1}^{-1} \partial_{1} \sigma_{1} V . \tag{10.56}
\end{equation*}
$$

To conclude this section, it has been shown that the field components $F_{1}, \partial_{3} F_{1}$ and $F_{2}$ can be expressed exactly in a lamellar layer: $F_{1}$ and $\partial_{3} F_{1}$ are provided by equations (10.42) and (10.43) (where $F \rightarrow F_{1}, L F \rightarrow-i \partial_{3} F_{1}$, and $K=V \lambda^{2} V^{-1}$ ); $F_{2}$ is given by equation (10.56). These expressions are based on the knowledge of the eigenvectors $V$ of $K$, the eigenvalues $\lambda^{2}$ of $K$ and the operators $U=\sigma_{1} V$ and $W=\sigma_{1}^{-1} \partial_{1} \sigma_{1} V$. All these quantities can be calculated exactly, as shown in appendix.

Finally, it is stressed that the propagation constants inside the lamellar layer are the square roots of the eigenvalues of the operator $K$, i.e. $\pm \sqrt{\lambda^{2}}$. If there is some absorption, i.e. $\operatorname{Im} \sigma>0$, or if a small positive imaginary part is added to the frequency $\omega$ (the limit $\omega=\lim _{\eta \downarrow 0}(\omega+i \eta)$ is considered [12,13]), the eigenvalues $\lambda^{2}$ of $K$ cannot be purely real. It follows that the propagation constants $\pm \sqrt{\lambda^{2}}$ have non zero imaginary parts and, moreover, the half of them $\left(+\sqrt{\lambda^{2}}\right)$ have strictly positive imaginary part and the second half of them $\left(-\sqrt{\lambda^{2}}\right)$ have strictly negative imaginary part. In this case, the assumption (10.16) on the eigenvalues of $L$ is well justified.

### 10.5 Numerical algorithm

A general solution for numerical algorithm has been proposed by $\mathrm{L} . \mathrm{Li}$ in the case of modal methods of gratings [3]. This solution is based on the definition of $S$ or $R$ matrices which are well-conditioned.

### 10.5.1 $R$ matrix for a single lamellar layer

In this section, the expression of a $R$ matrix associated with a lamellar layer is established: it is defined by the relationship

$$
\left[\begin{array}{l}
F_{1}(0)  \tag{10.57}\\
F_{1}(h)
\end{array}\right]=R\left[\begin{array}{l}
F_{2}(0) \\
F_{2}(h)
\end{array}\right] .
$$

First, equation (10.43) is used to provide an expression of the solution of (10.49):

$$
\begin{align*}
& F_{1}(0)-V \exp \left[i \sqrt{\lambda^{2}} h\right] V^{-1} F_{1}(h)=-i V \frac{1}{\sqrt{\lambda^{2}}} V^{-1} \partial_{3} F_{1}(0)+i V \frac{1}{\sqrt{\lambda^{2}}} \exp \left[i \sqrt{\lambda^{2}} h\right] V^{-1} \partial_{3} F_{1}(h), \\
& F_{1}(h)-V \exp \left[i \sqrt{\lambda^{2}} h\right] V^{-1} F_{1}(0)=-i V \frac{1}{\sqrt{\lambda^{2}}} \exp \left[i \sqrt{\lambda^{2}} h\right] V^{-1} \partial_{3} F_{1}(0)+i V \frac{1}{\sqrt{\lambda^{2}}} V^{-1} \partial_{3} F_{1}(h) \tag{10.58}
\end{align*}
$$

This set of equations is written using $2 \times 2$ matrices:

$$
A\left[\begin{array}{l}
F_{1}(0)  \tag{10.59}\\
F_{1}(h)
\end{array}\right]=B\left[\begin{array}{l}
\partial_{3} F_{1}(0) \\
\partial_{3} F_{1}(h)
\end{array}\right]
$$

where

$$
A=\left[\begin{array}{lr}
1 & -V \exp \left[i \sqrt{\lambda^{2}} h\right] V^{-1}  \tag{10.60}\\
-V \exp \left[i \sqrt{\lambda^{2}} h\right] V^{-1} & 1
\end{array}\right]
$$

and

$$
B=\left[\begin{array}{lr}
-i V \frac{1}{\sqrt{\lambda^{2}}} V^{-1} & i V \frac{1}{\sqrt{\lambda^{2}}} \exp \left[i \sqrt{\lambda^{2}} h\right] V^{-1}  \tag{10.61}\\
-i V \frac{1}{\sqrt{\lambda^{2}}} \exp \left[i \sqrt{\lambda^{2}} h\right] V^{-1} & \\
i V \frac{1}{\sqrt{\lambda^{2}}} V^{-1}
\end{array}\right]
$$

Next, from (10.56), the field $\partial_{3} F_{1}$ is related to $F_{1}$ and $F_{2}$ from

$$
\begin{equation*}
\partial_{3} F_{1}=-i k_{2} V\left[\lambda^{2}+k_{2}^{2}\right] U^{-1} W \frac{1}{\lambda^{2}+k_{2}^{2}} V^{-1} F_{1}+V\left[\lambda^{2}+k_{2}^{2}\right] U^{-1} F_{2} \tag{10.62}
\end{equation*}
$$

Defining the two matrices

$$
C=\left[\begin{array}{llr}
-i k_{2} V\left[\lambda^{2}+k_{2}^{2}\right] U^{-1} W \frac{1}{\lambda^{2}+k_{2}^{2}} V^{-1} & 0  \tag{10.63}\\
0 & -i k_{2} V\left[\lambda^{2}+k_{2}^{2}\right] U^{-1} W \frac{1}{\lambda^{2}+k_{2}^{2}} V^{-1}
\end{array}\right]
$$

and

$$
D=\left[\begin{array}{lr}
V\left[\lambda^{2}+k_{2}^{2}\right] U^{-1} & 0  \tag{10.64}\\
0 & V\left[\lambda^{2}+k_{2}^{2}\right] U^{-1}
\end{array}\right]
$$

the relationship (10.59) becomes

$$
A\left[\begin{array}{l}
F_{1}(0)  \tag{10.65}\\
F_{1}(h)
\end{array}\right]=B C\left[\begin{array}{l}
F_{1}(0) \\
F_{1}(h)
\end{array}\right]+B D\left[\begin{array}{l}
F_{2}(0) \\
F_{2}(h)
\end{array}\right] .
$$

Finally, according to the definition (10.57), the $R$ matrix is given by

$$
\begin{equation*}
R=\frac{1}{A-B C} B D . \tag{10.66}
\end{equation*}
$$

It is stressed that the $R$ matrix is numerically stable. Indeed, exponential functions always have arguments such that they decrease and, when inverted, they are always added to wellconditioned functions. For example, in the matrix $[A-B C]$, one can check that the diagonal blocs are well-conditioned since they do not contain any exponential function, while the offdiagonal blocs decrease exponentially. When inverted, this matrix will have the same behaviour with well-conditioned diagonal blocs and exponentially decreasing off-diagonal blocs.

### 10.5.2 $R$ matrix for a stack of lamellar layers

Here, a system made of two lamellar layers is considered. The first layer is located between the planes $x_{3}=-h_{1}$ and $x_{3}=0$, and the second layer between the planes $x_{3}=0$ and $x_{3}=h_{2}$. Let $R_{1}$ and $R_{2}$ be the $R$ matrices associated with these layers:

$$
\left[\begin{array}{c}
F_{1}\left(-h_{1}\right)  \tag{10.67}\\
F_{1}(0)
\end{array}\right]=R_{1}\left[\begin{array}{c}
F_{2}\left(-h_{1}\right) \\
F_{2}(0)
\end{array}\right], \quad\left[\begin{array}{c}
F_{1}(0) \\
F_{1}\left(h_{2}\right)
\end{array}\right]=R_{2}\left[\begin{array}{c}
F_{2}(0) \\
F_{2}\left(h_{2}\right)
\end{array}\right] .
$$

Then, the $R$ matrix associated with the stack of the two layers is determined by eliminating the components $F_{1}(0)$ and $F_{2}(0)$ in the equations above. Denoting by $R_{1, i j}$ and $R_{2, i j}(i, j=1,2)$ the blocs of $R_{1}$ and $R_{2}$,

$$
R_{1}=\left[\begin{array}{ll}
R_{1,11} & R_{1,12}  \tag{10.68}\\
R_{1,21} & R_{1,22}
\end{array}\right], \quad R_{2}=\left[\begin{array}{ll}
R_{2,11} & R_{2,12} \\
R_{2,21} & R_{2,22}
\end{array}\right],
$$

the expression of $R$ is given by

$$
R=\left[\begin{array}{cc}
R_{1,11}-R_{1,12} \frac{1}{R_{1,22}-R_{2,11}} R_{1,21} & R_{1,12} \frac{1}{R_{1,22}-R_{2,11}} R_{2,12}  \tag{10.69}\\
-R_{2,21} \frac{1}{R_{1,22}-R_{2,11}} R_{1,21} & R_{2,22}-R_{2,21} \frac{1}{R_{1,22}-R_{2,11}} R_{2,12}
\end{array}\right] .
$$

Again, one can check that the algorithm is stable since the only inverted blocs are the diagonal ones, which are well-conditioned.

### 10.6 Numerical application

A simple numerical example is considered to put the exact modal method to the test. The structure is made of a set of rectangular rods with dielectric constant $\varepsilon_{1,1} / \varepsilon_{0}=12.96$ (corresponding to the index 3.6 of Si at optical wavelengths), width $w_{1,1}=0.28 d$ and height $h=d /(2 \sqrt{2})$, where $d$ is the spatial period of the grating (see figure 10.4). This lamellar grating


Figure 10.4: The considered structure for the numerical example: a single layer made of rectangular rods.
is illuminated by a plane wave with an incident angle $\theta^{i}=45^{\circ}$. The oscillating frequency is $\omega=2 \pi /\left(d \sqrt{\varepsilon_{0} \mu_{0}}\right)$, which corresponds to a wavelength equal to the spatial period $d$. The efficiency diffracted in the order zero, i.e. at the reflected angle $\theta^{r}=\theta^{i}=45^{\circ}$ is calculated for both $s$ and $p$-polarizations which correspond respectively to the electric and magnetic fields reduced to a single component along the invariance axis $x_{2}$. Each component of the electromagnetic field is described by a finite number $(2 n+1)$ of exact modes. Reflected efficiency in the order zero


Figure 10.5: Efficiency in the zero order for s-polarization (left panel) and p-polarization (right panel).
for different values of the number of exact modes $(2 n+1)$ is represented on figure 10.5. These curves show that these efficiencies differ from their converged value with less than one percent from $(2 n+1)=7$ in $s$-polarization and $(2 n+1)=5$ in $p$-polarization (the converged values are respectively 0.7323 and 0.9487 in $s$ and $p$ polarizations). This convergence is found to be faster than in the case of the modal method with Fourier basis [2] where an error smaller than one percent is obtained from $(2 n+1)=17$ and $(2 n+1)=27$ in $s$ and $p$ polarizations respectively (the method [2] contains all the techniques to improve the convergence of the truncated Fourier series [15]). This improvement of the convergence resulting from the use of the exact modes
becomes a significant advantage when three dimensional woodpile structures are considered [9] (the total number of modes $(2 n+1)^{2}$ can be reduced by a factor of 10 ).

### 10.7 Appendix. Calculation of the exact modes and eigenvalues

It is shown here how to determine exactly the eigenvalues and the eigenfunctions of the operator $K$ associated with a lamellar layer in a very general case. An analogous reasoning provides the ones of the operator associated with the others lamellar layers.

From the expression (10.50), every eigenvalue $\Lambda_{n}$ of the operator $K$ is either an eigenvalue of $K_{\varepsilon_{1}}$ or $K_{\mu_{1}}$. So, it is sufficient to determine the set of eigenvalues $\left\{\Lambda_{v_{1}, n} \mid n \in \mathbb{N}\right\}$ associated with the set of eigenfunctions $\left\{\phi_{v_{1}, n} \mid n \in \mathbb{N}\right\}$ of the scalar operator $K_{v_{1}}$, with $v_{1}=\varepsilon_{1}$ or $v_{1}=\mu_{1}$.

### 10.7.1 The equation satisfied by the exact eigenvalues

From the expression (10.50), the operator $K_{v_{1}}$ is the sum of $\omega^{2} \varepsilon_{1} \mu_{1}+\partial_{1} v_{1}^{-1} \partial_{1} v_{1}$ and $\partial_{2}^{2}$ : the first part is an operator of the single variable $x_{1}$ and the second part is an operator of the single variable $x_{2}$. Thus, we can perform a variable separation: every eigenfunction of $L_{v_{1}}$ can be written

$$
\begin{equation*}
\phi_{v_{1}, n}\left(x_{1}, x_{2}\right)=\phi_{n_{1}}^{(1)}\left(x_{1}\right) \phi_{n_{2}}^{(2)}\left(x_{2}\right) \quad n_{1}, n_{2} \in \mathbb{N}, \tag{10.70}
\end{equation*}
$$

where $\phi_{n_{1}}^{(1)}$ and $\phi_{n_{2}}^{(2)}$ are respectively eigenfunctions of the first and second operators which constitute $L_{v_{1}}$.

In the case of a lamellar grating which is invariant in the direction $x_{2}$, this direction of invariance is considered using the Fourier decomposition (10.5). Thus the eigenfunction of $\partial_{2}^{2}$ is just

$$
\begin{equation*}
\phi_{n_{2}}^{(2)}\left(x_{2}\right)=\exp \left[i k_{2} x_{2}\right] \quad k_{2} \in \mathbb{R}, \tag{10.71}
\end{equation*}
$$

and the integer $n_{2}$ plays no role (integer $n$ will be simply $n_{1}$ ). In the case of woodpile crystals, it is easy to verify that the plane-wave

$$
\begin{equation*}
\phi_{n_{2}}^{(2)}\left(x_{2}\right)=\exp \left\{i\left[k_{2}+2 \pi p\left(n_{2}\right) / d_{2}\right] x_{2} / d_{2}\right\} \quad p\left(n_{2}\right) \in \mathbb{Z} \tag{10.72}
\end{equation*}
$$

is an eigenfunction of the operator $\partial_{2}^{2}$ and satisfies the partial Bloch boundary condition (10.8) adapted for the variable $x_{2}$. Let $\Lambda_{n_{2}}^{(2)}$ be the associated eigenvalue. Then, from $(10.71,10.72)$,

$$
\begin{equation*}
\Lambda_{n_{2}}^{(2)}=-\left[k_{2}+2 \pi q\left(n_{2}\right) / d_{2}\right]^{2} \tag{10.73}
\end{equation*}
$$

Note that, for lamellar gratings and eigenfunctions (10.71), the integer $q\left(n_{2}\right)$ is set to zero.
The $x_{1}$-dependency of the eigenfunction (10.70) is determined using the usual transfer matrix $[16,17,18]$. Let $\lambda_{n_{1}}^{(1)}$ be the eigenvalue associated with $\phi_{n_{1}}^{(1)}$ :

$$
\begin{equation*}
\left[\omega^{2} \varepsilon_{1} \mu_{1}+\partial_{1} v_{1}^{-1} \partial_{1} v_{1}\right] \phi_{n_{1}}^{(1)}=\Lambda_{n_{1}}^{(1)} \phi_{n_{1}}^{(1)} . \tag{10.74}
\end{equation*}
$$

In order to obtain a set of first order differential equations, the following column vector is introduced

$$
F_{n_{1}}=\left[\begin{array}{c}
v_{1} \phi_{n_{1}}^{(1)}  \tag{10.75}\\
v_{1}^{-1} \partial_{1} v_{1} \phi_{n_{1}}^{(1)}
\end{array}\right] .
$$

Note that, from equation (10.74), the two components of this vector are continuous functions. Now, suppose that the unit cell of the considered lamellar layer is made of $J$ rods of width $w_{1, j}$, permittivity $\varepsilon_{1, j}$ and permeability $\mu_{1, j}, j=1,2, \cdots, J$ (figure 10.6): we denote by $v_{1, j}$ the value


Figure 10.6: A layer made of three rods per unit cell $(J=3)$ : the three rods have width $w_{1, j}$, permittivity $\varepsilon_{1, j}$ and permeability $\mu_{1, j}, j=1,2,3$.
of the function $v_{1}$ in the $\operatorname{rod} j, j=1,2, \cdots, J$. Then, from equation (10.74), the vector (10.75) satisfies [18]

$$
\begin{equation*}
F_{n_{1}}(d)=T_{1}\left(\Lambda_{n_{1}}^{(1)}\right) F_{n_{1}}(0), \tag{10.76}
\end{equation*}
$$

where

$$
\begin{gather*}
T_{1}(\Lambda)=T_{1, J}(\Lambda) T_{1, J-1}(\Lambda) \cdots T_{1,1}(\Lambda),  \tag{10.77}\\
T_{1, j}(\Lambda)=P_{1, j}\left(\Lambda, w_{1, j}\right),  \tag{10.78}\\
P_{1, j}(\Lambda, w)=\left[\begin{array}{lr}
\cos \left(\beta_{1, j} w\right) & v_{1, j} \beta_{1, j}^{-1} \sin \left(\beta_{1, j} w\right) \\
-v_{1, j}^{-1} \beta_{1, j} \sin \left(\beta_{1, j} w\right) & \cos \left(\beta_{1, j} w\right)
\end{array}\right],  \tag{10.79}\\
\beta_{1, j}=\sqrt{\omega^{2} \varepsilon_{1, j} \mu_{1, j}-\Lambda} \quad j=1,2, \cdots, J . \tag{10.80}
\end{gather*}
$$

Note that the four elements of each matrix $T_{1, j}$ only depend on $\beta_{1, j}^{2}$ : the expression (10.79) is independent of the definition of the square root (10.80). In addition to (10.76), the vector (10.75) has to satisfy the partial Bloch boundary condition (10.8) for the variable $x_{1}$ :

$$
\begin{equation*}
F_{n_{1}}(d)=\exp \left[i k_{1} d\right] F_{n_{1}}(0) . \tag{10.81}
\end{equation*}
$$

The combination of (10.76) and (10.81) implies that $\exp \left[i k_{1} d\right]$ is an eigenvalue of the matrix $T_{1}\left(\Lambda_{n_{1}}^{(1)}\right)$ : the equation

$$
\begin{equation*}
\operatorname{det}\left\{T_{1}\left(\Lambda_{n_{1}}^{(1)}\right)-\exp \left[i k_{1} d\right]\right\}=0 \tag{10.82}
\end{equation*}
$$

determines the eigenvalues $\Lambda_{n_{1}}^{(1)}$. This last equation can be simplified using the fact that $\operatorname{det} T_{1}=$ 1 (since, from (10.78), $\left.\operatorname{det} T_{1, j}=1, j=1,2, \cdots, J\right)$ : if $\exp \left[i k_{1} d\right]$ is an eigenvalue of $T_{1}$, then $\exp \left[-i k_{1} d\right]$ is also. Thus, the equation (10.82) is equivalent to

$$
\begin{equation*}
\operatorname{tr} T_{1}\left(\Lambda_{n_{1}}^{(1)}\right)-2 \cos \left[k_{1} d\right]=0, \tag{10.83}
\end{equation*}
$$

where $\operatorname{tr} T_{1}$ is the trace of matrix $T_{1}$. Once the eigenvalues $\Lambda_{n_{1}}^{(1)}$ are determined from (10.83), the associated eigenvectors $\phi_{n_{1}}^{(1)}$ are also obtained using the transfer matrix [17]: firstly, the eigenvector $F_{n_{1}}(0)$ in $\mathbb{C}^{2}$ (associated with the eigenvalue $\exp \left[i k_{1} d\right]$ ) of the matrix $T_{1}\left(\Lambda_{n_{1}}^{(1)}\right)$ is determined; secondly, the expression of $\phi_{n_{1}}^{(1)}$ in the rod $j$ can be deduced from

$$
\begin{equation*}
F_{n_{1}}\left(x_{1}\right)=P_{1, j}\left(\Lambda_{n_{1}}^{(1)}, x_{1}-x_{1, j}\right) F_{n_{1}}\left(x_{1, j-1}\right), \tag{10.84}
\end{equation*}
$$

where

$$
\begin{equation*}
x_{1,0}=0, \quad x_{1, j}=\sum_{q=1}^{j} w_{1, q} \quad j=1,2, \cdots, J . \tag{10.85}
\end{equation*}
$$

Finally the eigenvalues of the operator $K_{V_{1}}$ are

$$
\begin{equation*}
\lambda_{v_{1}, n}=\lambda_{n_{1}}^{(1)}+\lambda_{n_{2}}^{(2)}, \tag{10.86}
\end{equation*}
$$

whose the two parts are respectively given by (10.83) and (10.73), and the expression of associated eigenvectors is (10.70) whose the two parts are respectively given by (10.84) and (10.72). Concerning the functions of the operators $U=\sigma_{1} V$ and $W=\sigma_{1}^{-1} \partial_{1} \sigma_{1} V$ used in section 10.4 [see equation (10.56)], they are equal to the functions

$$
\begin{equation*}
\left(v_{1} \phi_{n_{1}}^{(1)}\right)\left(x_{1}\right) \phi_{n_{2}}^{(2)}\left(x_{2}\right), \quad\left(v_{1}^{-1} \partial_{1} v_{1} \phi_{n_{1}}^{(1)}\right)\left(x_{1}\right) \phi_{n_{2}}^{(2)}\left(x_{2}\right), \tag{10.87}
\end{equation*}
$$

where $n_{1}$ and $n_{2}$ are in $\mathbb{N}$, the expression of $v_{1} \phi_{n_{1}}^{(1)}$ and $v_{1}^{-1} \partial_{1} v_{1} \phi_{n_{1}}^{(1)}$ in the rod $j$ can be deduced from $(10.75,10.84)$ and the expression of $\phi_{n_{2}}^{(2)}$ is given by (10.72).

### 10.7.2 Real eigenvalues

Here, we suppose that the permittivity and permeability are real positive functions:

$$
\begin{equation*}
\varepsilon_{1}\left(x_{1}\right) \in \mathbb{R}, \quad \varepsilon_{+}>\varepsilon_{1}\left(x_{1}\right)>0 ; \quad \mu_{1}\left(x_{1}\right) \in \mathbb{R}, \quad \mu_{+}>\mu_{1}\left(x_{1}\right)>0 . \tag{10.88}
\end{equation*}
$$

Under these conditions, the operator $K_{V_{1}}$ is selfadjoint and its eigenvalues are real when the following inner product is used:

$$
\begin{equation*}
(\phi, \psi) \longrightarrow \frac{1}{d} \int_{0}^{d} v_{1}\left(x_{1}\right) \overline{\phi\left(x_{1}\right)} \boldsymbol{\psi}\left(x_{1}\right) d x_{1} . \tag{10.89}
\end{equation*}
$$

The only difficulty in the numerical determination of the eigenvalues (10.86) is to find the real numbers $\lambda_{n_{1}}^{(1)}$ which satisfy the transcendental equation (10.83).

Since the numbers $\Lambda_{n_{1}}^{(1)}$ are eigenvalues of the operator $\omega^{2} \varepsilon_{1} \mu_{1}+\partial_{1} v_{1}^{-1} \partial_{1} v_{1} \leq \omega^{2} \varepsilon_{+} \mu_{+}$, these numbers are on the semi-axis $\left(-\infty, \omega^{2} \varepsilon_{+} \mu_{+}\right]$. This property makes their numerical determination easier and provides a way to number them:

$$
\begin{equation*}
\omega^{2} \varepsilon_{+} \mu_{+} \geq \Lambda_{v_{1}, 1} \geq \Lambda_{v_{1}, 2} \cdots \geq \Lambda_{v_{1}, n} \geq \cdots \tag{10.90}
\end{equation*}
$$

However, two difficulties can occur in this numerical determination. We give herein the solutions we have adopted.

The first difficulty comes from the possibility for two consecutive numbers $\Lambda_{n_{1}}^{(1)}$ to be very close to each other. Our solution is to use an algorithm which determines the zeros of the function $\operatorname{tr} T_{1}(\Lambda)-2 \cos \left[k_{1} d\right]$ on the left side of equation (10.83) by taking into account this function together with its derivative with respect to $\Lambda$. If two numbers $\Lambda_{n_{1}}^{(1)}$ are very close one to each other, then the derivative is close to zero. Thus such algorithm needs to determine the function

$$
\begin{equation*}
\frac{d}{d \Lambda}\left\{\operatorname{tr} T_{1}(\Lambda)-2 \cos \left[k_{1} d\right]\right\}=\operatorname{tr} \frac{d T_{1}}{d \Lambda}(\lambda) \tag{10.91}
\end{equation*}
$$

The expression of the derivative of the matrix $T_{1}$ can be deduced from $(10.77,10.78)$ :

$$
\begin{equation*}
\frac{d T_{1}}{d \Lambda}=\frac{d T_{1, J}}{d \Lambda} T_{1, J-1} \cdots T_{1,1}+T_{1, J} \frac{d T_{1, J-1}}{d \Lambda} \cdots T_{1,1}+\cdots+T_{1, J} T_{1, J-1} \cdots \frac{d T_{1,1}}{d \Lambda} \tag{10.92}
\end{equation*}
$$

where, for $j=1,2, \cdots, J$, the derivative of matrices

$$
\frac{d T_{1, j}}{d \Lambda}=\frac{1}{2}\left[\begin{array}{ll}
a_{1, j} & b_{1, j}  \tag{10.93}\\
c_{1, j} & d_{1, j}
\end{array}\right]
$$

is given by

$$
\begin{align*}
& a_{1, j}=w_{1, j} \beta_{1, j}^{-1} \sin \left[\beta_{1, j} w_{1, j}\right], \\
& b_{1, j}=v_{1, j} \beta_{1, j}^{-3} \sin \left[\beta_{1, j} w_{1, j}\right]-v_{1, j} w_{1, j} \beta_{1, j}^{-2} \cos \left[\beta_{1, j} w_{1, j}\right], \\
& c_{1, j}=v_{1, j}^{-1} \beta_{1, j}^{-1} \sin \left[\beta_{1, j} w_{1, j}\right]+v_{1, j}^{-1} w_{1, j} \cos \left[\beta_{1, j} w_{1, j}\right],  \tag{10.94}\\
& d_{1, j}=w_{1, j} \beta_{1, j}^{-1} \sin \left[\beta_{1, j} w_{1, j}\right] .
\end{align*}
$$

The second difficulty comes from the possibility of numerical instabilities in the expressions ( $10.79,10.93$ ) since the numbers $\beta_{1, j}(10.80)$ can have non-vanishing imaginary part. A solution is to multiply the four coefficients of matrices $T_{1, j}$ and their derivative (10.93) by the number

$$
\begin{equation*}
N_{j}=\exp \left[-\left|\operatorname{Im}\left(\beta_{1, j}\right)\right| w_{1, j}\right] \quad j=1,2, \cdots, J, \tag{10.95}
\end{equation*}
$$

and the term $2 \cos \left[k_{1} d\right]$ wich appears in (10.83) by the product

$$
\begin{equation*}
N=N_{J} N_{J-1} \cdots N_{1} . \tag{10.96}
\end{equation*}
$$

### 10.7.3 Complex eigenvalues

Here, the permittivity and permeability can take any complex value: $v_{1, j}$ is in $\mathbb{C}$, where $v_{1}=$ $\varepsilon_{1}, \mu_{1}$ and $j=1,2, \cdots, J$. The operator $K_{v_{1}}$ is not selfadjoint and then, its eigenvalues are, in general, in the complex plane. The determination of these complex eigenvalues $\lambda_{n_{1}}^{(1)}$ which satisfy the equation (10.83) has been intensively studied using different methods [6, 7, 19].

We present here a method similar to the one presented in [7]: the complex eigenvalues are deduced from the real eigenvalues by an analytic continuation. However, our method differs from the one presented in [7] since we make varying the phase of the numbers $v_{1, j}$ instead of their imaginary part. We think that it is better to make varying the phase since, from that we have observed, it leaves invariant the generalization to the complex case

$$
\begin{equation*}
\operatorname{Re}\left(\Lambda_{v_{1}, 1}\right) \geq \operatorname{Re}\left(\Lambda_{v_{1}, 2}\right) \cdots \geq \operatorname{Re}\left(\Lambda_{v_{1}, p}\right) \cdots \tag{10.97}
\end{equation*}
$$

of the numbering used when the eigenvalues are real (10.90).
We define for all $t$ in $[0,1]$ the functions

$$
\begin{equation*}
\widetilde{v}_{1, j}(t)=\left|v_{1, j}\right| \exp \left[i t \arg \left(v_{1, j}\right)\right], \tag{10.98}
\end{equation*}
$$

where $\arg \left(v_{1, j}\right)$ is the phase of the complex number $v_{1, j}, v_{1}=\varepsilon_{1}, \mu_{1}$ and $j=1,2, \cdots, J$. Substituting the numbers $v_{1, j}$ (where $v_{1}=\varepsilon_{1}, \mu_{1}$ ) for $\widetilde{v}_{1, j}(t)$ in equations (10.77,10.78), we obtain the matrix $\widetilde{T}_{1}(\Lambda, t)$. For each value of $t$, we define the numbers $\widetilde{\Lambda}_{n_{1}}^{(1)}(t)$ which satisfy

$$
\begin{equation*}
\operatorname{tr} \widetilde{T}_{1}\left[\widetilde{\Lambda}_{n_{1}}^{(1)}(t), t\right]-2 \cos \left[k_{1} d\right]=0 . \tag{10.99}
\end{equation*}
$$

Then, the numbers $\widetilde{\Lambda}_{n_{1}}^{(1)}(1)$ are the desired complex eigenvalues $\Lambda_{p_{1}}^{(1)}$ and the numbers $\widetilde{\Lambda}_{n_{1}}^{(1)}(0)$ are real eigenvalues which can be determined using the method presented in the previous section 10.7.2. Assuming that $\widetilde{\Lambda}_{n_{1}}^{(1)}(t)$ are continuous and differentiable functions of $t$, the complex numbers $\widetilde{\Lambda}_{n_{1}}^{(1)}(1)$ can be estimated from the numbers $\widetilde{\Lambda}_{n_{1}}^{(1)}(0)$ by a numerical integration [7] of

$$
\begin{equation*}
\frac{d \widetilde{\Lambda}_{n_{1}}^{(1)}}{d t}(t)=-\frac{\operatorname{tr}\left(\partial \widetilde{T}_{1} / \partial \Lambda\right)\left[\widetilde{\Lambda}_{n_{1}}^{(1)}(t), t\right]}{\operatorname{tr}\left(\partial \widetilde{T}_{1} / \partial t\right)\left[\widetilde{\Lambda}_{n_{1}}^{(1)}(t), t\right]}, \tag{10.100}
\end{equation*}
$$

where $\partial \widetilde{T}_{1} / \partial \Lambda$ is given by substituting the numbers $v_{1, j}$ for $\widetilde{v}_{1, j}(t)$ in equations $(10.94,10.93)$ and $\partial \widetilde{T}_{1} / \partial t$ is determined similarly. Finally the obtained estimates of numbers $\widetilde{\Lambda}_{n_{1}}^{(1)}(1)$ are used to initiate any of the classical methods for the numerical solution of equations [7]. Then, one obtains the desired complex eigenvalues.

In order to eliminate the numerical instabilities, one has to multiply each matrix $T_{1, j}$ and their derivatives by the numbers $N_{j}(10.95)$ as in the previous section 10.7.2.

### 10.7.4 Eigenfunctions

From (10.84), the expression of each eigenfunction $\phi_{n_{1}}^{(1)}$ is given by the coefficients of the column vectors $F_{n_{1}}\left(x_{1, j}\right), j=0,1, \cdots, J$. On the numerical side, the only difficulty comes from the fact that numerical instabilities in the expression of the transfer matrices $(10.78,10.79)$. A solution based on the $R$-matrix algorithm (or $S$-matrix) should consist in using the algorithm presented in [20] to obtain the vector $F_{n_{1}}\left(x_{1,0}\right)$ (and the vector $F_{n_{1}}\left(x_{1, J}\right)=\exp \left[k_{1} d\right] F_{n_{1}}\left(x_{1,0}\right)$ ) and then, the algorithm presented in [21, section V$]$ to obtain the vectors $F_{n_{1}}\left(x_{1, j}\right), j=1,2, \cdots, J-1$. However, we propose to use another solution which benefits of the fact that we deal with $2 \times 2$ matrices.

We define the following complex coefficients:

$$
\begin{gather*}
{\left[\begin{array}{cc}
\mathscr{T}_{11}^{j} & \mathscr{T}_{12}^{j} \\
\mathscr{T}_{21}^{j} & \mathscr{T}_{22}^{j}
\end{array}\right]=T_{1, J}\left(\Lambda_{n_{1}}^{(1)}\right) T_{1, J-1}\left(\Lambda_{n_{1}}^{(1)}\right) \cdots T_{1, j}\left(\Lambda_{n_{1}}^{(1)}\right),}  \tag{10.101}\\
{\left[\begin{array}{cc}
\tau_{11}^{j} & \tau_{12}^{j} \\
\tau_{21}^{j} & \tau_{22}^{j}
\end{array}\right]=T_{1, j}\left(\Lambda_{n_{1}}^{(1)}\right) T_{1, j-1}\left(\Lambda_{n_{1}}^{(1)}\right) \cdots T_{1,1}\left(\Lambda_{n_{1}}^{(1)}\right),}  \tag{10.102}\\
{\left[\begin{array}{c}
\mathscr{F}_{1}^{j} \\
\mathscr{F}_{2}^{j}
\end{array}\right]=F_{n_{1}}\left(x_{1, j}\right) \quad j=0,1, \cdots, \mathrm{~J} .} \tag{10.103}
\end{gather*}
$$

Since $F_{n_{1}}\left(x_{1,0}\right)$ is an eigenvector of the matrix $T_{1}\left(\Lambda_{n_{1}}^{(1)}\right)$ associated with the eigenvalue $\left.\exp \left[k_{1} d\right]\right)$, its coefficients satisfy

$$
\begin{equation*}
\mathscr{F}_{2}^{0}=-\frac{\mathscr{T}_{11}^{J} N-\exp \left[k_{1} d\right] N}{\mathscr{T}_{12}^{J} N} \mathscr{F}_{1}^{0}, \tag{10.104}
\end{equation*}
$$

where the numbers $\mathscr{T}_{11}^{J} N$ and $\mathscr{T}_{12}^{J} N$ are obtained by multiplying each coefficient of matrices $T_{1, j}\left(\Lambda_{n_{1}}^{(1)}\right)$ by the number $N_{j}$. The coefficients $\mathscr{F}_{1}^{J}$ and $\mathscr{F}_{2}^{J}$ are deduced from $(10.81,10.104)$ and then, one can obtain the other coefficients for $j=1,2, \cdots, J-1$ :

$$
\begin{align*}
& \mathscr{F}_{1}^{j}=\frac{\mathscr{T}_{22}^{j+1} \tau_{11}^{j} N}{\mathscr{T}_{21}^{j+1} \tau_{11}^{j} N+\tau_{21}^{j} \mathscr{T}_{22}^{j+1} N}\left(\frac{\mathscr{F}_{2}^{J}}{\mathscr{T}_{22}^{j+1}}-\frac{\mathscr{F}_{2}^{0}}{\tau_{11}^{j}}\right), \\
& \mathscr{F}_{2}^{j}=\frac{\mathscr{T}_{11}^{j+1} \tau_{22}^{j} N}{\mathscr{T}_{11}^{j+1} \tau_{12}^{j} N+\tau_{22}^{j} \mathscr{T}_{12}^{j+1} N}\left(\frac{\mathscr{F}_{1}^{J}}{\mathscr{T}_{11}^{j+1}}-\frac{\mathscr{F}_{1}^{0}}{\tau_{22}^{j}}\right), \tag{10.105}
\end{align*}
$$

where, as in (10.104), the multiplication by the number $N$ consists in multiplying each coefficient of matrices $T_{1, j}\left(\Lambda_{n_{1}}^{(1)}\right)$ by the number $N_{j}$.

Finally these functions have to be normalized. From the definition (10.89) of the inner product, one has to compute

$$
\begin{equation*}
\left\|\phi_{n_{1}}^{(1)}\right\|_{v_{1}}^{2}=\frac{1}{d} \int_{0}^{d}\left|\phi_{n_{1}}^{(1)}\left(x_{1}\right)\right|^{2} v_{1}\left(x_{1}\right) d x_{1} \tag{10.106}
\end{equation*}
$$

when the functions $\varepsilon_{1}$ and $\mu_{1}$ have the property (10.88). In the general case (where $\varepsilon$ and $\mu$ are complex valued functions), one has to use the formalism presented in [8, section 2.3]. It is possible to compute analytically the expression (10.106):

$$
\begin{align*}
\left\|\phi_{n_{1}}^{(1)}\right\|_{v_{1}}^{2}=\frac{1}{2 d} \sum_{j=1}^{J} & \frac{w_{1, j}}{v_{1, j}}\left(\left|\mathscr{F}_{1}^{j-1}\right|^{2}+\beta_{1, j}^{-2} v_{1, j}^{-2}\left|\mathscr{F}_{2}^{j-1}\right|^{2}\right)  \tag{10.107}\\
& -\beta_{1, j}^{-2} \operatorname{Re}\left(i \overline{\mathscr{F}}_{1}^{j-1} \mathscr{F}_{2}^{j-1}-i \overline{\mathscr{F}}_{1}^{j} \mathscr{F}_{2}^{j}\right) .
\end{align*}
$$

This expression allows to eliminate the numerical instabilities which can occur from the exponential functions. Note that all the coefficients of matrices defined in sections 10.4 and 10.5 (matrices $U, V$ and $W$ ) can be also computed analytically in order to eliminate the numerical instabilities.

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## Chapter 11:

Homogenization Techniques for Periodic Structures

Sebastien Guenneau,<br>Richard Craster,<br>Tryfon Antonakakis,<br>Kirill Cherednichenko, and<br>Shane Cooper

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## Chapter 11

# Homogenization Techniques for Periodic Structures 

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### 11.1 Introduction

In this chapter we describe a selection of mathematical techniques and results that suggest interesting links between the theory of gratings and the theory of homogenization, including a brief introduction to the latter. By no means do we purport to imply that homogenization theory is an exclusive method for studying gratings, neither do we hope to be exhaustive in our choice of topics within the subject of homogenization. Our preferences here are motivated most of all by our own latest research, and by our outlook to the future interactions between these two subjects. We have also attempted, in what follows, to contrast the "classical" homogenization (Section Ш.।2), which is well suited for the description of composites as we have known them since their advent until about a decade ago, and the "non-standard" approaches, high-frequency homogenization (Section $\amalg .2$ ) and high-contrast homogenization (Section 凹.3), which have been developing in close relation to the study of photonic crystals and metamaterials, which exhibit properties unseen in conventional composite media, such as negative refraction allowing for super-lensing through a flat heterogeneous lens, and cloaking, which considerably reduces the scattering by finite size objects (invisibility) in certain frequency range. These novel electromagnetic paradigms have renewed the interest of physicists and applied mathematicians alike in the theory of gratings [四].

### 11.1.1 Historical survey on homogenization theory

The development of theoretical physics and continuum mechanics in the second half of the 19th and first half of the 20th century has motivated the question of justifying the macrosopic view of physical phenomena (at the scales visible to the human eye) by "upscaling" the implied microscopic rules for particle interaction at the atomic level through the phenomena at the intermediate, "mesoscopic", level (from tenths to hundreds of microns). This ambition has led to an extensive worldwide programme of research, which is still far from being complete as of now. Trying to give a very crude, but more or less universally applicable, approximation of the aim of this extensive activity, one could say that it has to do with developing approaches to averaging out in some way material properties at one level with the aim of getting a less detailed, but almost equally precise, description of the material response. Almost every word in the last sentence needs to be clarified already, and this is essentially the point where one
could start giving an overview of the activities that took place during the years to follow the great physics advances of a century ago. Here we focus on the research that has been generally referred to as the theory of homogenization, starting from the early 1970s. Of course, even at that point it was not, strictly speaking, the beginning of the subject, but we will use this period as a kind of reference point in this survey.

The question that a mathematician may pose in relation to the perceived concept of "averaging out" the detailed features of a heterogeneous structure in order to get a more homogeneous description of its behaviour is the following: suppose that we have the simplest possible linear elliptic partial differential equation (PDE) with periodic coefficients of period $\eta>0$. What is the asymptotic behaviour of the solutions to this PDE as $\eta \rightarrow 0$ ? Can a boundary-value problem be written that is satisfied by the leading term in the asymptotics, no matter what the data unrelated to material properties are? Several research groups became engaged in addressing this question about four decades ago, most notably those led by N. S. Bakhvalov, E. De Giorgi, J.-L. Lions, V. A. Marchenko, see [2], [3], [4], [5] for some of the key contributions of that period. The work of these groups has immediately led to a number of different perspectives on the apparently basic question asked above, which in part was due to the different contexts that these research groups had had exposure to prior to dealing with the issue of averaging. Among these are the method of multiscale asymptotic expansions (also discussed later in this chapter), the ideas of compensated compactness (where the contribution by L. Tartar and F. Murat [6], [7] has to be mentioned specifically), the variational method (also known as the " $\Gamma$ convergence"). These approaches were subsequently applied to various contexts, both across a range of mathematical setups (minimisation problems, hyperbolic equations, problems with singular boundaries) and across a number of physical contexts (elasticity, electromagnetism, heat conduction). Some new approaches to homogenization appeared later on, too, such as the method of two-scale convergence by G. Nguetseng [ 8 ] and the periodic unfolding technique by D. Cioranescu, A. Damlamian and G. Griso [9]. Established textbooks that summarise these developments in different time periods, include, in addition to the already cited book [4], the monographs [10], [IT], [12], and more recently [13]]. The area that is perhaps worth a separate mention is that of stochastic homogenization, where some pioneering contributions were made by S. M. Kozlov [14]], G. C. Papanicolaou and S. R. S. Varadhan [15], and which has in recent years been approached with renewed interest.

A specific area of interest within the subject of homogenization that has been rapidly developing during the last decade or so is the study of the behaviour of "non-classical" periodic structures, which we understand here as those for which compactness of bounded-energy solution sequences fails to hold as $\eta \rightarrow 0$. The related mathematical research has been strongly linked to, and indeed influenced by, the parallel development of the area of metamaterials and their application in physics, in particular for electromagnetic phenomena. Metamaterials can be roughly defined as those whose properties at the macroscale are affected by higher-order behaviour as $\eta \rightarrow 0$. For example, in classical homogenization for elliptic second-order PDE one requires the leading ("homogenised solution") and the first-order ("corrector") terms in the $\eta$ -power-series expansion of the solution in order to determine the macroscopic properties, which results in a limit of the same type as the original problem, where the solution flux ("stress" in elasticity, "induction" in electromagnetics, "current" in electric conductivity, "heat flux" in heat conduction) depends on the solution gradient only ("strain" in elasticity, "field" in electromagnetics, "voltage" in electric conductivity, "temperature gradient" in heat condiction). If, however, one decides for some reason, or is forced by the specific problem setup, to include higher-order terms as well, they are likely to have to deal with an asymptotic limit of a different
type for small $\eta$, which may, say, include second gradients of the solution in its constitutive law. One possible reason for the need to include such unusual effects is the non-uniform (in $\eta$ ) ellipticity of the original problems or, using the language of materials science, the high-contrast in the material properties of the given periodic structure. Perhaps the earliest mathematical example of such degeneration is the so-called "double-porosity model", which was first considered by G. Allaire [16] and T. Arbogast, J. Douglas, U. Hornung [[17] in the early 1990s. A detailed analysis of the properties of double-porosity models, including their striking spectral behaviour did not appear until the work [[18] by V. V. Zhikov. We discuss the double-porosity model and its properties in more detail in Section [1.3.

Before moving on to the next section, it is important to mention one line of research within the homogenization area that has had a significant rôle in terms of application of mathematical analysis to materials, namely the subject of periodic singular structures (or "multi-structures", see [14]). While this subject is clearly linked to the general analysis of differential operators on singular domains (see [20]), there has been a series of works that develop specifically homogenization techniques for periodic structures of this kind (also referred to as "thin structures" in this context), e.g. [21]], [22]. It turns out that overall properties of such materials are similar to those of materials with high contrast. In the same vein, it is not difficult to see that compactness of bounded-energy sequences for problems on periodic thin structures does not hold (unless the sequence in question is suitably rescaled), which leads to the need for non-classical, higher-order, techniques in their analysis.

### 11.1.2 Multiple scale method: Homogenization of microstructured fibers



Figure 11.1: A diagram of the homogenization process: when the parameter $\eta$ gets smaller $\left(\eta<\eta^{\prime}\right)$, the number of cells inside the fixed domain $\Omega_{f}$ becomes larger. When $\eta \ll 1, \Omega_{f}$ is filled with a large number of small cells, and can thus be considered as an effective (or homogenized) medium. Such a medium is usually described by anisotropic parameters depending upon the resolution of auxiliary ("unit cell") problems set on the rescaled microcopic cell $Y$ which typically contains one inclusion D.

Let us consider a doubly periodic grating of pitch $\eta$ and finite extent such as shown in Fig.ㅍ.】. An interesting problem to look at is that of transverse electric (TE) modes- when the magnetic field has the form $(0,0, H)$ - propagating within a micro-structured fiber with infinite
conducting walls. Such an eigenvalue problem is known to have a discrete spectrum: we look for eigenfrequencies $\omega$ and associated eigenfields $H$ such that:

$$
\left(\mathscr{P}_{\eta}\right): \begin{cases}-\sum_{i, j=1}^{2} \frac{\partial}{\partial x_{i}}\left(\varepsilon_{i j}^{-1}\left(\frac{\mathbf{x}}{\eta}\right) \frac{\partial H(\mathbf{x})}{\partial x_{j}}\right)=\omega^{2} \mu_{0} \varepsilon_{0} H(\mathbf{x}) & \text { in } \Omega_{f} \\ \varepsilon_{i j}^{-1}\left(\frac{\mathbf{x}}{\eta}\right) \frac{\partial H(\mathbf{x})}{\partial x_{i}} n_{j}=0 & \text { on } \partial \Omega_{f}\end{cases}
$$

where we use the convention $\mathbf{x}=\left(x_{1}, x_{2}\right), \partial \Omega_{f}$ denotes the boundary $\Omega_{f}$, and $\mathbf{n}=\left(n_{1}, n_{2}\right)$ is the normal to the boundary. Here, $\varepsilon_{0} \mu_{0}=c^{-2}$ where $c$ is the speed of light in vacuum and we assume that matrix coefficients of relative permittivity $\varepsilon_{i j}(\mathbf{y})$, with $i, j=1,2$, are real, symmetric (with the convention $\mathbf{y}=\left(y_{1}, y_{2}\right)$ ), of period 1 (in $y_{1}$ et $y_{2}$ ) and satisfy:

$$
\begin{equation*}
M|\boldsymbol{\xi}|^{2} \geq \varepsilon_{i j}(\mathbf{y}) \xi_{i} \xi_{j} \geq m|\boldsymbol{\xi}|^{2}, \forall \boldsymbol{\xi} \in \mathbb{R}^{2}, \forall \mathbf{y} \in Y=[0,1]^{2}, \tag{11.1}
\end{equation*}
$$

where $|\boldsymbol{\xi}|^{2}=\left(\xi_{1}^{2}+\xi_{2}^{2}\right)$, for given strictly positive constants $M$ and $m$. This condition is met for all conventional dielectric media ${ }^{\text {WI }}$.
We can recast $\left(\mathscr{P}_{\eta}\right)$ as follows:

$$
-\frac{\partial}{\partial x_{i}} \sigma^{i}(H(\mathbf{x}))=\frac{\omega^{2}}{c^{2}} H(\mathbf{x})
$$

with

$$
\sigma^{i}(H(\mathbf{x}))=\varepsilon_{i j}^{-1}\left(\frac{\mathbf{x}}{\eta}\right) \frac{\partial H(\mathbf{x})}{\partial x_{j}} .
$$

The multiscale method relies upon the following ansatz:

$$
\begin{equation*}
H=H_{0}(\mathbf{x})+\eta H_{1}(\mathbf{x}, \mathbf{y})+\eta^{2} H_{2}(\mathbf{x}, \mathbf{y})+\ldots \tag{11.2}
\end{equation*}
$$

where $H_{i}(\mathbf{x}, \mathbf{y}), i=1,2, \ldots$ is a periodic function of period $Y$ in $\mathbf{y}$.
In order to proceed with the asymptotic algorithm, one needs to rescale the differential operator as follows

$$
\begin{equation*}
\frac{\partial H}{\partial x_{i}}=\left(\frac{\partial H_{0}}{\partial z_{i}}+\frac{\partial H_{1}}{\partial y_{i}}\right)+\eta\left(\frac{\partial H_{1}}{\partial z_{i}}+\frac{\partial H_{2}}{\partial y_{i}}\right)+\ldots \tag{11.3}
\end{equation*}
$$

where $\partial / \partial z_{i}$ stands for the partial derivative with respect to the $i$ th component of the macroscopic variable $\mathbf{x}$.
It is useful to set

$$
\sigma^{i}(H)=\sigma_{0}^{i}+\eta \sigma_{1}^{i}+\eta^{2} \sigma_{2}^{i}+\ldots
$$

what makes ([L.3) more compact.
Collecting coefficients sitting in front of the same powers of $\eta$, we obtain:

$$
\sigma_{0}^{i}(H)=\varepsilon_{i j}^{-1}(\mathbf{y})\left(\frac{\partial H_{0}}{\partial z_{i}}+\frac{\partial H_{1}}{\partial y_{i}}\right)
$$

[^4]$$
\sigma_{1}^{i}(H)=\varepsilon_{i j}^{-1}(\mathbf{y})\left(\frac{\partial H_{1}}{\partial z_{i}}+\frac{\partial H_{2}}{\partial y_{i}}\right)
$$
and so forth, all terms being periodic in $\mathbf{y}$ of period 1 .
Upon inspection of problem $\left(\mathscr{P}_{\eta}\right)$, we gather that
$$
-\left(\frac{1}{\eta} \frac{\partial}{\partial y_{i}}+\frac{\partial}{\partial z_{i}}\right)\left(\sigma_{0}^{i}+\eta \sigma_{1}^{i}+\ldots\right)=\frac{\omega^{2}}{c^{2}} H(\mathbf{x})+\ldots
$$
so that at order $\eta^{-1}$
$$
(\mathscr{A}):-\frac{\partial}{\partial y_{i}} \sigma_{0}^{i}=0
$$
and at order $\eta^{0}$
$$
(\mathscr{H}):-\frac{\partial}{\partial z_{i}} \sigma_{0}^{i}-\frac{\partial}{\partial y_{i}} \sigma_{1}^{i}=\frac{\omega^{2}}{c^{2}} H_{0}
$$
(the equations corresponding to higher orders in $\eta$ will not be used here).
Let us show that ( $\mathscr{H}$ ) provides us with an equation (known as the homogenized equation) associated with the macroscopic behaviour of the microstructured fiber. Its coefficients will be obtained thanks to $(\mathscr{A})$ which is an auxiliary problem related to the microscopic scale. We will therefore be able to compute $H_{0}$ and $H_{1}$ thus, in particular, the first terms of $H$ and $\sigma^{i}$.
In order to do so, let us introduce the mean on $Y$, which we denote $<$.$\rangle , which is an operator$ acting on the function $g$ of the variable $\mathbf{y}$ :
$$
<g>=\frac{1}{|Y|} \iint_{Y} g\left(y_{1}, y_{2}\right) d y_{1} d y_{2}
$$
where $|Y|$ is the area of the cell $Y$.
Applying the mean to both sides of $(\mathscr{H})$, we obtain:
$$
<(\mathscr{H})>:-\frac{\partial}{\partial z_{i}}<\sigma_{0}^{i}>-<\frac{\partial}{\partial y_{i}} \sigma_{1}^{i}>=\frac{\omega^{2}}{c^{2}} H_{0}<1>
$$
where we have used the fact that $<.>$ commutes with $\partial / \partial z_{i}$.
Moreover, invoking the divergence theorem, we observe that
$$
<\frac{\partial}{\partial y_{i}} \sigma_{1}^{i}>=\frac{1}{|Y|} \int_{Y} \frac{\partial}{\partial y_{i}} \sigma_{1}^{i}(\mathbf{y}) d \mathbf{y}=\frac{1}{|Y|} \int_{\partial Y} \sigma_{1}^{i}(\mathbf{y}) n_{i} d s
$$
where $\mathbf{n}=\left(n_{1}, n_{2}\right)$ is the unit outside normal to $\partial Y$ of $Y$. This normal takes opposite values on opposite sides of $Y$, hence the integral over $\partial Y$ vanishes.
Altogether, we obtain:
$$
<(\mathscr{H})>:-\frac{\partial}{\partial z_{i}}<\sigma_{0}^{i}>=\frac{\omega^{2}}{c^{2}} H_{0}
$$
which only involves the macroscopic variable $x$ and partial derivatives $\partial / \partial z_{i}$ with respect to the macroscopic variable. We now want to find a relation between $\left.<\sigma_{0}\right\rangle$ and the gradient in $\mathbf{x}$ of $H_{0}$. Indeed, we have seen that
$$
\sigma_{0}^{i}(H)=\varepsilon_{i j}^{-1}(\mathbf{y})\left(\frac{\partial H_{0}}{\partial z_{j}}+\frac{\partial H_{1}}{\partial y_{j}}\right)
$$
which from $(\mathscr{A})$ leads to
$$
(\mathscr{A} 1):-\frac{\partial}{\partial y_{i}}\left(\varepsilon_{i j}^{-1}(\mathbf{y}) \frac{\partial H_{1}}{\partial y_{j}}\right)=\left(\frac{\partial H_{0}}{\partial z_{j}}\right)\left(\frac{\partial}{\partial y_{i}} \varepsilon_{i j}^{-1}(\mathbf{y})\right) .
$$

We can look at $(\mathscr{A} 1)$ as an equation for the unknown $H_{1}(\mathbf{x}, \mathbf{y})$, periodic of period $Y$ in $\mathbf{y}$ and parametrized by $\mathbf{x}$. Such an equation is solved up to an additive constant. In addition to that, the parameter $\mathbf{x}$ is only involved via the factor $\partial H_{0} / \partial z_{j}$. Hence, by linearity, we can write the solution $H_{1}(\mathbf{x}, \mathbf{y})$ as follows:

$$
H_{1}(\mathbf{x}, \mathbf{y})=\frac{\partial H_{0}(\mathbf{x})}{\partial z_{j}} w^{j}(\mathbf{y}),
$$

where the two functions $w^{j}(\mathbf{y}), j=1,2$ are solutions to $(\mathscr{A} 1)$ corresponding to $\partial H_{0} / \partial z_{j}(\mathbf{x})$, $j=1,2$ equal to unity with the other ones being zero, that is solutions to:

$$
(\mathscr{A} 2):-\frac{\partial}{\partial y_{i}}\left(\varepsilon_{i j}^{-1}(\mathbf{y}) \frac{\partial w^{k}}{\partial y_{j}}\right)=\delta_{j k}\left(\frac{\partial}{\partial y_{i}} \varepsilon_{i j}^{-1}(\mathbf{y})\right),
$$

with $w^{k}(\mathbf{y}), k=1,2$ periodic functions in $\mathbf{y}$ of period $Y^{\square}$.
Since the functions $w^{k}(\mathbf{y})$ are known, we note that

$$
\sigma_{i}^{0}(\mathbf{x}, \mathbf{y})=\varepsilon_{i j}^{-1}(\mathbf{y})\left(\frac{\partial H_{0}}{\partial z_{j}}+\frac{\partial H_{1}}{\partial y_{j}}\right)=\varepsilon_{i j}^{-1}(\mathbf{y})\left(\frac{\partial H_{0}}{\partial z_{j}}+\frac{\partial H_{0}}{\partial z_{k}} \frac{\partial w^{k}(\mathbf{y})}{\partial y_{j}}\right),
$$

which can be written as

$$
\sigma_{0}^{i}(\mathbf{x}, \mathbf{y})=\left(\varepsilon_{i k}^{-1}(\mathbf{y})+\varepsilon_{i j}^{-1}(\mathbf{y}) \frac{\partial w^{k}(\mathbf{y})}{\partial y_{j}}\right) \frac{\partial H_{0}(\mathbf{x})}{\partial z_{k}} .
$$

Lets us now apply the mean to both sides of this equation. We obtain:

$$
<\sigma_{0}^{i}>(\mathbf{x})=\varepsilon_{\mathrm{hom}, i k}^{-1} \frac{\partial H_{0}(\mathbf{x})}{\partial z_{k}},
$$

which can be recast as the following homogenized problem:

$$
\left(\mathscr{P}_{0}\right): \begin{cases}-\sum_{i, k=1}^{2} \frac{\partial}{\partial x_{i}}\left(\varepsilon_{\text {hom }, i k}^{-1} \frac{\partial H_{0}(\mathbf{x})}{\partial x_{k}}\right)=\omega^{2} \mu_{0} \varepsilon_{0} H_{0}(\mathbf{x}) & \text {, in } \Omega_{f} \\ \varepsilon_{\text {hom }, i k}^{-1}\left(\frac{\mathbf{x}}{\eta}\right) \frac{\partial H_{0}(\mathbf{x})}{\partial x_{i}} n_{k}=0 & \text {,on } \partial \Omega_{f},\end{cases}
$$

where $\varepsilon_{\text {hom }, i k}^{-1}$ denote the coefficients of the homogenized matrix of permittivity given by:

$$
\begin{equation*}
\varepsilon_{\mathrm{hom}, i k}^{-1}=\frac{1}{|Y|} \iint_{Y}\left(\varepsilon_{i k}^{-1}(\mathbf{y})+\varepsilon_{i j}^{-1}(\mathbf{y}) \frac{\partial w^{k}(\mathbf{y})}{\partial y_{j}}\right) d \mathbf{y} . \tag{11.4}
\end{equation*}
$$

As an illustrative example for this homogenized problem, we consider a microstructured waveguide consisting of a medium with relative permittivity $\varepsilon=1.25$ with elliptic inclusions (of minor

[^5]

Figure 11.2: Potentials $V_{x}$ (left) and $V_{y}$ (right): The unit cell contains an elliptic inclusion of relative permittivity $(\varepsilon=4.0+3 i)$ with minor and major axis $a=0.3$ and $b=0.4$ in silica $(\varepsilon=1.25)$.
and major axes 0.3 cm and 0.4 cm respectively) with center to center spacing $d=0.1 \mathrm{~cm}$ with an infinite conducting boundary i.e. Neumann boundary conditions in the TE polarization.

We use the COMSOL MULTIPHYSICS finite element package to solve the annex problem and we find that $\left[\varepsilon_{\text {hom }}\right]$ from (IL.4) writes as [26]

$$
\left(\begin{array}{cc}
1.9296204 & -1.053308310^{-16} \\
-44.41744410^{-18} & 2.1127643
\end{array}\right)
$$

with $<\varepsilon>_{Y}=2.2867255$. The off diagonal terms can be neglected.
If we assume that the transverse propagating modes in the metallic waveguide have a small propagation constant $\gamma \ll 1$, the above mathematical model describes accurately the physics. We show in Fig. [.. 3 a comparison between two TE modes of the microstructured waveguide and its associated anisotropic homogenized counterpart. Both eigenfrequencies and eigenfields match well (note that we use the waveguide terminology wavenumber $k=\sqrt{\omega^{2} / c^{2}-\gamma^{2}}$ ).

### 11.1.3 The case of one-dimensional gratings: Application to invisibility cloaks

There is a case of particular importance for applications in grating theory: that of a periodic multilayered structure. Let us assume that the permittivity of this medium is $\varepsilon=\alpha$ in white layers and $\beta$ in yellow layers, as shown in Fig. Ш.4.
Equation $(\mathscr{A} 2)$ takes the form:

$$
(\mathscr{A} 3):-\frac{d}{d y}\left(\varepsilon^{-1}(y) \frac{d w}{\partial y}\right)=\left(\frac{d}{d y} \varepsilon^{-1}(y)\right)
$$

with $w(y)$, periodic function in $y$ of period 1 .
We deduce that

$$
-\frac{d w}{d y}=1+C \varepsilon(y)
$$

Noting that $\int_{Y} \frac{d w}{d y}=w(1)-w(0)=0$, this leads to

$$
\int_{Y}(1+C \varepsilon(y)) d y=0
$$

Since $|Y|=1$, we conclude that

$$
C=-\langle\varepsilon\rangle^{-1}
$$



Figure 11.3: Comparison between transverse electric fields $T E_{21}$ and $T E_{31}$ of a microstructured metallic waveguide for a propagation constant $\gamma=0.1 \mathrm{~cm}^{-1}$ (wavenumbers $k=0.7707 \mathrm{~cm}^{-1}$ and $k=0.5478 \mathrm{~cm}^{-1}$ respectively), see left panel, with the $T E_{21}$ and $T E_{31}$ modes of the corresponding homogenized anisotropic metallic waveguide for $\gamma=0.1 \mathrm{~cm}^{-1}\left(k=0.7607 \mathrm{~cm}^{-1}\right.$ and $k=0.5201 \mathrm{~cm}^{-1}$, where $k=\sqrt{\omega^{2} / c^{2}-\gamma^{2}}=\sqrt{\omega^{2} \varepsilon_{0} \mu_{0}-\gamma^{2}}$ were obtained from the computation of eigenvalues $\omega$ of homogenized problem $\left(\mathscr{P}_{0}\right)$ ), see right panel.

The homogenized permittivity takes the form:

$$
\begin{aligned}
\varepsilon_{\text {hom }}^{-1} & =\frac{1}{|Y|} \int_{Y}\left(\varepsilon^{-1}(y)+\varepsilon^{-1}(y) \frac{d w(y)}{d y}\right) d y \\
& =<\varepsilon^{-1}(y)>-<\varepsilon^{-1}(y)+C> \\
& =<\varepsilon^{-1}(y)>-<\varepsilon^{-1}(y)>+\ll \varepsilon(y)>^{-1}>=<\varepsilon(y)>^{-1}
\end{aligned}
$$

We note that if we now consider the full operator i.e. we include partial derivatives in $y_{1}$ and $y_{2}$, the anisotropic homogenized permittivity takes the form:

$$
\varepsilon_{\mathrm{hom}}^{-1}=\left(\begin{array}{cc}
<\varepsilon(y)^{-1}> & 0 \\
0 & <\varepsilon(y)>^{-1}
\end{array}\right),
$$

as the only contribution for $\varepsilon_{\text {hom, } 11}^{-1}$ is $1 /|Y| \int_{Y} \varepsilon^{-1}(y) d y$.
As an illustrative example of what artificial anisotropy can achieve, we propose the design of an invisibility cloak. For this, let us assume that we have a multilayered grating with periodicity along the radial axis. In the coordinate system $(r, \theta)$, the homogenized permittivity clearly has the same form as above. If we want to design an invisibility cloak with an alternation of two homogeneous isotropic layers of thicknesses $d_{A}$ and $d_{B}$ and permittivities $\alpha, \beta$, we then need to use the formula

$$
\frac{1}{\varepsilon_{r}}=\frac{1}{1+\eta}\left(\frac{1}{\alpha}+\frac{\eta}{\beta}\right), \quad \varepsilon_{\theta}=\frac{\alpha+\eta \beta}{1+\eta}
$$





Figure 11.4: Schematic of homogenization process for a one-dimensional grating with homogeneous dielectric layers of permittivity $\alpha$ and $\beta$ in white and yellow regions. When $\eta$ tends to zero the number of layers tends to infinity, and their thicknesses vanish, in such a way that the width of the overall stack remains constant.
where $\eta=d_{B} / d_{A}$ is the ratio of thicknesses for layers $A$ and $B$ and $d_{A}+d_{B}=1$.
We now note that the coordinate transformation $r^{\prime}=R_{1}+r \frac{R_{2}-R_{1}}{R_{2}}$ can compress a disc $r<R_{2}$ into a shell $R_{1}<r<R_{2}$, provided that the shell is described by the following anisotropic heterogeneous permittivity [27] $\underline{\underline{c}}^{\text {cloak }}$ (written in its diagonal basis):

$$
\begin{equation*}
\varepsilon_{r}^{\text {cloak }}=\left(\frac{R_{2}}{R_{2}-R_{1}}\right)^{2}\left(\frac{r^{\prime}-R_{1}}{r^{\prime}}\right)^{2}, \quad \varepsilon_{\theta}^{\text {cloak }}=\left(\frac{R_{2}}{R_{2}-R_{1}}\right)^{2}, \tag{11.5}
\end{equation*}
$$

where $R_{1}$ and $R_{2}$ are the interior and the exterior radii of the cloak. Such a metamaterial can be approximated using the formula (Ш.5), as first proposed in [28], which leads to the multilayered cloak shown in Fig. [Ш.5.

### 11.2 High-frequency homogenization

Many of the features of interest in photonic crystals [44, 45], or other periodic structures, such as all-angle negative refraction [46, 47, 48, 49] or ultrarefraction [50, 51] occur at high frequencies where the wavelength and microstructure dimension are of similar orders. Therefore the conventional low-frequency classical homogenisation clearly fails to capture the essential physics and a different approach to distill the physics into an effective model is required. Fortunately a high frequency homogenisation (HFH) theory as developed in [37] is capable of capturing features such as AANR and ultra-refraction [52] for some model structures. Somewhat tangentially, there is an existing literature in the analysis community on Bloch homogenisation [53, 54, 55, 56], that is related to what we call high frequency homogenisation. There is also a flourishing literature on developing homogenised elastic media, with frequency dependent effective parameters, based upon periodic media [38]. There is therefore considerable interest in creating effective continuum models of microstructured media that break free from the conventional low frequency homogenisation limitations.
(a)

(b)


Figure 11.5: Propagation of a plane wave of wavelength $710^{-7} m$ (red in the visible spectrum) from the left on a multilayered cloak of inner radius $R_{1}=1.510^{-8} m$ and outer radius $R_{2}=310^{-8} m$, consisting of 20 homogeneous layers of equal thickness and of respective relative permittivities $1680.70,0.25,80.75,0.25,29.39,0.25,16.37,0.25,10.99,0.25,8.18,0.25,6.50,0.25,5.40$, $0.25,4.63,0.25,4.06,0.25$ in vacuum. Importantly, one layer in two has the same permittivity.

### 11.2.1 High Frequency Homogenization for Scalar Waves

Waves propagating through photonic crystals and metamaterials have proven to show different effects depending on their frequency. The homogenization of a periodic material is not unique. The effective properties of a periodic medium change depending on the vibration modes within its cells. The dispersion diagram structure can be considered to be the identity of such a material and provides the most important information regarding group velocities, band-gaps of disallowed propagation frequency bands, Dirac cones and many other interesting effects. The goal of a homogenization theory is to provide an effective homogeneous medium that is equivalent, in the long scale, to the initial non-homogeneous medium composed of a short-scale periodic, or other microscale, structure. This was achieved initially using the classical theory of homogenization [4, 34, [1, 35, 36] and yields an intuitively obvious result that the effective medium's properties consist of simple averages of the original medium's properties. This is valid so long as the wavelength is very large compared to the size of the cells (here we focus on periodic media created by repeating cells). For shorter wavelengths of the order of a cell's length a more general theory has been developed [37] that also recovers the results of the classical homogenization theory. For clarity we present high frequency homogeniaztion (HFH) by means of an illustrative example and consider a two-dimensional lattice geometry for TE or TM polarised electromagnetic waves. With harmonic time dependence, $\exp (-i \Omega t)$ (assumed understood and henceforth suppressed), the governing equation is the scalar Helmholtz equation,

$$
\begin{equation*}
\nabla^{2} u+\Omega^{2} u=0, \tag{11.6}
\end{equation*}
$$

where $u$ represent $E_{Z}$ and $H_{Z}$, for TM and TE polarised electromagnetic waves respectively, and $\Omega^{2}=n^{2} \omega^{2} / c^{2}$. In our example the cells are square and each square cell of length $2 l$ contains a circular hole and the filled part of the cell has constant non-dimensionalized properties. The boundary conditions on the hole's surface, namely the boundary $\partial S_{2}$, depend on the polarisation and are taken to be either of Dirichlet or Neumann type. This approach assumes infinite conducting boundaries which is a good approximation for micro-waves. We adopt a multiscale approach where $l$ is the small length scale and $L$ is a large length scale and we set $\eta=l / L \ll 1$


Figure 11.6: Panel (a) An infinite square array of split ring resonators with the elementary cell shown as the dashed line inner square. Panel (b) shows the irreducible Brillouin zone, in wavenumber space, used for square arrays in perfectly periodic media based around the elementary cell shown of length $2 l(l=1$ in (b)). Figure reproduced from Proceedings of the Royal Society [40].
to be the ratio of these scales. The two length scales let us introduce the following two independent spatial variables, $\xi_{i}=x_{i} / l$ and $X_{i}=x_{i} / L$. The cell's reference coordinate system is then $-1<\xi<1$. By introducing the new variables in equation (Ш.6) we obtain,

$$
\begin{equation*}
u(\mathbf{X}, \boldsymbol{\xi}), \xi_{i} \xi_{i}+\Omega^{2} u(\mathbf{X}, \boldsymbol{\xi})+2 \eta u(\mathbf{X}, \boldsymbol{\xi}), \xi_{i} X_{i}+\eta^{2} u(\mathbf{X}, \boldsymbol{\xi}), x_{i} X_{i}=0 \tag{11.7}
\end{equation*}
$$

We now pose an ansatz for the field and the frequency,

$$
\begin{array}{r}
u(\mathbf{X}, \boldsymbol{\xi})=u_{0}(\mathbf{X}, \boldsymbol{\xi})+\eta u_{1}(\mathbf{X}, \boldsymbol{\xi})+\eta^{2} u_{2}(\mathbf{X}, \boldsymbol{\xi})+\ldots, \\
\Omega^{2}=\Omega_{0}^{2}+\eta \Omega_{1}^{2}+\eta^{2} \Omega_{2}^{2}+\ldots \tag{11.8}
\end{array}
$$

In this expansion we set $\Omega_{0}$ to be the frequency of standing waves that occur in the perfectly periodic setting. By substituting equations (Ш..8) into equation (Ш.7) and grouping equal powers of $\varepsilon$ through to second order, we obtain a hierarchy of three ordered equations:

$$
\begin{gather*}
u_{0, \xi_{i} \xi_{i}}+\Omega_{0}^{2} u_{0}=0,  \tag{11.9}\\
u_{1, \xi_{i} \xi_{i}}+\Omega_{0}^{2} u_{1}=-2 u_{0, \xi_{i} X_{i}}-\Omega_{1}^{2} u_{0},  \tag{11.10}\\
u_{2, \xi_{i} \xi_{i}}+\Omega_{0}^{2} u_{2}=-u_{0, X_{i} X_{i}}-2 u_{1, \xi_{i} X_{i}}-\Omega_{1}^{2} u_{1}-\Omega_{2}^{2} u_{0} . \tag{11.11}
\end{gather*}
$$

These equations are solved as in [40, 37] and hence the description is brief.
The asymptotic expansions are taken about the standing wave frequencies that occur at the corners of the irreducible Brillouin zone depicted in Fig. Ш.6. It should be noted that not all structured cells will have the usual symmetries of a square, as in Fig. [1.6(a) where there is no reflexion symmetry from the diagonals. As a consequence the usual triangular region $\Gamma Х M$ does not always represent the irreducible Brillouin zone and the square region $\Gamma M X N$ should be used instead. Also paths that cross the irreducible Brillouin zone have proven to yield interesting effects namely along the path $M X^{\prime}$ for large circular holes [39].

The subsequent asymptotic development considers small perturbations about the points $\Gamma, X$ and $M$ so that the boundary conditions of $u$ on the outer boundaries of the cell, namely $\partial S_{1}$, read,

$$
\begin{equation*}
\left.u\right|_{\xi_{i}=1}= \pm\left. u\right|_{\xi_{i}=-1} \quad \text { and }\left.\quad u_{, \xi_{i}}\right|_{\xi_{i}=1}= \pm\left. u_{, \xi_{i}}\right|_{\xi_{i}=-1} \tag{11.12}
\end{equation*}
$$

where the,+- stand for periodic and anti-periodic conditions respectively: the standing waves occur when these conditions are met. The conditions on $\partial S_{2}$ are either of Dirichlet or Neumann type. The theory that follows is similar for both boundary condition cases, but the latter one is illustrated herein. Neumann boudary condition on the hole's surface or equivalently electromagnetic waves in TE polarization yield,

$$
\begin{equation*}
\frac{\partial u}{\partial \mathbf{n}}=u,\left.x_{i} n_{i}\right|_{\partial S_{2}}=0 . \tag{11.13}
\end{equation*}
$$

which in terms of the two-scales and $u_{i}(\mathbf{X}, \boldsymbol{\xi})$ become

$$
\begin{equation*}
U_{0, \xi_{i}} n_{i}=0, \quad\left(U_{0} f_{0, X_{i}}+u_{1, \xi_{i}}\right) n_{i}=0, \quad\left(u_{1, X_{i}}+u_{2, \xi_{i}}\right) n_{i}=0 . \tag{11.14}
\end{equation*}
$$

The solution of the leading order equation is by introducing the following separation of variables $u_{0}=f_{0}(\mathbf{X}) U_{0}\left(\boldsymbol{\xi} ; \Omega_{0}\right)$. It is obvious that $f_{0}(\mathbf{X})$, which represents the behaviour of the solution in the long scale, is not set by the leading order equation and the resulting eigenvalue problem is solved on the short-scale for $\Omega_{0}$ and $U_{0}$ representing the standing wave frequencies and the associated cell's vibration modes respectively. To solve the first order equation ([..U) we take the integral over the cell of the product of equation ([.LI) with $U_{0}$ minus the product of equation (Ш.9) with $u_{1} / f_{0}$ and this yields $\Omega_{1}=0$. It then follows to solve for $u_{1}(\mathbf{X}, \boldsymbol{\xi})=f_{0, X_{i}}(\mathbf{X}) U_{1_{i}}(\boldsymbol{\xi})$ where the vector $\mathbf{U}_{1}$ is found as in [40]. By invoking a similar solvability condition for the second order equation we obtain a second order PDE for $f_{0}(\mathbf{X})$,

$$
\begin{array}{r}
T_{i j} f_{0, X_{i} X_{j}}+\Omega_{2}^{2} f_{0}=0 \quad \text { where, } \\
T_{i j}=\frac{t_{i j}}{\iint_{S} U_{0}^{2} d S} \text { for } i, j=1,2 \tag{11.15}
\end{array}
$$

entirely on the long scale with the coefficients $T_{i j}$ containing all the information of the cell's dynamical response and the tensor $t_{i j}$ represents dynamical averages of the properties of the medium. For Neumann boundary conditions on $\partial S_{2}$ its formulation reads,

$$
\begin{gather*}
t_{i i}=\iint_{S} U_{0}^{2} d S+\iint_{S}\left(U_{1_{i}, \xi_{i}} U_{0}-U_{1_{i}} U_{0, \xi_{i}}\right) d S \quad \text { for } \quad i=1 \text { or } 2,  \tag{11.16}\\
t_{i j}=\iint_{S}\left(U_{1_{j}, \xi_{i}} U_{0}-U_{1_{j}} U_{0, \xi_{i}}\right) d S \quad \text { for } \quad i \neq j . \tag{11.17}
\end{gather*}
$$

Note that there is no summation over repeated indexes for $t_{i i}$. The tensor depends on the boundary conditions of the holes and has a different form if Dirichlet type conditions are applied on $\partial S_{2}$.

The PDE for $f_{0}$ has several uses, and can be verified by re-creating asymptotically the dispersion curves for a perfect lattice system. One important result of equation ([1.15) is its use in the expansion of $\Omega$ namely in equation (Ш.区). In order to obtain $\Omega_{2}$ as a function of the Bloch wavenumbers we use the Bloch boundary conditions on the cell to solve for $f_{0}(\mathbf{X})=$ $\exp \left(i \kappa_{j} X_{j} / \eta\right)$, where $\kappa_{j}=K_{j}-d_{j}$ with $d_{j}=0, \pi / 2,-\pi / 2$ depending on the location in the Brillouin zone. The asymptotic dispersion relation now reads,

$$
\begin{equation*}
\Omega \sim \Omega_{0}+\frac{T_{i j}}{2 \Omega_{0}} \kappa_{i} \kappa_{j} . \tag{11.18}
\end{equation*}
$$



Figure 11.7: The dispersion diagram for a doubly periodic array of square cells with circular inclusions, of radius 0.4, free at their inner boundaries shown for the irreducible Brillouin zone of Fig. W.6. The dispersion curves are shown in solid lines and the asymptotic solutions from the high frequency homogenization theory are shown in dashed lines. Figure reproduced from Proceedings of the Royal Society [40].

Equation ([1.]8) yields the behaviour of the dispersion curves asymptotically around the standing wave frequencies that are naturally located at the edge points of the Brillouin zone. Fig. [. 8 illustrates the asymptotic dispersion curves for the first six dispersion bands of a square cell geometry with circular holes.

An assumption in the development of equation (Ш.L8) is that the standing wave frequencies are isolated. But one can clearly see in Fig. $\amalg .7$ that this is not the case for third standing wave frequency at point $\Gamma$ as well as for the second standing wave frequency at point $X$. A small alteration to the theory [40] enables the computation of the dispersion curves at such points by setting,

$$
\begin{equation*}
u_{0}=f_{0}^{(l)}(\mathbf{X}) U_{0}^{(l)}\left(\boldsymbol{\xi} ; \Omega_{0}\right) \tag{11.19}
\end{equation*}
$$

where we sum over the repeated superscripts $(l)$. Proceeding as before, we multiply equation (Ш.I(l) by $U_{0}^{(m)}$, substract $u_{1}\left(\left(U_{0, \xi_{i}}^{(m)}\right) \xi_{i}+\Omega_{0}^{2} U_{0}^{(m)}\right)$ then integrate over the cell to obtain,

$$
\begin{equation*}
\left(\frac{\partial}{\partial X_{j}} \mathbf{A}_{j m l}+\Omega_{1}^{2} \mathbf{B}_{m l}\right) \hat{f}_{0}^{(l)}=0, \quad \text { for } \quad m=1,2, \ldots, p \tag{11.20}
\end{equation*}
$$

$\Omega_{1}$ is not necessarily zero, and

$$
\begin{equation*}
\mathbf{A}_{j m l}=\iint_{S}\left(U_{0}^{(m)} U_{0, \xi_{j}}^{(l)}-U_{0, \xi_{j}}^{(m)} U_{0}^{(l)}\right) d S, \quad \mathbf{B}_{m l}=\iint U_{0}^{(l)} U_{0}^{(m)} d S \tag{11.21}
\end{equation*}
$$

There is now a system of coupled partial differential equations for the $f_{0}^{(l)}$ and, provided $\Omega_{1} \neq 0$, the leading order behaviour of the dispersion curves near the $\Omega_{0}$ is now linear (these then form Dirac cones).

For the perfect lattice, we set $f_{0}^{(l)}=\hat{f}_{0}^{(l)} \exp \left(i \kappa_{j} X_{j} / \eta\right)$ and obtain the following index equations,

$$
\begin{equation*}
\left(i \frac{\kappa_{j}}{\eta} \mathbf{A}_{j m l}+\Omega_{1}^{2} \mathbf{B}_{m l}\right) \hat{f}_{0}^{(l)}=0, \quad \text { for } \quad \mathrm{m}=1,2, \ldots, \mathrm{p} \tag{11.22}
\end{equation*}
$$

The system of equation (Ш.22) can be written simply as,

$$
\begin{equation*}
\mathbf{C} \hat{\mathbf{F}}_{0}=0 \tag{11.23}
\end{equation*}
$$

with $\mathbf{C}_{l l}=\Omega_{1}^{2} \mathbf{B}_{l l}$ and $\mathbf{C}_{m l}=i \kappa_{j} \mathbf{A}_{j m l} / \eta$ for $l \neq m$. One must then solve for $\Omega_{1}^{2}= \pm \sqrt{\alpha_{i j} \kappa_{i} \kappa_{j}} / \eta$ when the determinant of $\mathbf{C}$ vanishes and insert the result in,

$$
\begin{equation*}
\Omega \sim \Omega_{0} \pm \frac{1}{2 \Omega_{0}} \sqrt{\alpha_{i j} \kappa_{i} \kappa_{j}} . \tag{11.24}
\end{equation*}
$$

If the $\Omega_{1}$ are zero one must go to the next order.

### 11.2.1.1 Repeated eigenvalues: quadratic asymptotics

If $\Omega_{1}$ is zero, $u_{1}=f_{0, X_{k}}^{(l)} U_{1_{k}}^{(l)}$ (we again sum over all repeated ( $l$ ) superscripts) and we advance to second order using (Ш.Ш1). Taking the difference between the product of equation (Ш.Ш) with $U_{0}^{(m)}$ and $u_{2}\left(U_{0, \xi_{i} \xi_{i}}+\Omega_{0}^{2} U_{0}\right)$ and then integrating over the elementary cell gives

$$
\begin{gather*}
f_{0, X_{i} X_{i}}^{(l)} \iint_{S} U_{0}^{(m)} U_{0}^{(l)} d S+f_{0, X_{k} X_{j}}^{(l)} \iint_{S}\left(U_{0}^{(m)} U_{1_{k}, \xi_{j}}^{(l)}-U_{0, \xi_{j}}^{(m)} U_{1_{k}}^{(l)}\right) d S \\
+\Omega_{2}^{2} f_{0}^{(l)} \iint_{S} U_{0}^{(m)} U_{0}^{(l)} d S=0, \quad \text { for } \quad m=1,2, \ldots, p \tag{11.25}
\end{gather*}
$$

as a system of coupled PDEs. The above equation is presented more neatly as

$$
\begin{equation*}
f_{0, X_{i} X_{i}}^{(l)} \mathbf{A}_{m l}+f_{0, X_{k} X_{j}}^{(l)} \mathbf{D}_{k j m l}+\Omega_{2}^{2} f_{0} \mathbf{B}_{m l}=0, \quad \text { for } \quad m=1,2, \ldots, p \tag{11.26}
\end{equation*}
$$

For the Bloch wave setting, using $f_{0}^{(l)}(\mathbf{X})=\hat{f}_{0}^{(l)} \exp \left(i \kappa_{j} X_{j} / \eta\right)$ we obtain the following system,

$$
\begin{equation*}
\left(-\frac{\kappa_{i} \kappa_{i}}{\eta^{2}} \mathbf{A}_{m l}-\frac{\kappa_{k} \kappa_{j}}{\eta^{2}} \mathbf{D}_{k j m l}+\Omega_{2}^{2} \mathbf{B}_{m l}\right) \hat{f}_{0}^{(l)}=0, \quad \text { for } \quad m=1,2, \ldots, p \tag{11.27}
\end{equation*}
$$

and this determines the asymptotic dispersion curves.

### 11.2.1.2 The classical long wave zero frequency limit

The current theory simplifies if one enters the classical long wave, low frequency limit where $\Omega^{2} \sim O\left(\varepsilon^{2}\right)$ as $U_{0}$ becomes uniform, and without loss of generality is set to be unity, over the elementary cell. The final equation is again (1.15) where the tensor $t_{i j}$ simplifies to

$$
\begin{equation*}
t_{i i}=\iint_{S} d S+\iint_{S} U_{1_{i}, \xi_{i}} d S, \quad t_{i j}=\iint_{S} U_{1_{j}, \xi_{i}} d S \quad \text { for } \quad i \neq j \tag{11.28}
\end{equation*}
$$

(with no summation over repeated suffices in this equation) and $T_{i j}=t_{i j} / \iint_{S} d S$.


Figure 11.8: The dispersion diagrams for a doubly periodic array of square cells with ${ }^{\mathrm{M}} \underset{\mathrm{M}}{\mathrm{M}} \underset{\mathrm{M}}{\mathrm{M}} \underset{\mathrm{M}}{\mathrm{M}} \underset{\mathrm{M}}{\mathrm{M}} \underset{\mathrm{M}}{\mathrm{M}} \underset{\text { inclusions, free at }}{\mathrm{X}}$ their inner boundaries shown for the irreducible Brillouin zone of Fig. Ш.6 The dispersion curves are shown in solid lines and the asymptotic solutions from the high frequency homogenization theory are shown in dashed lines. Figure reproduced from Proceedings of the Royal Society [40].

### 11.2.2 Illustrations for Tranverse Electric Polarized Waves

Let us now turn to some illustrative examples. We present in Fig. Ш.8 the TE polarization waves for three types of SRR's (Split Ring Resonator's). Equation (11.15) represents the wave propagation in the effective medium. It is noticable that the $T_{i j}$ coefficients depend on the standing wave frequency and that $T_{11}$ is not necessarily equal to $T_{22}$ in order to yield an anisotropic effective medium for each separate frequency. Near some of the standing wave frequencies the anisotropy effects are very pronounced and well explained by the no longer elliptic equation (ㄴ.15).

In the above equations $U_{1_{i}}$ is a solution of,

$$
\begin{equation*}
U_{1_{j}, \xi_{i} \xi_{i}}=0 \tag{11.29}
\end{equation*}
$$

with boundary conditions $\left(f_{0, X_{i}}+u_{1, \xi_{i}}\right) n_{i}=0$ on the hole boundary. If the medium is homogeneous as it is in the illustrative examples herein, equation (Ш.29) is the same as that for $U_{0}$, but with different boundary conditions. The specific boundary conditions for $U_{1_{j}}$ are

$$
\begin{equation*}
U_{1_{j}, \xi_{i}} n_{i}=-n_{j} \quad \text { for } \quad j=1,2, \tag{11.30}
\end{equation*}
$$

where $n_{i}$ represent the normal vector components to the hole's surface. The role of $\mathbf{U}_{\mathbf{1}}$ is to ensure Neumann boundary conditions hold and the tensor contains simple averages of inverse permittivity and permeability supplemented by the correction term which takes into account the boundary conditions at $\partial S_{2}$. Equation ( $\left.\amalg .28\right)$ is the classical expression for the homogenised


Figure 11.9: Cloaking in square arrays of SRRs with four holes: A source at frequency $\Omega=2.8$, located in the center of a square metamaterial consisting of 64 SRRs shaped as in Fig. $W .8(b)$ produces a wave pattern reminiscent of (a) concentric spherical field, (b) cloaking of a rectangular inclusion inside a slab of a metamaterial consisting of 38 SRRs and (c) scattering of a plane wave from the same rectangular hole as the previous panel. (d) Zoom in dispersion diagram of Fig. Ш.8(b). Panels ( $e$ ), ( $f$ ) and ( $g$ ) present isofrequency plots of the respective the lower, middle and upper modes of the Dirac point. Figure reproduced from Proceedings of the Royal Society [40].
coefficient in a scalar wave equation with constant material properties; (Ш.29) is the wellknown annex problem of electrostatic type set on a periodic cell, see [4, [1] , and also holds for the homogenised vector Maxwell's system, where $\mathbf{U}_{\mathbf{1}}$ now has three components and $i, j=$ $1,2,3$ [41, 42, 43]].

### 11.2.2.1 Cloaking in metamaterials

SRRs with 4 holes are now used and the dispersion diagrams are in Fig. $[.8$ (b). The flat band along the $M \Gamma$ path is interesting for the fifth mode and we choose to illustrate cloaking effects that occur here. In Fig. Ш.Q(a), we set an harmonic source at the corresponding frequency $\Omega=2.8$ in an $8 \times 8$ array of SRRs and observe a wave pattern of concentric spherical modes. As can be seen in Figs. $\amalg .9(b)$ and $\amalg .9(\mathrm{c})$ a plane wave propagating at frequency $\Omega=2.8$ demonstrates perfect transmission through a slab composed of 38 SRRs but also cloaking of a rectangular inclusion where no scattering is seen before or after the metamaterial slab. Panel (d) of Fig. $\amalg .9$ shows the location in the band structure that is responsible for this effect. Note that the frequency of excitation is just below the Dirac cone point located at $\Omega=2.835$ where the group velocity is negative but also constant near that location of the Brillouin zone illustrated through an isofrequency plot of lower mode of the Dirac point in Fig. Ш.9(e). In constrast with
the isotropic features of panel (e), those of panels (f) and (g) show ultra-flattened isofrequency contours that relate to ultra-refraction, a regime more prone to omni-directivity than cloaking. The asymptotic system of equations (Ш.20) describing the effective medium at the Dirac point can be uncoupled to yield one same equation for all $f_{0}^{(j)}$, ,

$$
\begin{equation*}
f_{0, X_{i} X_{i}}^{(j)}+0.7191 \Omega_{1}^{4} f_{0}^{(3)}=0 \tag{11.31}
\end{equation*}
$$

After some further analysis, the PDE for $f_{0}^{(2)}$ is responsible for the effects at the frequency chosen $\Omega=2.8$.

### 11.2.2.2 Lensing via AANR and St Andrew's cross in metamaterials

We observe all-angle-negative-refraction effect in metamaterials with SRRs with 8 holes. The dispersion curves in Fig. Ш.8(c) are interesting, as the second curve displays the hallmark of an optical band for a photonic crystal (it has a negative group velocity around the $\Gamma$ point). However, this band is the upper edge of a low frequency stop band induced by the resonance of a SRR, whereas the optical band of a PC results from multiple scattering, which thus arises at higher frequencies. We are therefore in presence of a periodic structure behaving somewhat as a composite intermediate between a metamaterial and a photonic crystal. One of the most topical subjects in photonics is the so-called all-angle-negative- refraction (AANR), which was first described in [46]. AANR allows one to focus light emitted by a point, onto an image, even through a flat lens, provided that certain conditions for AANR are met, such as convex isofrequency contours shrinking with frequency about a point in the Brillouin zone [49]. In Fig. Ш.10, we show such an effect for a perfectly conducting photonic crystal (PC) in Fig. [.LO(a). In order to achieve AANR, we choose a frequency on the first dispersion curve (acoustic band) in Fig. Ш.8(c), and we take its intersection with the light line $\Omega=|\kappa|$ along the $X \Gamma$ path. This means that we achieve negative group velocity for waves propagating along the $X \Gamma$ direction of the array, hence the rotation by an angle $\pi / 4$ of every cell within the PC in panel (b) of Fig. [.I). This is a standard trick in optics that has the effect of moving the origin of the light-line dispersion to $X$ as, relative to the PC, the Bloch wavenumber is along $X \Gamma$. This then creates optical effects due to the interaction of the light-line with the acoustic branch, this would be absent if $\Gamma$ were the light-line origin.

The anisotropy of the effective material is reflected from coefficients $T_{11}=-5.53$ and $T_{22}=0.2946$. The same frequency of the first band is reachable at point $N$ of the Brillouin zone. By symmetry of the crystal, we would have $T_{11}=0.2946$ and $T_{22}=-5.53$. The resultant propagating waves would come from the superposition of the two effective media described above. Fig. ШІІ)(b) illustrates this anisotropy as the source wave only propagates at the prescribed directions.

### 11.2.3 Kirchoff Love Plates

HFH is by no means limited to the Helmholtz operator. HFH is here applied to flexural waves in two dimensions [59] for which the governing equation is a fourth order equation

$$
\begin{equation*}
\nabla^{4} u-\Omega^{2} u=0 \tag{11.32}
\end{equation*}
$$

assuming constant material parameters. Such a thin plate can be subject to point, or line, constraints and these are common place in structural engineering.


Figure 11.10: Lensing via AANR and St Andrew's cross in square arrays of SRRs with elght holes: (a) A line source at frequency $\Omega=1.1375$ located above a rectangular metamaterial consisting of of 90 SRRs as in Fig. $\square 1.8(c)$ displays an image underneath (lensing); (b) A line source at frequency $\Omega=1.25$ located inside a square metamaterial consisting of 49 SRRs as in Fig. DIS(c) displays the dynamically induced anisotropy of the effective medium; (c) Zoom in dispersion diagram of Fig. WI.8(c). Note that each cell in the arrays in (a) and (b) has been rotated through an angle $\pi / 4$. Figure reproduced from Proceedings of the Royal Society [40].
 plate (the dots represent the simple supports) in panel (a) with the elementary cell shown by the dotted lines and in (b) the irreducible Brillouin zone with the lettering for wavenumber positions shown. Figure reproduced from Proceedings of the Royal Society [59].


Figure 11.12: The dispersion diagram for a doubly periodic array of point simple supports shown for the irreducible Brillouin zone of Fig. ППП The figure shows the dispersion curves as solid lines. As dashed lines, the asymptotic solutions from the high frequency homogenization theory are shown. Figure reproduced from Proceedings of the Royal Society [59]

In two dimensions, only a few examples of constrained plates are available in the literature: a grillage of line constraints as in [60] that is effectively two coupled one dimensional problems, a periodic line array of point supports [61] raises the possibility of Rayleigh-Bloch modes and for doubly periodic point supports there are exact solutions by [62] (simply supported points) and by [63] (clamped points); the simply supported case is accessible via Fourier series and we choose this as an illustrative example that is of interest in its own right; it is shown in figure Ш.Ш1(a). In particular the simply supported plate has a zero-frequency stop-band and a non-trivial dispersion diagram. It is worth noting that classical homogenization is of no use in this setting with a zero frequency stop band. Naturally waves passing through periodically constrained plates have many similarities with those of photonics in optics.

We consider a double periodic array of points at $x_{1}=2 n_{1}, x_{2}=2 n_{2}$ where $u=0$ (with the first and second derivatives continuous) and so the elementary cell is one in $\left|x_{1}\right|<1,\left|x_{2}\right|<1$ with $u=0$ at the origin (see Figure Ш.П】); Floquet-Bloch conditions are applied at the edges of the cell.

Applying Bloch's theorem and Fourier series the displacement is readily found [62] as

$$
\begin{equation*}
u(\mathbf{x})=\exp (i \boldsymbol{\kappa} \cdot \mathbf{x}) \sum_{n_{1}, n_{2}} \frac{\exp (-i \pi \mathbf{N} \cdot \mathbf{x})}{\left[\left(\kappa_{1}-\pi n_{1}\right)^{2}+\left(\kappa_{2}-\pi n_{2}\right)^{2}\right]^{2}-\Omega^{2}}, \tag{11.33}
\end{equation*}
$$

where $\mathbf{N}=\left(n_{1}, n_{2}\right)$, and enforcing the condition at the origin gives the dispersion relation

$$
\begin{equation*}
D\left(\kappa_{1}, \kappa_{2}, \Omega\right)=\sum_{n_{1}, n_{2}} \frac{1}{\left[\left(\pi n_{1}-\kappa_{1}\right)^{2}+\left(\pi n_{2}-\kappa_{2}\right)^{2}\right]^{2}-\Omega^{2}}=0 \tag{11.34}
\end{equation*}
$$

In this two dimensional example a Bloch wavenumber vector $\boldsymbol{\kappa}=\left(\kappa_{1}, \kappa_{2}\right)$ is used and the dispersion relation can be characterised completely by considering the irreducible Brillouin zone $Г Х M$ shown in figure [..].

The dispersion diagram is shown in figure [..22; The singularities of the summand in equation ([1.34) correspond to solutions within the cell satisfying the Bloch conditions at the edges, in some cases these singular solutions also satisfy the conditions at the support and are therefore true solutions to the problem, a similar situation occurs in the clamped case considered using multipoles in [63]. Solid lines in figure [.12] label curves that are branches of the dispersion relation, notable features are the zero-frequency stop-band and also crossings of branches at the edges of the Brillouin zone. Branches of the dispersion relation that touch the edges of the Brillouin zone singly fall into two categories, those with multiple modes emerging at a same standing wave frequency (such as the lowest branch touching the left handside of the figure at M ) and those that are completely alone (such as the second lowest branch on the left at M).

The HFH theory can again be employed to find an effective PDE entirely upon the longscale that describes the behaviour local to the standing wave frequencies and the details are in [59], the asymptotics from the effective PDE are shown in Fig. [.12 as the dashed lines.

### 11.3 High-contrast homogenization

Periodic media offer a convenient tool in achieving control of electromagnetic waves, due to their relative simplicity from the point of view of the manufacturing process, and due to the possibility of using the Floquet-Bloch decomposition for the analysis of the spectrum of the wave equation in such media. The latter issue has received a considerable amount of interest in the mathematical community, in particular from the perspective of the inverse problem: how to achieve a given spectrum and/or density of states for the wave operator with periodic coefficients by designing an appropriate periodic structure? While the Floquet-Bloch decomposition provides a transparent procedure for answering the direct question, it does not yield a straightforward way of addressing the inverse question posed above.

One possibility for circumventing the difficulties associated with the inverse problem is by viewing the given periodic structure as a high-contrast one, if this is possible under the values of the material parameters used. The idea of considering high-contrast composites within the context of homogenization appeared first in the work by Allaire [16], which discussed the application of the two-scale convergence technique (Nguetseng [8]) to classical homogenization. A more detailed analysis of high-contrast composites, along with the derivation of an explicit formula for the related spectrum, was carried out in a major study by Zhikov [[18]. One of the obvious advantages in using high-contrast composites, or viewing a given composite as a high-contrast one, is in the mere existence of such formula for the spectrum. In the present section we focus on the results of the analysis of Zhikov, and on some more recent results for one-dimensional, layered, high-contrast periodic structures.

In order to get an as short as possible approach to the high-contrast theory, we consider the equation of electromagnetic wave propagation in the transverse electric (TE) polarisation, when the magnetic field has the form $(0,0, H)$, in the presence of sources with spatial density $f(\mathbf{x})$ :

$$
\begin{equation*}
-\operatorname{div}\left(\varepsilon^{\eta}\right)^{-1}(\mathbf{x} / \eta) \nabla H(\mathbf{x})=\omega^{2} H(\mathbf{x})+f(\mathbf{x}), \quad \mathbf{x} \in \Omega \subset \mathbb{R}^{2} \tag{11.35}
\end{equation*}
$$

where we normalise the speed of light $c$ to 1 for simplicity, which amounts to taking $\varepsilon_{0} \mu_{0}=1$ in section $\amalg .3$, and where the magnetic permeability is assumed to be equal to unity throughout the medium (i.e. $\mu=\mu_{0}$ ), and the function $f(\mathbf{x})$ is assumed to vanish outside some set that
has positive distance to the boundary of $\Omega$. The inverse dielectric permittivity tensor $\left(\varepsilon^{\eta}\right)^{-1}(\mathbf{y})$ is assumed in this section, for simplicity, to be a scalar, taking values $\eta^{\gamma} I$ and $I$, respectively, on $[0,1]^{2}$-periodic open sets $F_{0}$ and $F_{1}$, such that $\overline{F_{0}} \cup \overline{F_{1}}=\mathbb{R}^{2}$. Here $\gamma$ is a positive exponent representing a "contrast" between material properties of the two components of the structure that occupy the regions $F_{0}$ and $F_{1}$. In what follows we also assume that $F_{0} \cap[0,1]^{2}$ has a finite distance to the boundary of the unit cell $[0,1]^{2}$, so that the "soft" component $F_{0}$ consists of disjoint "inclusions", spaced $[0,1]^{2}$-periodically from each other, while the "stiff" component $F_{1}$ is a connected subset of $\mathbb{R}^{2}$. The matrix $\varepsilon^{\eta}$ represents the dielectric permittivity of the medium at a given point, however the analysis and conclusions of this section are equally applicable to acoustic wave propagation, which is the context we borrow the terms "soft" and "stiff" from. The assumed relation between the values of dielectric permittivity $\varepsilon^{\eta}$ (in acoustics, between the "stiffnesses") on the two components of the structure is close to the setting of what has been described as "arrow fibres" in the physics literature on electromagnetics, see e.g [64].

A simple dimensional analysis shows that if $\omega \sim 1$ then the soft inclusions are in resonance with the overall field if and only if $\gamma=2$, which is the case we focus on henceforth.

The above equation ([1.35) describes the wave profile for a TE-wave in the cylindrical domain $\Omega \times \mathbb{R}$ domain, and it is therefore supplied with the Neumann condition $\partial H / \partial n=0^{[1]}$ on the boundary of the domain and with the Sommerfeld radiation condition $\partial H / \partial|x|-\mathrm{i} \omega H=$ $o\left(|x|^{-1}\right)$ as $|x| \rightarrow \infty$.

In line with the previous sections, we apply the method of two-scale asymptotic expansions to the above problem, seeking the solution $H=H\left(x_{1}, x_{2}\right)=H(\mathbf{x})$ in the form (see also (Ш.2 in Section Ш.L.2)

$$
\begin{equation*}
H(\mathbf{x})=H_{0}(\mathbf{x}, \mathbf{x} / \eta)+\eta H_{1}(\mathbf{x}, \mathbf{x} / \eta)+\eta^{2} H_{2}(\mathbf{x}, \mathbf{x} / \eta)+\ldots \tag{11.36}
\end{equation*}
$$

where the functions involved are $[0,1]^{2}$-periodic with respect to the "fast" variable $y=x / \eta$. Substituting the expansion (Ш.36) into the equation (Ш.35) and rearranging the terms in the resulting expression in such a way that terms with equal powers of $\eta$ are grouped together, we obtain a sequence of recurrence relations for the functions $H_{k}, k=0,1, \ldots$, from which they are obtained sequentially. The first three of these equations can be transformed to the following system of equations for the leading-order term $H^{(0)}(\mathbf{x}, \mathbf{y})=u(\mathbf{x})+v(\mathbf{x}, \mathbf{y}), \mathbf{x} \in \Omega, \mathbf{y} \in[0,1]^{2}$ :

$$
\begin{gather*}
-\operatorname{div} \varepsilon_{\mathrm{hom}}^{-1} \nabla u(\mathbf{x})=\omega^{2}\left(u(\mathbf{x})+\int_{F_{0} \cap[0,1]^{2}} v(\mathbf{x}, \mathbf{y}) d \mathbf{y}\right)+f(\mathbf{x}), \quad \mathbf{x} \in \Omega,  \tag{11.37}\\
-\Delta_{\mathbf{y}} v(\mathbf{x}, \mathbf{y})=\omega^{2}(u(\mathbf{x})+v(\mathbf{x}, \mathbf{y}))+f(\mathbf{x}), \quad y \in F_{0} \cap[0,1]^{2}, \quad v(\mathbf{x}, \mathbf{y})=0, \quad y \in F_{1} \cap[0,1]^{2} . \tag{11.38}
\end{gather*}
$$

These equations are supplemented by the boundary conditions for the function $u$, of the same kind as in the problems with finite $\eta$. For the sake of simplifying the analysis, we assume that those inclusions that overlap with the boundary of $\Omega$ are substituted by the "main", "stiff" material, where $\left(\varepsilon^{\eta}\right)^{-1}=I$.

In the equation ([.37), the matrix $\varepsilon_{\text {hom }}$ is the classical homogenization matrix for the perforated medium $\varepsilon F_{1}$, see Section above. However, the properties of the system ([.37)([1.38) are rather different to those for the perforated-medium homogenised limit, described by

[^6]the equation $-\operatorname{div} \varepsilon_{\text {hom }}^{-1} \nabla u(\mathbf{x})=\omega^{2} u(\mathbf{x})+f(\mathbf{x})$. As we shall see next, the two-scale structure of (Ш1.37) - (Ш.38) means that the description of the spectra of the problems (Ш.35) in the limit as $\eta \rightarrow 0$ diverges dramatically from the usual moderate-contrast scenario.

The true value of the above limiting procedure is revealed by the statement of the convergence, as $\eta \rightarrow 0$, of the spectra of the original problems to the spectrum of the limit problem described above, see [[18] and by observing that the spectrum of the system (Ш.37)-(Ш.37) is evaluated easily as follows. We write an eigenfunction expansion for $v(\mathbf{x}, \mathbf{y})$ as a function of $y \in F_{0} \cap[0,1]^{2}$ :

$$
\begin{equation*}
v(\mathbf{x}, \mathbf{y})=\sum_{k=0}^{\infty} c_{k}(\mathbf{x}) \psi_{k}(\mathbf{y}), \tag{11.39}
\end{equation*}
$$

where $\psi_{k}$ are the (real-valued) eigenfunctions of the Dirichlet problem $-\Delta \psi_{k}=\lambda_{k} \psi_{k}, y \in F_{1} \cap$ $[0,1]^{2}$, arranged in the order of increasing eigenvalues $\lambda_{k}, k=0,1, \ldots$ and orthonormalised according to the conditions $\int_{F_{0} \cap[0,1]^{2}}\left|\psi_{k}(\mathbf{y})\right|^{2} d \mathbf{y}=1, k=0,1, \ldots$, and $\int_{F_{0} \cap[0,1]^{2}} \psi_{k}(\mathbf{y}) \psi_{l}(\mathbf{y}) d \mathbf{y}=0$, $k \neq l, k, l=0,1, \ldots$ Substituting (Ш.39) into (Ш.38), we find the values for the coefficients $c_{k}$, which yield an explicit expression for $v(\mathbf{x}, \mathbf{y})$ in terms of the function $u(\mathbf{x})$ :

$$
v(\mathbf{x}, \mathbf{y})=\left(\omega^{2} u(\mathbf{x})+f(\mathbf{x})\right) \sum_{k=0}^{\infty}\left(\int_{F_{0} \cap[0,1]^{2}} \psi_{k}(\mathbf{y}) d \mathbf{y}\right)\left(\lambda_{k}-\omega^{2}\right)^{-1} \psi_{k}(\mathbf{y}) .
$$

Finally, using the last expression in the first equation in ([1.37) yields an equation for the function $u$ only:

$$
\begin{equation*}
-\operatorname{div} \varepsilon_{\mathrm{hom}}^{-1} \nabla u(\mathbf{x})=\beta\left(\omega^{2}\right)\left(u(\mathbf{x})+\omega^{-2} f(\mathbf{x})\right), \quad \mathbf{x} \in \Omega, \tag{11.40}
\end{equation*}
$$

where the function $\beta$, which first appeared in the work [18], is given by

$$
\begin{equation*}
\beta\left(\omega^{2}\right)=\omega^{2}\left(1+\omega^{2} \sum_{k=0}^{\infty}\left(\int_{F_{0} \cap[0,1]^{2}} \psi_{k}(\mathbf{y}) d \mathbf{y}\right)^{2}\left(\lambda_{k}-\omega^{2}\right)^{-1}\right) \tag{11.41}
\end{equation*}
$$



Figure 11.13: The plot of the function $\beta$ describing the spectrum of the problem (Ш.37)-(Ш.38) subject to the boundary conditions. The stop bands for the problem in the whole space $\mathbb{R}^{2}$ are indicated by the red intervals of the horizontal axis. The spectra of the problems ([W.3ป) considered in the whole space converge, as $\eta \rightarrow 0$, to the closure of the complement of the union of the red intervals in the positive semiaxis.

The equation $(\mathbb{L} .40)$ is supplemented by appropriate boundary conditions and/or conditions at infinity, which are inherited from the $\eta$-dependent family, i.e. the Neumann condition at the boundary points $\mathbf{x} \in \partial \Omega$ and the radiation condition when $|\mathbf{x}| \rightarrow \infty$. Clearly, the spectrum of this limit problem consists of those values of $\omega^{2}$ for which $\beta\left(\omega^{2}\right)$ is in the spectrum of the operator generated by the differential expression $-\operatorname{div} \varepsilon_{\text {hom }}^{-1} \nabla$ subject to the same boundary conditions. For example, for the problem in the whole space $\mathbb{R}^{2}$ (describing the behaviour of TE-waves in a 3D periodic structure that is invariant in one specified direction) this procedure results in a band-gap spectrum shown in Fig. Ш.13. The end points of each pass band are found by a simple analysis of the formula (Ш.4]): the right ends of each pass band are given by those eigenvalues $\lambda_{k}$ of the Dirichlet Laplacian on the inclusion $F_{0} \cap[0,1]^{2}$ that possess at least one eigenfunction with non-zero integral over $F_{0} \cap[0,1]^{2}$ (otherwise the corresponding term in (Ш.4]) vanishes), while the left ends of the pass bands are given by solutions to the polynomial equation of infinite order $\beta\left(\omega^{2}\right)=0$. These points have a physical interpretation as eigenvalues of the so-called electrostatic problem on the inclusion, see [23].

As in the case of classical, moderate-contrast, periodic media, the fact of spectral convergence offers significant computational advantages over tackling the equations (Ш.35) directly: as $\eta \rightarrow 0$ the latter becomes increasingly demanding, while the former requires a single numerical procedure that serves all $\eta$ once the homogenised matrix $\varepsilon_{\text {hom }}$ and several eigenvalues $\lambda_{k}$ are calculated. A significant new feature, however, as compared to the classical case, is the fact of an infinite set of stop bands opening in the limit as $\eta \rightarrow$, which are easily controlled by the explicit description of the band endpoints. This immediately yields a host of applications of the above results for the design of band-gap devices with prescribed behaviour in the frequency interval of interest.

The theorem on spectral convergence for problems described by the equation ( $\amalg .35)$ is proved in [18]] under the assumption of connectedness of the domain $F_{1}$ occupied by the "stiff" component, via a variant of the extension procedure from $F_{1}$ to the whole of $\mathbb{R}^{2}$ for function sequences whose energy scales as $\eta^{-2}$ (or, equivalently, finite-energy sequences for the operator prior to the rescaling $\mathbf{x} / \eta=\mathbf{y}$ ). In the more recent works [24], [25], this assumption is dropped in a theorem about spectral convergence for a general class of high-contrast operators, via a version of the two-scale asymptotic analysis akin to ( $\amalg .36)$, for the Floquet-Bloch components of the resolvent of the original family of operators following the re-scaling $\mathbf{x} / \eta=\mathbf{y}$. In particular, in [24] a one-dimensional high-contrast model is analysed, which in 3D corresponds to a stack of dielectric layers aligned perpendicular to the direction of the magnetic field. Here the procedure described above for the 2D grating fails to yield a satisfactory limit description as $\eta \rightarrow 0$, i.e. a description where the spectra of problems for finite $\eta$ converge to the spectrum of the limit problem described by the system (Ш.37)-(Ш.38) as $\eta \rightarrow 0$. A more refined analysis of the structure of the related $\eta$-dependent family results in a statement of convergence to the set described by the inequalities

$$
\begin{equation*}
-1 \leq \frac{1}{2}(\alpha-\beta+1) \sqrt{\lambda} \sin (\sqrt{\lambda}(\alpha-\beta))+\cos (\sqrt{\lambda}(\alpha-\beta)) \leq 1 \tag{11.42}
\end{equation*}
$$

where $\alpha$ and $\beta$ denote the end-points of the inclusion in the unit cell, i.e. $F_{0} \cap[0,1]^{2}=(\alpha, \beta) \times$ $[0,1]$.

Similarly to the spectrum of the 2D high-contrast problem, described by the function $\beta$, the limit spectrum of the 1D problem has a band-gap structure, shown in Fig. [1.14, however the description of the location of the bands is different in that it is no longer obtained from the inequality $\beta>0$, where $\beta$ is the 1D analogue of (Ш.4D). Importantly, the asymptotic
behaviour of the density of states function as $\eta \rightarrow 0$ is also very different in the two cases. One can show that the family of resolvents for the problems ([L.35) converges, up to a suitable unitary transformation, to the resolvent of a certain operator whose spectrum is given exactly by (Ш.42), see [25]. The rate of convergence is rigorously shown to be $O(\eta)$, as is anticipated by the expansion (Ш.36).


Figure 11.14: The square root of the limit spectrum for a $1 D$ high-contrast periodic stack, in TE polarisation. The oscillating solid line is the graph of the function $f(\omega)=\cos (\omega / 2)-\omega \sin (\omega / 2) / 4$ in ([1.42) with $\alpha=1 / 4$, $\beta=3 / 4$. The square root of the spectrum is the union of the intervals indicated by bold lines.

The above 1D result is generalised to the case of an oblique incidence of an electromagnetic wave on the same 3D layered structure. Suppose that $x_{2}$ is the coordinate across the stack. Then, assuming for simplicity that the wave vector $(\varkappa, 0,0)$ is parallel to the direction $x_{1}$, it can be shown that all three components of the magnetic field are non-vanishing, with the magnetic component $H=H_{3}$ satisfying the equation

$$
-\left(\left(\varepsilon^{\eta}\right)^{-1}(x / \eta) H^{\prime}(x)\right)^{\prime}=\left(\omega^{2}-\left(\varepsilon^{\eta}\right)^{-1}(x / \eta) \varkappa^{2}\right) H(x)
$$

subject to the same boundary conditions as before. The modified limit spectrum for this family is given by those $\omega^{2}$ for which ( $c f$. (Ш.42))

$$
\begin{equation*}
-1 \leq \frac{1}{2}(\alpha-\beta+1)\left(\omega-\frac{\varkappa^{2}}{\omega}\right) \sin (\sqrt{\lambda}(\alpha-\beta))+\cos (\sqrt{\lambda}(\alpha-\beta)) \leq 1, \quad \omega>0 \tag{11.43}
\end{equation*}
$$

where, as before, $\alpha$ and $\beta$ describe the "soft" inclusion layer in the unit cell, see [24]. The set of $\omega$ described by the inequalities ([.43) is similar to that shown in Figure Ш.]4, the only significant difference between the two cases being a low-frequency gap opening near $\omega=0$ for (‥43).


Figure 11.15: Superlens application of grating: (a) A time harmonic source at frequency 0.473 displays an image through a square array of square inclusions; (b) Effective magnetism versus frequency using (Ш1.44) for square inclusions of relative permittivity 100 with sidelength $a=0.5 d$ in matrix of relative permittivity 1 (grating pitch $d=0.1$ ); Negative values of the effective magnetism are in the frequency region $[0.432,0.534]$.

### 11.4 Conclusion and further applications of grating theory

To conclude this chapter, we would like to stress that advances in homogenization theory over the past forty years have been fuelled by research in composites [36]. The philosophy of the necessity for rigour expressed by Lord Rayleigh in 1892 concerning the Lorentz-Lorenz equations (also known as Maxwell-Garnett formulae) can be viewed as the foundation act of homogenization: 'In the application of our results to the electric theory of light we contemplate a medium interrupted by spherical, or cylindrical, obstacles, whose inductive capacity is different from that of the undisturbed medium. On the other hand, the magnetic constant is supposed to retain its value unbroken. This being so, the kinetic energy of the electric currents for the same total flux is the same as if there were no obstacles, at least if we regard the wavelength as infinitely great.' In this paper, John William Strutt, the third Lord Rayleigh [29], was able to solve Laplace's equation in two dimensions for rectangular arrays of cylinders, and in threedimensions for cubic lattices of spheres. The original proof of Lord Rayleigh suffered from a conditionally convergent sum in order to compute the dipolar field in the array. Many authors in the theoretical physics and applied mathematics communities proposed extensions of Rayleigh's method to avoid this drawback. Another limit of Rayleigh's algorithm is that it does not hold when the volume fraction of inclusions increases. So-called multipole methods have been developed in conjunction with lattice sums in order to overcome such obstacles, see e.g. [30] for a comprehensive review of these methods. In parallel to these developments, the quasistatic limit for gratings has been the subject of intensive research, one might cite [31] and [32] for important contributions in the 1980s, and [33] for a comprehensive review of the modern theory of gratings, including a close inspection of homogenization limit.

Interestingly, in the pure mathematics community, Zhikov's work on high-contrast homogenization [18] has had important applications in metamaterials, with the interpretation of his homogenized equations in terms of effective magnetism first put forward by O'Brien and Pendry [65], and then by Bouchitté and Felbacq [66], although these authors did not seem to be aware at that time of Zhikov's seminal paper [[18]. In order to grasp the physical importance
of (Ш.40))-(Ш.4I), we consider the case of square inclusions of sidelength $a=d / 2$, where $d$ is the pitch of a bi-periodic grating. The eigenfunctions are $\psi_{n m}(\mathbf{y})=2 \sin \left(n \pi y_{1}\right) \sin \left(n \pi y_{2}\right)$ in (Ш.4II) and the corresponding eigenvalues are $k_{n m}^{2}=\pi^{2}\left(n^{2}+m^{2}\right)$. The right-hand side in the homogenized equation (Ш.40) can then be interpreted in terms of effective magnetism:

$$
\begin{equation*}
\mu_{\text {hom }}(k)=1+\frac{64 a^{2}}{\pi^{4}} \sum_{(n, m) o d d} \frac{k^{2}}{n^{2} m^{2}\left(k_{n m}^{2} / a^{2}-k^{2}\right)} . \tag{11.44}
\end{equation*}
$$

This function can be computed numerically for instance with Matlab and demonstrates that negative values can be achieved for $\mu_{\text {hom }}$ near resonances, see Fig. ㄸ.15(b). This allows for superlensing via negative refraction, as shown in Fig. ШШ.I5(a).

Finally, we would like to point out that high-order homogenization techniques [67] suggest that most gratings display some artificial magnetism and chirality when the wavelength is no longer much larger than the periodicity [68]. We hope we have convinced the reader that there is a whole new range of physical effects in gratings which classical, high-frequency and high-contrast homogenization theories can capture.

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[^0]:    ${ }^{1}$ Littrow mount means retrodiffusion, when the diffracted beams propagates in a direction opposite to the incident beam, i.e. $\sin \theta_{-1}=-\sin \theta_{i}$

[^1]:    ${ }^{1}$ The continuity of the tangential component of the magnetic field is valid for materials having bounded values of permittivity. When the permittivity of the grating material is infinite, as in the model of perfectly conducting material, this condition does not hold.

[^2]:    ${ }^{2}$ This property does not hold if the upper medium is not air but has an index $\tilde{v}$ greater than the index $v$ of the lower medium, provided that the incidence is chosen in such a way that the incident wave is totally reflected by a plane interface (Total Internal Reflection). In that case, $\sin (\theta)$ is replaced by $\tilde{v} \sin (\theta)$ in equation 2.60, in such a way that the zeroth order is evanescent if $\tilde{v} \sin (\theta)>v$.

[^3]:    ${ }^{3}$ It must be remembered that the conventions for the measurements of the angles of incidence and diffraction in $\mathscr{R}_{0}$ are opposite

[^4]:    ${ }^{1}$ When the periodic medium is assumed to be isotropic, $\varepsilon_{i j}(y)=\varepsilon(\mathbf{y}) \delta_{i j}$, with the Kronecker symbol $\delta_{i j}=1$ if $i=j$ and 0 otherwise. For instance, (Ш.ل) has typically the bounds $M=13$ and $m=1$ in optics. One class of problems where this condition (Ш.ل) is violated (the bound below, to be more precise) is considered in Section $\amalg .3$ on high-contrast homogenization.

[^5]:    ${ }^{2}$ We note that $(\mathscr{A} 2)$ are two equations which merely depend upon $\varepsilon_{i j}^{-1}(\mathbf{y})$, that is on the microscopic properties of the periodic medium. The two functions $w^{k}$ (defined up to an additive constant) can be computed once for all, independently of $\Omega_{f}$.

[^6]:    ${ }^{3}$ Neumann boundary conditions i.e. infinite conducting walls is a good model for metals in microwaves, but much less so in the visible range of frequencies wherein absorption by metals need be taken into account. Note also that in the TM polarization case, when the electric field takes the form $(0,0, E)$, our analysis applies mutatis mutandis by interchanging the roles of $\varepsilon$ and $\mu, H$ and $E$, and Neumann boundary conditions by Dirichlet ones.

