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## To cite this version:

Hao Chen, Martin Traizet. STACKING DISORDER IN PERIODIC MINIMAL SURFACES. SIAM
Journal on Mathematical Analysis, 2021, 53 (1), pp.855-887. hal-03897973

## HAL Id: hal-03897973

## https://hal.science/hal-03897973

Submitted on 14 Dec 2022

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# STACKING DISORDER IN PERIODIC MINIMAL SURFACES 

HAO CHEN AND MARTIN TRAIZET


#### Abstract

We construct 1-parameter families of non-periodic embedded minimal surfaces of infinite genus in $T \times \mathbb{R}$, where $T$ denotes a flat 2-tori. Each of our families converges to a foliation of $T \times \mathbb{R}$ by $T$. These surfaces then lift to minimal surfaces in $\mathbb{R}^{3}$ that are periodic in horizontal directions but not periodic in the vertical direction. In the language of crystallography, our construction can be interpreted as disordered stacking of layers of periodically arranged catenoid necks. Our work is motivated by experimental observations of twinning defects in periodic minimal surfaces, which we reproduce as special cases of stacking disorder.


## 1. Introduction

1.1. Background. Triply periodic minimal surfaces (TPMSs) is a topic of trans-disciplinary interest. On the one hand, the mathematical notion has been employed to model many structures in nature (e.g. biological membrane) and in laboratory (e.g. lyotropic liquid crystals); we refer the readers to the book $\left[\mathrm{HBL}^{+} 96\right]$ for more information. On the other hand, natural scientists have been contributing with important mathematical discoveries, many long precede the rigorous mathematical treatment. Examples include the famous gyroid discovered in [Sch70] and proved in [GBW96], as well as its deformations discovered in [FHL93, FH99] and recently proved in [Che19].

The current paper is another example in which mathematics is inspired by natural sciences. In [HXBC11], mesoporous crystals exhibiting the structure of Schwarz' D surface are synthesized. Remarkably, a twinning structure, which looks like two copies of Schwarz' D surface glued along a reflection plane, is observed. In other word, the periodicity is broken in the direction orthogonal to the reflection plane. Thereafter, many other crystal defects are experimentally observed in more TPMS structures, leading to a growing demand of mathematical understanding.

Recently, the first named author [Che18] responded to this demand with numerical experiments in Surface Evolver [Bra92]. More specifically, periodic twinning defects are numerically introduced into rPD surfaces (see Figure 1) and the gyroid. Success of these experiments provide strong evidences for the existence of single twinning defects.

Moreover, he also became aware of the node-opening techniques developed by the second named author [Tra02]. The idea is to glue catenoid necks among horizontal planes. When the planes are infinitesimally close, the necks degenerate to singular points termed nodes. It

[^0]is proved that, if the limit positions of the nodes satisfy a balancing condition and a nondegeneracy condition, then it is possible to push the planes a little bit away from each other, giving a 1-parameter family of minimal surfaces along the way. The technique has been used to construct TPMSs [Tra08] by gluing necks among finitely many flat tori, and non-periodic minimal surfaces with infinitely many planar ends [MT12] by gluing necks among infinitely many Riemann spheres.

In this paper, we combine the techniques in [Tra08] and [MT12] to glue necks among infinitely many flat tori. Then each balanced and non-degenerate arrangement of nodes gives rise to a 1-parameter families of minimal surfaces. Seen in $T \times \mathbb{R}$, each of these family converges to a foliation of $T \times \mathbb{R}$ by $T$. Seen in $\mathbb{R}^{3}$, the minimal surfaces are periodic in two independent horizontal directions but not periodic in any other independent direction.

Our motivation is to rigorously construct twinning defects, but the examples produced by our construction is far richer. In the language of crystallography, our construction can be seen as stacking layers of periodically arranged catenoid necks. In the case that $T$ is the 60 -degree torus, for example, we will see that any bi-infinite sequence of 5 stacking patterns gives arise to a 1-parameter family of minimal surfaces. These are then uncountably many families. In particular, a twinning defect arises from a stacking fault, which is not periodic but still quite ordered from a physics point of view. But most of our examples does not exhibit any order, hence should be considered as stacking disorders.

Back to the twinning, experiments and simulations have shown that TPMSs with twinning defects decay exponentially to the standard TPMSs. We will provide mathematical proof to this physics phenomenon, hence finally justify the term "TPMS twinning". More specifically, we will prove that if a configuration is eventually periodic, then the corresponding minimal surface is asymptotic to a TPMS. The proof uses weighted Banach space as in [Tra13].

Our construction uses Implicit Function Theorem, hence only works near the degenerate limit of foliations, which is not physically plausible. However, physicists have proposed formation mechanisms for TPMSs in nature and in laboratory (e.g. [CF97, MBF94, CCM ${ }^{+}$06, $\left.\mathrm{TBC}^{+} 15\right]$ ) that are very similar to node-opening, some even with experiment evidences. Hence we may hope that some of the minimal surfaces constructed in this paper, including those with stacking disorders, would be one day observed in laboratory.
1.2. Mathematical setting. A doubly periodic minimal surface (DPMS) $M$ is invariant by two independent translations, which we may assume to be horizontal. Let $\Gamma$ be the twodimensional lattice generated by these translations, then $M$ projects to a minimal surface $M / \Gamma$ in $\mathbb{R}^{3} / \Gamma=T \times \mathbb{R}$, where $T=\mathbb{R}^{2} / \Gamma$ denotes a flat 2-torus. Immediate examples of infinite genus are given by triply periodic minimal surfaces (TPMSs), if one ignores one of their three periods. Motivated by experimental observations mentioned above, we are particularly interested in non-periodic DPMS with infinite genus.

A flat torus in $T \times \mathbb{R}$ is horizontal if it has the form $T \times\{h\}$ for some $h \in \mathbb{R}$; then $h$ is called the height of the torus. Informally speaking, we construct minimal surfaces that look like infinitely many horizontal flat tori in $T \times \mathbb{R}$, ordered by increasing height, with one catenoid neck between each adjacent pair. The tori are then labeled by $k \in \mathbb{Z}$ in the order of height. The catenoid necks are also labeled by $k \in \mathbb{Z}$, such that the $k$-th neck is between the $k$-th and the $(k+1)$-th tori.

Remark 1.1. Our construction can, in principle, handle finitely many necks between each adjacent pair of tori. But in view of the immediate interest from material sciences, we will


Figure 1. Twinning defects in an rPD surface near the catenoid limit. This is actually an approximation by a TPMS with large vertical period. The surface has a horizontal symmetry plane in the middle. The image was computed in Surface Evolver [Bra92] using the procedure in [Che18].
only glue one catenoid neck between each adjacent tori. This also eases the notations and facilitates the proofs, but still produces a rich variety of examples.

More formally, we say that a minimal surfaces $M \in T \times \mathbb{R}$ is stacked if there is an increasing sequence of real numbers $\left(h_{k}\right)_{k \in \mathbb{Z}}$ satisfying

- $M \cap\left(T \times\left\{h_{k}\right\}\right)$ has a single connected component that projects to a null-homotopic smooth simple closed curve in $T$;
- $M \cap\left(T \times\left(h_{k}, h_{k+1}\right)\right)$ is homeomorphic to $T$ with two disks removed.

Then the $k$-th neck can be interpreted as an annular neighborhood of $M \cap\left(T \times\left\{h_{k}\right\}\right)$.
Remark 1.2. The term "stacked" is borrowed from crystallography. Closed-packed structures are often described as a result of stacking layers of periodically arranged atoms, one on top of another. Analogously, a stacked minimal surface can be seen as obtained by stacking layers of periodically arranged catenoid necks.

We intend to construct 1-parameter families $M_{t}, t>0$, of stacked minimal surfaces such that, in the limit $t \rightarrow 0$, every neck converges to a catenoid after suitable rescaling. If we rescale to keep the size of the torus, then $M_{t}$ will converge to a foliation of $T \times \mathbb{R}$ by $T$, and necks converge to singular points, which we call nodes.

### 1.3. Definitions and main result.

Definition 1.1. A configuration is a sequence $\boldsymbol{p}=\left(p_{k}\right)_{k \in \mathbb{Z}}$ such that $p_{k} \in T$ for all $k \in \mathbb{Z}$.
We use $p_{k}$ to prescribe the limit position of the $k$-th node, and assume that

Hypothesis 1.2 (Uniform separation). There exists a constant $\varrho>0$ such that $p_{k}$ and $p_{\ell}$ are at least at distance $\varrho$ apart whenever $|k-\ell|=1$.

Assume that $2 \omega_{1}=1$ and $2 \omega_{2}=\tau$ generate the lattice $\Gamma$, so $T=T_{\tau}=\mathbb{C} /(\mathbb{Z}+\tau \mathbb{Z})$. Then for $p \in T$, we use $x(p ; \tau)$ and $y(p ; \tau)$ to denote its coordinates in the basis 1 and $\tau$, and define the function

$$
\xi(p ; \tau)=x(p ; \tau) \eta_{1}(\tau)+y(p ; \tau) \eta_{2}(\tau),
$$

where $\eta_{i}(\tau)=\zeta\left(z+2 \omega_{i} ; \tau\right)-\zeta(z ; \tau)=2 \zeta\left(\omega_{i} ; \tau\right)$ for $i=1,2$, and

$$
\zeta(z ; \tau)=\frac{1}{z} \sum_{0 \neq u \in \mathbb{Z}+\tau \mathbb{Z}}\left(\frac{1}{z-u}+\frac{1}{u}+\frac{z}{u^{2}}\right)
$$

is the Weierstrass zeta function associated to $\Gamma$.
The following definitions are borrowed from [Tra08] and [MT12]. Given a configuration $\boldsymbol{p}$, the force $F_{k}$ exerted on the node $p_{k}$ by other nodes in the configuration is

$$
\begin{equation*}
F_{k}:=\zeta\left(p_{k+1}-p_{k} ; \tau\right)+\zeta\left(p_{k-1}-p_{k} ; \tau\right)+2 \xi\left(p_{k} ; \tau\right)-\xi\left(p_{k+1} ; \tau\right)-\xi\left(p_{k-1} ; \tau\right) \tag{1}
\end{equation*}
$$

Definition 1.3. A configuration is said to be balanced if $F_{k}=0$ for all $k \in \mathbb{Z}$.
Since our construction uses the Implicit Function Theorem, we need the differential of the force to be invertible in some sense. As explained in [MT12], $\left(p_{k}\right)$ is not the right parameter to formulate non-degeneracy and one needs to introduce the sequence $\boldsymbol{q}=\left(q_{k}\right)$ defined by

$$
q_{k}=p_{k}-p_{k-1} .
$$

Note that the uniform separation hypothesis can be reformulated as $\left(q_{k}\right)$ being bounded away from 0 . From now on, the term "configuration" refers to an infinite sequence ( $q_{k}$ ) without any warning.

Under the new variables, (1) becomes

$$
F_{k}:=\zeta\left(q_{k+1} ; \tau\right)-\zeta\left(q_{k} ; \tau\right)+\xi\left(q_{k} ; \tau\right)-\xi\left(q_{k+1} ; \tau\right)=G_{k+1}-G_{k} .
$$

where

$$
G_{k}=G\left(q_{k} ; \tau\right)=\zeta\left(q_{k} ; \tau\right)-\xi\left(q_{k} ; \tau\right) .
$$

A configuration is then balanced if $\left(G_{k}\right)$ is a constant sequence, i.e. $G_{k}=G_{0}$ for all $k \in \mathbb{Z}$.
Definition 1.4. A configuration is said to be non-degenerate if the differential of $\left(G_{k}\right)_{k \in \mathbb{Z}}$ with respect to $\left(q_{k}\right)_{k \in \mathbb{Z}}$, as a map from $\ell^{\infty}$ to itself, is an isomorphism.

Note that $G(q ; \tau)$ is periodic in $q$ but not meromorphic. The function $G(q ; \tau)$ is called the Hecke form in [Lan95]. It was proved by Hecke [Hec27] that, if $q=\left(k_{1}+k_{2} \tau\right) / N$ with $\operatorname{gcd}\left(k_{1}, k_{2}, N\right)=1$, then $G(q ; \tau)$ is a modular form of weight 1 with respect to the congruence group $\Gamma(N)$. Recently, $G(q ; \tau)$ gained popularity for its deep connection with the mean field equation [LW10, CKLW18].

Now we are ready to state our main theorems
Theorem 1.5. If a configuration $\boldsymbol{p}$ is balanced, non-degenerate, and satisfies the uniform separation hypothesis, then there exists in $T \times \mathbb{R}$ a 1 -parameter family $\left(M_{t}\right)_{0<t<\epsilon}$ of embedded stacked minimal surfaces which, in the limit $t \rightarrow 0$, converges to a foliation of $T \times \mathbb{R}$ by $T$. Moreover, the necks have asymptotically catenoidal shape and their limiting positions in $T$ are prescribed by $\boldsymbol{p}$.

Theorem 1.6. Let $\boldsymbol{q}$ and $\boldsymbol{q}^{\prime}$ be two balanced and non-degenerate configurations that satisfy the uniform separation hypothesis. Assume that $\left(q_{k}\right)$ is periodic (in the sense $q_{k+N}=q_{k}$ ) and $q_{k}^{\prime}=q_{k}$ for all $k \geq 0$. Let $\left(M_{t}\right)$ and $\left(M_{t}^{\prime}\right)$ denote the corresponding 1-parameter families of minimal surfaces. Then $M_{t}$ is a TPMS and $M_{t}^{\prime}$ is asymptotic to $M_{t}$ as $x_{3} \rightarrow \infty$.

The paper is organized by increasing technicality. After reviewing some examples in Section 2, we prove our main theorems in Section 3 and Section 4, respectively. Technical ingredients of the proofs are delayed to later sections. In Section 5 we prove the existence and smooth dependence on parameters of a holomorphic 1-form $\omega$ that we used in the Weierstrass data in Section 3. In Section 6, we study the asymptotic behavior of $\omega$ and other parameters, which is crucial for proving the TPMS asymptotic behavior in Section 4.

## 2. Examples

Given an infinite sequence of planes, if only one node is opened between each adjacent pair, the result is necessarily a Riemann minimal example [MT12]. We now show that opening nodes among flat tori is a sharp contrast. Although we only open one node between each adjacent pair of tori, we still obtain a rich variety of balanced configurations. We produce configurations using the following

Proposition 2.1. Let $q_{0}, q_{1}, \cdots, q_{n-1}$ be $n$ solutions of the equation

$$
\begin{equation*}
G(q ; \tau)=C \tag{2}
\end{equation*}
$$

for the same complex constant $C$. Assume that the differential of $G$ with respect to $q$ at $q_{k}$, as a function from $\mathbb{R}^{2}$ to $\mathbb{R}^{2}$, is non-singular for each $k$. Then any bi-infinite sequence of $q_{0}, q_{1}, \cdots, q_{n-1}$ is a balanced, non-degenerate configuration satisfying the uniform separation hypothesis.

Proof. The configuration is balanced by definition. Since there is only a finite number of points $q_{k}$, the uniform separation hypothesis is satisfied and the differentials of $G$ with respect to $q$ at $q_{k}$, as well as their inverses, are uniformly bounded. Then the differential of $\left(G_{k}\right)$ with respect to $\left(q_{k}\right)$ is clearly an automorphism of $\ell^{\infty}$.

We can produce a rich variety of examples thanks to the fact that (2) often has several solutions, which we can combine in any arbitrary way to form stack disorders. The case $C=0$ is special, as in this case, the solutions to $G(q, \tau)=0$ are critical points of the Green function on a flat torus. The number of critical points and their non-degeneracy has been investigated in [LW10, CKLW18]. We will exploit their results to discuss the non-degeneracy of our configurations in the case $C=0$.
2.1. TPMSs of genus 3. In this section, we assume that the configuration is periodic with period $2\left(\right.$ namely $q_{2 k}=q_{0}$ and $q_{2 k+1}=q_{1}$ for all $\left.k \in \mathbb{Z}\right)$. The resulting TPMS have genus 3 in the quotient by their period lattice. In this section, we discuss how to recover classical families of TPMSs from such balanced configurations. This completes the discussion in [Tra08] which was incomplete.

If $q_{0}=q_{1}$ (so $\boldsymbol{q}$ is constant) the configuration is automatically balanced. These configurations give rise to TPMSs in Meeks' 5-parameter family. Some famous examples are

Example 1.

- $\operatorname{Re} \tau=0, q_{0}=(1+\tau) / 2$, gives an orthorhombic deformation family of Schwarz' P surface.
- $|\tau|=1, q_{0}=(1+\tau) / 2$, gives another orthorhombic deformation family of Schwarz' P surface.
- $\operatorname{Re} \tau=0, q_{0}=1 / 2$, gives an orthorhombic deformation family Schwarz' CLP surface.
- $\tau=\exp (\mathrm{i} \pi / 3), q_{0}=(1+\tau) / 3$, gives a rhombohedral deformation family of Schwarz' D and P surfaces (rPD family).

Some examples with $q_{0} \neq q_{1}$ are:
Example 2. For arbitrary $\tau, q_{0}=1 / 2$ and $q_{1}=\tau / 2$ is balanced and generically nondegenerate [Che18]. In particular, when $\operatorname{Re} \tau=0$, we recover the newly discovered o $\Delta$ surfaces [CW18a].
Example 3. Examples with $q_{0}=-q_{1}$ were studied in [CW18b]. In particular

- $\tau=\exp (\mathrm{i} \pi / 3), q_{0}=(1+\tau) / 3$, gives hexagonal Schwarz' H family.
- There exists a real number $\pi / 2>\theta^{*}>\pi / 3$ such that whenever $\tau=\exp (i \theta)$ with $\theta<\theta^{*}$, the configuration is balanced with $q_{0}=c(1+\tau)$ for a unique $c<1 / 2$. This leads to the orthorhombic deformations of Schwarz' H surfaces in [CW18b]. Existence and uniqueness of $\theta^{*}$ was essentially proved [Web02, WHW09], and independently in [LW10]. Its value was computed in [CW18b] explicitly as

$$
\begin{equation*}
\theta^{*}=2 \arctan \frac{K^{\prime}(m)}{K(m)} \approx 1.23409 \tag{3}
\end{equation*}
$$

where $m$ is the unique solution of $2 E(m)=K(m)$, and $K(m), K^{\prime}(m)$ and $E(m)$ are elliptic integrals of the first kind, associated first kind, and second kind, respectively.

- For general $\tau$, non-trivial $q_{0}$ 's that give balanced configurations are studied in [Web02, LW10, CKLW18] and numerically in [CW18b].
Remark 2.1. For crystallographers, the rPD and the H surfaces are analogous to, respectively, the cubic and hexagonal close-packing.
Remark 2.2. Interestingly, configurations in Examples 1 and 3 give rise to TPMSs no matter their degeneracy. Those in Example 1 form a 4-parameter family, and they are limits of Meeks 5-parameter family. Those in Example 3 are degenerate only if $q_{0}$ is a 2 -division point, hence reduces to Example 1. In particular, the degenerate configuration with $\tau=\exp \left(i \theta^{*}\right)$ and $q_{0}=q_{1}=(1+\tau) / 2$ is considered in [CW18b]. It is the limit of a 1-parameter family of TPMSs that are degenerate in the sense that the same deformation of the lattices may lead to different deformations of the TPMSs. It is not clear to what extent does this phenomenon generalize.
2.2. Disordered stacking. In the TPMS examples, we have seen a lot balanced configurations with period 2 and $q_{0} \neq q_{1}$. Each of them provides a pair of solutions, namely $q=q_{0}$ and $q=q_{1}$, to the equation $G(q ; \tau)=C$, so we can apply proposition 2.1.

Example 4 (Twinning). Consider a balanced configuration $\boldsymbol{q}$ of period 2 with $q_{0}=-q_{1}$. Since $G(q ; \tau)$ is odd in $q, q_{0}$ and $q_{1}$ must solve the equation $G(q ; \tau)=0$. Then the configuration $\widetilde{\boldsymbol{q}}$ with $\widetilde{q}_{k}=q_{0}$ for $k \leq 0$ and $\widetilde{q}_{k}=q_{1}$ for $k>0$ is also balanced. If $\boldsymbol{q}$ is non-degenerate, then the constant sequence $\left(q_{k}=q_{0}\right)$ gives rise to a family $\left(M_{t}\right)$ of Meeks' TPMSs, and $\widetilde{\boldsymbol{q}}$ gives rise to a family $\left(\widetilde{M}_{t}\right)$ that are twinnings of $\left(M_{t}\right)$. The experimentally observed D-twin corresponds to $\tau=\exp (\mathrm{i} \pi / 3)$ and $q_{0}=(1+\tau) / 3$.

Example 5. For any $\tau$, the 2 -division points $1 / 2, \tau / 2$ and $(1+\tau) / 2$ are trivial solutions to $G(q ; \tau)=0$. Hence any bi-infinite sequences of the three 2 -division points is balanced. Moreover, it is recently proved [CKLW18] that all three trivial solutions are non-degenerate for a generic $\tau$, and at least two of them are non-degenerate for any $\tau$. Hence any flat 2torus admits uncountably many balanced and non-degenerate configuration, giving rise to uncountably many 1 -parameter families of non-periodic minimal surfaces.
Example 6. Since $G(q ; \tau)$ is odd in $q$, non-trivial solutions of $G(q ; \tau)=0$ must appear in pairs. Using a deep connection with the mean field equation, Lin and Wang proved that $G(q ; \tau)=0$ has at most one non-trivial solution pair for a fixed $\tau$ [LW10, Theorem 1.2]. In other words, $G(q ; \tau)=0$ has either three or five solutions. Moreover, in the case of five solutions, all solutions are non-degenerate [LW17]. Hence any bi-infinite sequence of the five solutions is a balanced and non-degenerate configurations. This again gives us uncountably many 1-parameter families of non-periodic minimal surfaces. An example with five solutions is $\tau=\exp (\mathrm{i} \pi / 3)$, with the non-trivial solutions $\pm(1+\tau) / 3$.
2.3. Remarks. The Hecke form $G(q ; \tau)$ has been studied independently by the PDE and minimal surface communities. Hence we would like to point out some connections between their approaches.

Solutions to $G(q ; \tau)=0$ are particularly interesting as they are the critical points of the Green function on a flat torus $T_{\tau}$ [LW10, CKLW18]. This is no surprise in the context of node-opening construction of TPMSs. In [Tra08], the forces between nodes are compared to electrostatic forces between electric charges. The Green function is nothing but the potential function of the electric field generated by periodically arranged charges. The balancing condition asks that all charges are in equilibrium, hence at a critical point of the potential.

A 2-division point $\omega$ is degenerate if

$$
\frac{\tau \wp(\omega ; \tau)+\eta_{2}(\tau)}{\wp(\omega ; \tau)+\eta_{1}(\tau)}
$$

is real. This is the quotient of periods of the elliptic function $\wp(z ; \tau)-\wp(\omega ; \tau)$. So if $\omega$ is degenerate, the torus $T_{\tau}$ admits a meromorphic 1-form with a double pole, a double zero, and only real periods. Such tori are no stranger to the minimal surface theory. In particular, the unique rhombic torus with period quotient -1 was used to construct helicoids with handles [Web02, WHW09], and its angle has an explicit expression as given in (3) (see [CW18b]). On the PDE side, existence of this torus was independently proved in [LW10].

Following [WHW09], we propose a simple construction for the torus and the 1 -form: Slit the complex plane along the real segment $[-1,1]$. Identify the top edge of $[-1,-x]$ (resp. $[-x, 1])$ with the bottom edge of $[x, 1]$ (resp. $[-1, x]$ ), where $x \in[0,1)$ and $(x+1) /(x-1)$ is the quotient of periods. The result is a torus carrying a cone metric with two cone singularities, one of cone angle $6 \pi$ at the point identified with $\pm 1$ and $\pm x$, the other of cone angle $-2 \pi$ at $\infty$. Its periods are obviously real. The same torus with flat metric is $T_{\tau}$.

It follows easily from [Web02] that there exists a unique torus for each real period quotient. This essentially proves Theorem 6.1(1) in [CKLW18]. In fact, all ingredients to prove Example 5 were essentially known in the minimal surface theory.

The equation $G(q ; \tau)=C$ contains the key for understanding catenoid limits of periodic minimal surfaces. On the PDE side, following the electrostatic interpretation above, solutions to $G(q ; \tau)=C \neq 0$ are critical points of the periodic electric potential in the presence of a constant background field. We propose the following conjecture on the number of solutions.

Conjecture 2.2. $G(q ; \tau)=C$ has at most five solutions for any fixed $\tau$ and $C$.

## 3. Construction

3.1. Parameters. The parameters of the construction are a real number $t$ in a neighborhood of 0 and four sequences of complex numbers $\boldsymbol{a}=\left(a_{k}\right), \boldsymbol{b}=\left(b_{k}\right), \boldsymbol{v}=\left(v_{k}\right)$ and $\boldsymbol{\tau}=\left(\tau_{k}\right)$ in $\ell^{\infty}$. Each parameter is in a small $\ell^{\infty}$-neighborhood of a central value denoted with an underscore. We will calculate that the central value of the parameters are:

$$
\begin{equation*}
\underline{a}_{k}=-\frac{1}{2}, \quad \underline{b}_{k}=\frac{1}{2} \xi\left(\underline{v}_{k} ; \underline{\tau}_{k}\right), \quad \underline{v}_{k}=(-\operatorname{conj})^{k} q_{k}, \quad \underline{\tau}_{k}=(-\operatorname{conj})^{k} \tau \tag{4}
\end{equation*}
$$

where conj denotes conjugation, $\tau$ and $\left(q_{k}\right)$ prescribe the flat 2-torus and the configuration as in the introduction. We use $\boldsymbol{x}=(\boldsymbol{a}, \boldsymbol{b}, \boldsymbol{v}, \boldsymbol{\tau})$ to denote the vector of all parameters but $t$. When required, the dependence of objects on parameters will be denoted with a bracket as in $\Sigma[t, \boldsymbol{x}]$, but will be omitted most of the time.
3.2. Opening nodes and the Gauss map. We denote by $T_{k}=T_{k}[\boldsymbol{x}]$ the torus $\mathbb{C} /\left(\mathbb{Z}+\tau_{k} \mathbb{Z}\right)$. The point $z=0$ in $T_{k}$ is denoted by $0_{k}$. We define the elliptic function $g_{k}=g_{k}[\boldsymbol{x}]$ on $T_{k}$ by

$$
g_{k}(z)=a_{k}\left(\zeta\left(z ; \tau_{k}\right)-\zeta\left(z-v_{k} ; \tau_{k}\right)\right)+b_{k}
$$

It has two simple poles at $0_{k}$ and $v_{k}$, with residues $a_{k}$ and $-a_{k}$, respectively. Observe that $1 / g_{k}$ is a local complex coordinate in a neighborhood of $0_{k}$ and $v_{k}$. Hence for a sufficiently small $\varepsilon>0,1 / g_{k}$ gives a diffeomorphism $z_{k}^{+}$from a neighborhood of $v_{k}$ in $T_{k}$ to the disk $D(0,2 \varepsilon) \subset \mathbb{C}$, and a diffeomorphism $z_{k}^{-}$from a neighborhood of $0_{k}$ in $T_{k}$ to the disk $D(0,2 \varepsilon)$. Provided that $\boldsymbol{x}$ is sufficiently close to $\underline{\boldsymbol{x}}$ and by Hypothesis $1.2, \varepsilon$ can be chosen independent of $k$ and $\boldsymbol{x}$. Let $D_{k}^{ \pm}$be the disk $\left|z_{k}^{ \pm}\right|<\varepsilon$ in $T_{k}$.

Consider the disjoint union of all $T_{k}$ for $k \in \mathbb{Z}$. If $t=0$, identify $v_{k} \in T_{k}$ and $0_{k+1} \in T_{k+1}$ to create a node. The resulting Riemann surface with nodes is denoted $\Sigma[0, \boldsymbol{x}]$. If $t \neq 0$ and $|t|<\varepsilon$, then for each $k \in \mathbb{Z}$, remove the disks $\left|z_{k}^{ \pm}\right| \leq t^{2} / \varepsilon$ from $T_{k}$, and let $A_{k}^{ \pm}$be the annuli $t^{2} / \varepsilon<\left|z_{k}^{ \pm}\right|<\varepsilon$. Identify $A_{k}^{+}$and $A_{k+1}^{-}$by $z_{k}^{+} z_{k+1}^{-}=t^{2}$. This opens nodes and creates a neck between $T_{k}$ and $T_{k+1}$. The resulting Riemann surface is denoted $\Sigma=\Sigma[t, \boldsymbol{x}]$.

If $t \neq 0$, we define the Gauss map $g=g[t, \boldsymbol{x}]$ explicitly on $\Sigma[t, \boldsymbol{x}]$ by

$$
g(z)=\left(\operatorname{tg}_{k}(z)\right)^{(-1)^{k+1}} \quad \text { in } T_{k}
$$

Then $g$ takes the same value at the points that are identified when defining $\Sigma$. So $g$ is a well-defined meromorphic function on $\Sigma$.

### 3.3. Height differential. We define

$$
\Omega_{k}:=T_{k} \backslash\left(\overline{D_{k}^{+}} \cup \overline{D_{k}^{-}}\right) \quad \text { and } \quad \Omega:=\bigcup_{k \in \mathbb{Z}} \Omega_{k} \subset \Sigma
$$

All circles $\partial D_{k}^{-}$are homologous in $\Sigma$. This homology class is denoted $\gamma$. We denote by $\alpha_{k}$ and $\beta_{k}$ the standard generators of the homology of $T_{k}$, namely the homology classes of $[0,1]$ and $\left[0, \tau_{k}\right]$ modulo $\mathbb{Z}+\tau_{k} \mathbb{Z}$. We choose representatives of $\alpha_{k}$ and $\beta_{k}$ within $\Omega_{k}$, so they can be seen as curves on $\Sigma$.

By Proposition 5.1 in Section 5, for $t$ small enough, there exists a holomorphic 1-form $\omega=\omega[t, \boldsymbol{x}]$ on $\Sigma[t, \boldsymbol{x}]$ with imaginary periods on $\alpha_{k}, \beta_{k}$ for all $k \in \mathbb{Z}$ and $\int_{\gamma} \omega=2 \pi \mathrm{i}$. We define the height differential $d h$ by

$$
d h=t \omega
$$

If $t=0, \omega$ is allowed to have simple poles at the nodes (a so-called regular 1 -form on a Riemann surface with nodes). So $\omega$ has simple poles at $0_{k}$ and $v_{k}$, with residues 1 and -1 respectively, and imaginary periods on $\alpha_{k}$ and $\beta_{k}$. By Proposition 5.1, we have explicitly

$$
\begin{equation*}
\omega[0, \boldsymbol{x}]=\left(\zeta\left(z ; \tau_{k}\right)-\zeta\left(z-v_{k} ; \tau_{k}\right)-\xi\left(v_{k} ; \tau_{k}\right)\right) d z \quad \text { in } T_{k} . \tag{5}
\end{equation*}
$$

Finally, $\omega[t, \boldsymbol{x}]$ restricted to $\Omega$ depends smoothly on $(t, \boldsymbol{x})$ in a sense which we now explain.
Some care is required because the domain $\Omega$ depends on the parameters. To formulate the smooth dependence, we pullback $\omega$ to a fixed domain as follows. Let $\mathbb{T}=\mathbb{C} /(\mathbb{Z}+\mathrm{i} \mathbb{Z})$ be the standard square torus. Let $\psi_{k}[\boldsymbol{x}]$ be the diffeomorphism defined by

$$
\psi_{k}[\boldsymbol{x}]: \mathbb{T} \rightarrow T_{k}[\boldsymbol{x}], \quad \psi_{k}(x+\mathrm{i} y)=x+\tau_{k} y .
$$

Let $\widetilde{v}_{k}=\psi_{k}^{-1}\left(v_{k}\right)$. Fix a small $\varepsilon^{\prime}>0$ and define

$$
\widetilde{\Omega}_{k}=\mathbb{T} \backslash\left(\overline{D\left(0, \varepsilon^{\prime}\right)} \cup \overline{D\left(\widetilde{v}_{k}, \varepsilon^{\prime}\right)}\right) \quad \text { and } \quad \widetilde{\Omega}=\bigsqcup_{k \in \mathbb{Z}} \widetilde{\Omega}_{k}
$$

If $\varepsilon^{\prime}$ is small enough and $\boldsymbol{x}$ is close enough to $\underline{\boldsymbol{x}}$, we have $\Omega_{k} \subset \psi_{k}\left(\widetilde{\Omega}_{k}\right)$. Moreover, if $t$ is small enough, the disks $\left|z_{k}^{ \pm}\right|<t^{2} / \varepsilon$ which were removed when opening nodes are outside $\psi_{k}\left(\widetilde{\Omega}_{k}\right)$, so we can see $\psi_{k}\left(\widetilde{\Omega}_{k}\right)$ as a domain in $\Sigma$. We define $\psi: \widetilde{\Omega} \rightarrow \Sigma$ by $\psi=\psi_{k}$ on $\widetilde{\Omega}_{k}$. Then $\psi^{*} \omega$ is a smooth 1 -form on $\widetilde{\Omega}$. (Note that $\psi^{*} \omega$ is not holomorphic, because $\psi_{k}$ is not conformal.)

We define the pointwise norm of a (not necessarily holomorphic) 1-form $\eta$ on a domain $U$ in $\mathbb{C}$ or a torus $\mathbb{C} / \Gamma$ by $|\eta(z)|=\sup _{X \in \mathbb{C}^{*}} \frac{|\eta(z) X|}{|X|}$. We denote $C^{0}(U)$ the Banach space of 1-forms $\eta=f_{1} d x+f_{2} d y$ with $f_{1}, f_{2}$ bounded continuous functions on $U$, with the sup norm. We can now state:

Proposition 3.1. The map $(t, \boldsymbol{x}) \mapsto \psi^{*} \omega[t, \boldsymbol{x}]$ is smooth from a neighborhood of $(0, \underline{\boldsymbol{x}})$ to $C^{0}(\widetilde{\Omega})$.

This is the content of Proposition 5.1(c).
Remark 3.1. The point here is that $\widetilde{\Omega}$ is a fixed domain, independent of the parameters. The domain $\psi(\widetilde{\Omega}) \subset \Sigma$ depends on $\left(\tau_{k}\right)$ but not on $t$ nor the other parameters. It contains the domain $\Omega$ which fully depends on $\boldsymbol{x}$. Since $\tau_{k}$ is in a neighborhood of $\tau_{k}$, we have $c^{-1} \leq\left|\psi_{k}^{*} d z\right| \leq c$ for some uniform constant $c$. Hence for any holomorphic 1-form $\eta$ on $\Sigma$,

$$
c^{-1}\left\|\psi_{k}^{*} \eta\right\|_{C^{0}(\tilde{\Omega})} \leq\|\eta\|_{C^{0}(\psi(\widetilde{\Omega}))} \leq c\left\|\psi_{k}^{*} \eta\right\|_{C^{0}(\tilde{\Omega})}
$$

3.4. Zeros of the height differential. We use the Weierstrass parametrization

$$
\Sigma[t, \boldsymbol{x}] \ni z \mapsto \operatorname{Re} \int_{z_{0}}^{z}\left(\Phi_{1}, \Phi_{2}, \Phi_{3}\right) \in \mathbb{R}^{3},
$$

where

$$
\left(\Phi_{1}, \Phi_{2}, \Phi_{3}\right):=\left(\frac{1}{2}\left(g^{-1}-g\right), \frac{\mathrm{i}}{2}\left(g^{-1}+g\right), 1\right) d h
$$

are holomorphic differentials, so that the Weierstrass parametrization is an immersion. So we need to solve the Regularity Problem, which asks that $d h$ has a zero at each zero or pole of $g$, with the same multiplicity, and no other zeros.

The elliptic function $g_{k}[\boldsymbol{x}]$ has degree 2 so it has two zeros in $T_{k}$ which we denote $Z_{k, 1}[\boldsymbol{x}]$ and $Z_{k, 2}[\boldsymbol{x}]$. Since $g_{k}[\boldsymbol{x}]$ has poles at $0_{k}$ and $v_{k}$, its zeros are in $\Omega_{k}$ provided that $\varepsilon$ is small enough. Note that we cannot rule out the possibility of a double zero $Z_{k, 1}=Z_{k, 2}$ at $\underline{\boldsymbol{x}}$,
in which case $Z_{k, 1}$ and $Z_{k, 2}$ are not smooth functions of $\boldsymbol{x}$. Nevertheless, by Weierstrass Preparation Theorem, in a neighborhood of $\underline{\boldsymbol{x}}, Z_{k, 1}+Z_{k, 2}$ and $Z_{k, 1} Z_{k, 2}$ are smooth functions of $\boldsymbol{x}$. The gauss map $g$, by definition, has zeros (resp. poles) at $Z_{k, 1}$ and $Z_{k, 2}$ for $k$ odd (resp. $k$ even).

Proposition 3.2. For $(t, \boldsymbol{x})$ close to $(0, \underline{\boldsymbol{x}})$, $\omega$ has two zeros in $\Omega_{k}$ for $k \in \mathbb{Z}$ (counting multiplicity) and no other zeros. The Regularity Problem is equivalent to

$$
\begin{equation*}
\frac{\omega}{d z}\left(Z_{k, 1}\right)+\frac{\omega}{d z}\left(Z_{k, 2}\right)=0 \tag{6}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{\partial \Omega_{k}} g_{k}^{-1} \omega=0 \tag{7}
\end{equation*}
$$

for $k \in \mathbb{Z}$.
Proof. By (5), $\omega[0, \boldsymbol{x}] / d z$ is an elliptic function of degree 2 in $T_{k}$, with simple poles at $0_{k}$ and $v_{k}$. So it has two zeros in $\Omega_{k}$ as long as $\varepsilon$ is small enough. Hence $\psi^{*} \omega[0, \boldsymbol{x}]$ has two zeros in $\widetilde{\Omega}_{k}$, counting multiplicity. By Proposition 3.1 and Weierstrass Preparation Theorem, for $t$ small enough, $\psi^{*} \omega[t, \boldsymbol{x}]$ has two zeros in $\widetilde{\Omega}_{k}$, so $\omega[t, \boldsymbol{x}]$ has two zeros in $\psi\left(\widetilde{\Omega}_{k}\right)$. By the same proof of Corollary 2 in [Tra13], $\omega$ has no zero in the annuli $A_{k}^{ \pm}$, so it has two zeros in $\Omega_{k}$ for each $k \in \mathbb{Z}$, and no further zero.

If $Z_{k, 1} \neq Z_{k, 2}$, then $d z / g_{k}$ is a meromorphic 1-form on $T_{k}$ with two simple poles at $Z_{k, 1}$ and $Z_{k, 2}$. Therefore, by the Residue Theorem,

$$
\begin{gathered}
\frac{1}{g_{k}^{\prime}\left(Z_{k, 1}\right)}+\frac{1}{g_{k}^{\prime}\left(Z_{k, 2}\right)}=0 \\
\int_{\partial \Omega_{k}} g_{k}^{-1} \omega= \\
=2 \pi \mathrm{i}\left(\frac{\omega}{g_{k}^{\prime} d z}\left(Z_{k, 1}\right)+\frac{\omega}{g_{k}^{\prime} d z}\left(Z_{k, 2}\right)\right) \\
=\frac{2 \pi \mathrm{i}}{g_{k}^{\prime}\left(Z_{k, 1}\right)}\left(\frac{\omega}{d z}\left(Z_{k, 1}\right)-\frac{\omega}{d z}\left(Z_{k, 2}\right)\right)
\end{gathered}
$$

Hence (6) and (7) are equivalent to $\frac{\omega}{d z}\left(Z_{k, 1}\right)=\frac{\omega}{d z}\left(Z_{k, 2}\right)=0$.
If $Z_{k, 1}=Z_{k, 2}, g_{k}$ has a double zero at $Z_{k, 1}$. It is then easy to see, using the Residue Theorem, that (6) and (7) are equivalent to $\omega$ having a double zero at $Z_{k, 1}$.

Proposition 3.3. For $t$ in a neighborhood of 0 , there exists a unique value of $\left(b_{k}\right) \in \ell^{\infty}$, depending smoothly on $t$ and the other parameters, such that (6) is solved for all $k \in \mathbb{Z}$. Moreover, at $t=0$,

$$
b_{k}=-a_{k} \xi\left(v_{k} ; \tau_{k}\right)
$$

so

$$
\begin{equation*}
g_{k} d z=a_{k} \omega \quad \text { in } T_{k} \tag{8}
\end{equation*}
$$

Proof. Define

$$
\mathcal{E}_{k}(t, \boldsymbol{x})=\frac{\omega}{d z}\left(Z_{k, 1}\right)+\frac{\omega}{d z}\left(Z_{k, 2}\right) .
$$

Then by the Residue Theorem,

$$
\mathcal{E}_{k}(t, \boldsymbol{x})=\frac{1}{2 \pi \mathrm{i}} \int_{\partial \Omega_{k}} \frac{\omega}{z-Z_{k, 1}}+\frac{\omega}{z-Z_{k, 2}}=\frac{1}{2 \pi \mathrm{i}} \int_{\partial \Omega_{k}} \frac{\left(2 z-\left(Z_{k, 1}+Z_{k, 2}\right)\right) \omega}{z^{2}-\left(Z_{k, 1}+Z_{k, 2}\right) z+Z_{k, 1} Z_{k, 2}} .
$$

Using Proposition 3.1, $\left(\mathcal{E}_{k}\right)_{k \in \mathbb{Z}}$ is a smooth map with value in $\ell^{\infty}$. At $t=0$, we have by (5)

$$
\frac{\omega}{d z}=\frac{1}{a_{k}}\left(g_{k}(z)-b_{k}\right)-\xi\left(v_{k} ; \tau_{k}\right)
$$

Since $g_{k}\left(Z_{k, i}\right)=0$,

$$
\mathcal{E}_{k}(0, \boldsymbol{x})=-2\left(\frac{b_{k}}{a_{k}}+\xi\left(v_{k} ; \tau_{k}\right)\right)
$$

Since $\underline{a}_{k}=-1 / 2$, the partial differential of $\left(\mathcal{E}_{k}\right)$ with respect to $\left(b_{k}\right)$ is an automorphism of $\ell^{\infty}$. Proposition 3.3 then follows from the Implicit Function Theorem.

If $b_{k}$ is given by Proposition 3.3, then at $t=0, \omega$ and $g_{k}$ are proportional in $T_{k}$ hence have the same zeros. Then (7) is satisfied at $t=0$ disregard of the value of the other parameters. So we cannot easily solve (7) for $t \neq 0$ using the Implicit Function Theorem. We will solve (7) in Section 3.6.
3.5. The Period Problem. From now on, we assume that $\left(b_{k}\right)$ is given by Proposition 3.3 and $\boldsymbol{x}=\left(a_{k}, v_{k}, \tau_{k}\right)_{k \in \mathbb{Z}}$ denotes the remaining parameters. The height differential has imaginary periods by definition. It remains to solve the following Period Problems for all $k \in \mathbb{Z}$ :

$$
\begin{align*}
\operatorname{Re} \int_{\alpha_{k}} \Phi_{1}=(-1)^{k}, & \operatorname{Re} \int_{\alpha_{k}} \Phi_{2}=0  \tag{9}\\
\operatorname{Re} \int_{\beta_{k}} \Phi_{1}=\operatorname{Re} \tau, & \operatorname{Re} \int_{\beta_{k}} \Phi_{2}=\operatorname{Im} \tau \\
\operatorname{Re} \int_{\gamma} \Phi_{1}=0, & \operatorname{Re} \int_{\gamma} \Phi_{2}=0
\end{align*}
$$

Proposition 3.4. For $t$ small enough, there exists unique values for the parameters $\left(a_{k}\right)$ and $\left(\tau_{k}\right)$ in $\ell^{\infty}$, depending smoothly on $t$ and $\left(v_{k}\right)$, such that (9) and (10) are satisfied for all $k \in \mathbb{Z}$. Moreover, at $t=0, a_{k}=-1 / 2$ and $\tau_{k}=(-\operatorname{conj})^{k}(\tau)$, disregard of the value of $\left(v_{k}\right)$.
Proof. we define

$$
\begin{aligned}
& \mathcal{P}_{k, 1}(t, \boldsymbol{x})=\operatorname{conj}\left(\int_{\alpha_{k}} g^{-1} d h\right)-\int_{\alpha_{k}} g d h \\
& \mathcal{P}_{k, 2}(t, \boldsymbol{x})=\operatorname{conj}\left(\int_{\beta_{k}} g^{-1} d h\right)-\int_{\beta_{k}} g d h
\end{aligned}
$$

Equations (9) and (10) are equivalent to

$$
\left\{\begin{array}{l}
\mathcal{P}_{k, 1}(t, \boldsymbol{x})=2(-1)^{k}  \tag{12}\\
\mathcal{P}_{k, 2}(t, \boldsymbol{x})=2 \tau
\end{array}\right.
$$

We have in $T_{k}$ :

$$
g d h=\left\{\begin{array}{ll}
g_{k}^{-1} \omega & \text { if } k \text { even } \\
t^{2} g_{k} \omega & \text { if } k \text { odd }
\end{array} \quad \text { and } \quad g^{-1} d h= \begin{cases}t^{2} g_{k} \omega & \text { if } k \text { even } \\
g_{k}^{-1} \omega & \text { if } k \text { odd }\end{cases}\right.
$$

We can take $\alpha_{k}=\psi_{k}\left(\widetilde{\alpha}_{k}\right)$ and $\beta_{k}=\psi_{k}\left(\widetilde{\beta}_{k}\right)$ where $\widetilde{\alpha}_{k}$ and $\widetilde{\beta}_{k}$ are fixed curves in $\widetilde{\Omega}_{k}$. By Proposition 3.1, $\left(\mathcal{P}_{k, 1}\right)_{k \in \mathbb{Z}}$ and $\left(\mathcal{P}_{k, 2}\right)_{k \in \mathbb{Z}}$ are smooth maps with value in $\ell^{\infty}$. At $t=0$, we have by (8) that: For $k$ even,

$$
\mathcal{P}_{k, 1}(0, \boldsymbol{x})=-\int_{\alpha_{k}} a_{k}^{-1} d z=\frac{-1}{a_{k}} \quad \text { and } \quad \mathcal{P}_{k, 2}(0, \boldsymbol{x})=-\int_{\beta_{k}} a_{k}^{-1} d z=\frac{-\tau_{k}}{a_{k}} .
$$

The solution to (12) is then $a_{k}=-1 / 2$ and $\tau_{k}=\tau$. For $k$ odd,

$$
\mathcal{P}_{k, 1}(0, \boldsymbol{x})=\operatorname{conj} \int_{\alpha_{k}} a_{k}^{-1} d z=\operatorname{conj} \frac{1}{a_{k}} \quad \text { and } \quad \mathcal{P}_{k, 2}(0, \boldsymbol{x})=\operatorname{conj} \int_{\beta_{k}} a_{k}^{-1} d z=\operatorname{conj} \frac{\tau_{k}}{a_{k}} .
$$

The solution to (12) is then $a_{k}=-1 / 2$ and $\tau_{k}=-\operatorname{conj} \tau$. The partial differential of $\left(\left(\mathcal{P}_{k, 1}\right),\left(\mathcal{P}_{k, 2}\right)\right)$ with respect to $\left(\left(a_{k}\right),\left(\tau_{k}\right)\right)$ is clearly an automorphism of $\ell^{\infty} \times \ell^{\infty}$. Proposition 3.4 then follows from the Implicit Function Theorem.
3.6. Balancing. From now on, we assume that the parameters $\left(a_{k}\right)$ and $\left(\tau_{k}\right)$ are given by Proposition 3.4. So the only remaining parameters are $t$ and $\boldsymbol{v}=\left(v_{k}\right)$. It remains to solve (7) and (11). We define

$$
\mathcal{G}_{k}(t, \boldsymbol{v})=\operatorname{conj}^{k} \int_{\partial D_{k}^{-}} g_{k} \omega
$$

Proposition 3.5. For $t \neq 0$, (7) and (11) are equivalent to $\mathcal{G}_{k}(t, \boldsymbol{v})=\mathcal{G}_{0}(t, \boldsymbol{v})$ for all $k \in \mathbb{Z}$.
Proof. We have for $k \in \mathbb{Z}$ :

$$
\begin{aligned}
\int_{\partial \Omega_{k}} g_{k}^{-1} \omega & =-\int_{\partial D_{k}^{-}} z_{k}^{-} \omega-\int_{\partial D_{k}^{+}} z_{k}^{+} \omega \\
& =\int_{\partial D_{k-1}^{+}} \frac{t^{2}}{z_{k-1}^{+}} \omega+\int_{\partial D_{k+1}^{-}} \frac{t^{2}}{z_{k+1}^{-}} \omega \\
& =t^{2} \int_{\partial D_{k-1}^{+}} g_{k-1} \omega+t^{2} \int_{\partial D_{k+1}^{-}} g_{k+1} \omega \\
& \left.=-t^{2} \int_{\partial D_{k-1}^{-}} g_{k-1} \omega+t^{2} \int_{\partial D_{k+1}^{-}} g_{k+1} \omega \quad \text { (because } g_{k-1} \omega \text { is holomorphic in } \Omega_{k-1}\right) \\
& =t^{2} \operatorname{conj}^{k+1}\left(\mathcal{G}_{k+1}-\mathcal{G}_{k-1}\right) .
\end{aligned}
$$

Hence (7) is equivalent to $\mathcal{G}_{k+1}=\mathcal{G}_{k-1}$ for $k \in \mathbb{Z}$.
We chose $\partial D_{1}^{-}$as a representative of $\gamma$. Equation (11) is equivalent to

$$
\int_{\partial D_{1}^{-}} g^{-1} d h-\operatorname{conj}\left(\int_{\partial D_{1}^{-}} g d h\right)=0 .
$$

We have

$$
\begin{gathered}
\int_{\partial D_{1}^{-}} \overline{g d h}=t^{2} \int_{\partial D_{1}^{-}} \overline{g_{1} \omega}=t^{2} \mathcal{G}_{1} \\
\int_{\partial D_{1}^{-}} g^{-1} d h=\int_{\partial D_{1}^{-}} g_{1}^{-1} \omega=-t^{2} \int_{\partial D_{0}^{+}} g_{0} \omega=t^{2} \int_{\partial D_{0}^{-}} g_{0} \omega=t^{2} \mathcal{G}_{0}
\end{gathered}
$$

Hence (11) is equivalent to $\mathcal{G}_{1}=\mathcal{G}_{0}$.
Proposition 3.6. Assume that the configuration $\boldsymbol{q}=\left(q_{k}\right)$ is balanced and non-degenerate. For $t$ in a neighborhood of 0 , there exists a unique $\boldsymbol{v}(t)$ in $\ell^{\infty}$, depending smoothly on $t$, such that $v_{k}(0)=(-\operatorname{conj})^{k} q_{k}$ and

$$
\begin{equation*}
\mathcal{G}_{k}(t, \boldsymbol{v}(t))=-2 \pi \mathrm{i} G\left(q_{0} ; \tau\right) \tag{13}
\end{equation*}
$$

for all $k \in \mathbb{Z}$, so (7) and (11) are solved.

Proof. First of all, $\left(\mathcal{G}_{k}(t, \boldsymbol{v})\right)$ is a smooth function of $(t, \boldsymbol{v})$ with value in $\ell^{\infty}$ by Proposition 3.1. At $t=0$, we have in $T_{k}$ :

$$
\begin{aligned}
g_{k} \frac{\omega}{d z} & =a_{k}\left(\frac{\omega}{d z}\right)^{2}=\frac{-1}{2}\left(\zeta\left(z ; \tau_{k}\right)-\zeta\left(z-v_{k} ; \tau_{k}\right)-\xi\left(v_{k} ; \tau_{k}\right)\right)^{2} \\
\operatorname{Res}_{0_{k}}\left(g_{k} \omega\right) & =\frac{-1}{2} \operatorname{Res}_{0_{k}}\left[\zeta\left(z ; \tau_{k}\right)^{2}-2 \zeta\left(z ; \tau_{k}\right)\left(\zeta\left(z-v_{k} ; \tau_{k}\right)+\xi\left(v_{k} ; \tau_{k}\right)\right)\right] \\
& =-\zeta\left(v_{k} ; \tau_{k}\right)+\xi\left(v_{k} ; \tau_{k}\right)=-G\left(v_{k} ; \tau_{k}\right)
\end{aligned}
$$

Here we used the fact that $\zeta$ is odd, hence $\zeta^{2}$ is even and has no residue at 0 . Then by the Residue Theorem,

$$
\mathcal{G}_{k}(0, \boldsymbol{v})=\operatorname{conj}^{k}\left(-2 \pi \mathrm{i} G\left(v_{k} ; \tau_{k}\right)\right)=-2 \pi \mathrm{i}(-\operatorname{conj})^{k} G\left(v_{k} ; \tau_{k}\right)
$$

Recall that $\tau_{k}=(-\operatorname{conj})^{k}(\tau)$ at $t=0$. Now change to the variable $u_{k}=(-\operatorname{conj})^{k} v_{k}$ with central value $\underline{u}_{k}=(-\operatorname{conj})^{k} \underline{v}_{k}=q_{k}$. Using the definition of $\zeta, \xi$ and $G$, one easily checks that

$$
G(-\bar{z} ;-\bar{\tau})=-\operatorname{conj} G(z ; \tau)
$$

Hence

$$
\mathcal{G}_{k}(0, \boldsymbol{v})=-2 \pi \mathrm{i} G\left(u_{k} ; \tau\right)
$$

Then Proposition 3.6 follows from the balance and non-degeneracy of $\left(q_{k}\right)$ and the Implicit Function Theorem.

Remark 3.2. The horizontal component of the flux of $\gamma$, identified with a complex number, is given by

$$
-\mathrm{i} \int_{\gamma} g d h=-\mathrm{i} \int_{\partial D_{1}^{-}} t^{2} g_{1} \omega=-\mathrm{i} t^{2} \overline{\mathcal{G}_{1}}
$$

So (13) for $k=1$ normalizes the horizontal part of the flux of $\gamma$ (which can also be taken as a free parameter).
3.7. Embeddedness. We denote by $\boldsymbol{x}(t)$ the value of the parameters given by Proposition $3.3,3.4$ and $3.6,\left(\Sigma_{t}, g_{t}, d h_{t}\right)$ the corresponding Weierstrass data, $f_{t}: \Sigma_{t} \rightarrow \mathbb{R}^{3} / \Gamma$ the immersion given by Weierstrass Representation and $M_{t}=f_{t}\left(\Sigma_{t}\right)$. Recall that $\Gamma$ is the 2-dimensional lattice generated by the horizontal vectors $(1,0,0)$ and $(\operatorname{Re} \tau, \operatorname{Im} \tau, 0)$. The goal of this section is to prove that $M_{t}$ is embedded and has the geometry described in Section 1.2. The argument is very similar to Section 4.10 of [MT12], so we will only sketch it.

We write $f_{t}=\left(X_{t}, h_{t}\right)$ with

$$
X_{t}(z)=\frac{1}{2} \operatorname{conj}\left(\int_{z_{0}}^{z} g_{t}^{-1} d h_{t}\right)-\frac{1}{2} \int_{z_{0}}^{z} g_{t} d h_{t} \quad \text { and } \quad h_{t}(z)=\operatorname{Re} \int_{z_{0}}^{z} d h_{t}
$$

We fix a base point $\widetilde{O}_{k}$ in $\widetilde{\Omega}_{k}$, away from the zeros of $g_{k} \circ \psi_{k}$, and let $O_{k}=\psi_{k}\left(\widetilde{O}_{k}\right) \in \Omega_{k}$. Let $w_{k}(t) \in \Sigma_{t}$ be the point $z_{k}^{+}=t$ which is identified with the point $z_{k+1}^{-}=t$ (the "middle" of the $k$-th neck). For $r>0$, we denote $\Omega_{k, r}$ the torus $T_{k}$ minus the two disks $\left|z_{k}^{ \pm}\right| \leq r$. By the computations in Section 3.5, we have in $\Omega_{k, r}$

$$
\begin{equation*}
\lim _{t \rightarrow 0} d X_{t}=(-\operatorname{conj})^{k} d z \tag{14}
\end{equation*}
$$

so the image of $\Omega_{k, r}$ is a graph for $t$ small enough and

$$
\lim _{t \rightarrow 0} X_{t}\left(w_{k}(t)\right)-X_{t}\left(w_{k-1}(t)\right)=(-\operatorname{conj})^{k} \int_{0_{k}}^{\underline{v}_{k}} d z=(-\operatorname{conj})^{k}\left(\underline{v}_{k}\right)=q_{k}
$$

Remark 3.3. This computation is not rigorous because the limit (14) only holds in $\Omega_{k, r}$. It can be made rigorous by expanding $\omega_{t}$ in Laurent series in the annuli $A_{k-1}$ and $A_{k}$, see details in Appendix A of [MT12].

Recall that $q_{k}=p_{k}-p_{k-1}$. We may translate $f_{t}$ horizontally so that $X_{t}\left(w_{0}(t)\right)=p_{0}$. Then

$$
\lim _{t \rightarrow 0} X_{t}\left(w_{k}(t)\right)=p_{k}
$$

In other words, $p_{k}$ is the limit position of the $k$-th neck. Note however that the convergence is not uniform with respect to $k$. By (5), we have for $z \in \Omega_{k, r}$

$$
\begin{aligned}
& \lim _{t \rightarrow 0} \frac{1}{t}\left(h_{t}(z)-h_{t}\left(O_{k}\right)\right)=\operatorname{Re} \lim _{t \rightarrow 0} \int_{O_{k}}^{z} \omega_{t} \\
&=\operatorname{Re} \int_{O_{k}}^{z}\left(\zeta\left(z ; \underline{\tau}_{k}\right)-\zeta\left(z-\underline{v}_{k} ; \underline{\tau}_{k}\right)-\xi\left(\underline{v}_{k} ; \underline{\tau}_{k}\right)\right) d z=: u_{k}(z)
\end{aligned}
$$

The function $u_{k}$ is bounded in $\Omega_{k, r}$ by a uniform constant $C(r)$ depending only on $r$. By Lemma A. 2 in [MT12],

$$
\int_{O_{k-1}}^{O_{k}} \omega_{t} \simeq-2 \log t \operatorname{Res}_{0_{k}}\left(\omega_{0}\right)=-2 \log t \quad \text { as } t \rightarrow 0
$$

Hence

$$
h_{t}\left(O_{k}\right)-h_{t}\left(O_{k-1}\right) \sim-2 t \log t
$$

This ensures that for $t$ small enough, $f_{t}\left(\Omega_{k, r}\right)$ lies strictly above $f_{t}\left(\Omega_{k-1, r}\right)$ so the images $f_{t}\left(\Omega_{k, r}\right)$ for $k \in \mathbb{Z}$ are disjoint.

Since the function $u_{k}$ has two logarithmic singularities at $0_{k}$ and $\underline{v}_{k}$, for $c$ large enough, the level lines $u_{k}= \pm c$ are convex curves, which are included in $\Omega_{k, r}$ provided $r$ is small enough. Define

$$
h_{k}^{ \pm}=h_{t}\left(O_{k}\right) \pm t c .
$$

Then for $t$ small enough, $M_{t}$ intersects the planes $x_{3}=h_{k}^{+}$and $x_{3}=h_{k}^{-}$in two convex curves denoted $\gamma_{k}^{+}$and $\gamma_{k}^{-}$, which are included in $f_{t}\left(\Omega_{k, r}\right)$. Then $M_{t} \cap\left\{h_{k}^{+}<x_{3}<h_{k+1}^{-}\right\}$is a minimal annulus bounded by two convex curves in parallel planes. Such an annulus is foliated by convex curves by a theorem of Shiffman [Shi56]. This proves that $M_{t}$ is embedded.

## 4. Convergence to TPMSs

Assume that the configuration $\left(q_{k}\right)$ is periodic. Then the corresponding minimal surface $M_{t}$ is a TPMS (as an easy consequence of uniqueness in the Implicit Function Theorem). Let $\left(q_{k}^{\prime}\right)$ be another balanced, non-degenerate configuration with the same horizontal lattice $\Gamma$, and assume that $q_{k}^{\prime}=q_{k}$ for all $k \geq 0$. Let $M_{t}^{\prime}$ be the family of minimal surfaces corresponding to the configuration $\left(q_{k}^{\prime}\right)$. In this section, we prove that $M_{t}^{\prime}$ is asymptotic to $M_{t}$ as the vertical coordinate $x_{3} \rightarrow+\infty$. This finally justify the term "TPMS twinning".

We use primes for all objects associated to the configuration $\left(q_{k}^{\prime}\right)$. So $\Sigma^{\prime}=\Sigma\left[t, \boldsymbol{x}^{\prime}(t)\right]$ and $f^{\prime}: \Sigma^{\prime} \rightarrow \mathbb{R}^{3} / \Gamma$ is the immersion given by the Weierstrass data $g^{\prime}=g\left[t, \boldsymbol{x}^{\prime}(t)\right]$ and $d h^{\prime}=$ $t \omega\left[t, \boldsymbol{x}^{\prime}(t)\right]$. We will omit the dependence on $t$, hence will write, for instance, $\omega=\omega[t, \boldsymbol{x}(t)]$
and $\omega^{\prime}=\omega\left[t, \boldsymbol{x}^{\prime}(t)\right]$, and in the same way $\psi_{k}=\psi_{k}[t, \boldsymbol{x}(t)]$ and $\psi_{k}^{\prime}=\psi_{k}\left[t, \boldsymbol{x}^{\prime}(t)\right]$. Otherwise notations are as in Section 3.7.

We use the same letter $C$ to denote any constant that is independent of $t>0$ and $k \in \mathbb{Z}$.
Fix an arbitrary $\delta>1$. We will prove in Proposition 6.5 that, for $t$ sufficiently small, $\boldsymbol{x}-\boldsymbol{x}^{\prime}$ decays like $\delta^{-k}$. That is

$$
\left|x_{k}-x_{k}^{\prime}\right| \leq \frac{C}{\delta^{k}} .
$$

Then Proposition 6.1 implies that, for $t$ sufficiently small, the difference between $\psi_{k}^{*} \omega$ and $\left(\psi_{k}^{\prime}\right)^{*} \omega^{\prime}$ in $\widetilde{\Omega}_{k}$ also decays like $\delta^{-k}$. That is

$$
\left\|\psi_{k}^{*} \omega-\left(\psi_{k}^{\prime}\right)^{*} \omega^{\prime}\right\|_{C^{0}\left(\widetilde{\Omega}_{k}\right)} \leq \frac{C}{\delta^{k}} .
$$

It is convenient to scale the third coordinate of $f$ and $f^{\prime}$ by $t^{-1}$. So let $S: \mathbb{R}^{3} \rightarrow \mathbb{R}^{3}$ be the linear map defined by $S\left(x_{1}, x_{2}, x_{3}\right)=\left(x_{1}, x_{2}, t^{-1} x_{3}\right)$ and define $\widehat{f}=S \circ f$ and $\widehat{f^{\prime}}=S \circ f^{\prime}$.
Proposition 4.1. For $k \in \mathbb{N}$,

$$
\left\|d\left(\widehat{f} \circ \psi_{k}\right)-d\left(\widehat{f}^{\prime} \circ \psi_{k}^{\prime}\right)\right\|_{C^{0}\left(\widetilde{\Omega}_{k}\right)} \leq \frac{C}{\delta^{k}}
$$

Proof. We have in $\Omega_{k}$

$$
d \widehat{f}= \begin{cases}\left(\frac{1}{2}\left(t^{2} \overline{g_{k} \omega}-g_{k}^{-1} \omega\right), \operatorname{Re} \omega\right) & \text { if } k \text { is even }  \tag{15}\\ \left(\frac{1}{2}\left(\overline{g_{k}^{-1} \omega}-t^{2} g_{k} \omega\right), \operatorname{Re} \omega\right) & \text { if } k \text { is odd }\end{cases}
$$

and similar formulas for $d \widehat{f}^{\prime}$. By Propositions 6.1 and 6.5 , we have

$$
\begin{equation*}
\left\|\psi_{k}^{*} \omega-\left(\psi_{k}^{\prime}\right)^{*} \omega^{\prime}\right\|_{C^{0}\left(\tilde{\Omega}_{k}\right)} \leq \frac{C}{\delta^{k}} \tag{16}
\end{equation*}
$$

Using Proposition 6.5 and the definition of $g_{k}$, we obtain

$$
\begin{equation*}
\left\|\psi_{k}^{*}\left(g_{k} \omega\right)-\left(\psi_{k}^{\prime}\right)^{*}\left(g_{k}^{\prime} \omega^{\prime}\right)\right\|_{C^{0}\left(\tilde{\Omega}_{k}\right)} \leq \frac{C}{\delta^{k}} . \tag{17}
\end{equation*}
$$

Let $\widetilde{D}_{k, 1}$ and $\widetilde{D}_{k, 2}$ be two disks containing the two zeros of $g_{k} \circ \psi_{k}$, so that $g_{k}^{-1} \circ \psi_{k}$ is uniformly bounded outside these disks. Then

$$
\begin{equation*}
\left\|\psi_{k}^{*}\left(g_{k}^{-1} \omega\right)-\left(\psi_{k}^{\prime}\right)^{*}\left(\left(g_{k}^{\prime}\right)^{-1} \omega^{\prime}\right)\right\|_{C^{0}\left(\widetilde{\Omega}_{k} \backslash \widetilde{D}_{k, 1} \cup \widetilde{D}_{k, 2}\right)} \leq \frac{C}{\delta^{k}} . \tag{18}
\end{equation*}
$$

Since the Regularity Problem is solved, $\left(g_{k}^{\prime}\right)^{-1} \omega^{\prime}$ extends holomorphically to the zeros of $g_{k}^{\prime}$. Since $\psi_{k}(x+\mathrm{i} y)-\psi_{k}^{\prime}(x+\mathrm{i} y)=\left(\tau_{k}-\tau_{k}^{\prime}\right) y$, we have

$$
\begin{equation*}
\left\|\psi_{k}^{*}\left(\left(g_{k}^{\prime}\right)^{-1} \omega^{\prime}\right)-\left(\psi_{k}^{\prime}\right)^{*}\left(\left(g_{k}^{\prime}\right)^{-1} \omega^{\prime}\right)\right\|_{C^{0}\left(\widetilde{\Omega}_{k}\right)} \leq \frac{C}{\delta^{k}} . \tag{19}
\end{equation*}
$$

From (18) and (19), we obtain by the Triangular Inequality

$$
\left\|\psi_{k}^{*}\left[g_{k}^{-1} \omega-\left(g_{k}^{\prime}\right)^{-1} \omega^{\prime}\right]\right\|_{C^{0}\left(\partial \widetilde{D}_{k, i}\right)} \leq \frac{C}{\delta^{k}} .
$$

Since $g_{k}^{-1} \omega-\left(g_{k}^{\prime}\right)^{-1} \omega^{\prime}$ is holomorphic, we obtain by the Maximum Principle (for the holomorphic structure on $\mathbb{T}$ induced by $\psi_{k}$ )

$$
\left\|\psi_{k}^{*}\left[g_{k}^{-1} \omega-\left(g_{k}^{\prime}\right)^{-1} \omega^{\prime}\right]\right\|_{C^{0}\left(\widetilde{D}_{k, i}\right)} \leq \frac{C}{\delta^{k}} .
$$

Using (19),

$$
\begin{equation*}
\left\|\psi_{k}^{*}\left(g_{k}^{-1} \omega\right)-\left(\psi_{k}^{\prime}\right)^{*}\left(\left(g_{k}^{\prime}\right)^{-1} \omega^{\prime}\right)\right\|_{C^{0}\left(\widetilde{D}_{k, i}\right)} \leq \frac{C}{\delta^{k}} . \tag{20}
\end{equation*}
$$

Proposition 4.1 follows from (16), (17), (18) and (20).
Recall that $O_{k}=\psi_{k}\left(\widetilde{O}_{k}\right)$ and $O_{k}^{\prime}=\psi_{k}^{\prime}\left(\widetilde{O}_{k}\right)$.
Proposition 4.2. For $k \in \mathbb{N}$,

$$
\left\|\widehat{f}\left(O_{k+1}\right)-\widehat{f}\left(O_{k}\right)-\widehat{f}^{\prime}\left(O_{k+1}^{\prime}\right)+\widehat{f}^{\prime}\left(O_{k}^{\prime}\right)\right\| \leq \frac{C}{\delta^{k}} .
$$

Hence the sequence $\left(\widehat{f}\left(O_{k}\right)-\widehat{f}^{\prime}\left(O_{k}^{\prime}\right)\right)_{k \in \mathbb{N}}$ is Cauchy. By translation, we may assume that its limit is zero, so

$$
\left\|\widehat{f}\left(O_{k}\right)-\widehat{f}^{\prime}\left(O_{k}^{\prime}\right)\right\| \leq \frac{C^{\prime}}{\delta^{k}} \quad \text { with } C^{\prime}=\frac{C \delta}{\delta-1} .
$$

Proof. Recall from Section 3.7 that $w_{k} \in \Sigma$ denotes the point $z_{k}^{+}=t$, identified with $z_{k+1}^{-}=t$. Similarly, we introduce $w_{k}^{\prime} \in \Sigma^{\prime}$ to denote the point $z_{k}^{\prime+}=t$, identified with $z_{k+1}^{\prime-1}=t$. Moreover, let $m_{k}^{ \pm} \in \Sigma$ denotes the point $z_{k}^{ \pm}=\varepsilon$, and $m_{k}^{\prime \pm} \in \Sigma^{\prime}$ denotes the point $z_{k}^{\prime \pm}=\varepsilon$.

Proposition 4.2 follows from the following three estimates: For $k \in \mathbb{N}$ :

$$
\begin{align*}
\left\|\widehat{f}\left(O_{k}\right)-\widehat{f}\left(m_{k}^{ \pm}\right)-\widehat{f}^{\prime}\left(O_{k}^{\prime}\right)+\widehat{f}^{\prime}\left(m_{k}^{\prime \pm}\right)\right\| & \leq \frac{C}{\delta^{k}},  \tag{21}\\
\left\|\widehat{f}\left(m_{k}^{+}\right)-\widehat{f}\left(w_{k}\right)-\widehat{f}^{\prime}\left(m_{k}^{\prime+}\right)+\widehat{f}^{\prime}\left(w_{k}^{\prime}\right)\right\| & \leq \frac{C}{\delta^{k}},  \tag{22}\\
\left\|\widehat{f}\left(m_{k+1}^{-}\right)-\widehat{f}\left(w_{k}\right)-\widehat{f}^{\prime}\left(m_{k+1}^{\prime-}\right)+\widehat{f}^{\prime}\left(w_{k}^{\prime}\right)\right\| & \leq \frac{C}{\delta^{k}} . \tag{23}
\end{align*}
$$

Inequality (21) follows from Proposition 4.1,

$$
\psi_{k}^{-1}\left(O_{k}\right)=\left(\psi_{k}^{\prime}\right)^{-1}\left(O_{k}^{\prime}\right)=\widetilde{O}_{k} \quad \text { and } \quad\left|\psi_{k}^{-1}\left(m_{k}^{ \pm}\right)-\left(\psi_{k}^{\prime}\right)^{-1}\left(m_{k}^{\prime \pm}\right)\right| \leq \frac{C}{\delta_{k}} .
$$

To prove (22), we follow the proof of Lemma A. 1 in [MT12]. We write the Laurent series of $\omega$ in the annulus $A_{k}$ in term of the complex coordinate $z=z_{k}^{+}$as

$$
\omega=\frac{-d z}{z}+\sum_{n \geq 1} c_{k, n}^{+} z^{n-1} d z+\sum_{n \geq 1} t^{2 n} c_{k, n}^{-} \frac{d z}{z^{n+1}}
$$

where

$$
\begin{gathered}
c_{k, n}^{+}=\frac{1}{2 \pi \mathrm{i}} \int_{|z|=\varepsilon} \frac{\omega}{z^{n}}=\frac{1}{2 \pi \mathrm{i}} \int_{\partial D_{k}^{+}} \frac{\omega}{\left(z_{k}^{+}\right)^{n}} \\
c_{k, n}^{-}=\frac{t^{-2 n}}{2 \pi \mathrm{i}} \int_{|z|=\varepsilon} z^{n} \omega=\frac{1}{2 \pi \mathrm{i}} \int_{\partial D_{k}^{+}}\left(t^{-2} z_{k}^{+}\right)^{n} \omega=\frac{-1}{2 \pi \mathrm{i}} \int_{\partial D_{k+1}^{-}} \frac{\omega}{\left(z_{k+1}^{-}\right)^{n}} .
\end{gathered}
$$

We expand $\omega^{\prime}$ in the annulus $A_{k}^{\prime}$ in the same way, with coefficients $c_{k, n}^{\prime \pm}$ given by similar formulas. Using Proposition 4.1, we obtain the following estimate

$$
\left|c_{k, n}^{ \pm}-c_{k, n}^{\prime}\right| \leq \frac{C n}{\delta^{k}(2 \varepsilon)^{n}} .
$$

By integration, we obtain the following estimates (see details in Appendix A of [MT12]):

$$
\begin{aligned}
\left|\int_{z=\varepsilon}^{t} \omega-\omega^{\prime}\right| & \leq \frac{C}{\delta^{k}} \\
\left|\int_{z=\varepsilon}^{t} z\left(\omega-\omega^{\prime}\right)\right| & \leq \frac{C}{\delta^{k}} \\
\left|\int_{z=\varepsilon}^{t} t^{2} z^{-1}\left(\omega-\omega^{\prime}\right)\right| & \leq \frac{C}{\delta^{k}}
\end{aligned}
$$

Then (22) follows from (15). (23) is proved in the same way, using $z=z_{k+1}^{-}$as a local coordinate.

Assume that the configuration $\left(q_{k}\right)$ is periodic with even period $N$, i.e. $q_{k+N}=q_{k}$. By uniqueness in the Implicit Function Theorem, the resulting immersion $f$ is periodic. More precisely, if we define $\sigma: \Sigma \rightarrow \Sigma$ by $z \in T_{k} \mapsto z \in T_{k+N}$ then $f \circ \sigma=f+\mathcal{T}$ (where the period $\mathcal{T} \in \mathbb{R}^{3}$ depends on $t$ ).
Proposition 4.3. Let $M=f(\Sigma)$ and $M^{\prime}=f^{\prime}\left(\Sigma^{\prime}\right)$. Define $M_{\ell}^{\prime}=M^{\prime}-\ell \mathcal{T}$ for $\ell \in \mathbb{N}$. Then

$$
\lim _{\ell \rightarrow \infty} M_{\ell}^{\prime}=M
$$

Here the limit is for the smooth convergence on compact subsets of $\mathbb{R}^{3} / \Gamma$.
Proof. By periodicity we have

$$
\widetilde{\Omega}_{k+N}=\widetilde{\Omega}_{k} \quad \text { and } \quad f \circ \psi_{k+N}=f \circ \psi_{k}+\mathcal{T} \quad \text { in } \widetilde{\Omega}_{k}
$$

By Propositions 4.1 and 4.2, we have for $\ell \in \mathbb{N}$

$$
\left\|f \circ \psi_{k+\ell N}-f^{\prime} \circ \psi_{k+\ell N}^{\prime}\right\|_{C^{0}\left(\widetilde{\Omega}_{k}\right)} \leq \frac{C}{\delta^{k+\ell N}}
$$

Define $f_{\ell}^{\prime}=f-\ell \mathcal{T}$. Then

$$
\left\|f \circ \psi_{k}-f_{\ell}^{\prime} \circ \psi_{k+\ell N}^{\prime}\right\|_{C^{0}\left(\widetilde{\Omega}_{k}\right)} \leq \frac{C}{\delta^{k+\ell N}}
$$

Hence

$$
\begin{equation*}
\lim _{\ell \rightarrow \infty} f_{\ell}^{\prime} \circ \psi_{k+\ell N}^{\prime}\left(\widetilde{\Omega}_{k}\right)=f \circ \psi_{k}\left(\widetilde{\Omega}_{k}\right) \tag{24}
\end{equation*}
$$

Recall from Section 3.7 that we have defined heights $h_{k}^{ \pm}$such that $M \cap\left\{h_{k}^{-}<x_{3}<h_{k}^{+}\right\}$is included in $f \circ \psi_{k}\left(\widetilde{\Omega}_{k}\right)$ and is bounded by two convex curves denoted $\gamma_{k}^{+}$and $\gamma_{k}^{-}$. Then for $\ell$ large enough, $M_{\ell}^{\prime} \cap\left\{h_{k}^{-}<x_{3}<h_{k}^{+}\right\}$is included in $f_{\ell}^{\prime} \circ \psi_{k+\ell N}^{\prime}\left(\widetilde{\Omega}_{k}\right)$ and is bounded by two convex curves denoted $\gamma_{k, \ell}^{\prime+}$ and $\gamma_{k, \ell}^{\prime-}$. By (24) we have

$$
\lim _{\ell \rightarrow \infty} M_{\ell}^{\prime} \cap\left\{h_{k}^{-}<x_{3}<h_{k}^{+}\right\}=M \cap\left\{h_{k}^{-}<x_{3}<h_{k}^{+}\right\}
$$

Let $A_{k, \ell}^{\prime}=M_{\ell}^{\prime} \cap\left\{h_{k}^{+}<x_{3}<h_{k+1}^{-}\right\}$. Then $A_{k, \ell}^{\prime}$ is an unstable minimal annulus bounded by $\gamma_{k, \ell}^{\prime+}$ and $\gamma_{k+1, \ell}^{\prime-}$. By Theorem 2(a) in [Tra10], $A_{k, \ell}^{\prime}$ converges subsequentially as $\ell \rightarrow \infty$ to an unstable annulus bounded by $\gamma_{k}^{+}$and $\gamma_{k+1}^{-}$. By [MW91], this annulus is unique so the whole sequence converges and

$$
\lim _{\ell \rightarrow \infty} A_{k, \ell}^{\prime}=M \cap\left\{h_{k}^{+}<x_{3}<h_{k+1}^{-}\right\}
$$

Hence for any integers $k_{1}<k_{2}$, we have

$$
\lim _{\ell \rightarrow \infty} M_{\ell}^{\prime} \cap\left\{h_{k_{1}}^{+}<x_{3}<h_{k_{2}}^{+}\right\}=M \cap\left\{h_{k_{1}}^{+}<x_{3}<h_{k_{2}}^{+}\right\}
$$

which proves Proposition 4.3 since $\lim _{k \rightarrow \pm \infty} h_{k}^{+}= \pm \infty$.

## 5. The holomorphic 1-FORMS $\omega$

The goal of this section is to prove:

## Proposition 5.1.

(a) For $(t, \boldsymbol{x})$ in a neighborhood of $(0, \underline{\boldsymbol{x}})$, there exists a holomorphic (regular if $t=0$ ) 1-form $\omega[t, \boldsymbol{x}]$ on $\Sigma[t, \boldsymbol{x}]$ with imaginary periods on $\alpha_{k}$ and $\beta_{k}$ for all $k \in \mathbb{Z}$ and $\int_{\gamma} \omega=2 \pi \mathrm{i}$.
(b) At $t=0$, we have for all $k \in \mathbb{Z}$ :

$$
\omega[0, \boldsymbol{x}]=\left(\zeta\left(z ; \tau_{k}\right)-\zeta\left(z-v_{k} ; \tau_{k}\right)-\xi\left(v_{k} ; \tau_{k}\right)\right) d z \quad \text { in } T_{k} .
$$

(c) The pullback $\psi^{*} \omega$ is in $C^{0}(\widetilde{\Omega})$ and depends smoothly on $(t, \boldsymbol{x})$ in a neighborhood of $(0, \underline{x})$.

Remark 5.1. If the configuration is periodic with period $2 N$ (namely, $q_{k+2 N}=q_{k}$ ), the quotient of $\Sigma$ by its period is a compact Riemann surface, obtained by opening $2 N$ nodes between $2 N$ tori. In this case, the existence of $\omega$ follows from the standard theory of Opening Nodes [Fay73]. To prove the existence of $\omega$ in the non-periodic case, we adapt the argument in [Tra13] which allows for infinitely many nodes. The difference is that Riemann spheres are replaced by tori.

In the following, we use the same letter $C$ to designate all uniform constants.

### 5.1. Preliminaries.

## Definition 5.2.

(a) For $p, q \in T_{k}, p \neq q$, we denote by $\omega_{k, p, q}$ the unique meromorphic 1-form on $T_{k}$ with simple poles at $p, q$ with residues 1 and -1 , and imaginary periods on $\alpha_{k}$ and $\beta_{k}$. So $\omega_{k, p, q}$ is an abelian differential of the third kind with real normalisation.
(b) For $n \geq 2$, we denote $\omega_{k, n}^{+}$(resp. $\omega_{k, n}^{-}$) the unique meromorphic 1-form on $T_{k}$ with a pole of multiplicity $n$ at $v_{k}$ (resp. $0_{k}$ ) with principal part

$$
\frac{d z_{k}^{ \pm}}{\left(z_{k}^{ \pm}\right)^{n}}
$$

and imaginary periods on $\alpha_{k}$ and $\beta_{k}$. So $\omega_{k, n}^{ \pm}$are abelian differentials of the second kind with real normalisation. Recall that it depends on the choice of the local coordinate $z_{k}^{ \pm}$used to define the principal part.

Lemma 5.3. The abelian differential $\omega_{k, p, q}$ is explicitely given by

$$
\begin{equation*}
\omega_{k, p, q}=\left(\zeta\left(z-p ; \tau_{k}\right)-\zeta\left(z-q ; \tau_{k}\right)-\xi\left(q-p ; \tau_{k}\right)\right) d z \tag{25}
\end{equation*}
$$

Proof. Let $\omega_{k, p, q}^{\prime}$ be the right-hand side of (25). Then $\omega_{k, p, q}^{\prime}$ has simple poles at $p, q$ with residues 1 and -1 . Using the quasi-periodicity of $\zeta, \omega_{k, p, q}^{\prime}$ is independent of the choice of the representatives of $p$ and $q$ modulo $\mathbb{Z}+\tau_{k} \mathbb{Z}$. We take these representatives in the fundamental parallelogram spanned by 1 and $\tau_{k}$. We may represent $\alpha_{k}$ and $\beta_{k}$ by curves which do not
intersect the segment $[p, q]$. Using the quasi-periodicity of $\zeta$ (and omitting the argument $\tau_{k}$ everywhere)

$$
\frac{\partial}{\partial q} \int_{\alpha_{k}}(\zeta(z-p)-\zeta(z-q)) d z=\int_{\alpha_{k}} \frac{d \zeta}{d z}(z-q) d z=[\zeta(z-q)]_{\alpha_{k}(0)}^{\alpha_{k}(1)}=\eta_{1}
$$

Hence

$$
\int_{\alpha_{k}}(\zeta(z-p)-\zeta(z-q)) d z=(q-p) \eta_{1}
$$

Write $q-p=x+y \tau_{k}$ and recall the definition of $\xi$ and Legendre Relation $\eta_{1} \tau_{k}-\eta_{2}=2 \pi \mathrm{i}$. Then

$$
\int_{\alpha_{k}} \omega_{k, p, q}^{\prime}=\left(x+y \tau_{k}\right) \eta_{1}-\left(x \eta_{1}+y \eta_{2}\right)=y\left(\eta_{1} \tau_{k}-\eta_{2}\right)=2 \pi \mathrm{i} y
$$

In the same way,

$$
\int_{\beta_{k}} \omega_{k, p, q}^{\prime}=\left(x+y \tau_{k}\right) \eta_{2}-\left(x \eta_{1}+y \eta_{2}\right) \tau_{k}=x\left(\eta_{2}-\eta_{1} \tau_{k}\right)=-2 \pi \mathrm{i} x
$$

Hence $\omega_{k, p, q}^{\prime}$ has imaginary periods so is equal to $\omega_{k, p, q}$ by uniqueness.
Lemma 5.4. There exists a uniform constant $C$ such that for $(t, \boldsymbol{x})$ in a neighborhood of $(0, \underline{\boldsymbol{x}})$ and $k \in \mathbb{Z}$ :

$$
\left\|\omega_{k, n}^{ \pm}\right\|_{C^{0}\left(\Omega_{k}\right)} \leq C\left(\frac{2}{\varepsilon}\right)^{n-1}
$$

Proof. We only prove the + case. The - case follows similarly.
Let $\eta_{k, n}^{+}$be the unique meromorphic 1-form on $T_{k}$ with a pole at $v_{k}$ with the same principal part as $\omega_{k, n}^{+}$and normalized by $\int_{\alpha_{k}} \eta_{k, n}^{+}=0$. These two differentials are related by

$$
\begin{equation*}
\omega_{k, n}^{+}=\eta_{k, n}^{+}+\frac{\mathrm{i}}{\operatorname{Im}\left(\tau_{k}\right)}\left(\operatorname{Re} \int_{\beta_{k}} \eta_{k, n}^{+}\right) d z \tag{26}
\end{equation*}
$$

Indeed, the right-hand side has imaginary periods.
Write $\eta_{k, n}^{+}=f_{k, n}^{+} d z$. For $p, q$ in $T_{k} \backslash D_{k}^{+}$, we have by the Residue Theorem in $T_{k} \backslash D_{k}^{+}$:

$$
\int_{\partial D_{k}^{+}} f_{k, n}^{+} \omega_{k, p, q}=-2 \pi \mathrm{i}\left(f_{k, n}^{+}(p)-f_{k, n}^{+}(q)\right)
$$

On the other hand, since $\omega_{k, p, q}$ is holomorphic in $D_{k}^{+}$and by definition of the principal part of $f_{k, n}^{+}$at $v_{k}$ :

$$
\int_{\partial D_{k}^{+}} f_{k, n}^{+} \omega_{k, p, q}=\int_{\partial D_{k}^{+}} \frac{d z_{k}^{+}}{d z}\left(z_{k}^{+}\right)^{-n} \omega_{k, p, q}
$$

Hence

$$
f_{k, n}^{+}(p)-f_{k, n}^{+}(q)=-\chi_{k, n}^{+}(p, q) \quad \text { with } \quad \chi_{k, n}^{+}(p, q)=\frac{1}{2 \pi \mathrm{i}} \int_{\partial D_{k}^{+}} \frac{d z_{k}^{+}}{d z}\left(z_{k}^{+}\right)^{-n} \omega_{k, p, q}
$$

Integrating with respect to $q$ on $\alpha_{k}$, we obtain the following integral representation of $f_{k, n}^{+}(p)$ for $p \in T_{k} \backslash D_{k}^{+}$:

$$
\begin{equation*}
f_{k, n}^{+}(p)=\int_{\alpha_{k}}\left(f_{k, n}^{+}(p)-f_{k, n}^{+}(q)\right) d q=-\int_{\alpha_{k}} \chi_{k, n}^{+}(p, q) d q \tag{27}
\end{equation*}
$$

By Cauchy Theorem, we may replace the circle $\partial D_{k}^{+}$by the circle $\left|z_{k}^{+}\right|=\frac{\varepsilon}{2}$ in the definition of $\chi_{k, n}^{+}$. Using Lemma 5.3, there exists a uniform constant $C$ (independent of $k \in \mathbb{Z}$ and $\boldsymbol{x}$ in a neighborhood of $\underline{\boldsymbol{x}}$ ) such that for all $p, q \in \Omega_{k}$ and $z$ on the circle $\left|z_{k}^{+}\right|=\frac{\varepsilon}{2}$ :

$$
\left|\omega_{k, p, q}(z)\right| \leq C .
$$

Then for $p, q \in \Omega_{k}$ :

$$
\left|\chi_{k, n}^{+}(p, q)\right| \leq \frac{C}{2 \pi} \int_{\left|z_{k}^{+}\right|=\varepsilon / 2} \frac{\left|d z_{k}^{+}\right|}{\left|z_{k}^{+}\right|^{n}}=\frac{C}{2 \pi}\left(\frac{2}{\varepsilon}\right)^{n} 2 \pi \frac{\varepsilon}{2}=C\left(\frac{2}{\varepsilon}\right)^{n-1}
$$

By (27),

$$
\begin{equation*}
\left\|\eta_{k, n}^{+}\right\|_{C^{0}\left(\Omega_{k}\right)} \leq C\left(\frac{2}{\varepsilon}\right)^{n-1} \tag{28}
\end{equation*}
$$

Lemma 5.4 then follows from (26).
5.2. Existence. Let $\Lambda$ denotes the set $\mathbb{Z} \times\{n \in \mathbb{N}: n \geq 2\} \times\{+,-\}$. We look for $\omega$ in the form $\omega[t, \boldsymbol{x}]=\widetilde{\omega}(\boldsymbol{x}, \boldsymbol{\lambda}(t, \boldsymbol{x}))$ where

$$
\begin{equation*}
\widetilde{\omega}(\boldsymbol{x}, \boldsymbol{\lambda})=\omega_{k, 0_{k}, v_{k}}+\sum_{n=2}^{\infty} \rho^{n-1}\left(\lambda_{k, n}^{+} \omega_{k, n}^{+}+\lambda_{k, n}^{-} \omega_{k, n}^{-}\right) \quad \text { in } T_{k}, \tag{29}
\end{equation*}
$$

where $\rho \leq \frac{\varepsilon}{4}$ is a fixed positive number, and $\boldsymbol{\lambda}=\left(\lambda_{k, n}^{s}\right)_{(k, n, s) \in \Lambda} \in \ell^{\infty}$ is a sequence of complex numbers to be determined as a function of $(t, \boldsymbol{x})$. Observe that, formally, $\widetilde{\omega}(\boldsymbol{x}, \boldsymbol{\lambda})$ has the desired periods. Regarding convergence, by Lemma 5.4, we have in $\Omega_{k}$ :

$$
\begin{equation*}
\sum_{n=2}^{\infty} \rho^{n-1}\left|\left(\lambda_{k, n}^{+} \omega_{k, n}^{+}+\lambda_{k, n}^{-} \omega_{k, n}^{-}\right)\right| \leq 2 C\|\boldsymbol{\lambda}\|_{\infty} \sum_{n=2}^{\infty}\left(\frac{2 \rho}{\varepsilon}\right)^{n-1} \leq 2 C\|\boldsymbol{\lambda}\|_{\infty} \tag{30}
\end{equation*}
$$

so the series (29) converges absolutely in $\Omega_{k}$ and by Lemma 5.3,

$$
\begin{equation*}
\|\widetilde{\omega}(\boldsymbol{x}, \boldsymbol{\lambda})\|_{C^{0}\left(\Omega_{k}\right)} \leq C\left(1+\|\boldsymbol{\lambda}\|_{\infty}\right) \tag{31}
\end{equation*}
$$

We assume for now the convergence outside $\Omega_{k}$.
Lemma 5.5. Assume that the series (29) converges in the annulus $A_{k}^{ \pm}$for all $k \in \mathbb{Z}$. For $t \neq 0, \widetilde{\omega}(\boldsymbol{x}, \boldsymbol{\lambda})$ is a well-defined 1 -form on $\Sigma[t, \boldsymbol{x}]$ if and only if for all $k \in \mathbb{Z}$ and $n \geq 2$,

$$
\begin{aligned}
& \lambda_{k, n}^{+}=\frac{-1}{2 \pi \mathrm{i}} \int_{\partial D_{k+1}^{-}}\left(\frac{t^{2}}{\rho z_{k+1}^{-}}\right)^{n-1} \widetilde{\omega}(\boldsymbol{x}, \boldsymbol{\lambda}) \quad \text { and } \\
& \lambda_{k, n}^{-}=\frac{-1}{2 \pi \mathrm{i}} \int_{\partial D_{k-1}^{+}}\left(\frac{t^{2}}{\rho z_{k-1}^{+}}\right)^{n-1} \widetilde{\omega}(\boldsymbol{x}, \boldsymbol{\lambda})
\end{aligned}
$$

Proof. Fix $t \neq 0$ and define a diffeomorphism

$$
\varphi_{k}=\left(z_{k+1}^{-}\right)^{-1} \circ \frac{t^{2}}{z_{k}^{+}}: A_{k}^{+} \rightarrow A_{k+1}^{-}
$$

so $z \in A_{k}^{+}$is identified with $\varphi_{k}(z) \in A_{k+1}^{-}$when opening nodes. Then $\widetilde{\omega}$ is well-defined on $\Sigma$ if and only if $\varphi_{k}^{*} \widetilde{\omega}=\widetilde{\omega}$ in the annulus $A_{k}^{+}$for all $k \in \mathbb{Z}$. Using the theorem on Laurent series,
this is equivalent to

$$
\begin{equation*}
\int_{\partial D_{k}^{+}}\left(z_{k}^{+}\right)^{n}\left(\varphi_{k}^{*} \widetilde{\omega}-\widetilde{\omega}\right)=0 \quad \text { for all } k \in \mathbb{Z} \text { and } n \in \mathbb{Z} \tag{32}
\end{equation*}
$$

We have

$$
\begin{aligned}
\int_{\partial D_{k}^{+}}\left(z_{k}^{+}\right)^{n} \varphi_{k}^{*} \widetilde{\omega} & =\int_{\partial D_{k}^{+}} \varphi_{k}^{*}\left[\left(\frac{t^{2}}{z_{k+1}^{-}}\right)^{n} \widetilde{\omega}\right] & & \text { by definition of } \varphi_{k} \\
& =\int_{\varphi_{k}\left(\partial D_{k}^{+}\right)}\left(\frac{t^{2}}{z_{k+1}^{-}}\right)^{n} \widetilde{\omega} & & \text { by a change of variable } \\
& =-\int_{\partial D_{k+1}^{-}}\left(\frac{t^{2}}{z_{k+1}^{-}}\right)^{n} \widetilde{\omega} & & \text { because } \widetilde{\omega} \text { is holomorphic in } A_{k+1}^{-}
\end{aligned}
$$

For $n=0,(32)$ is always satisfied, because

$$
-\int_{\partial D_{k+1}^{-}} \widetilde{\omega}-\int_{\partial D_{k}^{+}} \widetilde{\omega}=-2 \pi \mathrm{i} \operatorname{Res}_{0_{k+1}}\left(\omega_{k+1,0_{k+1}, v_{k+1}}\right)-2 \pi \mathrm{i} \operatorname{Res}_{v_{k}}\left(\omega_{k, 0_{k}, v_{k}}\right)=0
$$

For $n \geq 1$, we have by the Residue Theorem and definition of $\widetilde{\omega}(\boldsymbol{x}, \boldsymbol{\lambda})$

$$
\int_{\partial D_{k}^{+}}\left(z_{k}^{+}\right)^{n} \widetilde{\omega}=2 \pi \mathrm{i} \rho^{n} \lambda_{k, n+1}^{+}
$$

so (32) is equivalent to

$$
2 \pi \mathrm{i} \rho^{n} \lambda_{k, n+1}^{+}=-\int_{\partial D_{k+1}^{-}}\left(\frac{t^{2}}{z_{k+1}^{-}}\right)^{n} \widetilde{\omega} .
$$

For $n \leq-1$, we have by the Residue Theorem

$$
\int_{\partial D_{k+1}^{-}}\left(\frac{t^{2}}{z_{k+1}^{-}}\right)^{n} \widetilde{\omega}=2 \pi \mathrm{i} t^{2 n} \rho^{-n} \lambda_{k+1,-n+1}^{-}
$$

so (32) is equivalent to

$$
2 \pi \mathrm{i} t^{2 n} \rho^{-n} \lambda_{k+1,-n+1}^{-}=-\int_{\partial D_{k}^{+}}\left(z_{k}^{+}\right)^{n} \widetilde{\omega}
$$

which, after replacing $n$ by $-n$ and $k$ by $k-1$, becomes

$$
2 \pi \mathrm{i} \rho^{n} \lambda_{k, n+1}^{-}=-\int_{\partial D_{k-1}^{+}}\left(\frac{t^{2}}{z_{k-1}^{+}}\right)^{n} \widetilde{\omega}
$$

for all $n \geq 1$. Collecting all results gives Lemma 5.5.

In view of Lemma 5.5, we define

$$
L_{k, n}^{ \pm}(t, \boldsymbol{x}, \boldsymbol{\lambda})=\frac{-1}{2 \pi \mathrm{i}} \int_{\partial D_{k \pm 1}^{\mp}}\left(\frac{t^{2}}{\rho z_{k \pm 1}^{\mp}}\right)^{n-1} \widetilde{\omega}(\boldsymbol{x}, \boldsymbol{\lambda})
$$

and $\boldsymbol{L}=\left(L_{k, n}^{s}\right)_{(k, n, s) \in \Lambda}$, so $\widetilde{\omega}(\boldsymbol{x}, \boldsymbol{\lambda})$ is well-defined on $\Sigma[t, \boldsymbol{x}]$ if and only if $\boldsymbol{\lambda}=\boldsymbol{L}(t, \boldsymbol{x}, \boldsymbol{\lambda})$. Observe that $\boldsymbol{L}(t, \boldsymbol{x}, \boldsymbol{\lambda})$ is defined for all $\boldsymbol{\lambda} \in \ell^{\infty}$. In particular, we do not need the convergence in $A_{k}^{ \pm}$to define $\boldsymbol{L}$. Also, $\omega$ and $\boldsymbol{L}$ are affine with respect to $\boldsymbol{\lambda}$.

Lemma 5.6. For $(t, \boldsymbol{x})$ in a neighborhood of $(0, \underline{\boldsymbol{x}}), \boldsymbol{\lambda} \mapsto \boldsymbol{L}(t, \boldsymbol{x}, \boldsymbol{\lambda})$ is contracting from $\ell^{\infty}$ to itself, hence has a fixed point $\boldsymbol{\lambda}(t, \boldsymbol{x})$ by the Fixed Point Theorem.

Proof. We can bound the length of $\partial D_{k}^{ \pm}$by a uniform constant $\ell$. By Estimate (30):

$$
\left|L_{k, n}^{ \pm}(t, \boldsymbol{x}, \boldsymbol{\lambda})-L_{k, n}^{ \pm}(t, \boldsymbol{x}, \mathbf{0})\right| \leq \frac{\ell}{2 \pi}\left(\frac{t^{2}}{\rho \varepsilon}\right)^{n-1} 2 C\|\boldsymbol{\lambda}\|_{\infty}
$$

Hence if $t^{2} \leq \rho \varepsilon$,

$$
\|\boldsymbol{L}(t, \boldsymbol{x}, \boldsymbol{\lambda})-\boldsymbol{L}(t, \boldsymbol{x}, \mathbf{0})\|_{\infty} \leq \frac{C \ell t^{2}}{\pi \rho \varepsilon}\|\boldsymbol{\lambda}\|_{\infty}
$$

so $\boldsymbol{L}$ is contracting for $t$ sufficiently small.

Now we verify the convergence of (29) outside $\Omega_{k}$.
Lemma 5.7. If $\boldsymbol{\lambda}=\boldsymbol{L}(t, \boldsymbol{x}, \boldsymbol{\lambda})$ and $t$ is sufficiently small, the series (29) converges absolutely in the annulus $A_{k}^{ \pm}$for all $k \in \mathbb{Z}$.

Proof. We only deal with the convergence of $\sum_{n \geq 2} \rho^{n-1} \lambda_{k, n}^{+} \omega_{k, n}^{+}$. Its convergence in $A_{k}^{-}$is straightforward because $\omega_{k, n}^{+}$is holomorphic in $D_{k}^{-}$and we already know the convergence on $\partial D_{k}^{-}$. It remains to prove the convergence in $A_{k}^{+}$.

By Cauchy Theorem, we can replace the circle $\partial D_{k+1}^{-}$by the circle $\left|z_{k+1}^{-}\right|=2 \varepsilon$ in the definition of $L_{k, n}^{+}(t, \boldsymbol{x}, \boldsymbol{\lambda})$. Using Estimate (31), this gives

$$
\begin{equation*}
\left|\lambda_{k, n}^{+}\right|=\left|L_{k, n}^{+}(t, \boldsymbol{x}, \boldsymbol{\lambda})\right| \leq C\left(1+\|\boldsymbol{\lambda}\|_{\infty}\right)\left(\frac{t^{2}}{2 \rho \varepsilon}\right)^{n-1} \tag{33}
\end{equation*}
$$

By definition, the function

$$
\frac{\omega_{k, n}^{+}}{d z}-\frac{\left(z_{k}^{+}\right)^{\prime}}{\left(z_{k}^{+}\right)^{n}}
$$

extends holomorphically to the disk $D_{k}^{+}$. By the maximum principle and Lemma 5.4

$$
\sup _{D_{k}^{+}}\left|\frac{\omega_{k, n}^{+}}{d z}-\frac{\left(z_{k}^{+}\right)^{\prime}}{\left(z_{k}^{+}\right)^{n}}\right|=\max _{\partial D_{k}^{+}}\left|\frac{\omega_{k, n}^{+}}{d z}-\frac{\left(z_{k}^{+}\right)^{\prime}}{\left(z_{k}^{+}\right)^{n}}\right| \leq C\left(\frac{2}{\varepsilon}\right)^{n-1}+\frac{C}{\varepsilon^{n}} \leq C\left(\frac{2}{\varepsilon}\right)^{n}
$$

Hence recalling the definition of $A_{k}^{+}$, and provided $2 t^{2} \leq \varepsilon^{2}$ :

$$
\sup _{A_{k}^{+}}\left|\omega_{k, n}^{+}\right| \leq C\left(\frac{2}{\varepsilon}\right)^{n}+C\left(\frac{\varepsilon}{t^{2}}\right)^{n} \leq C\left(\frac{\varepsilon}{t^{2}}\right)^{n}
$$

Using Estimate (33), we obtain

$$
\sup _{A_{k}^{+}} \rho^{n-1}\left|\lambda_{k, n}^{+} \omega_{k, n}^{+}\right| \leq \frac{C \varepsilon}{t^{2} 2^{n}}\left(1+\|\boldsymbol{\lambda}\|_{\infty}\right)
$$

Hence the series $\sum_{n \geq 2} \rho^{n-1} \lambda_{k, n}^{+} \omega_{k, n}^{+}$converges absolutely in $A_{k}^{+}$.
Convergence of $\sum_{n \geq 2} \rho^{n-1} \lambda_{k, n}^{-} \omega_{k, n}^{-}$follows similarly.

We define $\omega[t, \boldsymbol{x}]=\widetilde{\omega}(t, \boldsymbol{\lambda}(t, \boldsymbol{x}))$. By Lemmas 5.5 and 5.7, $\omega[t, \boldsymbol{x}]$ is a well-defined holomorphic 1-form on $\Sigma[t, \boldsymbol{x}]$ and has the desired periods by definition. This proves Proposition 5.1(a). At $t=0, \boldsymbol{L}=\mathbf{0}$ so $\boldsymbol{\lambda}=\mathbf{0}$ and $\omega=\omega_{k, 0_{k}, v_{k}}$ in $T_{k}$. Proposition 5.1(b) follows from Lemma 5.3.
5.3. Smooth dependence on parameters. We denote $\widetilde{\gamma}_{k}^{-}$the circle $|z|=2 \varepsilon^{\prime}$ in $\mathbb{T}=$ $\mathbb{C} /(\mathbb{Z}+\mathrm{i} \mathbb{Z})$ and $\widetilde{\gamma}_{k}^{+}$the circle $\left|z-\widetilde{v}_{k}\right|=2 \varepsilon^{\prime}$. These two circles are fixed and included in the domain $\widetilde{\Omega}_{k}$. Define for $\eta \in C^{0}(\widetilde{\Omega})$

$$
\widetilde{L}_{k, n}^{ \pm}(t, \boldsymbol{x}, \eta)=\frac{-1}{2 \pi \mathrm{i}} \int_{\widetilde{\gamma}_{k \pm 1}^{\mp}}\left(\frac{t^{2}}{\rho z_{k \pm 1}^{\mp} \circ \psi_{k \pm 1}}\right)^{n-1} \eta
$$

and let $\widetilde{\boldsymbol{L}}=\left(\widetilde{L}_{k, n}^{s}\right)_{(k, n, s) \in \Lambda}$. Using a change of variable

$$
\widetilde{L}_{k, n}^{ \pm}\left(t, \boldsymbol{x},\left(\psi^{*} \widetilde{\omega}\right)(\boldsymbol{x}, \boldsymbol{\lambda})\right)=\frac{-1}{2 \pi \mathrm{i}} \int_{\psi_{k \pm 1}\left(\widetilde{\gamma}_{ \pm 1}\right)}\left(\frac{t^{2}}{\rho z_{k \pm 1}^{\mp}}\right)^{n-1} \widetilde{\omega}(\boldsymbol{x}, \boldsymbol{\lambda})=L_{k, n}^{ \pm}(t, \boldsymbol{x}, \boldsymbol{\lambda}) .
$$

Hence

$$
\begin{equation*}
\boldsymbol{L}(t, \boldsymbol{x}, \boldsymbol{\lambda})=\widetilde{\boldsymbol{L}}\left(t, \boldsymbol{x},\left(\psi^{*} \widetilde{\omega}\right)(\boldsymbol{x}, \boldsymbol{\lambda})\right) . \tag{34}
\end{equation*}
$$

By Lemma 5.8 below and composition, $\boldsymbol{L}$ is smooth so its fixed point $\boldsymbol{\lambda}(t, \boldsymbol{x})$ depends smoothly on $(t, \boldsymbol{x})$. By the first point of Lemma 5.8, $\left(\psi^{*} \omega\right)[t, \boldsymbol{x}]=\left(\psi^{*} \widetilde{\omega}\right)(\boldsymbol{x}, \boldsymbol{\lambda}(t, \boldsymbol{x}))$ depends smoothly on $(t, \boldsymbol{x})$. This proves Proposition 5.1(c).

## Lemma 5.8.

(a) $\psi^{*} \widetilde{\omega}$ is a smooth function of $\boldsymbol{x}$ in an $\ell^{\infty}$-neighborhood of $\underline{\boldsymbol{x}}$ and $\boldsymbol{\lambda} \in \ell^{\infty}$, with value in $C^{0}(\widetilde{\Omega})$.
(b) $\widetilde{\boldsymbol{L}}$ is a smooth function of $t$ in a neighborhood of $0, \boldsymbol{x}$ in an $\ell^{\infty}$-neighborhood of $\underline{\boldsymbol{x}}$ and $\eta \in C^{0}(\widetilde{\Omega})$, with value in $\ell^{\infty}$.

For any infinite set $K$, if $\left(V_{k}\right)_{k \in K}$ is a sequence of normed spaces, we denote

$$
\left(\bigoplus_{k \in K} V_{k}\right)_{\infty}=\left\{x \in \prod_{k \in K} V_{k}:\|x\|_{\infty}=\sup _{k \in K}\left\|x_{k}\right\|<\infty\right\}
$$

To prove Lemma 5.8, we use the following elementary fact.
Proposition 5.9. For $k \in K$, let $f_{k}: B(0, r) \subset U_{k} \rightarrow V_{k}$ be a smooth function between normed spaces. Assume that there exists uniform constants $C(m)$ such that

$$
\forall m \in \mathbb{N}, \quad \forall k \in K, \quad \forall x_{k} \in B(0, r), \quad\left\|d^{m} f_{k}\left(x_{k}\right)\right\| \leq C(m)
$$

where $d^{m} f_{k}$ denotes the $m$-th order differential of $f_{k}$. Let $U^{\infty}=\left(\bigoplus_{k \in K} U_{k}\right)_{\infty}$ and $V^{\infty}=$ $\left(\oplus_{k \in K} V_{k}\right)_{\infty}$. Define $\boldsymbol{f}: B(0, r) \subset U^{\infty} \rightarrow V^{\infty}$ by $\boldsymbol{f}(\boldsymbol{x})=\left(f_{k}\left(x_{k}\right)\right)_{k \in K}$. Then $\boldsymbol{f}$ is smooth and $d \boldsymbol{f}(\boldsymbol{x}) \boldsymbol{h}=\left(d f_{k}\left(x_{k}\right) h_{k}\right)_{k \in K}$.

We summarize the hypothesis of Proposition 5.9 by saying that the functions $f_{k}$ have uniformly bounded derivatives.

Proof. It is straightforward to prove that $f$ is differentiable (with the indicated differential) using Taylor Formula with integral remainder. Smoothness follows by induction.

Recall that $\Omega_{k, r}$ denotes the torus $T_{k}$ minus the disks $\left|z_{k}^{ \pm}\right| \leq r$. We fix a uniform positive $\varepsilon^{\prime \prime}<\varepsilon$ so that for all $\boldsymbol{x}$ in a neighborhood of $\underline{\boldsymbol{x}}$ and $k \in \mathbb{Z}, \psi_{k}\left(\widetilde{\Omega}_{k}\right) \subset \Omega_{k, \varepsilon^{\prime \prime}}$. We choose $\rho$ such that $\rho \leq \varepsilon^{\prime \prime} / 4$.

## Claim 5.10.

(a) For $k \in \mathbb{Z}, \psi_{k}^{*} \omega_{k, 0_{k}, v_{k}}$ is a smooth function of $\left(\tau_{k}, v_{k}\right)$ in a neighborhood of $\left(\underline{\tau}_{k}, \underline{v}_{k}\right)$, with value in $C^{0}\left(\widetilde{\Omega}_{k}\right)$ and has uniformly bounded derivatives.
(b) For $k \in \mathbb{Z}$ and $n \geq 2$, $\left(\frac{\varepsilon^{\prime \prime}}{2}\right)^{n-1} \psi_{k}^{*} \omega_{k, n}^{ \pm}$is a smooth function of $x_{k}=\left(a_{k}, b_{k}, v_{k}, \tau_{k}\right)$ in a neighborhood of $\underline{x}_{k}$, with value in $C^{0}\left(\widetilde{\Omega}_{k}\right)$, and has uniformly (with respect to $k$ and n) bounded derivatives.
(c) If $t^{2}<\rho \varepsilon^{\prime \prime}$, then for $k \in \mathbb{Z}$ and $n \geq 2,\left(\frac{t^{2}}{\rho z_{k}^{ \pm} \circ \psi_{k}}\right)^{n-1}$ restricted to $\widetilde{\gamma}_{k}^{ \pm}$is a smooth function of $t$ in a neighborhood of 0 and $x_{k}$ in a neighborhood of $\underline{x}_{k}$, with value in $C^{0}\left(\widetilde{\gamma}_{k}^{ \pm}\right)$, and has uniformly bounded derivatives.

Proof.
(a) follows from the explicit formula in Lemma 5.3 and Hypothesis 1.2.
(b) By (28), we have, for some uniform constant $C$

$$
\left\|\eta_{k, n}^{ \pm}\right\|_{C^{0}\left(\Omega_{k, \varepsilon^{\prime \prime}}\right)} \leq C\left(\frac{2}{\varepsilon^{\prime \prime}}\right)^{n-1}
$$

Hence

$$
\left\|\psi_{k}^{*} \eta_{k, n}^{ \pm}\right\|_{C^{0}\left(\widetilde{\Omega}_{k}\right)} \leq 2 C\left(\frac{2}{\varepsilon^{\prime \prime}}\right)^{n-1}
$$

Observe that $\psi_{k}$ depends holomorphically on $\tau_{k}$. Also, $\eta_{k, n}^{ \pm}$depends holomorphically on $x_{k}\left(\omega_{k, n}^{ \pm}\right.$does not). Hence for $z \in \widetilde{\Omega}_{k}, \psi_{k}^{*} \eta_{k, n}^{ \pm}(z)$ depends holomorphically on $x_{k}$. By Cauchy Estimate, restricting the parameter $\boldsymbol{x}$ to a smaller neighborhood of $\underline{\boldsymbol{x}}$, we have

$$
\left\|d^{m} \psi_{k}^{*} \eta_{k, n}^{ \pm}\right\|_{C^{0}\left(\widetilde{\Omega}_{k}\right)} \leq C(m)\left(\frac{2}{\varepsilon^{\prime \prime}}\right)^{n-1}
$$

for some uniform constants $C(m)$, where $d^{m}$ denotes the $m$-th order differential with respect to $x_{k}$. (b) follows from (26).
(c) Since $\psi_{k}\left(\widetilde{\gamma}_{k}^{ \pm}\right) \subset \psi_{k}\left(\widetilde{\Omega}_{k}\right) \subset \Omega_{k, \varepsilon^{\prime \prime}}$, we have $\left|z_{k}^{ \pm} \circ \psi_{k}\right| \geq \varepsilon^{\prime \prime}$ on $\widetilde{\gamma}_{k}^{ \pm}$. Hence

$$
\left\|\left(\frac{t^{2}}{\rho z_{k}^{ \pm} \circ \psi_{k}}\right)^{n-1}\right\|_{C^{0}\left(\widetilde{\gamma}_{k}^{ \pm}\right)} \leq 1
$$

Since $z_{k}^{ \pm} \circ \psi_{k}(z)$ depends holomorphically on $x_{k}$, we obtain uniform estimates of the derivatives by Cauchy Estimate.

Proof of Lemma 5.8. We write $\widetilde{\omega}_{1}$ for the first term in the definition of $\widetilde{\omega}$ and $\widetilde{\omega}_{2}$ for the second term (the sum for $n \geq 2$ ).
(a) $\left(\psi^{*} \widetilde{\omega}_{1}\right)(\boldsymbol{x}) \in C^{0}(\widetilde{\Omega})$ and is a smooth function of $\boldsymbol{x}$. This follows from Claim 5.10(a) and Proposition 5.9. $\left(\psi^{*} \widetilde{\omega}_{2}\right)(\boldsymbol{x}, \boldsymbol{\lambda}) \in C^{0}(\widetilde{\Omega})$ and is a smooth function of $\boldsymbol{x}$ and $\boldsymbol{\lambda}$. This follows from Claim 5.10(b), Proposition 5.9 and the fact that the bilinear operator

$$
\begin{equation*}
(\eta, \boldsymbol{\lambda}) \mapsto\left(\sum_{n=2}^{\infty}\left(\frac{2 \rho}{\varepsilon^{\prime \prime}}\right)^{n-1}\left(\lambda_{k, n}^{+} \eta_{k, n}^{+}+\lambda_{k, n}^{-} \eta_{k, n}^{-}\right)\right)_{k \in \mathbb{Z}} \tag{35}
\end{equation*}
$$

is bounded from $\left(\bigoplus_{(k, n, s) \in \Lambda} C^{0}\left(\widetilde{\Omega}_{k}\right)\right)_{\infty} \times \ell^{\infty}$ to $C^{0}(\widetilde{\Omega})$. (Since $\rho \leq \varepsilon^{\prime \prime} / 4$, it has norm at most 2). This proves Lemma 5.8(a).
(b) Lemma 5.8(b) follows from Claim 5.10(c), Proposition 5.9 and the fact that the bilinear operator

$$
\begin{equation*}
(\boldsymbol{f}, \eta) \mapsto\left(\frac{-1}{2 \pi \mathrm{i}} \int_{\widetilde{\gamma}_{k+s}^{-s}} f_{k+s, n}^{-s} \eta\right)_{(k, n, s) \in \Lambda} \tag{36}
\end{equation*}
$$

is bounded from $\left(\bigoplus_{(k, n, s) \in \Lambda} C^{0}\left(\widetilde{\gamma}_{k}^{s}\right)\right)_{\infty} \times C^{0}(\widetilde{\Omega})$ to $\ell^{\infty}$.

## 6. Asymptotic Behavior

Assume that we are given two configurations $\left(q_{k}\right)$ and $\left(q_{k}^{\prime}\right)$ such that $q_{k}=q_{k}^{\prime}$ for all $k \geq 0$. In this section we prove that the holomorphic 1 -form $\omega^{\prime}$ and the parameters $\boldsymbol{x}^{\prime}$ are asymptotic to $\omega$ and $\boldsymbol{x}$. These results were used in Section 4 to prove the asymptotic behaviors of the minimal surfaces.
6.1. Asymptotic behavior of $\omega$. For any infinite set $K$ equiped with a weight function $\sigma: K \rightarrow[1, \infty)$, we define the Banach space $\ell^{\infty, \sigma}(K)$ as the space of sequences of complex numbers $\boldsymbol{u}=\left(u_{k}\right)_{k \in K}$ such that the norm

$$
\|\boldsymbol{u}\|_{\infty, \sigma}=\sup _{k \in K} \sigma(k)\left|u_{k}\right|
$$

is finite. The argument $K$ will be omitted if it is clear in the context.
If $\left(V_{k}\right)_{k \in K}$ is a sequence of normed spaces, we define

$$
\left(\bigoplus_{k \in K} V_{k}\right)_{\infty, \sigma}=\left\{\boldsymbol{x} \in \prod_{k \in K} V_{k}:\|\boldsymbol{x}\|_{\infty, \sigma}=\sup _{k \in K} \sigma(k)\left\|x_{k}\right\|<\infty\right\}
$$

Let $\underline{\boldsymbol{x}}$ and $\underline{\boldsymbol{x}}^{\prime}$ be the central values corresponding to the configurations $\left(q_{k}\right)$ and $\left(q_{k}^{\prime}\right)$, as given by (4). Our goal is to compare $\omega[t, \boldsymbol{x}]$ to $\omega\left[t, \boldsymbol{x}^{\prime}\right]$ for $\boldsymbol{x}, \boldsymbol{x}^{\prime}$ in a neighborhood of $\underline{\boldsymbol{x}}, \underline{\boldsymbol{x}}^{\prime}$. For this purpose, we replace the definition of $\widetilde{\Omega}_{k}$ and $\widetilde{\Omega}$ in Section 3.3 by

$$
\widetilde{\Omega}_{k}=\mathbb{T} \backslash\left(\bar{D}\left(0, \varepsilon^{\prime}\right) \cup \bar{D}\left(\underline{\widetilde{v}}_{k}, \varepsilon^{\prime}\right) \cup \bar{D}\left(\underline{v}_{k}^{\prime}, \varepsilon^{\prime}\right)\right) \quad \text { and } \quad \widetilde{\Omega}=\bigsqcup_{k \in \mathbb{Z}} \widetilde{\Omega}_{k}
$$

Now $\left(\psi^{*} \omega\right)[t, \boldsymbol{x}]$ and $\left(\psi^{*} \omega\right)\left[t, \boldsymbol{x}^{\prime}\right]$ are both defined on the same fixed domain $\widetilde{\Omega}$ so we can compare them.

Fix $\delta>1$ and define the weight $\sigma: \mathbb{Z} \rightarrow[1, \infty)$ by $\sigma(k)=1$ if $k \leq 0$ and $\sigma(k)=\delta^{k}$ if $k \geq 0$. We extend this weight to $\Lambda$ by $\sigma(k, n, s)=\sigma(k)$ for all $(k, n, s) \in \Lambda$. We have $\underline{x}_{k}=\underline{x}_{k}^{\prime}$ for all $k \geq 0$ so $\underline{\boldsymbol{x}}^{\prime}-\underline{\boldsymbol{x}} \in \ell^{\infty, \sigma}\left(\mathbb{Z}, \mathbb{C}^{4}\right)$. It will be convenient to write

$$
\Delta \boldsymbol{x}=\boldsymbol{x}^{\prime}-\boldsymbol{x}
$$

We define the weighted space $C^{0, \sigma}(\widetilde{\Omega})$ (not to be confused with a Hölder space) by

$$
C^{0, \sigma}(\widetilde{\Omega})=\left(\bigoplus_{k \in \mathbb{Z}} C^{0}\left(\widetilde{\Omega}_{k}\right)\right)_{\infty, \sigma}
$$

So functions in $C^{0, \sigma}(\widetilde{\Omega})$ decay like $\delta^{-k}$ in $\widetilde{\Omega}_{k}$ as $k \rightarrow+\infty$.
Proposition 6.1. For $t$ small enough, $\boldsymbol{x}$ in an $\ell^{\infty}$-neighborhood of $\underline{\boldsymbol{x}}$ and $\Delta \boldsymbol{x}$ in an $\ell^{\infty, \sigma}$ neighborhood of $\Delta \underline{\boldsymbol{x}}=\underline{\boldsymbol{x}}^{\prime}-\underline{\boldsymbol{x}}$,

$$
\left(\psi^{*} \omega\right)[t, \boldsymbol{x}+\Delta \boldsymbol{x}]-\left(\psi^{*} \omega\right)[t, \boldsymbol{x}] \in C^{0, \sigma}(\widetilde{\Omega})
$$

and depends smoothly on $t, \boldsymbol{x} \in \ell^{\infty}\left(\mathbb{Z}, \mathbb{C}^{4}\right)$ and $\Delta \boldsymbol{x} \in \ell^{\infty, \sigma}\left(\mathbb{Z}, \mathbb{C}^{4}\right)$.
Proof. Recall that $\boldsymbol{\lambda}(t, \boldsymbol{x})$ is given by Lemma 5.6 and define

$$
\boldsymbol{\mu}\left(t, \boldsymbol{x}, \boldsymbol{x}^{\prime}\right)=\boldsymbol{\lambda}(t, \boldsymbol{x}+\Delta \boldsymbol{x})-\boldsymbol{\lambda}(t, \boldsymbol{x}) \in \ell^{\infty}(\Lambda)
$$

Our goal is to prove that $\boldsymbol{\mu}(t, \boldsymbol{x}, \Delta \boldsymbol{x}) \in \ell^{\infty, \sigma}(\Lambda)$.
Let

$$
\Delta \boldsymbol{L}(t, \boldsymbol{x}, \Delta \boldsymbol{x}, \boldsymbol{\lambda}, \Delta \boldsymbol{\lambda})=\boldsymbol{L}(t, \boldsymbol{x}+\Delta \boldsymbol{x}, \boldsymbol{\lambda}+\Delta \boldsymbol{\lambda})-\boldsymbol{L}(t, \boldsymbol{x}, \boldsymbol{\lambda})
$$

Then
$\Delta \boldsymbol{L}(t, \boldsymbol{x}, \Delta \boldsymbol{x}, \boldsymbol{\lambda}(t, \boldsymbol{x}), \boldsymbol{\mu}(t, \boldsymbol{x}, \Delta \boldsymbol{x}))=\boldsymbol{L}(t, \boldsymbol{x}+\Delta \boldsymbol{x}, \boldsymbol{\lambda}(t, \boldsymbol{x}+\Delta \boldsymbol{x}))-\boldsymbol{L}(t, \boldsymbol{x}, \boldsymbol{\lambda}(t, \boldsymbol{x}))=\boldsymbol{\mu}(t, \boldsymbol{x}, \Delta \boldsymbol{x})$.
In other words, $\boldsymbol{\mu}(t, \boldsymbol{x}, \Delta \boldsymbol{x})$ is a fixed point of $\Delta \boldsymbol{L}$ with respect to the $\Delta \boldsymbol{\lambda}$ variable.
Define

$$
\Delta \widetilde{\boldsymbol{L}}(t, \boldsymbol{x}, \Delta \boldsymbol{x}, \eta, \Delta \eta)=\widetilde{\boldsymbol{L}}(t, \boldsymbol{x}+\Delta \boldsymbol{x}, \eta+\Delta \eta)-\widetilde{\boldsymbol{L}}(t, \boldsymbol{x}, \eta)
$$

and

$$
\boldsymbol{H}(\boldsymbol{x}, \Delta \boldsymbol{x}, \boldsymbol{\lambda}, \Delta \boldsymbol{\lambda})=\left(\psi^{*} \widetilde{\omega}\right)(\boldsymbol{x}+\Delta \boldsymbol{x}, \boldsymbol{\lambda}+\Delta \boldsymbol{\lambda})-\left(\psi^{*} \widetilde{\omega}\right)(\boldsymbol{x}, \boldsymbol{\lambda})
$$

Recalling (34), we have

$$
\Delta \boldsymbol{L}(t, \boldsymbol{x}, \Delta \boldsymbol{x}, \boldsymbol{\lambda}, \Delta \boldsymbol{\lambda})=\Delta \widetilde{\boldsymbol{L}}\left(t, \boldsymbol{x}, \Delta \boldsymbol{x},\left(\psi^{*} \widetilde{\omega}\right)(\boldsymbol{x}, \boldsymbol{\lambda}), \boldsymbol{H}(\boldsymbol{x}, \Delta \boldsymbol{x}, \boldsymbol{\lambda}, \Delta \boldsymbol{\lambda})\right)
$$

By Lemma 6.2 below, $\Delta \boldsymbol{L} \in \ell^{\infty, \sigma}$ is smooth with respect to $\boldsymbol{x}, \boldsymbol{\lambda} \in \ell^{\infty}$ and $\Delta \boldsymbol{x}, \Delta \boldsymbol{\lambda} \in \ell^{\infty, \sigma}$, and is contracting with respect to $\Delta \boldsymbol{\lambda}$, so its unique fixed point $\boldsymbol{\mu}(t, \boldsymbol{x}, \Delta \boldsymbol{x})$ is in $\ell^{\infty, \sigma}$ and depends smoothly on $t, \boldsymbol{x}, \Delta \boldsymbol{x}$ in their respective spaces. Finally, we have by definition

$$
\begin{aligned}
\left(\psi^{*} \omega\right)[t, \boldsymbol{x}+\Delta \boldsymbol{x}]-\left(\psi^{*} \omega\right)[t, \boldsymbol{x}] & =\left(\psi^{*} \widetilde{\omega}\right)(\boldsymbol{x}+\Delta \boldsymbol{x}, \boldsymbol{\lambda}(t, \boldsymbol{x}+\Delta \boldsymbol{x}))-\left(\psi^{*} \widetilde{\omega}\right)(\boldsymbol{x}, \boldsymbol{\lambda}(t, \boldsymbol{x})) \\
& =\boldsymbol{H}(\boldsymbol{x}, \Delta \boldsymbol{x}, \boldsymbol{\lambda}(t, \boldsymbol{x}), \boldsymbol{\mu}(t, \boldsymbol{x}, \Delta \boldsymbol{x}))
\end{aligned}
$$

so Proposition 6.1 follows from Lemma 6.2(a).

## Lemma 6.2.

 $\Delta \boldsymbol{\lambda} \in \ell^{\infty, \sigma}$,

$$
\boldsymbol{H}(\boldsymbol{x}, \Delta \boldsymbol{x}, \boldsymbol{\lambda}, \Delta \boldsymbol{\lambda}) \in C^{0, \sigma}(\widetilde{\Omega})
$$

and depends smoothly on $\boldsymbol{x}, \Delta \boldsymbol{x}, \boldsymbol{\lambda}, \Delta \boldsymbol{\lambda}$.
(b) Fort in a neighborhood of $0, \boldsymbol{x}$ in an $\ell^{\infty}$-neighborhood of $\underline{\boldsymbol{x}}, \Delta \boldsymbol{x}$ in an $\ell^{\infty, \sigma}$-neighborhood of $\Delta \underline{\boldsymbol{x}}, \eta \in C^{0}(\widetilde{\Omega})$ and $\Delta \eta \in C^{0, \sigma}(\widetilde{\Omega})$,

$$
\Delta \widetilde{\boldsymbol{L}}(t, \boldsymbol{x}, \Delta \boldsymbol{x}, \eta, \Delta \eta) \in \ell^{\infty, \sigma}
$$

and depends smoothly on $t, \boldsymbol{x}, \Delta \boldsymbol{x}, \eta, \Delta \eta$.
(c) For $t$ small enough, $\Delta \boldsymbol{L}$ is contracting with respect to $\Delta \boldsymbol{\lambda}$, as a map from $\ell^{\infty, \sigma}$ to itself.

We need the following
Proposition 6.3. Under the same notations and hypothesis as in Proposition 5.9, let $\sigma$ : $K \rightarrow[1, \infty)$ be an arbitrary weight. Let $U^{\infty, \sigma}=\left(\bigoplus_{k \in K} U_{k}\right)_{\infty, \sigma}$ and $V^{\infty, \sigma}=\left(\bigoplus_{k \in K} V_{k}\right)_{\infty, \sigma}$. Define for $\boldsymbol{x} \in B(\mathbf{0}, r / 2) \subset U^{\infty}$ and $\Delta \boldsymbol{x} \in B(\mathbf{0}, r / 2) \subset U^{\infty, \sigma}$

$$
\Delta \boldsymbol{f}(\boldsymbol{x}, \Delta \boldsymbol{x})=\left(f_{k}\left(x_{k}+\Delta x_{k}\right)-f_{k}\left(x_{k}\right)\right)_{k \in K}
$$

Then $\Delta \boldsymbol{f}(\boldsymbol{x}, \Delta \boldsymbol{x}) \in V^{\infty, \sigma}, \Delta \boldsymbol{f}$ is smooth and

$$
d(\Delta \boldsymbol{f})(\boldsymbol{x}, \Delta \boldsymbol{x})(\boldsymbol{h}, \Delta \boldsymbol{h})=\left(d f_{k}\left(x_{k}+\Delta x_{k}\right)\left(h_{k}+\Delta h_{k}\right)-d f_{k}\left(x_{k}\right) h_{k}\right)_{k \in K}
$$

Proof. By the Mean Value Inequality

$$
\sigma(k)\left\|f_{k}\left(x_{k}+\Delta x_{k}\right)-f_{k}\left(x_{k}\right)\right\| \leq C \sigma(k)\left\|\Delta x_{k}\right\|
$$

Hence $\Delta \boldsymbol{f}(\boldsymbol{x}, \Delta \boldsymbol{x}) \in V^{\infty, \sigma}$. Define

$$
l_{k}\left(h_{k}, \Delta h_{k}\right)=d f_{k}\left(x_{k}+\Delta x_{k}\right)\left(h_{k}+\Delta h_{k}\right)-d f_{k}\left(x_{k}\right) h_{k} \quad \text { and } \quad \boldsymbol{l}(\boldsymbol{h}, \Delta \boldsymbol{h})=\left(l_{k}\left(h_{k}, \Delta h_{k}\right)\right)_{k \in K}
$$

Using the Mean Value Inequality, one easily obtains

$$
\sigma(k)\left\|l_{k}\left(h_{k}, \Delta h_{k}\right)\right\| \leq C \sigma(k)\left(\left\|\Delta h_{k}\right\|+\left\|\Delta x_{k}\right\|\left\|h_{k}\right\|\right)
$$

Hence $\boldsymbol{l}$ is a bounded operator from $U^{\infty} \times U^{\infty, \sigma}$ to $V^{\infty, \sigma}$. Using Taylor Formula with integral remainder, we have

$$
\begin{aligned}
& \Delta f_{k}\left(x_{k}+h_{k}, \Delta x_{k}+\Delta h_{k}\right)-\Delta f_{k}\left(x_{k}, \Delta x_{k}\right)-l_{k}\left(h_{k}, \Delta h_{k}\right) \\
= & f_{k}\left(x_{k}+\Delta x_{k}+h_{k}+\Delta h_{k}\right)-f_{k}\left(x_{k}+\Delta x_{k}\right) \\
& \quad-d f_{k}\left(x_{k}+\Delta x_{k}\right)\left(h_{k}+\Delta h_{k}\right)-\left[f_{k}\left(x_{k}+h_{k}\right)-f_{k}\left(x_{k}\right)-d f_{k}\left(x_{k}\right) h_{k}\right] \\
= & \int_{0}^{1}(1-t)\left[d^{2} f_{k}\left(x_{k}+\Delta x_{k}+t\left(h_{k}+\Delta h_{k}\right)\right)\left(h_{k}+\Delta h_{k}\right)^{2}-d^{2} f_{k}\left(x_{k}+t h_{k}\right) h_{k}^{2}\right] d t
\end{aligned}
$$

By the Mean Value Inequality

$$
\begin{aligned}
& \sigma(k)\left\|\Delta f_{k}\left(x_{k}+h_{k}, \Delta x_{k}+\Delta h_{k}\right)-\Delta f_{k}\left(x_{k}, \Delta x_{k}\right)-l_{k}\left(h_{k}, \Delta h_{k}\right)\right\| \\
& \leq C \sigma(k)\left[2\left\|h_{k}\right\|\left\|\Delta h_{k}\right\|+\left\|\Delta h_{k}\right\|^{2}+\left(\left\|\Delta x_{k}\right\|+\left\|\Delta h_{k}\right\|\right)\left\|h_{k}\right\|^{2}\right] .
\end{aligned}
$$

Hence

$$
\|\Delta \boldsymbol{f}(\boldsymbol{x}+\boldsymbol{h}, \Delta \boldsymbol{x}+\Delta \boldsymbol{h})-\Delta \boldsymbol{f}(\boldsymbol{x}, \Delta \boldsymbol{x})-\boldsymbol{l}(\boldsymbol{h}, \Delta \boldsymbol{h})\|_{\infty, \sigma}=O\left(\left(\|\boldsymbol{h}\|_{\infty}+\|\Delta \boldsymbol{h}\|_{\infty, \sigma}\right)^{2}\right)
$$

so $\Delta \boldsymbol{f}$ is differentiable with $d(\Delta \boldsymbol{f})(\boldsymbol{x}, \Delta \boldsymbol{x})=\boldsymbol{l}$. Smoothness follows by induction.

We shall use the following corollary with $K=\mathbb{Z}, K^{+}=\mathbb{N}$ and $K^{-}=\mathbb{Z} \backslash \mathbb{N}$ :
Corollary 6.4. With the same notation as in Proposition 6.3, assume that for $k \in K, f_{k}$ is defined in $B\left(\underline{x}_{k}, r\right) \cup B\left(\underline{x}_{k}^{\prime}, r\right)$ in $U_{k}$ and has uniformly bounded derivatives. Assume that $K$ admits a partition $\left(K^{+}, K^{-}\right)$such that for all $k \in K^{+}, \underline{x}_{k}=\underline{x}_{k}^{\prime}$, and for all $k \in K^{-}, \sigma(k)=1$. Then for $\boldsymbol{x} \in B(\underline{\boldsymbol{x}}, r / 2) \subset U^{\infty}$ and $\Delta \boldsymbol{x} \in B\left(\underline{\boldsymbol{x}}^{\prime}-\underline{\boldsymbol{x}}, r / 2\right) \subset U^{\infty, \sigma}, \Delta \boldsymbol{f}(\boldsymbol{x}, \Delta \boldsymbol{x}) \in V^{\infty, \sigma}$ and depends smoothly on $\boldsymbol{x}$ and $\Delta \boldsymbol{x}$.

Proof. We decompose a sequence $\boldsymbol{x}=\left(x_{k}\right)_{k \in K}$ as $\boldsymbol{x}=\boldsymbol{x}^{+}+\boldsymbol{x}^{-}$with $\boldsymbol{x}^{+}$supported on $K^{+}$ and $\boldsymbol{x}^{-}$supported on $K^{-}$. By Proposition 5.9, $\boldsymbol{f}^{-}\left(\boldsymbol{x}^{-}\right) \in V^{\infty}$ and $\boldsymbol{f}^{-}\left(\boldsymbol{x}^{-}+\Delta \boldsymbol{x}^{-}\right) \in V^{\infty}$ so $\Delta \boldsymbol{f}^{-}\left(\boldsymbol{x}^{-}, \Delta \boldsymbol{x}^{-}\right) \in V^{\infty}$. Since $\sigma=1$ on $K^{-}, \Delta \boldsymbol{f}^{-}\left(\boldsymbol{x}^{-}, \Delta \boldsymbol{x}^{-}\right) \in V^{\infty, \sigma}$. Since $\underline{\boldsymbol{x}}^{+}=\left(\underline{\boldsymbol{x}}^{\prime}\right)^{+}$, we may use the change of variable $\boldsymbol{x}^{+}=\underline{\boldsymbol{x}}^{+}+\boldsymbol{y}^{+}$and conclude that $\Delta \boldsymbol{f}^{+}\left(\boldsymbol{x}^{+}, \Delta \boldsymbol{x}^{+}\right) \in V^{\infty, \sigma}$ by Proposition 6.3.

Proof of Lemma 6.2. Lemma 6.2(a) follows from the following.

- $\left(\psi^{*} \widetilde{\omega}_{1}\right)(\boldsymbol{x}+\Delta \boldsymbol{x})-\left(\psi^{*} \widetilde{\omega}_{1}\right)(\boldsymbol{x}) \in C^{0, \sigma}(\widetilde{\Omega})$ and is a smooth function of $\boldsymbol{x}, \Delta \boldsymbol{x}$. This follows from Claim 5.10(a) and Corollary 6.4.
- $\left(\psi^{*} \widetilde{\omega}_{2}\right)(\boldsymbol{x}+\Delta \boldsymbol{x}, \boldsymbol{\lambda})-\left(\psi^{*} \widetilde{\omega}_{2}\right)(\boldsymbol{x}, \boldsymbol{\lambda}) \in C^{0, \sigma}(\widetilde{\Omega})$ and is a smooth function of $\boldsymbol{x}, \Delta \boldsymbol{x}$ and $\boldsymbol{\lambda}$. This follows from Claim 5.10 (b), Corollary 6.4 and the fact that the bilinear operator (35) is bounded from $\left(\bigoplus_{(k, n, s) \in \Lambda} C^{0}\left(\widetilde{\Omega}_{k}\right)\right)_{\infty, \sigma} \times \ell^{\infty}(\Lambda)$ to $C^{0, \sigma}(\widetilde{\Omega})$.
- $\left(\psi^{*} \widetilde{\omega}_{2}\right)(\boldsymbol{x}+\Delta \boldsymbol{x}, \Delta \boldsymbol{\lambda}) \in C^{0, \sigma}(\widetilde{\Omega})$ and is a smooth function of $\boldsymbol{x}, \Delta \boldsymbol{x}$ and $\Delta \boldsymbol{\lambda}$. This follows from Claim $5.10(\mathrm{~b})$, Proposition 5.9 and the fact that the bilinear operator (35) is bounded from $\left(\bigoplus_{(k, n, s) \in \Lambda} C^{0}\left(\widetilde{\Omega}_{k}\right)\right)_{\infty} \times \ell^{\infty, \sigma}(\Lambda)$ to $C^{0, \sigma}(\widetilde{\Omega})$.

Lemma $6.2(\mathrm{~b})$ follows from the following.

- $\widetilde{\boldsymbol{L}}(t, \boldsymbol{x}+\Delta \boldsymbol{x}, \eta)-\widetilde{\boldsymbol{L}}(t, \boldsymbol{x}, \eta) \in \ell^{\infty, \sigma}(\Lambda)$ and depends smoothly on $t, \boldsymbol{x}, \Delta \boldsymbol{x}$ and $\eta$. This follows from Claim 5.10 (c), Corollary 6.4 and the fact that the bilinear operator (36) is bounded from $\left(\bigoplus_{(k, n, s) \in \Lambda} C^{0}\left(\widetilde{\gamma}_{k}^{s}\right)\right)_{\infty, \sigma} \times C^{0}(\widetilde{\Omega})$ to $\ell^{\infty, \sigma}(\Lambda)$. This uses that $\frac{\sigma(k)}{\sigma(k \pm 1)} \leq \delta$ and explains our choice of the weight $\sigma$.
- $\widetilde{\boldsymbol{L}}(t, \boldsymbol{x}+\Delta \boldsymbol{x}, \Delta \eta) \in \ell^{\infty, \sigma}(\Lambda)$ and depends smoothly on $t, \boldsymbol{x}, \Delta \boldsymbol{x}$ and $\Delta \eta$. This follows from Claim 5.10(c), Proposition 5.9 and the fact that the bilinear operator (36) is bounded from $\left(\bigoplus_{(k, n, s) \in \Lambda} C^{0}\left(\widetilde{\gamma}_{k}^{s}\right)\right)_{\infty} \times C^{0, \sigma}(\widetilde{\Omega})$ to $\ell^{\infty, \sigma}(\Lambda)$.
Finally, we have

$$
\Delta \boldsymbol{L}(t, \boldsymbol{x}, \Delta \boldsymbol{x}, \boldsymbol{\lambda}, \Delta \boldsymbol{\lambda})-\Delta \boldsymbol{L}(t, \boldsymbol{x}, \Delta \boldsymbol{x}, \boldsymbol{\lambda}, 0)=\boldsymbol{L}(t, \boldsymbol{x}+\Delta \boldsymbol{x}, \boldsymbol{\lambda}+\Delta \boldsymbol{\lambda})-\boldsymbol{L}(t, \boldsymbol{x}+\Delta \boldsymbol{x}, \boldsymbol{\lambda})
$$

Using Estimate (30) as in the proof of Lemma 5.6,

$$
\|\boldsymbol{L}(t, \boldsymbol{x}+\Delta \boldsymbol{x}, \boldsymbol{\lambda}+\Delta \boldsymbol{\lambda})-\boldsymbol{L}(t, \boldsymbol{x}+\Delta \boldsymbol{x}, \boldsymbol{\lambda})\|_{\infty, \sigma} \leq \frac{C \ell \delta t^{2}}{\pi \rho \varepsilon}\|\Delta \boldsymbol{\lambda}\|_{\infty, \sigma}
$$

so $\Delta \boldsymbol{L}$ is contracting with respect to $\Delta \boldsymbol{\lambda}$ for $t$ small enough.
6.2. Asymptotic behavior of the parameters. Let $\boldsymbol{x}(t)$ and $\boldsymbol{x}^{\prime}(t)$ be the solutions obtained in Section 3 from the configurations $\left(q_{k}\right)$ and $\left(q_{k}^{\prime}\right)$, respectively.
Proposition 6.5. For $t$ small enough, $\boldsymbol{x}^{\prime}(t)-\boldsymbol{x}(t) \in \ell^{\infty, \sigma}$.
Proof. Recall the definition of $\mathcal{E}_{k}$ in Section 3.4, $\mathcal{P}_{k, 1}$ and $\mathcal{P}_{k, 2}$ in Section 3.5 and $\mathcal{G}_{k}$ in Section 3.6. We have solved equations by three consecutive applications of the Implicit Function Theorem. But we could have solved all of them by one single application. Indeed, consider the change of parameter

$$
b_{k}=-a_{k} \xi\left(v_{k}, \tau_{k}\right)+\widehat{b}_{k}
$$

By the computations in Sections 3.4, 3.5 and 3.6 , the jacobian of $\left(\mathcal{E}_{k}, \mathcal{P}_{k, 1}, \mathcal{P}_{k, 2}, \mathcal{G}_{k}\right)$ with respect to $\left(\widehat{b}_{k}, a_{k}, \tau_{k}, v_{k}\right)$ has upper-triangular form with $\mathbb{R}$-linear automorphisms of $\mathbb{C}$ on the diagonal, whose inverses are uniformly bounded with respect to $k$. Define

$$
\mathcal{F}_{k}(t, \boldsymbol{x})=\left(\mathcal{E}_{k}(t, \boldsymbol{x}), \mathcal{P}_{k, 1}(t, \boldsymbol{x}), \mathcal{P}_{k, 2}(t, \boldsymbol{x}), \mathcal{G}_{k}(t, \boldsymbol{x})+2 \pi \mathrm{i} G\left(q_{0} ; T\right)\right) \quad \text { and } \quad \mathcal{F}=\left(\mathcal{F}_{k}\right)_{k \in \mathbb{Z}} .
$$

Then $d_{\boldsymbol{x}} \mathcal{F}\left(0, \underline{\boldsymbol{x}}^{\prime}\right)$ is an automorphism of $\ell^{\infty}$, and restricts to an automorphism of $\ell^{\infty, \sigma}$. Define, for $\boldsymbol{x}$ in an $\ell^{\infty}$-neighborhood of $\underline{\boldsymbol{x}}$ and $\Delta \boldsymbol{x}$ in an $\ell^{\infty, \sigma}$-neighborhood of $\Delta \underline{\boldsymbol{x}}=\underline{\boldsymbol{x}}^{\prime}-\underline{\boldsymbol{x}}$

$$
\Delta \mathcal{F}(t, \boldsymbol{x}, \Delta \boldsymbol{x})=\mathcal{F}(t, \boldsymbol{x}+\Delta \boldsymbol{x})-\mathcal{F}(t, \boldsymbol{x})
$$

By Lemma 6.6 below, $\Delta \boldsymbol{\mathcal { F }}(t, \boldsymbol{x}, \Delta \boldsymbol{x}) \in \ell^{\infty, \sigma}$. We have

$$
\Delta \mathcal{F}(0, \underline{\boldsymbol{x}}, \Delta \underline{\boldsymbol{x}})=\mathcal{F}\left(t, \underline{\boldsymbol{x}}^{\prime}\right)-\mathcal{F}(t, \underline{\boldsymbol{x}})=0 \quad \text { and } \quad d_{\Delta \boldsymbol{x}}(\Delta \mathcal{F})(0, \underline{\boldsymbol{x}}, \Delta \underline{\boldsymbol{x}})=d_{\boldsymbol{x}} \mathcal{F}\left(0, \underline{\boldsymbol{x}}^{\prime}\right)
$$

By the Implicit Function Theorem, for $t$ small enough and $\boldsymbol{x}$ in a neighborhood of $\boldsymbol{x}$, there exists $\Delta \boldsymbol{x}(t, \boldsymbol{x}) \in \ell^{\infty, \sigma}$ such that $\Delta \mathcal{F}(t, \boldsymbol{x}, \Delta \boldsymbol{x}(t, \boldsymbol{x}))=0$. We substitute $\boldsymbol{x}=\boldsymbol{x}(t)$ and obtain $\mathcal{F}(t, \boldsymbol{x}(t)+\Delta \boldsymbol{x}(t, \boldsymbol{x}(t)))=\mathcal{F}\left(t, \boldsymbol{x}^{\prime}(t)\right)=0$. By uniqueness, $\boldsymbol{x}^{\prime}(t)=\boldsymbol{x}(t)+\Delta \boldsymbol{x}(t, \boldsymbol{x}(t))$, which proves Proposition 6.5.
Lemma 6.6. For $t$ in a neighborhood of $0, \boldsymbol{x}$ in an $\ell^{\infty}$-neighborhood of $\underline{\boldsymbol{x}}$ and $\Delta \boldsymbol{x}$ in an $\ell^{\infty, \sigma}$-neighborhood of $\Delta \underline{\boldsymbol{x}}, \Delta \boldsymbol{\mathcal { F }}(t, \boldsymbol{x}, \Delta \boldsymbol{x}) \in \ell^{\infty, \sigma}$ and depends smoothly on $t, \boldsymbol{x}$ and $\Delta \boldsymbol{x}$.
Proof. Define for $z \in \mathbb{T}$ :

$$
f_{k}(z)=\frac{\psi_{k}(z)-\left(Z_{k, 1}+Z_{k, 2}\right)}{\psi_{k}(z)^{2}-\left(Z_{k, 1}+Z_{k, 2}\right) \psi_{k}(z)+Z_{k, 1} Z_{k, 2}} \quad \text { and } \quad \boldsymbol{f}[\boldsymbol{x}]=\left(f_{k}[\boldsymbol{x}]\right)_{k \in \mathbb{Z}}
$$

By Cauchy Theorem and a change of variable,

$$
\mathcal{E}_{k}(t, \boldsymbol{x})=\frac{1}{2 \pi \mathrm{i}} \int_{\partial \widetilde{\Omega}_{k, 2 \varepsilon^{\prime}}} f_{k} \psi_{k}^{*} \omega[t, \boldsymbol{x}]
$$

Hence we can write

$$
\mathcal{E}(t, \boldsymbol{x})=\boldsymbol{B}\left(\boldsymbol{f}[\boldsymbol{x}],\left(\psi^{*} \omega\right)[t, \boldsymbol{x}]\right) \quad \text { where } \quad \boldsymbol{B}(\boldsymbol{f}, \eta)=\left(\frac{1}{2 \pi \mathrm{i}} \int_{\partial \widetilde{\Omega}_{k, 2 \varepsilon^{\prime}}} f_{k} \eta_{k}\right)_{k \in \mathbb{Z}}
$$

Using Weierstrass Preparation Theorem, the symmetric functions of $Z_{k, 1}$ and $Z_{k, 2}$ are holomorphic functions of $x_{k}$. Hence $f_{k}$ is a smooth function of $x_{k}$ with value in $C^{0}\left(\partial \widetilde{\Omega}_{k, 2 \varepsilon^{\prime}}\right)$. Using Corollary 6.4 and that the bilinear operator

$$
\boldsymbol{B}:\left(\bigoplus_{k \in \mathbb{Z}} C^{0}\left(\partial \widetilde{\Omega}_{k, 2 \varepsilon^{\prime}}\right)\right)_{\infty, \sigma} \times C^{0}(\widetilde{\Omega}) \rightarrow \ell^{\infty, \sigma}
$$

is bounded, we conclude that

$$
\begin{equation*}
\boldsymbol{B}\left(\boldsymbol{f}[\boldsymbol{x}+\Delta \boldsymbol{x}]-\boldsymbol{f}[\boldsymbol{x}],\left(\psi^{*} \omega\right)[t, \boldsymbol{x}]\right) \in \ell^{\infty, \sigma} \tag{37}
\end{equation*}
$$

and depends smoothly on $t, \boldsymbol{x}$ and $\Delta \boldsymbol{x}$. Using Proposition 6.1 and that the bilinear operator

$$
\boldsymbol{B}:\left(\bigoplus_{k \in \mathbb{Z}} C^{0}\left(\partial \widetilde{\Omega}_{k, 2 \varepsilon^{\prime}}\right)\right)_{\infty} \times C^{0, \sigma}(\widetilde{\Omega}) \rightarrow \ell^{\infty, \sigma}
$$

is bounded, we obtain

$$
\begin{equation*}
\boldsymbol{B}\left(\boldsymbol{f}[\boldsymbol{x}+\Delta \boldsymbol{x}],\left(\psi^{*} \omega\right)[t, \boldsymbol{x}+\Delta \boldsymbol{x}]-\left(\psi^{*} \omega\right)[t, \boldsymbol{x}]\right) \in \ell^{\infty, \sigma} \tag{38}
\end{equation*}
$$

and depends smoothly on $t, \boldsymbol{x}$ and $\Delta \boldsymbol{x}$. Adding (37) and (38), we conclude that

$$
\mathcal{E}(t, \boldsymbol{x}+\Delta \boldsymbol{x})-\boldsymbol{\mathcal { E }}(t, \boldsymbol{x}) \in \ell^{\infty, \sigma}
$$

and depends smoothly on $t, \boldsymbol{x}$ and $\Delta \boldsymbol{x}$.
Finally, $\mathcal{P}_{k, 1}, \mathcal{P}_{k, 2}$ and $\mathcal{G}_{k}$ are defined as integrals of $\omega$ times powers of $g_{k}$ on certain curves in $T_{k}$, so we can deal with them in the same way as $\mathcal{E}_{k}$.

## References

[Bra92] Kenneth A. Brakke. The Surface Evolver. Experiment. Math., 1(2):141-165, 1992.
$\left[\mathrm{CCM}^{+} 06\right]$ Charlotte E. Conn, Oscar Ces, Xavier Mulet, Stephanie Finet, Roland Winter, John M. Seddon, and Richard H. Templer. Dynamics of structural transformations between lamellar and inverse bicontinuous cubic lyotropic phases. Physical review letters, 96(10):108102, 2006.
[CF97] Thierry Charitat and Bertrand Fourcade. Lattice of passages connecting membranes. Journal de Physique II, 7(1):15-35, 1997.
[Che18] Hao Chen. Minimal twin surfaces. Exp. Math., 2018. online first.
[Che19] Hao Chen. Existence of the rhombohedral and tetragonal deformation families of the gyroid. page 21 pp., 2019. preprint, arXiv:1901.04006.
[CKLW18] Zhijie Chen, Ting-Jung Kuo, Chang-Shou Lin, and Chin-Lung Wang. Green function, Painlevé VI equation, and Eisenstein series of weight one. J. Differential Geom., 108(2):185-241, 2018.
[CW18a] Hao Chen and Matthias Weber. A new deformation family of Schwarz' D surface. page 15 pp., 2018. preprint, arXiv:1804.01442.
[CW18b] Hao Chen and Matthias Weber. An orthorhombic deformation family of Schwarz' H surfaces. page 17 pp., 2018. preprint, arXiv:1807. 10631.
[Fay73] John D. Fay. Theta Functions on Riemann Surfaces, volume 352 of Lecture Notes in Mathematics. 1973.
[FH99] Andrew Fogden and Stephen T. Hyde. Continuous transformations of cubic minimal surfaces. The European Physical Journal B-Condensed Matter and Complex Systems, 7(1):91-104, 1999.
[FHL93] Andrew Fogden, Markus Haeberlein, and Sven Lidin. Generalizations of the gyroid surface. J. Phys. I, 3(12):2371-2385, 1993.
[GBW96] Karsten Große-Brauckmann and Meinhard Wohlgemuth. The gyroid is embedded and has constant mean curvature companions. Calc. Var. Partial Differential Equations, 4(6):499-523, 1996.
$\left[\mathrm{HBL}^{+} 96\right]$ Stephen T. Hyde, Zoltan Blum, Tomas Landh, Sven Lidin, Barry W. Ninham, Sten Andersson, and Kåre Larsson. The Language of Shape: The Role of Curvature in Condensed Matter: Physics, Chemistry and Biology. Elsevier Science, 1996.
[Hec27] Erich Hecke. Zur Theorie der elliptischen Modulfunktionen. Math. Ann., 97(1):210-242, 1927.
[HXBC11] Lu Han, Ping Xiong, Jingfeng Bai, and Shunai Che. Spontaneous formation and characterization of silica mesoporous crystal spheres with reverse multiply twinned polyhedral hollows. Journal of the American Chemical Society, 133(16):6106-6109, 2011.
[Lan95] Serge Lang. Introduction to modular forms, volume 222 of Grundlehren der Mathematischen Wissenschaften. Springer-Verlag, Berlin, 1995. With appendixes by D. Zagier and Walter Feit, Corrected reprint of the 1976 original.
[LW10] Chang-Shou Lin and Chin-Lung Wang. Elliptic functions, Green functions and the mean field equations on tori. Ann. of Math. (2), 172(2):911-954, 2010.
[LW17] Chang-Shou Lin and Chin-Lung Wang. On the minimality of extra critical points of Green functions on flat tori. Int. Math. Res. Not. IMRN, (18):5591-5608, 2017.
[MBF94] Xavier Michalet, David Bensimon, and Bertrand Fourcade. Fluctuating vesicles of nonspherical topology. Physical review letters, 72(1):168, 1994.
[MT12] Filippo Morabito and Martin Traizet. Non-periodic Riemann examples with handles. Adv. Math., 229(1):26-53, 2012.
[MW91] William H. Meeks and Brian White. Minimal surfaces bounded by convex curves in parallel planes. Comment. Math. Helvetici, 66:263-278, 1991.
[Sch70] Alan H. Schoen. Infinite periodic minimal surfaces without self-intersections. Technical Note D5541, NASA, Cambridge, Mass., May 1970.
[Shi56] Max Shiffman. On surfaces of stationary area bounded by two circles, or convex curves, in parallel planes. Ann. of Math. (2), 63:77-90, 1956.
$\left[\mathrm{TBC}^{+} 15\right]$ T.-Y. Dora Tang, Nicholas J. Brooks, Oscar Ces, John M. Seddon, and Richard H. Templer. Structural studies of the lamellar to bicontinuous gyroid cubic (Q GII) phase transitions under limited hydration conditions. Soft matter, 11(10):1991-1997, 2015.
[Tra02] Martin Traizet. An embedded minimal surface with no symmetries. J. Diff. Geom., 60:103-153, 2002.
[Tra08] Martin Traizet. On the genus of triply periodic minimal surfaces. J. Diff. Geom., 79(2):243-275, 2008.
[Tra10] Martin Traizet. On minimal surfaces bounded by two convex curves in parallel planes. Comment. Math. Helv., 85:39-71, 2010.
[Tra13] Martin Traizet. Opening infinitely many nodes. J. Reine Angew. Math., 684:165-186, 2013.
[Web02] Matthias Weber. Period quotient maps of meromorphic 1-forms and minimal surfaces on tori. J. Geom. Anal., 12(2):325-354, 2002.
[WHW09] Matthias Weber, David Hoffman, and Michael Wolf. An embedded genus-one helicoid. Ann. of Math. (2), 169(2):347-448, 2009.
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[^0]:    Date: August 20, 2019.
    2010 Mathematics Subject Classification. Primary 53A10.
    Key words and phrases. minimal surfaces.
    H. Chen is supported by Individual Research Grant from Deutsche Forschungsgemeinschaft within the project "Defects in Triply Periodic Minimal Surfaces", Projektnummer 398759432.

