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# Spectral measures of empirical autocovariance matrices of high dimensional Gaussian stationary processes 

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#### Abstract

Consider the empirical autocovariance matrix at a given non-zero time lag based on observations from a multivariate complex Gaussian stationary time series. The spectral analysis of these autocovariance matrices can be useful in certain statistical problems, such as those related to testing for white noise. We study the behavior of their spectral measures in the asymptotic regime where the time series dimension and the observation window length both grow to infinity, and at the same rate. Following a general framework in the field of the spectral analysis of large random non-Hermitian matrices, at first the probabilistic behavior of the small singular values of the shifted versions of the autocovariance matrix are obtained. This is then used to infer about the large sample behaviour of the empirical spectral measure of the autocovariance matrices at any lag. Matrix orthogonal polynomials on the unit circle play a crucial role in our study.


Keywords: High-dimensional times series analysis, Large non-Hermitian matrix theory, Limit spectral distribution, Matrix orthogonal polynomials, Multivariate stationary processes, Small singular values.

AMS 2020 Subject Classification: Primary 60B20; Secondary 60G57, 62M10, 33C47.

## 1 Background, assumptions and results

### 1.1 Background

Consider a multivariate time-series sequence $\left(\boldsymbol{x}^{(N)}\right)_{N=1,2, \ldots}$ where for each $N, \boldsymbol{x}^{(N)}=\left(\boldsymbol{x}_{k}^{(N)}\right)_{k \in \mathbb{Z}}$ is a $\mathbb{C}^{N}$-valued centered Gaussian stationary process in the discrete time parameter $k$. Let

$$
R_{L}^{(N)}=\mathbb{E} \boldsymbol{x}_{L}^{(N)}\left(\boldsymbol{x}_{0}^{(N)}\right)^{*}
$$

be the autocovariance matrix of $\boldsymbol{x}^{(N)}$ at lag $L$ (throughout this article, * stands for the complex conjugate). Let $\left(n_{N}\right)_{N}$ be an increasing sequence of positive integers such that

$$
0<\liminf _{N \rightarrow \infty} \frac{N}{n_{N}} \leqslant \limsup _{N \rightarrow \infty} \frac{N}{n_{N}}<\infty .
$$

Assume that for each $N$, we have the samples $\boldsymbol{x}_{0}^{(N)}, \ldots, \boldsymbol{x}_{n_{N}-1}^{(N)}$ of the process $\boldsymbol{x}^{(N)}$. Fixing an integer $L \geqslant 0$, the empirical autocovariance matrix of order $L$ is given by

$$
\widehat{R}_{L}^{(N)}=\frac{1}{n_{N}} \sum_{\ell=0}^{n_{N}-1} \boldsymbol{x}_{\ell+L}^{(N)} \boldsymbol{x}_{\ell}^{(N)^{*}} \in \mathbb{C}^{N \times N}
$$

[^0]where the sum $\ell+L$ will be taken modulo $-n_{N}$. For any matrix $M \in \mathbb{C}^{m \times m}$, let $\left\{\lambda_{0}(M), \ldots, \lambda_{m-1}(M)\right\}$ be its eigenvalues. The spectral measure of $\widehat{R}_{L}^{(N)}$ is then defined as
$$
\mu_{N}=\frac{1}{N} \sum_{\ell=0}^{N-1} \delta_{\lambda_{\ell}\left(\hat{R}_{L}^{(N)}\right)}
$$

We are interested in studying this measure as $N \rightarrow \infty$. In the field of mutlivariate time series analysis, it is classically assumed that $N$ is fixed while the observation window length converges to infinity, in which case, under standard sets of assumptions, $\mu_{N}$ converges weakly to the spectral measure of $R_{L}^{(N)}$ in the almost sure sense. This is no more true in the asymptotic regime that we consider in this paper, where the time series dimension and the window length are both large and of the same order of magnitude.

Note that for $L=0, \widehat{R}_{L}^{(N)}$ is Hermitian and several results are known for this case under different assumptions. Our aim is to consider the cases $L \geqslant 1$ in which case $\widehat{R}_{L}^{(N)}$ are nonHermitian. Generally speaking, the study of the spectral measure of non-Hermitian matrices is much harder than that for Hermitian matrices. See for example [18, 26, 14]. As an example, it took a tremendous amount of effort from researchers over a long period of time to establish the limit of the empirical spectral measure of the matrix all whose elements are real-valued iid with mean zero and variance 1 (see [33]).

In [4], using the ideas from [18, 26, 35], we identified the limit spectral measure of $\hat{R}_{L}^{(N)}$ in the particular case when the time series is a white noise process. The same setting is considered in [37], but they relax the modulo $-n_{N}$ summation when constructing $\hat{R}_{L}^{(N)}$. However, when the time series is not a white noise, no such result appears to exists in the literature. The nature of the empirical spectral measure of $\widehat{R}_{L}^{(N)}$ for different values of $L$ reflect the degree of dependence in the underlying process. Thus for example, the plots of the empirical spectral measure can serve as graphical tests for white noise or for the order of dependence in the time series. Theoretical support for these tests rests on the asymptotic behavior of the empirical spectral measure. Moreover, once this behavior is identified, it can potentially be used to develop significance tests for other statistical hypothesis.

For any matrix $M \in \mathbb{C}^{m \times m}$, let $s_{0}(M) \geqslant s_{1}(M) \geqslant \cdots \geqslant s_{m-1}(M)$ be its singular values arranged in a non-increasing order. It is well-known that for a high dimensional non-Hermitian matrices, say $M_{N}$, the behavior its spectral measure is connected to the probabilistic behavior of the small singular values of the related matrix $M_{N}-z I_{N}$ for $z \in \mathbb{C} \backslash\{0\}$. With this in mind, we shall seek solutions to the following problems. Below $I_{k}$ denotes the identity matrix of order $k$.

- The behavior of the smallest singular value $s_{N-1}\left(\widehat{R}_{L}^{(N)}-z I_{N}\right)$ for an arbitrary $z \in \mathbb{C} \backslash\{0\}$.
- For an arbitrary $\beta \in(0,1)$, the behavior of the "small" singular values $s_{N-\ell}\left(\widehat{R}_{L}^{(N)}-z I_{N}\right)$ for $\ell \in\left\{\left\lfloor N^{\beta}\right\rfloor, \ldots,\lfloor N / 2\rfloor\right\}$ and $z \in \mathbb{C} \backslash\{0\}$. Later this will help in controlling the magnitude of the singular values $s_{N-\ell}$ when the indices $\ell$ are close to $N^{\beta}$.
- The behavior of $\mu_{N}$ as $N \rightarrow \infty$ via the existence of a deterministic equivalent.


### 1.2 Assumptions

We shall assume that for every $N,\left(\boldsymbol{x}_{k}^{(N)}\right)_{k}$ is a stationary Gaussian process whose spectral density exists and satisfies some reasonable regularity conditions. As we shall see, this provides the opportunity to use a variety of technical tools. If the time series are not Gaussian, then the situation is way more involved technically but results similar to those in this paper are expected to hold under suitable restrictions on the time series. Let $\mathbb{T}$ denote the unit circle of the complex plane. Let $\mathcal{H}_{+}^{N}$ be the set of $N \times N$ Hermitian non-negative matrices. Suppose that for each positive integer $N$, there is an integrable function $S^{(N)}: \mathbb{T} \rightarrow \mathcal{H}_{+}^{N}$ such that for each $\ell \in \mathbb{Z}$,

$$
\begin{equation*}
R_{L}^{(N)}=\frac{1}{2 \pi} \int_{0}^{2 \pi} e^{-\imath L \theta} S^{(N)}\left(e^{\imath \theta}\right) d \theta \tag{1}
\end{equation*}
$$

This $S^{(N)}$ is called the spectral density of $\left\{R_{L}^{(N)}, L \geqslant 0\right\}$ or of the corresponding stationary process. Note that

$$
s_{N-1}\left(R_{0}^{(N)}\right)=s_{N-1}\left((2 \pi)^{-1} \int_{0}^{2 \pi} S^{(N)}\left(e^{\imath \theta}\right) d \theta\right)
$$

We further assume that for each $N, S^{(N)}$ is non-trivial in the sense that for each non-zero $\mathbb{C}^{N_{-}}$ valued polynomial $p(z)$,

$$
\int_{0}^{2 \pi} p\left(e^{\imath \theta}\right)^{*} S^{(N)}\left(e^{\imath \theta}\right) p\left(e^{\imath \theta}\right) d \theta>0
$$

We now turn to more substantial assumptions on $S^{(N)}$.
The first assumption is akin to uniform equicontinuity of $\left\{S^{(N)}\right\}$. For $h>0$, let

$$
\boldsymbol{w}\left(S^{(N)}, h\right)=\sup _{\theta} \sup _{|\psi| \leqslant h}\left\|S^{(N)}\left(e^{\imath(\theta+\psi)}\right)-S^{(N)}\left(e^{\imath \theta}\right)\right\|
$$

be the modulus of continuity of $S^{(N)}$ with respect to the spectral norm $\|\cdot\|$.
Assumption 1. 1. For any $\varepsilon>0$, there exists $h>0$ such that for all $N \in \mathbb{N}^{*}$,

$$
\sup _{N \in \mathbb{N}} \boldsymbol{w}\left(S^{(N)}, h\right)<\varepsilon
$$

2. With $\left\|S^{(N)}\right\|_{\infty}^{\mathbb{T}}=\max _{\theta}\left\|S^{(N)}\left(e^{\imath \theta}\right)\right\|$,

$$
\boldsymbol{M}:=\sup _{N}\left\|S^{(N)}\right\|_{\infty}^{\mathbb{T}}<\infty .
$$

3. $\boldsymbol{m}:=\inf _{N} s_{N-1}\left(R_{0}^{(N)}\right)>0$.

We allow the spectral density $S^{(N)}\left(e^{\imath \theta}\right)$ to be singular or close to singular at certain points of $\mathbb{T}$, but within the restrictions provided by the two following assumptions. Assumption 2 stipulates that $S^{(N)}\left(e^{\imath \theta}\right)$ can be close to being singular only on a set of frequencies with a small Lebesgue measure and it implies Assumption 1-3. Assumption 3 puts additional constraint on $s_{N-1}\left(S^{(N)}\left(e^{\imath \theta}\right)\right)$. Examples where Assumptions 2 and 3 are satisfied are provided in Section 2 below. Let Leb(•) denote the Lebesgue measure.
Assumption 2. Suppose that for any $\kappa \in(0,1)$, there exists $\delta>0$ such that

$$
\sup _{N} \operatorname{Leb}\left\{z \in \mathbb{T}: s_{N-1}\left(S^{(N)}(z)\right) \leqslant \delta\right\} \leqslant \kappa .
$$

## Assumption 3.

$$
\frac{1}{N} \int_{0}^{2 \pi} \log \left\|S^{(N)}\left(e^{2 \theta}\right)^{-1}\right\| d \theta \xrightarrow[N \rightarrow \infty]{ } 0
$$

### 1.3 Results

Before we state our results, we wish to recall that we have used modulo- $n_{N}$ summation to construct $\widehat{R}_{L}^{(N)}$. This is for convenience and helps in the details of the proofs. The results we establish continue to hold for the sample autocovariance matrices when we define them via the usual summation over all indices that maintain a lag $L$. See for instance [37] who relax the modulo- $n_{N}$ summation in the context of the model with i.i.d. processes considered in [4].

The first result is on a probabilistic bound on the magnitude of the smallest singular value of $\left(\widehat{R}_{L}^{(N)}-z\right)$. In problems similar to ours, often the optimal factor at the left hand side of the bound given below is $N^{-1} t$ instead of our $N^{-3 / 2} t$ when $\varepsilon>0$ is fixed. Our weaker bound will serve our purpose. We shall elaborate on this issue during the course of the proof.

Theorem 1. Suppose Assumptions 1, 2 and 3 hold. Then, for each $z \neq 0$ and arbitrarily small $\varepsilon>0$, there exists a constant $c_{1}$ such that for all small $t>0$ and for all large $N$,

$$
\mathbb{P}\left[s_{N-1}\left(\widehat{R}_{L}^{(N)}-z\right) \leqslant N^{-3 / 2} t\right] \leqslant \varepsilon t+\exp \left(-c_{1} \varepsilon^{2} N\right)
$$

The behavior of the small singular values is handled by Theorem 2. The behavior of $s_{N-k}\left(\hat{R}_{L}^{(N)}-\right.$ $z$ ) for values of $k$ which are close to $N^{\beta}$ is more important. The theorem implies that $s_{N-N^{\beta}}\left(\widehat{R}_{L}^{(N)}-\right.$ $z) \gtrsim N^{\beta / 2-1}$ with large probability. Again, though the rate is not optimal, it will be sufficient for our needs. Further comments on this issue will be provided in the course of the proof, see Remark 1 below.

Theorem 2. Suppose Assumptions 1, 2 and 3 hold. Let $\beta \in(0,1)$. Then, for each $z \neq 0$, there exist two positive constants $c_{2}$ and $C_{2}$ such that for all $k \in\left[\left\lfloor N^{\beta}\right\rfloor,[N / 2\rfloor\right]$, and for $N \geqslant N_{0}$, where $N_{0}$ is independent of $k$,

$$
\mathbb{P}\left[s_{N-k-1}\left(\widehat{R}_{L}^{(N)}-z\right) \leqslant C_{2} \sqrt{k} / N\right] \leqslant \exp \left(-c_{2} k\right)
$$

We now turn to the large- $N$ behavior of $\mu_{N}$. For this we rely on the well-known Hermitization technique due to Girko [12] (see [3] for a comprehensive exposition). Let $\mu$ be a probability measure on $\mathbb{C}$ that integrates $\log |\cdot|$ near infinity. Then its log-potential $U_{\mu}(\cdot): \mathbb{C} \rightarrow(-\infty, \infty]$ is defined below. The measure $\mu$ can be recovered from $U_{\mu}(z)$.

$$
U_{\mu}(z)=-\int_{\mathbb{C}} \log |w-z| \mu(d w)
$$

For the empirical spectral measure $\mu_{N}$, we can write

$$
U_{\mu_{N}}(z)=-\frac{1}{N} \sum_{\ell=0}^{N-1} \log \left|\lambda_{\ell}\left(\widehat{R}_{L}^{(N)}\right)-z\right|=-\frac{1}{2 N} \log \operatorname{det}\left(\widehat{R}_{L}^{(N)}-z\right)\left(\widehat{R}_{L}^{(N)}-z\right)^{*}=-\int \log t \nu_{z, N}(d t)
$$

where $\nu_{z, N}$, the empirical measure (on $\mathbb{R}$ ) of the singular values of the matrix $\hat{R}_{L}^{(N)}-z$ given by

$$
\nu_{z, N}=\frac{1}{N} \sum_{\ell=0}^{N-1} \delta_{s_{\ell}\left(\hat{R}_{L}^{(N)}-z\right)}
$$

Given a matrix $M \in \mathbb{C}^{N \times N}$, denote its so-called Hermitized version as

$$
\boldsymbol{H}(M)=\left[\begin{array}{ll}
M^{*} & M
\end{array}\right]
$$

As is well-known, the spectral measure $\check{\nu}_{z, N}($ on $\mathbb{R})$ of $\boldsymbol{H}\left(\widehat{R}_{L}^{(N)}-z\right)$ is given by

$$
\check{\nu}_{z, N}=\frac{1}{2 N} \sum_{\ell=0}^{N-1}\left(\delta_{s_{\ell}\left(\hat{R}_{L}^{(N)}-z\right)}+\delta_{-s_{\ell}\left(\hat{R}_{L}^{(N)}-z\right)}\right)
$$

The measure $\check{\nu}_{z, N}$ is clearly symmetric in the sense that $\check{\nu}_{z, N}(B)=\check{\nu}_{z, N}(-B)$ for each Borel set $B \subset \mathbb{R}$, and it is the symmetrized version of $\nu_{z, N}$. Observe that $\int \log t \nu_{z, N}(d t)=\int \log |t| \check{\nu}_{z, N}(d t)$.

Note that all of these are random measures. To study the behavior of $\mu_{N}$, we will find it more convenient to study $\check{\nu}_{z, N}$ instead of $\nu_{z, N}$, since the matrix $\boldsymbol{H}\left(\widehat{R}_{L}^{(N)}-z\right)$ is Hermitian. This approach is formalized in the following general proposition which provides conditions under which the sequence of random measures $\left(\mu_{N}\right)$ can be approximated by some sequence $\left(\boldsymbol{\mu}_{N}\right)$ of deterministic measures. Its proof can be found in [3] (see also [7]). The uniform integrability condition needed in the proposition is ensured by Theorems 1 and 2 (again, see [3] for the proof details).

Proposition 3. Assume that for almost every $z \in \mathbb{C}$, the following two conditions hold:

1. With probability one, $\log |\cdot|$ is uniformly integrable with respect to $\left\{\check{\nu}_{z, N}\right\}$.
2. There exists a tight sequence of deterministic symmetric probability measures $\left(\check{\boldsymbol{\nu}}_{z, N}\right)_{N}$ on $\mathbb{R}$ such that for each bounded and continuous function $f: \mathbb{R} \rightarrow \mathbb{R}$,

$$
\int f d \check{\nu}_{z, N}-\int f d \check{\boldsymbol{\nu}}_{z, N} \xrightarrow[N \rightarrow \infty]{\text { a.s. }} 0 .
$$

Then, there exists a tight sequence of deterministic probability measure $\left(\boldsymbol{\mu}_{N}\right)$ on $\mathbb{C}$ such that for each bounded and continuous function $f: \mathbb{C} \rightarrow \mathbb{C}$, we have

$$
\int f d \mu_{N}-\int f d \boldsymbol{\mu}_{N} \xrightarrow[N \rightarrow \infty]{\text { a.s. }} 0
$$

Moreover, the logarithmic potential of $\boldsymbol{\mu}_{N}$ is

$$
U_{\boldsymbol{\mu}_{N}}(z)=-\int \log |t| \check{\boldsymbol{\nu}}_{z, N}(d t) .
$$

We are thus left to consider the asymptotic properties of $\check{\nu}_{z, N}$ with a goal to comply with Condition 2 of Proposition 3. Classically, the central object that is used for this is the resolvent (in what follows $z \in \mathbb{C}$ is arbitrary)

$$
Q^{(N)}(z, \eta)=\left(\boldsymbol{H}\left(\widehat{R}_{L}^{(N)}-z\right)-\eta I_{2 N}\right)^{-1}
$$

of $\boldsymbol{H}\left(\widehat{R}_{L}^{(N)}-z\right)$ in the complex variable $\eta \in \mathbb{C}_{+} \triangleq\{w \in \mathbb{C}: \Im w>0\}$. The resolvent $Q^{(N)}(z, \cdot)$ is a typical example of a so-called matrix Stieltjes transform. Before we recall its nature, let us recall that the Stieltjes transform of any probability measure $\nu$ on $\mathbb{R}$ is given by

$$
g_{\nu}(\eta)=\int_{\mathbb{R}} \frac{1}{\lambda-\eta} \nu(d \lambda), \quad \eta \in \mathbb{C} \backslash \mathbb{R}
$$

The real and imaginary parts of a square matrix $M$ are respectively given as

$$
\Re M=\frac{M+M^{*}}{2} \quad \text { and } \quad \Im M=\frac{M-M^{*}}{2 \imath} .
$$

Given an integer $m>0$, we let $\mathcal{M}_{+}^{m}$ denote the set of matrices $M \in \mathbb{C}^{m \times m}$ such that $\Im M>0$.
The following result can be found in the literature dealing with the moment problem and related topics, see, e.g., $[2,11]$.

Proposition 4 (matrix Stieltjes transform). Let $F: \mathbb{C}_{+} \rightarrow \mathbb{C}^{m \times m}$ be a matrix-valued function. Then, the following facts are equivalent.

1. $F$ is the Stieltjes transform of an $\mathcal{H}_{+}^{m}$-valued measure $\mu$ on $\mathbb{R}$ such that $\mu(\mathbb{R})=I_{m}$.
2. $F$ is analytic, $F(\eta) \in \mathcal{M}_{+}^{m}$ for $\eta \in \mathbb{C}_{+}$, and $-\imath t F(\imath t)$ converges to $I_{m}$ as $t \rightarrow \infty$.

Such an $F$ is called a matrix Stieltjes transform. Let $\mathcal{S}^{m}$ denote the set of all such $m \times m$ matrix Stieltjes transforms. It is known that if $F \in \mathcal{S}^{m}$, then
(1) $\|F(\eta)\| \leqslant 1 / \Im \eta$.
(2) $m^{-1} \operatorname{tr} F(\eta)=g_{\nu}(\eta)$ for some probability measure $\nu$ on $\mathbb{R}$.

An illustration of (2) is provided by the resolvent we just defined: it can be checked that

$$
Q^{(N)}(z, \cdot) \in \mathcal{S}^{2 N} \text { and } g_{\check{\nu}_{z, N}}(\eta)=(2 N)^{-1} \operatorname{tr} Q^{(N)}(z, \eta)
$$

In the remainder of this paper, whenever we write $M \in \mathbb{C}^{2 N \times 2 N}$ as $M=\left[\begin{array}{ll}M_{00} & M_{01} \\ M_{10} & M_{11}\end{array}\right]$, it is understood that the blocks $M_{u v}$ belong to $\mathbb{C}^{N \times N}$. With this notation, we define the linear operator $\mathscr{T}: \mathbb{C}^{2 N \times 2 N} \rightarrow \mathbb{C}^{2 \times 2}$ as

$$
\mathscr{T}\left(\left[\begin{array}{ll}
M_{00} & M_{01} \\
M_{10} & M_{11}
\end{array}\right]\right)=\left[\begin{array}{ll}
\operatorname{tr} M_{00} / n_{N} & \operatorname{tr} M_{01} / n_{N} \\
\operatorname{tr} M_{10} / n_{N} & \operatorname{tr} M_{11} / n_{N}
\end{array}\right] .
$$

Given an integer $L>0$, we also define the $2 \times 2$ Hermitian unitary matrix function $U_{L}$ on $\mathbb{T}$ as

$$
U_{L}\left(e^{\imath \theta}\right)=\left[\begin{array}{ll}
e^{\imath L \theta} & e^{-\imath L \theta}
\end{array}\right]
$$

The next technical result is on the existence of a unique solution for a functional equation. It will be used to show that $\left(\check{\boldsymbol{\nu}}_{z, N}\right)_{N}$ approximates $\left(\check{\nu}_{z, N}\right)_{N}$.

Theorem 5. Let $\Sigma: \mathbb{T} \rightarrow \mathcal{H}_{+}^{N}$ be a continuous function, and let $z \in \mathbb{C}$. Given a function $M(\eta) \in \mathcal{S}^{2 N}$, the function displayed below is well-defined and belongs to $\mathcal{S}^{2 N}$ as a function of $\eta$.

$$
\mathcal{F}_{\Sigma, z}(M(\eta), \eta)=\left(\frac{1}{2 \pi} \int_{0}^{2 \pi}\left(\mathscr{T}\left(\left(I_{2} \otimes \Sigma\left(e^{\imath \theta}\right)\right) M(\eta)\right)+U_{L}\left(e^{\imath \theta}\right)\right)^{-1} \otimes \Sigma\left(e^{\imath \theta}\right) d \theta-\left[\begin{array}{ll}
\eta & z \\
\bar{z} & \eta
\end{array}\right] \otimes I_{N}\right)^{-1}
$$

Moreover, the functional equation in the parameter $\eta \in \mathbb{C}_{+}$

$$
\begin{equation*}
P(z, \eta)=\mathcal{F}_{\Sigma, z}(P(z, \eta), \eta) \tag{2}
\end{equation*}
$$

admits a unique solution in the class $P(z, \cdot) \in \mathcal{S}^{2 N}$. Write

$$
P(z, \cdot)=\left[\begin{array}{ll}
P_{00}(z, \cdot) & P_{01}(z, \cdot) \\
P_{10}(z, \cdot) & P_{11}(z, \cdot)
\end{array}\right] \quad \text { and } \quad \Lambda(d t)=\left[\begin{array}{ll}
\Lambda_{00}(d t) & \Lambda_{01}(d t) \\
\Lambda_{10}(d t) & \Lambda_{11}(d t)
\end{array}\right]
$$

where $\Lambda(d t)$ is the matrix measure whose matrix Stieltjes transform is $P(z, \cdot)$, where $P_{i i}(z, \cdot)$ is the Stieltjes transform of $\Lambda_{i i}(d t)$. The positive matrix measures $\Lambda_{00}$ and $\Lambda_{11}$ are symmetric.

One consequence of this theorem is that the probability measure $\zeta$ on $\mathbb{R}$ such that $g_{\zeta}(\eta)=$ $(2 N)^{-1} \operatorname{tr} P(z, \eta)$ is symmetric. In fact (we omit the proof),

$$
g_{\zeta}(\eta)=N^{-1} \operatorname{tr} P_{00}(z, \eta)=N^{-1} \operatorname{tr} P_{11}(z, \eta)
$$

The first part of the next theorem claims the tightness of $\left(\check{\boldsymbol{\nu}}_{z, N}\right)_{N}$. Later we shall need the behavior of the difference $\left(Q^{(N)}(z, \eta)-G^{(N)}(z, \eta)\right)$. This is captured in the second part. The extra generality obtained by including the matrix $D$ is not necessary for the purposes of the present work in which we use only the specific choice of $D=I$. However this will be useful to derive convergence and related properties of $\left(\mu_{N}\right)_{N}$ and $\left(\boldsymbol{\mu}_{N}\right)_{N}$ that we wish to pursue in future.

Theorem 6. Suppose Assumption 1 holds. Let $G^{(N)}(z, \eta)$ be the solution in $\mathcal{S}^{2 N}$, of the equation

$$
G^{(N)}(z, \eta)=\mathcal{F}_{S^{(N)}, z}\left(G^{(N)}(z, \eta), \eta\right)
$$

as specified by Theorem 5 . Let $\check{\boldsymbol{\nu}}_{z, N}$ be the symmetric probability measure on $\mathbb{R}$ whose Stieltjes transform is

$$
g_{\check{\boldsymbol{\nu}}_{z, N}}(\eta)=\frac{1}{2 N} \operatorname{tr} G^{(N)}(z, \eta)
$$

Then, the sequence $\left(\check{\boldsymbol{\nu}}_{z, N}\right)_{N}$ is tight.
Let $D^{(N)} \in \mathbb{C}^{2 N \times 2 N}$ be an arbitrary deterministic matrix such that $\left\|D^{(N)}\right\|=1$. Then,

$$
\begin{equation*}
\forall \eta \in \mathbb{C}_{+}, \quad \frac{1}{2 N} \operatorname{tr} D^{(N)}\left(Q^{(N)}(z, \eta)-G^{(N)}(z, \eta)\right) \xrightarrow[N \rightarrow \infty]{\text { a.s. }} 0 \tag{3}
\end{equation*}
$$

Taking $D^{(N)}=I$, this theorem shows in particular that $g_{\check{\nu}_{z, N}}(\eta)-g_{\tilde{\nu}_{z, N}}(\eta) \xrightarrow[N \rightarrow \infty]{\text { a.s. }} 0$. We also have the following corollary, whose proof employs well-known facts about Stieltjes transforms (see, e.g., [22]), and is omitted.

Corollary 7. The sequence $\left(\check{\nu}_{z, N}\right)_{N}$ of probability measures is tight with probability one. Furthermore, for each bounded and continuous function $f: \mathbb{R} \rightarrow \mathbb{R}$, it holds that

$$
\int f d \check{\nu}_{z, N}-\int f d \check{\boldsymbol{\nu}}_{z, N} \xrightarrow[N \rightarrow \infty]{\text { a.s. }} 0
$$

Finally, with the help of Proposition 3, Theorems 1, 2, and 6 lead at once to the following result that characterizes the asymptotic behavior of $\left(\mu_{N}\right)$.

Theorem 8. Let Assumptions 1 to 3 hold true. Then, there exists a tight sequence of deterministic probability measures $\left(\boldsymbol{\mu}_{N}\right)$ on $\mathbb{C}$ such that, for each bounded and continuous function $f: \mathbb{C} \rightarrow \mathbb{C}$,

$$
\int f d \mu_{N}-\int f d \boldsymbol{\mu}_{N} \xrightarrow[N \rightarrow \infty]{\text { a.s. }} 0
$$

The function $\log |\cdot|$ is integrable with respect to $\check{\boldsymbol{\nu}}_{z, N}$ for each $z \neq 0$, and the measure $\boldsymbol{\mu}_{N}$ is defined by its logarithmic potential through the identity

$$
U_{\boldsymbol{\mu}_{N}}(z)=-\int \log |t| \check{\boldsymbol{\nu}}_{z, N}(d t)
$$

A natural issue is if can we say anything further about the behavior of $\left(\boldsymbol{\mu}_{N}\right)$. In particular, its weak convergence in turn would ensure the weak convergence of $\left(\mu_{N}\right)$. An idea that goes back to [9] and which has been frequently used in the literature of non-Hermitian matrices, connects the behaviour of $\left(\boldsymbol{\mu}_{N}\right)$ to that of $N^{-1} \operatorname{tr} G_{01}^{(N)}(z, \imath t)$ as $t \searrow 0$. In the present context the details need some significant efforts. Since this is outside the focus of this work, it will be pursued separately.

The remainder of the paper is organized as follows. Examples where Assumptions 2 and 3 hold, are provided in Section 2. Theorems 1 and 2 are proved in Section 3. The results related to the behavior of the singular values, i.e., Theorems 5 and 6, are proved in Section 4.

## Notation

The indices of the elements of a vector or a matrix start from zero. Given two integers $k$ and $m$, we write $[k: m]=\{k, k+1, \ldots, m-1\}$, this set being empty if $k \geqslant m$. We also write $[m]=[0: m]$. Assume $m>0$. For $k \in[m]$, we denote as $e_{m, k}$ the $k^{\text {th }}$ canonical basis vector of $\mathbb{C}^{m}$. When there is no ambiguity, we write $e_{k}$ for $e_{m, k}$. Given a matrix $M \in \mathbb{C}^{m \times n}$, and two sets $\mathcal{I} \subset[m]$ and $\mathcal{J} \subset[n]$, we denote as $M_{\mathcal{I}, \mathcal{J}}$ the $|\mathcal{I}| \times|\mathcal{J}|$ sub-matrix of $M$ that is obtained by retaining the rows of $M$ whose indices belong to $\mathcal{I}$ and the columns whose indices belong to $\mathcal{J}$. We also write $M_{\cdot, \mathcal{J}}=M_{[m], \mathcal{J}}$ and $M_{\mathcal{I}, \cdot}=M_{\mathcal{I},[n]}$. Given a vector $v \in \mathbb{C}^{m}$, we denote by $v_{\mathcal{I}}$ the $\mathbb{C}^{|\mathcal{I}|}$ sub-vector obtained by keeping the elements of $v$ whose indices are in $\mathcal{I}$.

For any matrix $M, M>0$ means that it is positive definite. The column span of a matrix $M$ is denoted by $\operatorname{span}(M)$. The orthogonal projection matrix onto $\operatorname{span}(M)$ (respectively, onto the
subspace orthogonal to $\operatorname{span}(M)$ ) is denoted by $\Pi_{M}$ (respectively, $\Pi_{M}^{\perp}$ ). The spectral norm of a matrix and the Euclidean norm of a vector are denoted by $\|\cdot\|$. The Hilbert-Schmidt norm of an operator will be denoted $\|\cdot\|_{\text {HS }}$. The notation $M \geqslant G$ where $M$ and $G$ are Hermitian matrices refers to the semi-definite ordering of such matrices.

We write $\mathbb{R}_{+}=[0, \infty)$. Suppose $B$ is a Borel set of $\mathbb{R}^{d}$. Then $\operatorname{Leb}(B)$ denotes its Lebesgue measure. For any $x \in \mathbb{R}^{d}, \operatorname{dist}(x, B)=\inf _{y \in B}\|x-y\|$. For $B$ in a metric space $E, \mathcal{V}_{\rho}^{E}(B)$ denotes its closed $\rho$-neighborhood. If the underlying space $E$ is clear, we simply write $\mathcal{V}_{\rho}(B)$.

The unit-sphere of $\mathbb{C}^{m}$ will be denoted as $\mathbb{S}^{m-1}$. The set of vectors of $\mathbb{S}^{m-1}$ that are supported by the (index) set $\mathcal{I} \subset[m]$ will be denoted by $\mathbb{S}_{\mathcal{I}}^{m-1}$.

The probability and the expectation with respect to the law of the vector $x$ will be denoted by $\mathbb{P}_{x}$ and $\mathbb{E}_{x}$. The centered and circularly symmetric complex Gaussian distribution with covariance matrix $\Sigma$ will be denoted by $\mathcal{N}_{\mathbb{C}}(0, \Sigma)$. Finally, $\otimes$ denotes the Kronecker product.

## 2 Examples where Assumptions 2 and 3 are satisfied

In this section, we provide examples where the spectral density satisfies Assumptions 2 and 3. Consider the moving average $M A(\infty)$ model:

$$
\begin{equation*}
\boldsymbol{x}_{k}^{(N)}=\boldsymbol{\xi}_{k}^{(N)}+\sum_{\ell \geqslant 1} A_{\ell}^{(N)} \boldsymbol{\xi}_{k-\ell}^{(N)} \tag{4}
\end{equation*}
$$

where $\boldsymbol{\xi}^{(N)}=\left(\boldsymbol{\xi}_{k}^{(N)}\right)_{k \in \mathbb{Z}}$ is an i.i.d. process with $\boldsymbol{\xi}_{k}^{(N)} \sim \mathcal{N}_{\mathbb{C}}\left(0, I_{N}\right)$ and $\left(A_{\ell}^{(N)}\right)_{\ell \geqslant 1}$ is a sequence of deterministic matrices which satisfy the minimal requirement $\sum_{\ell \geqslant 1}\left\|A_{\ell}^{(N)}\right\|<\infty$. Let us look at some particular cases:
(i) First suppose that

$$
\limsup _{N} \sum_{\ell \geqslant 1}\left\|A_{\ell}^{(N)}\right\|<1 .
$$

Then it is obvious that both Assumptions 2 and 3 are satisfied and $s_{N-1}\left(S^{(N)}(z)\right)$ remains bounded away from zero when $z$ runs through $\mathbb{T}$.
(ii) Now consider the $\mathrm{MA}(1)$ case, so that $A_{\ell}^{(N)}=0$ for $\ell \geqslant 2$, and suppose that we only have $\sup _{N}\left\|A_{1}^{(N)}\right\|<\infty$. In this case, Assumptions 2 and 3 , written in terms of $A_{1}^{(N)}$ are, respectively (the log in the second expression is always integrable),

$$
\begin{gather*}
\forall \kappa \in(0,1), \exists \delta>0, \sup _{N} \operatorname{Leb}\left\{z \in \mathbb{T}:\left\|\left(z-A_{1}^{(N)}\right)^{-1}\right\| \geqslant 1 / \delta\right\} \leqslant \kappa, \quad \text { and }  \tag{5a}\\
\frac{1}{N} \int_{0}^{2 \pi} \log \left\|\left(e^{\imath \theta}-A_{1}^{(N)}\right)^{-1}\right\| d \theta \underset{N \rightarrow \infty}{ } 0 . \tag{5b}
\end{gather*}
$$

These conditions are closely connected with the pseudospectrum of $A_{1}^{(N)}$ [34]. The Toeplitz matrices are among the matrices for which the pseudospectra are well-understood. Suppose that $A_{1}^{(N)}=\left[a_{k-\ell}^{(N)}\right]_{0 \leqslant k, \ell \leqslant N-1}$ is a Toeplitz matrix with the so-called symbol

$$
f_{N}(z)=\sum_{k=1-N}^{N-1} a_{k}^{(N)} z^{k}
$$

Then, by the Brown and Halmos theorem [5, Th. 4.29], Assumptions 2 and 3 are respectively satisfied if

$$
\begin{aligned}
\forall \kappa \in(0,1), \exists \delta>0 & , \limsup _{N} \operatorname{Leb}\left\{z \in \mathbb{T}: \operatorname{dist}\left(e^{\imath \theta}, \operatorname{conv} f_{N}(\mathbb{T})\right) \leqslant \delta\right\} \leqslant \kappa, \quad \text { and } \\
& \frac{1}{N} \int_{0}^{2 \pi} \log \operatorname{dist}\left(e^{\imath \theta}, \operatorname{conv} f_{N}(\mathbb{T})\right) d \theta \underset{N \rightarrow \infty}{ } 0
\end{aligned}
$$

where conv denotes the convex hull of a set.
(iii) Another particular case of the model in (4) is the following. Assume that $N=M K$ where $M$ and $K$ are two positive integers. Assume that the matrices $A_{\ell}^{(N)}$ are block-diagonal matrices of the form

$$
A_{\ell}^{(N)}=I_{M} \otimes B_{\ell}^{(K)}
$$

where $B_{\ell}^{(K)} \in \mathbb{C}^{K \times K}$. We also assume that $\mathfrak{b}=\sup _{K} \sum_{\ell \geqslant 1}\left\|B_{\ell}^{(K)}\right\|<\infty$. Thus the process $\boldsymbol{\xi}^{(N)}$ consists of $M$ i.i.d. streams of stationary processes each of dimension $K$. Let us show that if $M \rightarrow \infty$, then Assumption 3 is satisfied for this process. Further if $K$ is upper bounded, then Assumption 2 is satisfied.

The spectral density of each stream is $\Sigma^{(K)}\left(e^{\imath \theta}\right)=C^{(K)}\left(e^{\imath \theta}\right) C^{(K)}\left(e^{\imath \theta}\right)^{*}$, with $C^{(K)}(z)=$ $I_{K}+\sum_{\ell \geqslant 1} z^{\ell} B_{\ell}^{(K)}$ on the unit-disk. It is well known that (see, e.g., [25, Th. 6.1]),

$$
\begin{aligned}
\left.\left.\frac{1}{2 \pi} \int_{0}^{2 \pi}|\log | \operatorname{det} C^{(K)}\left(e^{\imath \theta}\right)\right|^{\frac{1}{K}} \right\rvert\, d \theta & \leqslant \frac{1}{2 \pi} \int_{0}^{2 \pi}\left(\operatorname{det} \Sigma^{(K)}\left(e^{\imath \theta}\right)\right)^{\frac{1}{K}} d \theta-\log \left|\operatorname{det} C^{(K)}(0)\right|^{\frac{1}{K}} \\
& =\frac{1}{2 \pi} \int_{0}^{2 \pi}\left(\operatorname{det} \Sigma^{(K)}\left(e^{\imath \theta}\right)\right)^{\frac{1}{K}} d \theta
\end{aligned}
$$

By Hadamard's inequality, the right hand side is bounded by $(1+\mathfrak{b})^{2}$. Thus,

$$
\begin{aligned}
\left|\log s_{K-1}\left(\Sigma^{(K)}\left(e^{\imath \theta}\right)\right)\right| & \leqslant\left|\sum_{\ell=0}^{K-1} \log s_{\ell}\left(\Sigma^{(K)}\left(e^{\imath \theta}\right)\right) \mathbb{1}_{s_{\ell}\left(\Sigma^{(K)}\left(e^{2 \theta}\right)\right) \leqslant 1}\right|+\log (1+\mathfrak{b})^{2} \\
& \leqslant\left|\sum_{\ell=0}^{K-1} \log s_{\ell}\left(\Sigma^{(K)}\left(e^{\imath \theta}\right)\right)\right|+(K+1) \log (1+\mathfrak{b})^{2} \\
& =2|\log | \operatorname{det} C^{(K)}\left(e^{\imath \theta}\right)| |+(K+1) \log (1+\mathfrak{b})^{2},
\end{aligned}
$$

and we deduce from the next to last display that

$$
\frac{1}{2 \pi} \int_{0}^{2 \pi}\left|\log s_{K-1}\left(\Sigma^{(K)}\left(e^{\imath \theta}\right)\right)\right| d \theta \leqslant(3 K+1) \log (1+\mathfrak{b})^{2}
$$

Since $S^{(N)}\left(e^{\imath \theta}\right)=I_{M} \otimes \Sigma^{(K)}\left(e^{\imath \theta}\right)$, Assumption 3 is satisfied when $M \rightarrow \infty$. Furthermore, given a small $\delta>0$, we have

$$
\frac{1}{2 \pi} \int_{0}^{2 \pi} \mathbb{1}_{s_{K-1}\left(\Sigma^{(K)}\left(e^{\imath \theta}\right)\right) \leqslant \delta} d \theta \leqslant \frac{1}{|\log \delta|} \frac{1}{2 \pi} \int_{0}^{2 \pi}\left|\log s_{K-1}\left(\Sigma^{(K)}\left(e^{\imath \theta}\right)\right)\right| d \theta \leqslant \frac{(3 K+1) \log (1+\mathfrak{b})^{2}}{|\log \delta|}
$$

This shows that Assumption 2 is satisfied if $K$ is bounded.

## 3 Small singular values: Proofs of Theorems 1 and 2

### 3.1 Outline of the proofs

In the sequel we shall most often omit the dependency on $N$ in the notations for simplicity. Writing $n=n_{N}$, define the matrix

$$
X=\left[\begin{array}{lll}
x_{0} & \cdots & x_{n-1}
\end{array}\right]=n^{-1 / 2}\left[\begin{array}{lll}
\boldsymbol{x}_{0} & \cdots & \boldsymbol{x}_{n-1}
\end{array}\right] \in \mathbb{C}^{N \times n}
$$

with $x_{k} \in \mathbb{C}^{N}$ being the $k^{\text {th }}$ column of $X$, and consider the $n \times n$ circulant matrix

$$
J=\left[\begin{array}{cccc}
0 & & & 1 \\
1 & \ddots & & \\
& \ddots & \ddots & \\
& & 1 & 0
\end{array}\right]
$$

Then, the sample autocovariance matrix $\widehat{R}_{L}$ can be rewritten as $\widehat{R}_{L}=X J^{L} X^{*}$. Let the so-called $n \times n$ Fourier matrix F be defined as

$$
\begin{equation*}
\mathrm{F}=\frac{1}{\sqrt{n}}[\exp (2 \imath \pi k \ell / n)]_{k, \ell=0}^{n-1} \tag{6}
\end{equation*}
$$

Then $J=\mathrm{F} \Omega \mathrm{F}^{*}$, where

$$
\Omega=\operatorname{diag}\left(\omega^{k}\right)_{k=0}^{n-1} \text { and } \omega=\exp (-2 \imath \pi / n)
$$

Furthermore, for $k \in[n]$, let

$$
y_{k}=\frac{1}{\sqrt{n}} \sum_{\ell=0}^{n-1} e^{2 \imath \pi k \ell} x_{\ell}
$$

be the discrete Fourier transform of the finite sequence $\left(x_{0}, \ldots, x_{n-1}\right)$, and define the $N \times n$ matrix

$$
Y=\left[\begin{array}{lll}
y_{0} & \cdots & y_{n-1}
\end{array}\right]=X F \in \mathbb{C}^{N \times n} .
$$

Then, we obviously have

$$
\widehat{R}_{L}=Y \Omega^{L} Y^{*}
$$

Note that the columns $y_{k}$ of $Y$ are "almost" independent, since they are the discrete Fourier transforms applied to a time window of a Gaussian stationary process. This is why we shall heavily rely on the above expression of $\widehat{R}_{L}$. Let us elaborate on this point before entering the core of the proof. Write

$$
y_{k}=\tilde{y}_{k}+\check{y}_{k}, \quad \check{y}_{k}=\mathbb{E}\left[y_{k} \mid Y_{k}\right],
$$

the latter being the conditional expectation with respect to the $\sigma$-field $\sigma\left(Y_{k}\right)$ generated by the elements of the matrix $Y_{k}=\left[y_{0} \cdots y_{k-1} y_{k+1} \cdots y_{n-1}\right] \in \mathbb{C}^{N \times(n-1)}$. Write $S_{k}=S\left(e^{2 \imath \pi k / n}\right)$ for brevity. Due to the Gaussianity, $\tilde{y}_{k}$ and $Y_{k}$ will be independent, and we shall prove that

$$
\mathbb{E}\left[\tilde{y}_{k} \tilde{y}_{k}^{*}\right] \simeq \mathbb{E}\left[y_{k} y_{k}^{*}\right] \simeq S_{k} .
$$

For this, as is frequently done in estimation theory for stationary processes, we shall make use of the orthogonal matrix polynomial theory with respect to the matrix measure $(2 \pi)^{-1} S\left(e^{\imath \theta}\right) d \theta$ on $\mathbb{T}$. Assumption 3, which is reminiscent of the notion of regular measures found in this literature [31, 30], will play a major role in our analysis that will be presented in Section 3.2.

We now consider the proof of Theorem 1, starting with a well-known linearization trick. Write

$$
H=\left[\begin{array}{cc}
\Omega^{-L} & Y^{*} \\
Y & z
\end{array}\right] \in \mathbb{C}^{(N+n) \times(N+n)}
$$

By using the well known formula for the inverse of a partitioned matrix (see [15, §0.7.3]), it follows that $\left\|\left(Y \Omega^{L} Y^{*}-z\right)^{-1}\right\| \leqslant\left\|H^{-1}\right\|$. Hence

$$
\begin{equation*}
s_{N+n-1}(H) \leqslant s_{n-1}\left(Y \Omega^{L} Y^{*}-z\right) \tag{7}
\end{equation*}
$$

and the problem then is to control $s_{N+n-1}(H)$. A similar problem, considered in [35], was that of the smallest singular value of a symmetric random matrix with iid elements above the diagonal, see also [20]. Even though our matrix is quite different from theirs, we borrow many of their ideas as well of its predecessors, such as [26].

First, it can be shown that for $C>0$ large enough, $\mathbb{P}[\|Y\| \geqslant C]$ is exponentially small (Lemma 16) by using some standard Gaussian calculations. Then, as in these articles, the control over $s_{N+n-1}(H)$ will be provided on the event $[\|Y\| \leqslant C]$. More specifically, we will prove that, for some constant $c>0$,

$$
\mathbb{P}\left[\left[s_{N+n-1}(H) \leqslant N^{-3 / 2} t\right] \cap[\|Y\| \leqslant C]\right] \leqslant \varepsilon t+\exp \left(-c \varepsilon^{2} n\right)
$$

The smallest singular value of $H$ can be obtained from the variational characterization

$$
s_{N+n-1}(H)=\inf _{u \in \mathbb{S}^{N+n-1}}\|H u\| .
$$

A well-established method to control the smallest singular value of a random matrix is to study the action of this matrix on the so-called compressible and incompressible vectors [18, 26]. Let $\theta, \rho \in(0,1)$ be fixed. A vector in $\mathbb{S}^{m-1}$ is called $\theta$-sparse if it does not have more than $\lfloor\theta m\rfloor$ non-zero elements. Given $\theta, \rho \in(0,1)$, we define the set of $(\theta, \rho)$-compressible vectors as

$$
\operatorname{comp}(\theta, \rho)=\mathbb{S}^{m-1} \cap \bigcup_{\substack{\mathcal{I} \subset[m] \\|\mathcal{I}|=\lfloor\theta m\rfloor}} \mathcal{V}_{\rho}^{\mathbb{C}^{m}}\left(\mathbb{S}_{\mathcal{I}}^{m-1}\right)
$$

In other words, this is the set of all unit vectors at a distance less or equal to $\rho$ from the set of the $\theta$-sparse unit vectors. The set incomp $(\theta, \rho)$ of $(\theta, \rho)$-incompressible vectors is the complementary set $\mathbb{S}^{m-1} \backslash \operatorname{comp}(\theta, \rho)$.

Getting back to the variational characterization, and writing $u=\left[v^{\top} w^{\top}\right]^{\top}$ with $v \in \mathbb{C}^{n}$, we (roughly) define the set

$$
\mathcal{S}=\left\{u \in \mathbb{S}^{N+n-1}:\|v\| \leqslant \text { a constant } \quad \text { or } \quad v /\|v\| \text { is compressible }\right\}
$$

and we write

$$
\begin{equation*}
s_{N+n-1}(H)=\inf _{u \in \mathcal{S}}\|H u\| \wedge \inf _{u \in \mathbb{S}^{N+n-1} \backslash \mathcal{S}}\|H u\| . \tag{8}
\end{equation*}
$$

For the first infimum, we focus on the component $v$ of the vector $u$ because $v$ impacts the first $n$ columns of $H$ which are nearly independent. Relying on the decomposition $y_{k}=\widetilde{y}_{k}+\check{y}_{k}$, we first show that $\mathbb{P}[\|H u\| \leqslant c]$ is exponentially small when $v /\|v\|$ is a sparse vector, and then we complete the analysis by an $\varepsilon$-net argument.

The second infimum requires other arguments. Let $h_{k}$ be the $k^{\text {th }}$ column of $H$, and so $H_{k}=$ $\left[h_{0} \cdots h_{k-1} h_{k+1} \cdots h_{N+n-1}\right] \in \mathbb{C}^{(N+n) \times(N+n-1)}$. Following a by now well-known idea of $[26$, Lem. 3.5], the infimum over the incompressible vectors can be controlled by managing the distances $\operatorname{dist}\left(h_{k}, H_{k}\right)$ between $h_{k}$ and $\operatorname{span}\left(H_{k}\right)$. Given $k \in[n]$, let $\Omega_{k}=\operatorname{diag}\left(\omega^{\ell}\right)_{\ell \neq k} \in \mathbb{C}^{(n-1) \times(n-1)}$. Let

$$
G_{k}=\left[\begin{array}{cc}
\Omega_{k}^{-L} & Y_{k}^{*} \\
Y_{k} & z
\end{array}\right] \in \mathbb{C}^{(N+n-1) \times(N+n-1)},
$$

and partition $G_{k}^{-1}$ as

$$
G_{k}^{-1}=\left[\begin{array}{ll}
E_{k} & F_{k} \\
P_{k} & D_{k}
\end{array}\right], \quad E_{k} \in \mathbb{C}^{(n-1) \times(n-1)}, \quad D_{k} \in \mathbb{C}^{N \times N} .
$$

Then, after some algebra, we get the following equation (similar to what was obtained in [35]):

$$
\operatorname{dist}\left(h_{k}, H_{k}\right)=\frac{\left|\omega^{-k L}-y_{k}^{*} D_{k} y_{k}\right|}{\sqrt{1+\left\|y_{k}^{*} P_{k}\right\|^{2}+\left\|y_{k}^{*} D_{k}\right\|^{2}}} .
$$

Let us assume temporarily that the $\check{y}_{k}$ are equal to zero. Restricting ourselves to the indices $k$ for which $s_{N-1}\left(S_{k}\right) \sim 1$ (which is allowed), we can show that for these indices

$$
\begin{equation*}
\sqrt{1+\left\|y_{k}^{*} P_{k}\right\|^{2}+\left\|y_{k}^{*} D_{k}\right\|^{2}} \sim\left\|S_{k}^{1 / 2} D_{k} S_{k}^{1 / 2}\right\|_{\mathrm{HS}} \tag{9}
\end{equation*}
$$

Consequently, the control of $\operatorname{dist}\left(h_{k}, H_{k}\right)$ can be reduced to the control of the probability (with respect to the law of $y_{k}$ ) that, $y_{k}^{*} D_{k} y_{k} /\left\|S_{k}^{1 / 2} D_{k} S_{k}^{1 / 2}\right\|_{\text {HS }}$ lies in a small ball. Due to the Gaussian nature of $y_{k}$, this task is easy and leads to a rate of $N^{-1} t$ in the statement of Theorem 1.

Now consider the term that involves $\check{y}_{k}$. Even though $\mathbb{E}\left\|\check{y}_{k}\right\|^{2}$ is small, its interdependence with $\left(P_{k}, D_{k}\right)$ prevents us from obtaining an approximation similar to (9). Its presence in fact is
responsible for the $N^{-3 / 2} t$ term (instead of $N^{-1} t$ ) in the statement of Theorem 1.
The proof of Theorem 2 is laid out on a similar canvas. Let $k \in\left[\left\lfloor N^{\beta}\right\rfloor,\lfloor N / 2\rfloor\right]$. In Lemma 23 below, we show that the smallest singular value of $H_{\cdot,[N+n-k]}$, is a lower bound for $s_{N-k-1}\left(Y \Omega^{L} Y^{*}-\right.$ $z)$. As in the proof of Theorem 1, this value can be characterized through the action of $H_{\cdot,[N+n-k]}$ on the compressible and incompressible vectors. The former can be handled exactly as for $H$. The latter can also be reduced to a distance problem, and indeed this term is easier to tackle than earlier, thanks to the rectangular nature of $H_{\cdot,[N+n-k]}$.

### 3.2 Statistical analysis of the process $\left(y_{k}\right)$

Recall the decomposition $y_{k}=\tilde{y}_{k}+\check{y}_{k}$, with $\check{y}_{k}=\mathbb{E}\left[y_{k} \mid Y_{k}\right]$. We now derive approximations for $n \mathbb{E}\left[y_{k} y_{k}^{*}\right]$ and $n \mathbb{E}\left[\tilde{y}_{k} \tilde{y}_{k}^{*}\right]$. Let $\mathcal{R}$ denote the $N n \times N n$ block-Toeplitz matrix

$$
\mathcal{R}=\mathbb{E} \operatorname{vec} X(\operatorname{vec} X)^{*}=\frac{1}{n}\left[\begin{array}{cccc}
R_{0} & R_{-1} & & R_{-n+1}  \tag{10}\\
R_{1} & \ddots & \ddots & \\
& \ddots & \ddots & R_{-1} \\
R_{n-1} & & R_{1} & R_{0}
\end{array}\right]=\frac{1}{n}\left[R_{k-\ell}\right]_{k, \ell=0}^{n-1}
$$

Since the spectral density $S\left(e^{\imath \theta}\right)$ is non-trivial, it is easy to show that this matrix is invertible [8]. Moreover, $\|n \mathcal{R}\| \leqslant \boldsymbol{M}$ from Assumption 1-2). We also write the Fourier matrix F as

$$
\mathrm{F}=\left[\begin{array}{c}
\mathrm{f}_{0} \\
\vdots \\
\mathrm{f}_{n-1}
\end{array}\right]
$$

with $\mathrm{f}_{k} \in \mathbb{C}^{1 \times n}$ being the $k^{\text {th }}$ row of F , and we define the function

$$
\mathrm{a}\left(e^{\imath \theta}\right)=\left[\begin{array}{llll}
1 & e^{-\imath \theta} & \cdots & e^{-\imath(n-1) \theta}
\end{array}\right]^{\top} \in \mathbb{C}^{n} .
$$

From Equation (1), we have

$$
\mathcal{R}=\frac{1}{2 \pi n} \int_{0}^{2 \pi}\left(\mathrm{a}\left(e^{\imath \theta}\right) \otimes I_{N}\right) S\left(e^{\imath \theta}\right)\left(\mathrm{a}\left(e^{\imath \theta}\right)^{*} \otimes I_{N}\right) d \theta
$$

Lemma 9. Let Assumption 1 hold true. Then,

$$
\max _{k \in[n]}\left\|n \mathbb{E}\left[y_{k} y_{k}^{*}\right]-S^{(N)}\left(e^{2 \imath \pi k / n}\right)\right\| \xrightarrow[N \rightarrow \infty]{ } 0
$$

Proof. Let $k \in[n]$. Since $\mathrm{F}=\mathrm{F}^{\top}$, we have $y_{k}=X \mathrm{~F} e_{k}=X \mathrm{f}_{k}^{\top}=\left(\mathrm{f}_{k} \otimes I_{N}\right)$ vec $X$. Hence,

$$
\begin{align*}
n \mathbb{E} y_{k} y_{k}^{*} & =n\left(\mathrm{f}_{k} \otimes I_{N}\right) \mathcal{R}\left(\mathrm{f}_{k}^{*} \otimes I_{N}\right)=\frac{1}{2 \pi} \int_{0}^{2 \pi}\left(\left(\mathrm{f}_{k} \mathrm{a}(\theta)\right) \otimes I_{N}\right) S\left(e^{\imath \theta}\right)\left(\left(\mathrm{a}(\theta)^{*} \mathrm{f}_{k}^{*}\right) \otimes I_{N}\right) d \theta \\
& =\frac{1}{2 \pi} \int_{0}^{2 \pi}\left|\mathrm{f}_{k} \mathrm{a}(\theta)\right|^{2} S\left(e^{\imath \theta}\right) d \theta=\frac{1}{2 \pi} \int_{0}^{2 \pi} F_{n-1}(\theta-2 \pi k / n) S\left(e^{\imath \theta}\right) d \theta \tag{11}
\end{align*}
$$

where $F_{n-1}$ is the Fejér kernel, defined as

$$
\begin{equation*}
F_{n-1}(\theta)=\frac{1}{n}\left|\sum_{\ell=0}^{n-1} e^{\imath \ell \theta}\right|^{2}=\frac{1}{n} \frac{\sin (n \theta / 2)^{2}}{\sin (\theta / 2)^{2}} \tag{12}
\end{equation*}
$$

This kernel satisfies

$$
\frac{1}{2 \pi} \int_{0}^{2 \pi} F_{n-1}(\theta) d \theta=1, \quad \text { and } \quad F_{n-1}(\theta) \leqslant n\left(\frac{\pi^{2}}{n^{2} \theta^{2}} \wedge \frac{\pi^{2}}{4}\right) \text { for } \theta \in[-\pi, \pi]
$$

see, e.g., [30, Page 136]. For an arbitrarily small number $\eta>0$, we know by Assumption 1-1) that there exists $\delta>0$ independent of $k$ and $N$ such that $\left\|S\left(e^{2 \theta}\right)-S\left(e^{2 \imath \pi k / n}\right)\right\| \leqslant \eta$ if $|\theta-2 \pi k / n| \leqslant \delta$. Splitting the integral at the right hand side of (11) as $\int_{|\theta-2 \pi k / n| \leqslant \delta}+\int_{|\theta-2 \pi k / n|>\delta}$ and using the properties of the Fejér kernel provided above along with Assumption 1-2), we obtain the result of the lemma after some standard derivations.

The handling of $n \mathbb{E}\left[\tilde{y}_{k} \tilde{y}_{k}^{*}\right]$ is more involved.
Proposition 10. Let Assumptions 1 and 3 hold true. Then,

$$
\forall \delta>0, \max _{k \in[n]}\left\|n \mathbb{E}\left[\tilde{y}_{k} \tilde{y}_{k}^{*}\right]-S^{(N)}\left(e^{2 \imath \pi k / n}\right)\right\| \mathbb{1}_{S^{(N)}}\left(e^{2 \imath \pi k / n}\right) \geqslant \delta I_{N} \xrightarrow[N \rightarrow \infty]{ } 0
$$

Proof of Proposition 10. The first step is to provide a single letter expression for $\mathbb{E}\left[\tilde{y}_{k} \tilde{y}_{k}^{*}\right]$. This expression is reminiscent of the formula to calculate the partial covariance of real valued random variables.
Lemma 11. For any given $k \in[n], \mathbb{E}\left[\tilde{y}_{k} \tilde{y}_{k}^{*}\right]=\left(\left(\mathrm{f}_{k} \otimes I\right) \mathcal{R}^{-1}\left(\mathrm{f}_{k}^{*} \otimes I\right)\right)^{-1}$.
Proof. The covariance matrix $\mathbb{E}\left[\tilde{y}_{k} \tilde{y}_{k}^{*}\right]$ coincides with the conditional covariance matrix $\operatorname{cov}\left(y_{k} \mid Y_{k}\right)$ of $y_{k}$ with respect to $\sigma\left(Y_{k}\right)$. Let $Z_{0}$ and $Z_{1}$ be two random square integrable vectors with arbitrary dimensions. Writing

$$
\mathbb{E}\left[\begin{array}{l}
Z_{0} \\
Z_{1}
\end{array}\right]\left[\begin{array}{ll}
Z_{0}^{*} & Z_{1}^{*}
\end{array}\right]=\left[\begin{array}{ll}
\Sigma_{00} & \Sigma_{01} \\
\Sigma_{10} & \Sigma_{11}
\end{array}\right]
$$

with the block dimensions at the right hand side being compatible with the dimensions of $Z_{0}$ and $Z_{1}$, it is well-known that $\operatorname{cov}\left(Z_{0} \mid Z_{1}\right)=\Sigma_{00}-\Sigma_{01} \Sigma_{11}^{-1} \Sigma_{10}$ when $\Sigma_{11}$ is invertible. Observe that

$$
\mathbb{E} \operatorname{vec} Y(\operatorname{vec} Y)^{*}=\mathbb{E} \operatorname{vec}(X \mathrm{~F})(\operatorname{vec}(X \mathrm{~F}))^{*}=\left(\mathrm{F} \otimes I_{N}\right) \mathcal{R}\left(\mathrm{F}^{*} \otimes I_{N}\right)
$$

Let $\mathrm{F}_{k} \in \mathbb{C}^{(n-1) \times n}$ be the matrix that remains after taking the row $\mathrm{f}_{k}$ out of F , and let $\mathcal{R}^{1 / 2}$ be the Hermitian square root of $\mathcal{R}$. Then

$$
\begin{aligned}
\mathbb{E}\left[\tilde{y}_{k} \tilde{y}_{k}^{*}\right]= & \left(\mathrm{f}_{k} \otimes I_{N}\right) \mathcal{R}^{1 / 2}\left(I_{(n-1) N}-\mathcal{R}^{1 / 2}\left(\mathrm{~F}_{k}^{*} \otimes I_{N}\right)\left(\left(\mathrm{F}_{k} \otimes I_{N}\right) \mathcal{R}\left(\mathrm{F}_{k}^{*} \otimes I_{N}\right)\right)^{-1}\left(\mathrm{~F}_{k} \otimes I_{N}\right) \mathcal{R}^{1 / 2}\right) \\
& \mathcal{R}^{1 / 2}\left(\mathrm{f}_{k}^{*} \otimes I_{N}\right) \\
= & \left(\mathrm{f}_{k} \otimes I_{N}\right) \mathcal{R}^{1 / 2} \Pi_{\mathcal{R}^{1 / 2}\left(\mathrm{~F}_{k}^{*} \otimes I_{N}\right)}^{\perp} \mathcal{R}^{1 / 2}\left(\mathrm{f}_{k}^{*} \otimes I_{N}\right) .
\end{aligned}
$$

Let $\mathcal{C}$ be a positively oriented circle with center zero and a radius small enough so that $\mathcal{R}$ has no eigenvalue in the closed disk delineated by $\mathcal{C}$. This implies that the positive definite matrix $\left(\mathrm{F}_{k} \otimes I\right) \mathcal{R}\left(\mathrm{F}_{k}^{*} \otimes I\right)$ does not have eigenvalues in this closed disk either. With this choice, the projection $\Pi_{\mathcal{R}^{1 / 2}\left(F_{k}^{*} \otimes I_{N}\right)}^{\perp}$ can be expressed as a contour integral, which leads us to write

$$
\begin{aligned}
\mathbb{E}\left[\tilde{y}_{k} \tilde{y}_{k}^{*}\right] & =\frac{-1}{2 \imath \pi} \oint_{\mathcal{C}}\left(\mathrm{f}_{k} \otimes I_{N}\right) \mathcal{R}^{1 / 2}\left(\mathcal{R}^{1 / 2}\left(\mathrm{~F}_{k}^{*} \otimes I_{N}\right)\left(\mathrm{F}_{k} \otimes I_{N}\right) \mathcal{R}^{1 / 2}-z\right)^{-1} \mathcal{R}^{1 / 2}\left(\mathrm{f}_{k}^{*} \otimes I_{N}\right) d z \\
& =\frac{-1}{2 \imath \pi} \oint_{\mathcal{C}}\left(\mathrm{f}_{k} \otimes I_{N}\right) \mathcal{R}^{1 / 2}\left(\mathcal{R}-\mathcal{R}^{1 / 2}\left(\mathrm{f}_{k}^{*} \otimes I_{N}\right)\left(\mathrm{f}_{k} \otimes I_{N}\right) \mathcal{R}^{1 / 2}-z\right)^{-1} \mathcal{R}^{1 / 2}\left(\mathrm{f}_{k}^{*} \otimes I_{N}\right) d z
\end{aligned}
$$

Let $\mathcal{Q}(z)=(\mathcal{R}-z)^{-1}$ when $z \in \mathbb{C}$ is not an eigenvalue of $\mathcal{R}$, and let

$$
\Sigma(z)=\left(\mathrm{f}_{k} \otimes I_{N}\right) \mathcal{R}^{1 / 2} \mathcal{Q}(z) \mathcal{R}^{1 / 2}\left(\mathrm{f}_{k}^{*} \otimes I_{N}\right)=I_{N}+z\left(\mathrm{f}_{k} \otimes I_{N}\right) \mathcal{Q}(z)\left(\mathrm{f}_{k}^{*} \otimes I_{N}\right)
$$

Now using the Sherman-Morrison-Woodbury formula, and a small residue calculation for the analytic function $\mathcal{Q}(z)$ in a neighborhood of the closed disk enclosed by $\mathcal{C}$, we get

$$
\begin{aligned}
\mathbb{E}\left[\tilde{y}_{k} \tilde{y}_{k}^{*}\right] & =\frac{-1}{2 \imath \pi} \oint_{\mathcal{C}}\left(\Sigma(z)+\Sigma(z)\left(I_{N}-\Sigma(z)\right)^{-1} \Sigma(z)\right) d z=\frac{1}{2 \imath \pi} \oint_{\mathcal{C}} z^{-1}\left(\left(\mathrm{f}_{k} \otimes I_{N}\right) \mathcal{Q}(z)\left(\mathrm{f}_{k}^{*} \otimes I_{N}\right)\right)^{-1} d z \\
& =\left(\left(\mathrm{f}_{k} \otimes I_{N}\right) \mathcal{R}^{-1}\left(\mathrm{f}_{k}^{*} \otimes I_{N}\right)\right)^{-1}
\end{aligned}
$$

We now reinterpret the expression provided by this lemma in the light of the matrix orthogonal polynomial theory. Let us quickly review some of the basic results of this theory. For the proofs of these results, the reader may consult [8].

For each $N$, consider the $N \times N$ matrix valued measure $\varrho_{N}$ defined on $\mathbb{T}$ as

$$
d \varrho_{N}(\theta)=\frac{1}{2 \pi} S^{(N)}\left(e^{\imath \theta}\right) d \theta
$$

Given two $\mathbb{C}^{N \times N}$ matrix-valued polynomials $F, G$, we define the $N \times N$ matrix sesquilinear function $\langle\langle F, G\rangle\rangle_{\varrho_{N}}$ with respect to $\varrho_{N}$ as

$$
\left\langle\langle F, G\rangle \varrho_{\varrho_{N}}=\int_{0}^{2 \pi} F\left(e^{\imath \theta}\right)^{*} d \varrho_{N}(\theta) G\left(e^{\imath \theta}\right)\right.
$$

We now define the sequence $\left(\Phi_{\ell}^{\varrho_{N}}\right)_{\ell=0,1,2, \ldots}$ of matrix orthogonal polynomials with respect to $\langle\langle\cdot, \cdot\rangle\rangle_{\varrho_{N}}$. The following conditions, which are analogous to a Gram-Schmidt orthogonalization, are enough to define this sequence:

- $\Phi_{\ell}^{\varrho_{N}}(z)$ is a $\mathbb{C}^{N \times N}$-valued monic matrix polynomial of degree $\ell$, and is thus written as $\Phi_{\ell}^{\varrho_{N}}(z)=z^{\ell} I_{N}+$ lower order coefficients.
- For each $\ell \geqslant 1$, the relation $\left\langle\left\langle\Phi_{\ell}^{\varrho_{N}}, z^{k}\right\rangle\right\rangle_{\varrho_{N}}=0$ holds for $k=0,1, \ldots, \ell-1$.

Since $S^{(N)}$ is non-trivial, the matrix $\left\langle\left\langle\Phi_{\ell}^{\varrho_{N}}, \Phi_{\ell}^{\varrho_{N}}\right\rangle \varrho_{N}\right.$ is positive definite for each $\ell \in \mathbb{N}$. Thus, one can define the sequence $\left(\varphi_{\ell}^{\varrho_{N}}\right)_{\ell \in \mathbb{N}}$ of the normalized versions of the polynomials $\Phi_{\ell}^{\varrho_{N}}$ as

$$
\varphi_{\ell}^{\varrho_{N}}(z)=\Phi_{\ell}^{\varrho_{N}}(z) \kappa_{\ell}^{\varrho_{N}}
$$

where the matrices $\left(\kappa_{\ell}^{\varrho_{N}}\right)_{\ell \in \mathbb{N}}$ are chosen in such a way that $\left\langle\left\langle\varphi_{\ell}^{\varrho_{N}}, \varphi_{k}^{\varrho_{N}}\right\rangle\right\rangle_{\varrho_{N}}=\mathbb{1}_{k=\ell} I_{N}$. This identity determines the matrices $\kappa_{\ell}^{\varrho_{N}}$ up to a right multiplication by a unitary matrix, the convenient choice of which is specified in $[8, \S 3.2]$ and is not relevant here.

The Christoffel-Darboux (CD) kernel of order $\ell$ for the measure $\varrho_{N}$ is defined for $z, u \in \mathbb{C}$ as

$$
K_{\ell}^{\varrho_{N}}(z, u)=\sum_{k=0}^{\ell} \varphi_{k}^{\varrho_{N}}(z) \varphi_{k}^{\varrho_{N}}(u)^{*}
$$

This function is a reproducing kernel, in the sense that for each matrix polynomial $P(z)$ with degree less or equal to $\ell$, the following equation holds and defines the kernel $K_{\ell}^{\varrho_{N}}$.

$$
\int_{0}^{2 \pi} K_{\ell}^{\varrho_{N}}\left(z, e^{\imath \theta}\right) d \varrho_{N}(\theta) P\left(e^{\imath \theta}\right) d \theta=P(z)
$$

The CD kernel $K_{\ell}^{Q_{N}}\left(e^{\imath \theta}, e^{\imath \theta}\right)$ is invertible, and it satisfies the following variational formula: For each $\theta \in[0,2 \pi)$ and for each $N \times N$ matrix polynomial $P(z)$ of degree at most $\ell$, with

$$
P\left(e^{\imath \theta}\right)=I_{N}
$$

it holds that

$$
\begin{equation*}
\langle\langle P, P\rangle\rangle_{\varrho_{N}} \geqslant K_{\ell}^{\varrho_{N}}\left(e^{\imath \theta}, e^{\imath \theta}\right)^{-1} \tag{13}
\end{equation*}
$$

with equality if and only if $P(z)=P_{\ell}^{\varrho_{N}, \theta}(z)$, with

$$
\begin{equation*}
P_{\ell}^{\varrho_{N}, \theta}(z)=K_{\ell}^{\varrho_{N}}\left(z, e^{\imath \theta}\right) K_{\ell}^{\varrho_{N}}\left(e^{\imath \theta}, e^{\imath \theta}\right)^{-1} . \tag{14}
\end{equation*}
$$

Getting back to our problem, it turns out that our covariance matrix $\mathbb{E}\left[\tilde{y}_{k} \tilde{y}_{k}^{*}\right]$ can be very simply expressed in terms of a CD kernel:
Lemma 12. $\mathbb{E}\left[\tilde{y}_{k} \tilde{y}_{k}^{*}\right]=K_{n-1}^{\varrho_{N}}\left(e^{2 \imath \pi k / n}, e^{2 \imath \pi k / n}\right)^{-1}$.
Proof. With Lemma 11 at hand, this result is an instance of a well-known result, see, e.g., [16] for the scalar measure case. We reproduce its proof for completeness. Let $P(z)=\sum_{k=0}^{n-1} z^{k} P_{k}$ be an arbitrary $N \times N$ matrix polynomial of degree at most $n-1$. Stack the coefficients of $P$ in the matrix $\mathrm{P}=\left[\begin{array}{lll}P_{0}^{\top} & \cdots & P_{n-1}^{T}\end{array}\right]^{\top}$, so that $P\left(e^{\imath \theta}\right)=\left(\mathrm{a}\left(e^{\imath \theta}\right)^{*} \otimes I_{N}\right) \mathrm{P}$. We have

$$
\begin{aligned}
& \frac{1}{2 \pi n} \int_{0}^{2 \pi}\left(\mathrm{a}\left(e^{\imath \theta}\right)^{*} \otimes I_{N}\right) \mathcal{R}^{-1}\left(\mathrm{a}\left(e^{\imath \psi}\right) \otimes I_{N}\right) S\left(e^{\imath \psi}\right) P\left(e^{\imath \psi}\right) d \psi \\
& =\left(\mathrm{a}\left(e^{\imath \theta}\right)^{*} \otimes I_{N}\right) \mathcal{R}^{-1}\left(\frac{1}{2 \pi n} \int_{0}^{2 \pi}\left(\mathrm{a}\left(e^{\imath \psi}\right) \otimes I_{N}\right) S\left(e^{\imath \psi}\right)\left(\mathrm{a}\left(e^{\imath \psi}\right)^{*} \otimes I_{N}\right) d \psi\right) \mathrm{P} \\
& =\left(\mathrm{a}\left(e^{\imath \theta}\right)^{*} \otimes I_{N}\right) \mathcal{R}^{-1} \mathcal{R} \mathrm{P} \\
& =P\left(e^{\imath \theta}\right) .
\end{aligned}
$$

From the uniqueness of the CD kernel as a reproducing kernel, we thus obtain that

$$
n^{-1}\left(\mathrm{a}\left(e^{\imath \theta}\right)^{*} \otimes I_{N}\right) \mathcal{R}^{-1}\left(\mathrm{a}\left(e^{\imath \psi}\right) \otimes I_{N}\right)=K_{n-1}^{\varrho_{N}}\left(e^{\imath \theta}, e^{\imath \psi}\right)
$$

The result follows upon observing that $\mathrm{f}_{k}=n^{-1 / 2} \mathrm{a}\left(e^{2 \pi \imath k / n}\right)^{*}$.
Now Proposition 10 will be obtained by studying the asymptotics of the CD kernel, a subject with a rich history in the literature of orthogonal polynomials. These asymptotics are well understood in the case where the underlying measure has a so-called regularity property à la Stahl and Totik [31] (see [30] for an extensive treatment of these ideas). We adapt the approach detailed in [ $30, \S 2.15$ and $\S 2.16]$, to the matrix measure case. This will consist of two main steps:

- We show that the polynomial $P_{n-1}^{\varrho_{N}, \theta}$, for which the variational inequality (13) for $\ell=n-1$ is an equality, satisfies

$$
\begin{equation*}
\max _{\theta, \psi \in[0,2 \pi)}\left\|P_{n-1}^{\varrho_{N}, \theta}\left(e^{\imath \psi}\right)\right\| \leqslant e^{n \varepsilon_{N}} \quad \text { with } \quad 0 \leqslant \varepsilon_{N} \rightarrow 0 \tag{15}
\end{equation*}
$$

This is where Assumption 3 comes into play.

- With the help of the variational characterization of the CD kernel, and by making use of the previous result coupled with a matrix version of the so-called Nevai's trial polynomial technique $[19,30]$, we show that

$$
\begin{align*}
\forall \varepsilon, \delta>0, \exists N_{0} \in \mathbb{N}^{*}, \forall N \geqslant N_{0}, \forall \theta_{0} \in[0,2 \pi), & S^{(N)}\left(e^{\imath \theta_{0}}\right) \geqslant \delta I_{N} \\
& \Rightarrow n K_{n-1}^{\varrho_{N}}\left(e^{2 \theta_{0}}, e^{\imath \theta_{0}}\right)^{-1} \geqslant(1-\varepsilon) S^{(N)}\left(e^{2 \theta_{0}}\right) . \tag{16}
\end{align*}
$$

In the following derivations, the sequence $\varepsilon_{N} \rightarrow 0$ can change from one display to another.
For a given $N \in \mathbb{N}^{*}$, consider the scalar measure $\zeta_{N}$ on $\mathbb{T}$ defined as

$$
d \zeta_{N}(\theta)=\frac{1}{2 \pi} s_{N-1}\left(S^{(N)}\left(e^{\imath \theta}\right)\right) d \theta
$$

Consider the sesquilinear form on the scalar polynomials

$$
\langle f, g\rangle_{\zeta_{N}}=\int_{0}^{2 \pi} f\left(e^{\imath \theta}\right)^{*} g\left(e^{\imath \theta}\right) d \zeta_{N}\left(e^{\imath \theta}\right)
$$

and denote as $\left(b_{\ell}^{\zeta_{N}}\right)_{\ell \in \mathbb{N}}$ the sequence of monic orthogonal polynomials with respect to $\langle\cdot, \cdot\rangle_{\zeta_{N}}$, which are the analogues of the $\Phi_{\ell}^{\varrho_{N}}$ above. Let $\tau_{\ell}^{\zeta_{N}}=\left\langle b_{\ell}^{\zeta_{N}}, b_{\ell}^{\zeta_{N}}\right\rangle_{\zeta_{N}}$, and consider the orthonormal polynomials $\left(\beta_{\ell}^{\zeta_{N}}\right)_{\ell \in \mathbb{N}}$ defined as $\beta_{\ell}^{\zeta_{N}}(z)=b_{\ell}^{\zeta_{N}}(z) / \sqrt{\tau_{\ell}^{\zeta_{N}}}$ (these are the analogues of the $\varphi_{\ell}^{\varrho_{N}}$ ). To establish (15), we first show that

$$
\begin{equation*}
\max _{\ell \in[n]} \max _{\theta \in[0,2 \pi)}\left|\beta_{\ell}^{\zeta_{N}}\left(e^{\imath \theta}\right)\right| \leqslant e^{n \varepsilon_{N}} \quad \text { with } \quad \varepsilon_{N} \geqslant 0, \varepsilon_{N} \rightarrow 0 \tag{17}
\end{equation*}
$$

Write

$$
\sigma_{N}^{2}=\tau_{0}^{\zeta_{N}}=\int_{0}^{2 \pi} d \zeta_{N}\left(e^{\imath \theta}\right)=\frac{1}{2 \pi} \int_{0}^{2 \pi} s_{N-1}\left(S^{(N)}\left(e^{\imath \theta}\right)\right) d \theta
$$

By the celebrated Szegö's theorem [30], $\left(\tau_{\ell}^{\zeta_{N}}\right)_{\ell}$ is a non-increasing sequence that converges as $\ell \rightarrow \infty$ towards $\exp \left((2 \pi)^{-1} \int_{0}^{2 \pi} \log s_{N-1}\left(S^{(N)}\left(e^{\imath \theta}\right)\right) d \theta\right)$. Thus, using Assumption 1-2), there is a constant $C>0$ such that

$$
C \geqslant \sigma_{N}^{2} \geqslant \tau_{n-1}^{\zeta_{N}} \geqslant \exp \left(\frac{1}{2 \pi} \int_{0}^{2 \pi} \log s_{N-1}\left(S^{(N)}\left(e^{\imath \theta}\right)\right) d \theta\right)
$$

By inspecting the proofs of [30, Prop. 2.15.12 and Th. 2.15.3], we furthermore get that for $|z|>1$,

$$
\max _{\ell \in[n]} \frac{\left|b_{\ell}^{\zeta_{N}}(z)\right|}{|z|^{\ell}} \leqslant \exp \left(n \sqrt{\frac{1}{n}\left(-\log \tau_{n-1}^{\zeta_{N}}+\log \sigma_{N}^{2}\right)}\right)
$$

(in [30], the analogue of our measure $\zeta_{N}$ is assumed to be a probability measure, hence our need to add the $\log \sigma_{N}^{2}$ factor in the expression above which is absent in [30]). Thus, since

$$
\frac{1}{n} \log C \geqslant \frac{1}{n} \log \sigma_{N}^{2} \geqslant \frac{1}{n} \log \tau_{n-1}^{\zeta_{N}} \geqslant \frac{1}{2 \pi n} \int_{0}^{2 \pi} \log s_{N-1}\left(S^{(N)}\left(e^{\imath \theta}\right)\right) d \theta
$$

we obtain from Assumption 3 that there exists a non-negative sequence $\varepsilon_{N} \rightarrow_{N} 0$ such that

$$
\max _{\ell \in[n]} \frac{\left|b_{\ell}^{\zeta_{N}}(z)\right|}{|z|^{\ell}} \leqslant \exp \left(n \varepsilon_{N}\right) .
$$

By the maximum principle, for any $\theta \in[0,2 \pi)$, any $\ell \in[n]$ and any $r>1$,

$$
\left|\beta_{\ell}^{\zeta_{N}}\left(e^{\imath \theta}\right)\right|=\left(\tau_{\ell}^{\zeta_{N}}\right)^{-1 / 2}\left|b_{\ell}^{\zeta_{N}}\left(e^{\imath \theta}\right)\right| \leqslant\left(\tau_{n-1}^{\zeta_{N}}\right)^{-1 / 2} \sup _{\psi \in[0,2 \pi)}\left|b_{\ell}^{\zeta_{N}}\left(r e^{\imath \psi}\right)\right| \leqslant r^{\ell}\left(\tau_{n-1}^{\zeta_{N}}\right)^{-1 / 2} \exp \left(n \epsilon_{N}\right)
$$

and by Szegö's theorem and Assumption 3 again, we obtain (17) by choosing $r$ as close to one as desired.

Proceeding, we now consider the matrix measure $d \zeta_{N}\left(e^{\imath \theta}\right) \otimes I_{N}=(2 \pi)^{-1} s_{N-1}\left(S^{(N)}\left(e^{\imath \theta}\right)\right) I_{N} d \theta$, equipped with the sesquilinear function

$$
\langle\langle F, G\rangle\rangle_{\zeta_{N} \otimes I_{N}}=\frac{1}{2 \pi} \int_{0}^{2 \pi} s_{N-1}\left(S^{(N)}\left(e^{\imath \theta}\right)\right) F^{*}\left(e^{\imath \theta}\right) G\left(e^{\imath \theta}\right) d \theta
$$

It is clear that the $\ell^{\text {th }}$ normalized orthogonal polynomial for $\langle\langle\cdot, \cdot\rangle\rangle \zeta_{N} \otimes I_{N}$ is $\beta_{\ell}^{\zeta_{N}}(z) I_{N}$, and the associated $\ell^{\text {th }}$ CD kernel is

$$
K_{\ell}^{\zeta_{N} \otimes I_{N}}(z, u)=\sum_{k=0}^{\ell} \beta_{k}^{\zeta_{N}}(z)\left(\beta_{k}^{\zeta_{N}}(u)\right)^{*} I_{N}
$$

Obviously, for each Borel set $A \subset[0,2 \pi), \zeta_{N}(A) I_{N} \leqslant \varrho_{N}(A)$ in the semi-definite ordering. Therefore, by the variational characterization of the CD kernels, we have

$$
K_{n-1}^{\zeta_{N} \otimes I_{N}}\left(e^{\imath \theta}, e^{\imath \theta}\right) \geqslant\left\langle\left\langle P_{n-1}^{\varrho_{N}, \theta}, P_{n-1}^{\varrho_{N}, \theta}\right\rangle\right\rangle_{\zeta_{N} \otimes I_{N}}^{-1} \geqslant\left\langle\left\langle P_{n-1}^{\varrho_{N}, \theta}, P_{n-1}^{\varrho_{N}, \theta}\right\rangle\right\rangle_{\varrho_{N}}^{-1}=K_{n-1}^{\varrho_{N}}\left(e^{\imath \theta}, e^{\imath \theta}\right)
$$

for all $\theta \in[0,2 \pi)$. From what precedes, this implies that

$$
\max _{\ell \in[n]} \max _{\theta \in[0,2 \pi)}\left\|\varphi_{\ell}^{\varrho_{N}}\left(e^{\imath \theta}\right)\right\| \leqslant e^{n \varepsilon_{N}} \quad \text { with } \quad 0 \leqslant \varepsilon_{N} \rightarrow 0
$$

Also notice that

$$
K_{n-1}^{\varrho_{N}}\left(e^{\imath \theta}, e^{\imath \theta}\right) \geqslant \varphi_{0}^{\varrho_{N}}\left(e^{\imath \theta}\right) \varphi_{0}^{\varrho_{N}}\left(e^{\imath \theta}\right)^{*}=\left(\int_{0}^{2 \pi} d \varrho_{N}(\theta)\right)^{-1}=\left(R_{0}^{(N)}\right)^{-1}
$$

Thus, using Assumption 1-3) and recalling Definition (14) of the polynomials $P_{\ell}^{\varrho_{N}, \theta}$, we obtain the bound (15).

We now prove (16). Let $\delta>0$ and $\theta_{0} \in(0,2 \pi)$ be such that $S^{(N)}\left(e^{\imath \theta_{0}}\right) \geqslant \delta I_{N}$. Consider the matrix measure

$$
d \varsigma_{N}(\theta)=\frac{1}{2 \pi} S^{(N)}\left(e^{\imath \theta_{0}}\right) d \theta
$$

The CD kernels for this measure are constant and satisfy the identity $\ell K_{\ell-1}^{\varsigma_{N}}\left(e^{\imath \theta}, e^{\imath \theta}\right)^{-1}=S^{(N)}\left(e^{\imath \theta_{0}}\right)$ for each integer $\ell>0$, as can be checked by, e.g., the technique of Lemma 12.

Let $\eta>0$ be an arbitrarily small number. By Assumption 1-1), we can choose $h>0$ small enough so that for each $N>0, \boldsymbol{w}\left(S^{(N)}, h\right)<\eta \delta$. Let $\theta \in\left[\theta_{0}-h, \theta_{0}+h\right]$. Then, for each vector $u \in \mathbb{C}^{N}$ with $\|u\|=1$, we have

$$
u^{*} S^{(N)}\left(e^{\imath \theta}\right) u=u^{*} S^{(N)}\left(e^{\imath \theta_{0}}\right) u\left(1+\frac{u^{*}\left(S^{(N)}\left(e^{\imath \theta}\right)-S^{(N)}\left(e^{\imath \theta_{0}}\right)\right) u}{u^{*} S^{(N)}\left(e^{\imath \theta_{0}}\right) u}\right)
$$

and hence,

$$
\forall \theta \in\left[\theta_{0}-h, \theta_{0}+h\right], S^{(N)}\left(e^{\imath \theta_{0}}\right) \leqslant(1+2 \eta) S^{(N)}\left(e^{\imath \theta}\right)
$$

The degree one polynomial $q(z)=0.5\left(z e^{-\imath \theta_{0}}+1\right)$ can be easily shown to satisfy $q\left(e^{\imath \theta_{0}}\right)=$ $\sup _{\theta}\left|q\left(e^{\imath \theta}\right)\right|=1$, and $\sup _{\theta-\theta_{0} \in[-\pi, \pi),\left|\theta-\theta_{0}\right| \geqslant h}\left|q\left(e^{2 \theta}\right)\right|=\cos (h / 2)<1$. Let $\tilde{\eta}>0$ be arbitrarily small, and let $m=n-1+\lfloor n \tilde{\eta}\rfloor$. Consider the "trial polynomial" $Q_{m}\left(e^{\imath \theta}\right)=P_{n-1}^{\varrho_{N}, \theta_{0}}\left(e^{\imath \theta}\right) q\left(e^{\imath \theta}\right)^{\lfloor n \tilde{\eta}\rfloor}$ with degree $m$. This polynomial has the following features:

- $Q_{m}\left(e^{2 \theta_{0}}\right)=I_{N}$.
- By Assertion (15) that we just proved, for $\theta-\theta_{0} \in[-\pi, \pi),\left|\theta-\theta_{0}\right| \geqslant h$, for all large $N$,

$$
\| Q_{n}\left(e^{\imath \theta}\left\|\leqslant|\cos (h / 2)|^{n \tilde{\eta}}\right\| P_{n-1}^{\varrho_{N}, \theta_{0}}\left(e^{\imath \theta}\right) \| \leqslant|\cos (h / 2)|^{n \tilde{\eta} / 2}\right.
$$

By the variational characterization of the CD kernels, we have

$$
\begin{aligned}
\frac{1}{m+1} S^{(N)}\left(e^{\imath \theta_{0}}\right) & =\left(K_{m}^{\varsigma_{N}}\left(e^{\imath \theta_{0}}, e^{\imath \theta_{0}}\right)\right)^{-1} \\
& =\left\langle\left\langle P_{m}^{\varsigma_{N}, \theta_{0}}, P_{m}^{\varsigma_{N}, \theta_{0}}\right\rangle\right\rangle_{\varsigma_{N}} \\
& \leqslant\left\langle\left\langle Q_{m}, Q_{m}\right\rangle\right\rangle_{\varsigma_{N}} \\
& \leqslant(1+2 \eta) \int_{\theta:\left|\theta-\theta_{0}\right| \leqslant h} Q_{m}\left(e^{\imath \theta}\right) d \varrho_{N}(\theta) Q_{m}\left(e^{\imath \theta}\right)^{*}+\int_{\theta:\left|\theta-\theta_{0}\right|>h} Q_{m}\left(e^{\imath \theta}\right) d \varsigma_{N}(\theta) Q_{m}\left(e^{\imath \theta}\right)^{*} \\
& \leqslant(1+2 \eta)\left\langle\left\langle P_{n-1}^{\varrho_{N}, \theta_{0}}, P_{n-1}^{\varrho_{N}, \theta_{0}}\right\rangle\right\rangle \varrho_{N}+|\cos (h / 2)|^{n \tilde{\eta}} \boldsymbol{M} \\
& =(1+2 \eta) K_{n-1}^{\varrho_{N}}\left(e^{\imath \theta_{0}}, e^{\imath \theta_{0}}\right)^{-1}+|\cos (h / 2)|^{n \tilde{\eta}} \boldsymbol{M}
\end{aligned}
$$

Since $\eta$ and $\tilde{\eta}$ are arbitrary, Assertion 16 holds true.
We now have all the elements to complete the proof of Proposition 12. Let $k \in[n]$ be such that $S^{(N)}\left(e^{2 \imath \pi k / n}\right) \geqslant \delta I_{N}$. Thanks to Lemma 12, we can replace $K_{n-1}^{\varrho_{N}}\left(e^{2 \theta_{0}}, e^{2 \theta_{0}}\right)^{-1}$ for $\theta_{0}=2 \pi k / n$ with $\mathbb{E} \tilde{y}_{k} \tilde{y}_{k}^{*}$ in Assertion 16. This provides a lower bound on $\mathbb{E} \tilde{y}_{k} \tilde{y}_{k}^{*}$. To obtain an upper bound on this matrix, we recall that $\tilde{y}_{k}$ and $\check{y}_{k}$ are independent, resulting in $\mathbb{E} \tilde{y}_{k} \tilde{y}_{k}^{*} \leqslant \mathbb{E} y_{k} y_{k}^{*}$, and use Lemma 9.

This completes the proof of Proposition 10. We note for completeness that we could have established the upper bound by using a Nevai's polynomial trial technique as well.

### 3.3 Technical results needed in the proofs of Theorems 1 and 2

The following lemma can be proven similarly to, e.g., [36, Corollary 4.2.13].
Lemma 13 (covering number). Let $\mathcal{U} \subset \mathbb{C}^{m}$ be a $k$-dimensional subspace. Given $\rho>0$, there exists a $\rho$-net of the Euclidean unit-ball of $\mathcal{U}$ with a cardinality bounded above by $(3 / \rho)^{2 k}$.

Lemma 14. Let $a_{0}, \ldots, a_{m-1}$ be non-negative real numbers such that there exist $0<c \leqslant C$ for which $c \leqslant m^{-1} \sum a_{k}$, and $\max a_{k} \leqslant C$. Given $x \in(0, c]$, the set $\mathcal{I}(x) \subset[m]$ defined as

$$
\mathcal{I}(x)=\left\{k \in[m]: a_{k}>x\right\}
$$

satisfies

$$
|\mathcal{I}| \geqslant \frac{c-x}{C-x} m
$$

Proof. We have $C|\mathcal{I}| \geqslant \sum_{k \in \mathcal{I}} a_{k} \geqslant m c-\sum_{k \in \mathcal{I}^{c}} a_{k} \geqslant m c-(m-|\mathcal{I}|) x$, hence the result.
For any random vector $\xi=\left[\xi_{0}, \ldots, \xi_{m-1}\right]^{\top} \in \mathbb{C}^{m}$, Lévy's anti-concentration function is defined for $t>0$ as

$$
\mathcal{L}(\xi, t)=\sup _{d \in \mathbb{C}^{m}} \mathbb{P}[\|\xi-d\| \leqslant t]
$$

(when the probability $\mathbb{P}$ above is taken with respect to some random vector $x$, we denote the associated Lévy's anti-concentration function as $\left.\mathcal{L}_{x}(\xi, t)\right)$. Assuming that the elements $\xi_{\ell}$ of $\xi$ are independent, letting $k \in[m]$, and defining $\xi^{(k)}=\left[\xi_{0}, \ldots, \xi_{k}\right]^{\top}$, the following restriction result is well-known [26, Lemma 2.1].

$$
\mathcal{L}(\xi, t) \leqslant \mathcal{L}\left(\xi^{(k)}, t\right)
$$

We shall need the following rather standard results on Gaussian vectors.
Lemma 15 (norm and anti-concentration results for Gaussian vectors). Let $\xi=\left[\xi_{0}, \ldots, \xi_{N-1}\right]^{\top} \sim$ $\mathcal{N}_{\mathbb{C}}\left(0, I_{N}\right)$, and let $\Sigma$ be an $N \times N$ covariance matrix. Then:

1. For $t>0$, it holds that $\mathbb{P}\left[\left\|\Sigma^{1 / 2} \xi\right\| \geqslant \sqrt{2 N t}\right] \leqslant \exp (-(t /\|\Sigma\|-1) N)$.

Assume that $s_{m-1}(\Sigma) \geqslant \alpha$ for some $m \in[N]$ and some $\alpha>0$. Then:
2. There exists a constant $c_{15,2}>0$ such that $\mathbb{P}\left[\left\|\Sigma^{1 / 2} \xi\right\| \leqslant \sqrt{\alpha m / 2}\right] \leqslant \exp \left(-c_{15,2} m\right)$.
3. There exists a constant $C_{15,3}>0$ such that for each deterministic non-zero matrix $M \in$ $\mathbb{C}^{N \times N}$ and each deterministic vector $a \in \mathbb{C}^{N}$,

$$
\mathcal{L}\left((\xi+a)^{*} M(\xi+a) /\|M\|_{\mathrm{HS}}, t\right) \leqslant C_{15,3} t
$$

4. There exists a constant $C_{15,4}=C_{15,4}(\alpha)>0$ such that $\mathcal{L}\left(\Sigma^{1 / 2} \xi, \sqrt{m} t\right) \leqslant\left(C_{15,4} t\right)^{m}$.
5. For each non-zero deterministic matrix $M \in \mathbb{C}^{N \times N}$,

$$
\mathbb{P}\left[\|M \xi\|^{2} \geqslant t\|M\|_{\mathrm{HS}}^{2}\right] \leqslant \exp (1-t / 2)
$$

This lemma is proven in the appendix.
The following results specifically concern our matrix model. The first one will be needed to bound the spectral norm of our matrix $Y$.

Lemma 16 (spectral norm of $Y$ ). Let Assumption 1-2) hold true. Then, there are two constants $c_{16}, c_{16}^{\prime}>0$ such that for each $t \geqslant c_{16}^{\prime}$,

$$
\mathbb{P}[\|Y\| \geqslant t] \leqslant \exp \left(-c_{16} N t^{2}\right)
$$

Proof. Since $Y=X F$, it is enough to prove the lemma for $X$. By a standard $\varepsilon$-net argument (see [32, Lm. 2.3.2]), we know that

$$
\mathbb{P}[\|X\| \geqslant t] \leqslant \sum_{u \in N} \mathbb{P}[\|X u\| \geqslant t / 2]
$$

where $\boldsymbol{N}$ is a $1 / 2-$ net of $\mathbb{S}^{n-1}$. Let $u \in \boldsymbol{N}$. Observing that $X u \sim \mathcal{N}_{\mathbb{C}}\left(0,\left(u^{\top} \otimes I_{N}\right) \mathcal{R}\left(\bar{u} \otimes I_{N}\right)\right)$, and recalling that $n\|\mathcal{R}\| \leqslant \boldsymbol{M}$, Lemma 15-1 shows that

$$
\left.\mathbb{P}[\|X u\| \geqslant t / 2] \leqslant \exp \left(-\left(\frac{t^{2} n}{8 \boldsymbol{M} N}\right)-1\right) N\right)
$$

By Lemma 13 and the union bound, we thus have that

$$
\mathbb{P}[\|X\| \geqslant t] \leqslant \exp \left(-\left(\frac{t^{2} n}{8 \boldsymbol{M} N}-1\right) N+(2 \log 6) n\right)
$$

Choosing $t$ large enough, we get the result.
We shall need to use a discrete frequency counterpart of Assumption 2:
Lemma 17. Let Assumptions $1-1$ and 2 hold true. Given $N \in \mathbb{N}^{*}$ and $\alpha>0$, let $\mathcal{C}_{\text {good }}(\alpha)$ be the subset of $[n]$ defined as

$$
\mathcal{C}_{\text {good }}(\alpha)=\left\{k \in[n]: S\left(e^{2 \imath \pi k / n}\right) \geqslant \alpha I_{N}\right\}
$$

Then,

$$
\forall \kappa \in(0,1), \exists \alpha>0,\left|\mathcal{C}_{\text {good }}(\alpha)\right| \geqslant(1-\kappa) n \quad \text { for all large } N .
$$

Proof. Let us identify a set $\mathcal{I} \subset[n]$ with the discretization of the unit-circle $\mathbb{T}$ given as $\{\exp (2 \imath \pi k / n)$ : $k \in \mathcal{I}\}$, and let us also denote any of these two sets as $\mathcal{I}$. Given a small real number $h>0$, let $\mathcal{V}_{h}^{\mathbb{T}}(\mathcal{I})$ be the closed $h$-neighborhood of $\mathcal{I}$ on $\mathbb{T}$ equipped with the curvilinear distance. We shall show that

$$
\begin{equation*}
|\mathcal{I}| \leqslant \frac{n}{2 \pi} \operatorname{Leb}\left(\mathcal{V}_{h}^{\mathbb{T}}(\mathcal{I})\right)+\frac{\pi}{h} \tag{18}
\end{equation*}
$$

We first observe that $\mathcal{V}_{h}^{\mathbb{T}}(\mathcal{I})$ consists of a finite number of disjoint closed arcs, the length of each arc being greater or equal to $2 h$. Given $\varphi_{0} \in[0,2 \pi)$ and $\varphi_{1} \in\left[\varphi_{0}+2 h, \varphi_{0}+2 h+2 \pi\right)$, let $\left\{e^{\imath \varphi}: \varphi_{0} \leqslant \varphi \leqslant \varphi_{1}\right\}$ be one such arc. Then, there are two integers $k_{\min } \leqslant k_{\max }$ in $\mathcal{I}$ such that $2 k_{\min } \pi / n=\varphi_{0}+h$ and $2 k_{\max } \pi / n=\varphi_{1}-h$, and all the elements of $\mathcal{I}$ in this arc belong to the set $\left\{k_{\min }, \ldots, k_{\max }\right\}$. Since there are at $\operatorname{most} k_{\max }-k_{\min }+1=n\left(\varphi_{1}-\varphi_{0}\right) /(2 \pi)-n h / \pi+1$ of these elements, and furthermore, since $\mathcal{V}_{h}^{\mathbb{T}}(\mathcal{I})$ consists of at most $\lfloor\pi / h\rfloor$ arcs, we obtain the inequality (18).

Fix an arbitrary $\kappa \in(0,1)$. By Assumption 2, there exists $\delta>0$ such that the set $B=$ $\left\{z \in \mathbb{T}: s_{N-1}(S(z)) \leqslant \delta\right\}$ satisfies $\operatorname{Leb}(B) \leqslant \kappa \pi$. Let $\alpha=\delta / 2$, and assume that the set $\mathcal{C}_{\text {bad }}(\alpha)=[n] \backslash \mathcal{C}_{\text {good }}(\alpha)$ is non-empty (otherwise the result of the lemma is true). Relying on Assumption 1-1), let $h>0$ be such that $\boldsymbol{w}(S, h) \leqslant \alpha$ for all $N$. Let $k \in \mathcal{C}_{\text {bad }}(\alpha)$. By the Wielandt-Hoffmann theorem and the triangle inequality, $s_{N-1}\left(S\left(e^{2 \theta}\right)\right) \leqslant \delta$ for each $\theta$ such that $|\theta-2 \pi k / n| \leqslant h$. In other words, $\mathcal{V}_{h}^{\mathbb{T}}\left(\mathcal{C}_{\text {bad }}(\alpha)\right) \subset B$. By inequality (18), we thus obtain that

$$
\left|\mathcal{C}_{\text {bad }}(\alpha)\right| \leqslant \frac{n}{2} \kappa+\frac{\pi}{h}
$$

which implies that $\left|\mathcal{C}_{\text {bad }}(\alpha)\right| \leqslant \kappa n$ for all large $n$.

We now enter the proofs of Theorems 1 and 2 . Given $C>0$, we denote as $\mathcal{E}_{\text {op }}(C)$ the event

$$
\mathcal{E}_{\mathrm{op}}(C)=[\|Y\| \leqslant C] .
$$

In the remainder of this section, the constants will be referred to by the letters $c, a$, or $C$, possibly with primes or numerical indices that refer generally to the statements (lemmas, propositions, ...) where these constants appear for the first time. These constants do not depend on $N$. The typical statements where they appear are of the type

$$
\mathbb{P}\left[[\cdots \leqslant a] \cap \mathcal{E}_{\mathrm{op}}(C)\right] \leqslant C^{\prime} a+\exp (-c N)
$$

or others in the same vein. Often, such inequalities hold true for all $N$ large enough. This detail will not always be mentioned.

### 3.4 Proof of Theorem 1

As mentioned in Section 3.1, our starting point for proving Theorem 1 is the variational characterization (8).

### 3.4.1 Compressible vectors

In the following statement and in the proof, the vectors $u \in \mathbb{S}^{N+n-1}$ will always take the form $u=\left[v^{\top} w^{\top}\right]^{\top}$ with $v \in \mathbb{C}^{n}$.

Proposition 18 (compressible vectors). There exist constants $\theta_{18}, \rho_{18}, C_{18} \in(0,1)$, and $a_{18}, c_{18}>$ 0 such that, for the set

$$
\mathcal{S}:=\left\{u \in \mathbb{S}^{N+n-1}:\|v\| \leqslant C_{18} \text { or } v /\|v\| \in \operatorname{comp}\left(\theta_{18}, \rho_{18}\right)\right\}
$$

(here $v /\|v\|=0$ if $v=0$ ), we have

$$
\mathbb{P}\left[\left[\inf _{u \in \mathcal{S}}\|H u\| \leqslant a_{18}\right] \cap \mathcal{E}_{\text {op }}(C)\right] \leqslant \exp \left(-c_{18} n\right) .
$$

Proof. For each $u \in \mathbb{S}^{N+n-1}$ the inequality $\|H u\| \leqslant a$ for $a>0$ implies

$$
\begin{align*}
\|Y v+z w\| & \leqslant a  \tag{19a}\\
\left\|\Omega^{-L} v+Y^{*} w\right\| & \leqslant a \tag{19b}
\end{align*}
$$

Take $a=|z| / 2$. On the event $\mathcal{E}_{\mathrm{op}}(C)$, we obtain from (19a) that

$$
|z| / 2 \geqslant|z|\|w\|-\|Y\|\|v\| \geqslant|z|(1-\|v\|)-C\|v\|
$$

since $\sqrt{1-x^{2}} \geqslant 1-x$ on $[0,1]$. Hence

$$
\|v\| \geqslant C_{18} \triangleq \frac{|z| / 2}{|z|+C}
$$

which trivially implies that

$$
\mathbb{P}\left[\left[\inf _{u \in \mathbb{S}^{N+n-1}:\|v\|<C_{18}}\|H u\| \leqslant|z| / 2\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right]=0 .
$$

Let $\mathcal{I} \in[n]$ be such that $|\mathcal{I}|=\left\lfloor\theta_{18} n\right\rfloor$, where $\theta_{18}$ will be chosen later. Let $V_{\mathcal{I}} \subset \mathbb{C}^{n}$ be the set of vectors $v$ such that $\|v\| \in\left[C_{18}, 1\right]$ and $v /\|v\| \in \mathbb{S}_{\mathcal{I}}^{n-1}$. Let $\boldsymbol{v}$ be a deterministic vector in $V_{\mathcal{I}}$, and
define the random vector $\boldsymbol{w}=-z^{-1} Y \boldsymbol{v}$. Writing $\boldsymbol{u}=\left[\boldsymbol{v}^{\boldsymbol{\top}} \boldsymbol{w}^{\top}\right]^{\top}$, we shall show that for $t>0$ small enough, there exists a constant $c>0$ such that

$$
\begin{equation*}
\mathbb{P}\left[[\|H \boldsymbol{u}\| \leqslant t] \cap \mathcal{E}_{\text {op }}(C)\right] \leqslant(c t)^{n} . \tag{20}
\end{equation*}
$$

On the event $\mathcal{E}_{\text {op }}(C)$, we have

$$
\begin{aligned}
\frac{C}{|z|} \geqslant\|\boldsymbol{w}\| \geqslant \frac{\left\|Y^{*} \boldsymbol{w}\right\|}{\|Y\|} & =\frac{\left\|Y^{*} \boldsymbol{w}+\Omega^{-L} \boldsymbol{v}-\Omega^{-L} \boldsymbol{v}\right\|}{\|Y\|} \\
& \geqslant \frac{\|\boldsymbol{v}\|-\left\|Y^{*} \boldsymbol{w}+\Omega^{-L} \boldsymbol{v}\right\|}{\|Y\|} \\
& \geqslant \frac{C_{18}-\left\|Y^{*} \boldsymbol{w}+\Omega^{-L} \boldsymbol{v}\right\|}{C} .
\end{aligned}
$$

By the choice of $\boldsymbol{w}$, we have $[\|H \boldsymbol{u}\| \leqslant t]=\left[\left\|\Omega^{-L} \boldsymbol{v}+Y^{*} \boldsymbol{w}\right\| \leqslant t\right]$. Assume that $t \leqslant C_{18} / 2$, and define the interval $\left.J=\left[C_{18} /(2 C), C /|z|\right]\right]$. Then,

$$
\left.[\|H \boldsymbol{u}\| \leqslant t] \cap \mathcal{E}_{\mathrm{op}}(C)\right] \subset\left[\|\boldsymbol{w}\| \in\left[\left(C_{18}-t\right) / C, C /|z|\right]\right] \subset[\|\boldsymbol{w}\| \in J] .
$$

If $J=\varnothing$, then the inequality (20) holds trivially. Otherwise, we write

$$
\mathbb{P}\left[[\|H \boldsymbol{u}\| \leqslant t] \cap \mathcal{E}_{\mathrm{op}}(C)\right] \leqslant \mathbb{P}\left[\left[\left\|\left(Y_{\cdot, \mathcal{I}^{\mathrm{c}}}\right)^{*} \boldsymbol{w}\right\| \leqslant t\right] \cap[\|\boldsymbol{w}\| \in J]\right],
$$

and for each $\ell \in \mathcal{I}^{\text {c }}$, we write

$$
y_{\ell}=\tilde{\boldsymbol{y}}_{\ell}+\check{\boldsymbol{y}}_{\ell} \text { where } \check{\boldsymbol{y}}_{\ell}=\mathbb{E}\left[y_{\ell} \mid Y_{\cdot, \mathcal{I}}\right] .
$$

Similarly, we write

$$
Y_{\cdot, \mathcal{I}^{\mathrm{c}}}=\tilde{Y}_{\cdot, \mathcal{I}^{\mathrm{c}}}+\check{Y}_{\cdot, \mathcal{I}^{\mathrm{c}}} \text { where } \check{Y}_{\cdot, \mathcal{I}^{\mathrm{c}}}=\mathbb{E}\left[Y_{\cdot}, \mathcal{I}^{\mathrm{c}} \mid Y_{\cdot, \mathcal{I}}\right]
$$

Since $\boldsymbol{w}$ is $\sigma\left(Y_{, \mathcal{I}}\right)$ measurable, the variables $\tilde{Y}_{, \mathcal{I}^{c}}$ and $\boldsymbol{w}$ are independent, and we get from the last inequality that

$$
\mathbb{P}\left[[\|H \boldsymbol{u}\| \leqslant t] \cap \mathcal{E}_{\mathrm{op}}(C)\right] \leqslant \sup _{w:\|w\| \in J} \mathcal{L}_{\tilde{Y}_{,}, \mathcal{I}^{\mathrm{c}}}\left(\left(\tilde{Y}_{, \mathcal{I}^{\mathrm{c}}}\right)^{*} w, t\right) .
$$

For each deterministic vector $w$ such that $\|w\| \in J$, we thus need to consider the probability law of the Gaussian vector $\left(\tilde{Y}_{., \mathcal{I}^{c}}\right)^{*} w$ by studying its covariance matrix $\Sigma=n \mathbb{E}\left(\widetilde{Y}_{., \mathcal{I}^{c}}\right)^{*} w w^{*} \tilde{Y}_{,}, \mathcal{I}^{c}$. Since $[\Sigma]_{i j}=n w^{*} \mathbb{E} y_{k} y_{\ell}^{*} w$ where $k$ (resp. $\ell$ ) is the $j^{\text {th }}$ (resp. $i^{\text {th }}$ ) element of $\mathcal{I}^{\text {c }}$, the matrix $\Sigma^{\top}$ is a $\left|\mathcal{I}^{c}\right| \times\left|\mathcal{I}^{c}\right|$ principal submatrix of the matrix

$$
n\left(I_{n} \otimes w^{*}\right) \mathbb{E} \operatorname{vec} Y(\operatorname{vec} Y)^{*}\left(I_{n} \otimes w\right)=n\left(\mathrm{~F} \otimes w^{*}\right) \mathcal{R}\left(\mathrm{F}^{*} \otimes w\right)
$$

Therefore,

$$
\|\Sigma\| \leqslant \boldsymbol{M}\|w\|^{2} \leqslant C_{\max }=\boldsymbol{M} C^{2} /|z|^{2}
$$

Recall that $y_{\ell}=\tilde{y}_{\ell}+\check{y}_{\ell}$, and observe that $\sigma\left(Y_{\cdot, \mathcal{I}}\right) \subset \sigma\left(Y_{\ell}\right)$ for each $\ell \in \mathcal{I}^{c}$. Therefore, for an arbitrary deterministic vector $u \in \mathbb{C}^{N}$, we have by Jensen's inequality

$$
\begin{aligned}
\mathbb{E}\left[\left|u^{*} \check{\boldsymbol{y}}_{\ell}\right|^{2}\right]=\mathbb{E}\left[\left|\mathbb{E}\left[u^{*} y_{\ell} \mid Y_{\cdot \mathcal{I}}\right]\right|^{2}\right] & =\mathbb{E}\left[\left|\mathbb{E}\left[\mathbb{E}\left[u^{*} y_{\ell} \mid Y_{\ell}\right] \mid Y_{\cdot, \mathcal{I}}\right]\right|^{2}\right] \\
& \leqslant \mathbb{E}\left[\left|\mathbb{E}\left[u^{*} y_{\ell} \mid Y_{\ell}\right]\right|^{2}\right] \\
& =\mathbb{E}\left[\left|u^{*} \check{y}_{\ell}\right|^{2}\right]
\end{aligned}
$$

for such $\ell$. Since

$$
\mathbb{E} y_{\ell} y_{\ell}^{*}=\mathbb{E} \tilde{y}_{\ell} \tilde{y}_{\ell}^{*}+\mathbb{E} \check{y}_{\ell} \check{y}_{\ell}^{*}=\mathbb{E} \tilde{\boldsymbol{y}}_{\ell} \tilde{\boldsymbol{y}}_{\ell}^{*}+\mathbb{E} \check{\boldsymbol{y}}_{\ell} \check{\boldsymbol{y}}_{\ell}^{*}
$$

we get from the previous inequality that

$$
\mathbb{E} \tilde{y}_{\ell} \tilde{y}_{\ell}^{*} \leqslant \mathbb{E} \tilde{\boldsymbol{y}}_{\ell} \tilde{\boldsymbol{y}}_{\ell}^{*}
$$

Choosing $\theta_{18} \leqslant 1 / 4$, we have that $\left|\mathcal{I}^{\mathrm{c}}\right| \geqslant 3 n / 4$. By Lemma 17 , there exists $\alpha>0$ such that $\left|\mathcal{I}^{\mathrm{c}} \cap \mathcal{C}_{\text {good }}(\alpha)\right| \geqslant n / 2$. With this, we have

$$
\operatorname{tr} \Sigma=n w^{*} \mathbb{E} \tilde{Y}_{., \mathcal{I}^{c}}\left(\tilde{Y}_{\cdot, \mathcal{I}^{\mathrm{c}}}\right)^{*} w=n \sum_{k \in \mathcal{I}^{\mathrm{c}}} w^{*} \mathbb{E} \tilde{\boldsymbol{y}}_{k} \tilde{\boldsymbol{y}}_{k}^{*} w \geqslant n \sum_{k \in \mathcal{I}^{\mathrm{c}} \cap \mathcal{C}_{\text {good }}} w^{*} \mathbb{E} \tilde{y}_{k} \tilde{y}_{k}^{*} w .
$$

By Proposition 10, we thus obtain that for all large $N$,

$$
\operatorname{tr} \Sigma \geqslant \frac{\alpha}{2}\left|\mathcal{I}^{\mathrm{c}} \cap \mathcal{C}_{\text {good }}\right|\|w\|^{2} \geqslant n \frac{\alpha}{4} \frac{C_{18}^{2}}{4 C^{2}} .
$$

Write $C_{\min }=\alpha C_{18}^{2} n /\left(16 C^{2}\left|\mathcal{I}^{c}\right|\right)$, and let $k$ be the largest integer in $\left[\left|\mathcal{I}^{c}\right|\right]$ such that $s_{k-1}(\Sigma) \geqslant$ $C_{\min } / 2$. Then, we obtain from Lemma 14 that

$$
k \geqslant \frac{C_{\min }}{2 C_{\max }-C_{\min }}\left|\mathcal{I}^{\mathrm{c}}\right| \geqslant \frac{C_{\min }}{2 C_{\max }-C_{\min }} \frac{3 n}{4} .
$$

With this at hand, we can deduce Inequality (20) from Lemma 15-4.
Now, still fixing $\mathcal{I}$ and $t$ as above, set $\rho_{18}$ as

$$
\rho_{18}=\frac{t}{3+C|z|^{-1}(2 C+1)} .
$$

Set $a_{18}=(|z| / 2) \wedge \rho_{18}$, and let $u=\left[v^{\top} w^{\top}\right]^{\top} \in \mathbb{S}^{N+n-1}$ be such that $v \in \mathcal{V}_{\rho_{18}}^{\mathbb{C}^{n}}\left(V_{\mathcal{I}}\right)$ and $\|H u\| \leqslant$ $a_{18}$. Let $\mathcal{K}_{\rho_{18}}$ be a $\rho_{18}-$ net of $V_{\mathcal{I}}$. By Lemma 13 , we can choose $\mathcal{K}_{\rho_{18}}$ in such a way that $\left|\mathcal{K}_{\rho_{18}}\right| \leqslant\left(3 / \rho_{18}\right)^{2 \theta_{18} n}$. By the triangle inequality, there is $\boldsymbol{v} \in \mathcal{K}_{\rho_{18}}$ such that $\|v-\boldsymbol{v}\| \leqslant 2 \rho_{18}$. Let $\boldsymbol{w}=-z^{-1} Y \boldsymbol{v}$, and write $\boldsymbol{u}=\left[\boldsymbol{v}^{\top} \boldsymbol{w}^{\top}\right]^{\top}$. Since $\|Y v+z w\| \leqslant \rho_{18}$ by (19a), we have

$$
\|Y(v-\boldsymbol{v})+z(w-\boldsymbol{w})\| \leqslant \rho_{18}
$$

Thus, on the event $\mathcal{E}_{\text {op }}(C)$,

$$
\|w-\boldsymbol{w}\| \leqslant|z|^{-1}\left(\rho_{18}+2 \rho_{18}\|Y\|\right) \leqslant|z|^{-1}\left(\rho_{18}+2 \rho_{18} C\right) .
$$

From the inequality $\left\|\Omega^{-L} v+Y^{*} w\right\| \leqslant \rho_{18}$ (see (19b)), we get

$$
\left\|\Omega^{-L} \boldsymbol{v}+Y^{*} \boldsymbol{w}+\Omega^{-L}(v-\boldsymbol{v})+Y^{*}(w-\boldsymbol{w})\right\| \leqslant \rho_{18}
$$

Thus,

$$
\left\|\Omega^{-L} \boldsymbol{v}+Y^{*} \boldsymbol{w}\right\| \leqslant\left(3+C|z|^{-1}(1+2 C)\right) \rho_{18}=t
$$

This implies that $\|H \boldsymbol{u}\| \leqslant t$. As a consequence,

$$
\left[\exists u \in \mathbb{S}^{N+n-1}: v \in \mathcal{V}_{\rho_{18}}\left(V_{\mathcal{I}}\right),\|H u\| \leqslant a_{18}\right] \subset\left[\exists \boldsymbol{v} \in \mathcal{K}_{\rho_{18}}:\left\|H\left[\boldsymbol{v}^{\top},-z^{-1}(Y \boldsymbol{v})^{\top}\right]^{\top}\right\| \leqslant t\right]
$$

Applying the union bound and using Inequality (20), we therefore get that

$$
\mathbb{P}\left[\left[\inf _{u \in \mathbb{S}^{N+n-1}:\|v\| \geqslant C_{18}, v \in \mathcal{V}_{\rho_{18}}\left(V_{\mathcal{I}}\right)}\|H u\| \leqslant a_{18}\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right] \leqslant\left(3 / \rho_{18}\right)^{2 \theta_{18} n}(c t)^{n}
$$

Now, considering all the sets $\mathcal{I} \subset[n]$ such that $|\mathcal{I}|=\left\lfloor\theta_{18} n\right\rfloor$, and using the bound $\binom{n}{m} \leqslant(e n / m)^{m}$ along with the union bound again, we get that

$$
\mathbb{P}\left[\left[_{u \in \mathbb{S}^{N+n-1}:\|v\| \geqslant C_{18}, v \in \operatorname{comp}\left(\theta_{18}, \rho_{18}\right)}\|H u\| \leqslant a_{18}\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right] \leqslant\left(e / \theta_{18}\right)^{\theta_{18} n}\left(3 / \rho_{18}\right)^{2 \theta_{18} n}(c t)^{n} .
$$

The proof is completed by choosing $\theta_{18}$ small enough.

### 3.4.2 Incompressible vectors

We now consider the action of $H$ on the vectors $u \in \mathbb{S}^{N+n-1}$ that belong to the complement of the set $\mathcal{S}$ of Proposition 18 in the unit-sphere.

Proposition 19 (incompressible vectors for the smallest singular value). There exists a constant $c_{19}>0$ such that for $\varepsilon>0$ arbitrarily small,

$$
\mathbb{P}\left[\left[\inf _{u \in \mathbb{S}^{N+n-1} \backslash \mathcal{S}}\|H u\| \leqslant n^{-3 / 2} t\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right] \leqslant \varepsilon t+\exp \left(-c_{19} \varepsilon^{2} N\right)
$$

The remainder of this subsection is devoted to the proof of this proposition. As in the proof of Proposition 18, a vector $u \in \mathbb{S}^{N+n-1}$ will be written $u=\left[v^{\top} w^{\top}\right]^{\top}$ with $v=\left[v_{0}, \ldots, v_{n-1}\right]^{\top} \in \mathbb{C}^{n}$. When $u \in \mathbb{S}^{N+n-1} \backslash \mathcal{S}$, the vector $\tilde{v}=v /\|v\|=\left[\tilde{v}_{0}, \ldots, \tilde{v}_{n-1}\right]^{\top}$ belongs now to incomp $\left(\theta_{18}, \rho_{18}\right)$. Associate with any such vector $u$, the set

$$
\mathcal{J}_{u}=\left\{k \in[n]:\left|v_{k}\right| \geqslant \rho_{18} C_{18} / \sqrt{2 n}\right\} .
$$

Then, $\left|\mathcal{J}_{u}\right|>\theta_{18} n$. Indeed, the set $\widetilde{\mathcal{J}}_{v}=\left\{k \in[n]:\left|\tilde{v}_{k}\right| \geqslant \rho_{18} / \sqrt{2 n}\right\}$ satisfies $\left|\widetilde{\mathcal{J}}_{v}\right|>\theta_{18} n$. To see this, denote as $\Pi_{\tilde{J}_{v}}$ the orthogonal projector on the vectors supported by $\widetilde{\mathcal{J}}_{v}$, and check that $\left\|\tilde{v}-\Pi_{\tilde{J}_{v}}(\tilde{v}) /\right\| \Pi_{\tilde{J}_{v}}(\tilde{v})\| \| \leqslant \rho_{18}$. If $\left|\widetilde{\mathcal{J}}_{v}\right| \leqslant \theta_{18} n$, we get a contradiction to the fact that $\tilde{v} \in$ $\operatorname{incomp}\left(\theta_{18}, \rho_{18}\right)$. It remains to check that if $u \in \mathbb{S}^{N+n-1} \backslash \mathcal{S}$, then $\tilde{\mathcal{J}}_{v} \subset \mathcal{J}_{u}$.

Now, choose $\alpha_{19}>0$ small enough so that the set $\mathcal{C}_{\text {good }}(\alpha)$ from Lemma 17 satisfies $\left|\mathcal{C}_{\text {good }}\left(\alpha_{19}\right)\right| \geqslant$ $\left(1-\theta_{18} / 2\right) n$. Observe that with this choice, $\left|\mathcal{J}_{u} \cap \mathcal{C}_{\text {good }}\left(\alpha_{19}\right)\right| \geqslant \theta_{18} n / 2$.

We now use the canvas of the proof of [26, Lem. 3.5] to relate the infimum over the incompressible vectors with the distance between column subspaces of $H$. Specifically, for each $u \in \mathbb{S}^{N+n-1} \backslash \mathcal{S}$, we have

$$
\|H u\| \geqslant \max _{k \in \mathcal{J}_{u} \cap \mathcal{C}_{\text {good }}\left(\alpha_{19}\right)}\left|v_{k}\right| \operatorname{dist}\left(h_{k}, H_{k}\right) \geqslant \frac{\rho_{18} C_{18}}{\sqrt{2 n}} \max _{k \in \mathcal{J}_{u} \cap \mathcal{C}_{\operatorname{good}}\left(\alpha_{19}\right)} \operatorname{dist}\left(h_{k}, H_{k}\right) .
$$

Thus,

$$
\begin{align*}
& \mathbb{P}\left[\left[\inf _{u \in \mathbb{S}^{N+n-1} \backslash \mathcal{S}}\|H u\| \leqslant \frac{\rho_{18} C_{18} t}{\sqrt{2} n^{3 / 2}}\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right] \\
& \leqslant \mathbb{P}\left[\left[\sum_{k \in \mathcal{J}_{u} \cap \mathcal{C}_{\operatorname{good}}\left(\alpha_{19}\right)} \operatorname{dist}\left(h_{k}, H_{k}\right) \leqslant t / n\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right] \\
& \leqslant \frac{1}{\left|\mathcal{J}_{u} \cap \mathcal{C}_{\operatorname{good}}\left(\alpha_{19}\right)\right|} \sum_{k \in \mathcal{J}_{u} \cap \mathcal{C}_{\text {good }}\left(\alpha_{19}\right)} \mathbb{P}\left[\left[\operatorname{dist}\left(h_{k}, H_{k}\right) \leqslant t / n\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right] \\
& \leqslant \frac{2}{\theta_{18} n} \sum_{k \in \mathcal{C}_{\operatorname{good}}\left(\alpha_{19}\right)} \mathbb{P}\left[\left[\operatorname{dist}\left(h_{k}, H_{k}\right) \leqslant t / n\right] \cap \mathcal{E}_{\text {op }}(C)\right] . \tag{21}
\end{align*}
$$

A workable formula for $\operatorname{dist}\left(h_{k}, H_{k}\right)$ is provided by the following lemma.
Lemma 20. Let $M \in \mathbb{C}^{n \times n}$, and partition this matrix as

$$
M=\left[\begin{array}{ll}
m_{0} & M_{-0}
\end{array}\right]=\left[\begin{array}{ll}
m_{00} & m_{01} \\
m_{10} & M_{11}
\end{array}\right]
$$

where $m_{0}$ is the first column of $M, M_{-0}$ is the submatrix that remains after extracting $m_{0}$, and $m_{00}$ is the first element of the vector $m_{0}$. Assume that $M_{11}$ is invertible. Then,

$$
\operatorname{dist}\left(m_{0}, M_{-0}\right)=\frac{\left|m_{00}-m_{01} M_{11}^{-1} m_{10}\right|}{\sqrt{1+\left\|m_{01} M_{11}^{-1}\right\|^{2}}} .
$$

Sketch of proof. Writing dist $\left(m_{0}, M_{-0}\right)^{2}=m_{0}^{*}\left(I-M_{-0}\left(M_{-0}^{*} M_{-0}\right)^{-1} M_{-0}^{*}\right) m_{0}$, and observing that $M_{-0}^{*} M_{-0}=M_{11}^{*} M_{11}+m_{01}^{*} m_{01}$, and thus, developing $\left(M_{-0}^{*} M_{-0}\right)^{-1}$ with the help of the Sherman-Morrison-Woodbury formula, we get the result after some algebra.

For convenience, let is recall the matrices $G_{k}$ and $G_{k}^{-1}$ from Section 3.1:

$$
G_{k}=\left[\begin{array}{cc}
\Omega_{k}^{-L} & Y_{k}^{*} \\
Y_{k} & z
\end{array}\right], \quad G_{k}^{-1}=\left[\begin{array}{cc}
E_{k} & F_{k} \\
P_{k} & D_{k}
\end{array}\right] .
$$

Using the above lemma, we get that

$$
\operatorname{dist}\left(h_{k}, H_{k}\right)=\frac{\left|\omega^{-k L}-y_{k}^{*} D_{k} y_{k}\right|}{\sqrt{1+\left\|\left[\begin{array}{ll}
0 & y_{k}^{*} \tag{22}
\end{array}\right] G_{k}^{-1}\right\|^{2}}}=\frac{\left|\omega^{-k L}-y_{k}^{*} D_{k} y_{k}\right|}{\sqrt{1+\left\|y_{k}^{*} P_{k}\right\|^{2}+\left\|y_{k}^{*} D_{k}\right\|^{2}}}
$$

We shall use the notation

$$
\operatorname{Num}_{k}=\left|\omega^{-k L}-y_{k}^{*} D_{k} y_{k}\right|, \quad \text { and } \quad \operatorname{Den}_{k}=\sqrt{1+\left\|y_{k}^{*} P_{k}\right\|^{2}+\left\|y_{k}^{*} D_{k}\right\|^{2}}
$$

Lemma 21. The following facts hold true.

1. Assume that $k \in \mathcal{C}_{\text {good }}\left(\alpha_{19}\right)$. Then, there exist $c_{21}, C_{21}>0$ such that

$$
\mathbb{P}\left[\left[\left\|\left[\begin{array}{ll}
0 & y_{k}^{*}
\end{array}\right] G_{k}^{-1}\right\| \leqslant C_{21}\right] \cap \mathcal{E}_{\text {op }}(C)\right] \leqslant \exp \left(-c_{21} N\right)
$$

2. On $\mathcal{E}_{\mathrm{op}}(C)$, for each $k \in[n]$, we have $\left\|P_{k}\right\|_{\mathrm{HS}} \leqslant C\left\|D_{k}\right\|_{\mathrm{HS}}$.

Proof. 1. Consider

$$
\left\|y_{k}\right\|=\left\|\left[\begin{array}{ll}
0 & y_{k}^{*}
\end{array}\right] G_{k}^{-1} G_{k}\right\| \leqslant\left\|\left[\begin{array}{ll}
0 & y_{k}^{*}
\end{array}\right] G_{k}^{-1}\right\|\left\|G_{k}\right\|
$$

On $\mathcal{E}_{\text {op }}(C)$, we have $\left\|G_{k}\right\| \leqslant C+|z|+1$. Moreover, since $k \in \mathcal{C}_{\text {good }}\left(\alpha_{19}\right)$, and since $\| n \mathbb{E} y_{k} y_{k}^{*}-$ $S\left(e^{2 \imath \pi k / n}\right) \|$ is small by Lemma 9 , there exist two constants $c^{\prime}, C^{\prime}>0$ such that

$$
\mathbb{P}\left[\left\|y_{k}\right\| \leqslant C^{\prime}\right] \leqslant \exp \left(-c^{\prime} N\right)
$$

by Lemma $15-2$. The proof of Item 1 is complete.
2. We show that for each constant non-zero vector $u \in \mathbb{C}^{N}$,

$$
\left\|u^{*} P_{k}\right\| \leqslant C\left\|u^{*} D_{k}\right\| .
$$

The result follows then from the identity $\|M\|_{\text {HS }}^{2}=\sum_{k}\left\|e_{k}^{*} M\right\|^{2}$, valid for each matrix $M$. The vectors $v^{*}=u^{*} P_{k}$ and $w^{*}=u^{*} D_{k}$ satisfy

$$
\left[\begin{array}{ll}
v^{*} & w^{*}
\end{array}\right]=\left[\begin{array}{ll}
0 & u^{*}
\end{array}\right] G_{k}^{-1}
$$

and hence,

$$
\left[\begin{array}{ll}
0 & u^{*}
\end{array}\right]=\left[\begin{array}{ll}
v^{*} & w^{*}
\end{array}\right]\left[\begin{array}{cc}
\Omega_{k}^{-L} & Y_{k}^{*} \\
Y_{k} & z
\end{array}\right]
$$

In particular, $v^{*} \Omega_{k}^{-L}+w^{*} Y_{k}=0$, and this shows that $\|v\| \leqslant C\|w\|$ on $\mathcal{E}_{\mathrm{op}}(C)$.

Lemma 21 will be used to control the value of $\operatorname{Den}_{k}$. Fix an arbitrary small real number $\varepsilon>0$, and define the event

$$
\mathcal{E}_{\mathrm{Den}}(\varepsilon)=\left[\operatorname{Den}_{k}^{2} \geqslant \varepsilon\left\|D_{k}\right\|_{\mathrm{HS}}^{2}\right]
$$

Lemma 22 (denominator not too large). Assume that $k \in \mathcal{C}_{\text {good }}\left(\alpha_{19}\right)$. Then, there exists $c_{22}>0$ such that for all large $n, \mathbb{P}\left[\mathcal{E}_{\text {Den }}(\varepsilon) \cap \mathcal{E}_{\text {op }}(C)\right] \leqslant \exp \left(-c_{22} \varepsilon n\right)$.
Proof. Defining the event $\mathcal{E}=\left[\operatorname{Den}_{k}^{2} \geqslant\left(C_{21}^{-2}+1\right)\left\|\left[\begin{array}{ll}0 & y_{k}^{*}\end{array}\right] G_{k}^{-1}\right\|^{2}\right]$, we have

$$
\mathbb{P}\left[\mathcal{E} \cap \mathcal{E}_{\mathrm{op}}(C)\right]=\mathbb{P}\left[\left[C_{21}^{2} \geqslant\left\|\left[\begin{array}{ll}
0 & y_{k}^{*}
\end{array}\right] G_{k}^{-1}\right\|^{2}\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right] \leqslant \exp \left(-c_{21} n\right)
$$

by Lemma 21-1. Writing $\mathcal{E}^{\prime}(\varepsilon)=\left[\left\|\left[\begin{array}{ll}0 & y_{k}^{*}\end{array}\right] G_{k}^{-1}\right\|^{2} \geqslant \varepsilon\left\|D_{k}\right\|_{\mathrm{HS}}^{2}\right]$, we have

$$
\begin{aligned}
\mathbb{P}\left[\mathcal{E}^{\prime}(\varepsilon) \cap \mathcal{E}_{\mathrm{op}}(C)\right] & =\mathbb{P}\left[\left[\left\|y_{k}^{*} P_{k}\right\|^{2}+\left\|y_{k}^{*} D_{k}\right\|^{2} \geqslant \varepsilon\left\|D_{k}\right\|_{\mathrm{HS}}^{2}\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right] \\
& \leqslant \mathbb{P}\left[\left\|y_{k}^{*} P_{k}\right\|^{2} \geqslant \varepsilon\left\|P_{k}\right\|_{\mathrm{HS}}^{2} /\left(2 C^{2}\right)\right]+\mathbb{P}\left[\left\|y_{k}^{*} D_{k}\right\|^{2} \geqslant \varepsilon\left\|D_{k}\right\|_{\mathrm{HS}}^{2} / 2\right]
\end{aligned}
$$

where we used Lemma 21-2. in the first inequality.
Note that

$$
\left\|y_{k}^{*} D_{k}\right\|^{2} \leqslant 2\left\|\tilde{y}_{k}^{*} D_{k}\right\|^{2}+2\left\|\check{y}_{k}^{*} D_{k}\right\|^{2} \text { and }\left\|\check{y}_{k}^{*} D_{k}\right\| \leqslant\left\|\check{y}_{k}\right\|^{2}\left\|D_{k}\right\|_{\mathrm{HS}}
$$

Hence there exists $c_{N} \rightarrow \infty$ such that

$$
\begin{aligned}
\mathbb{P}\left[\left\|y_{k}^{*} D_{k}\right\|^{2} \geqslant \varepsilon\left\|D_{k}\right\|_{\mathrm{HS}}^{2} / 2\right] & \leqslant \mathbb{P}\left[\left\|\tilde{y}_{k}^{*} D_{k}\right\|^{2} \geqslant \boldsymbol{\varepsilon}\left\|D_{k}\right\|_{\mathrm{HS}}^{2} / 8\right]+\mathbb{P}\left[\left\|\check{y}_{k}^{*}\right\|^{2} \geqslant \boldsymbol{\varepsilon} / 8\right] \\
& \leqslant \mathbb{P}\left[\left\|\tilde{y}_{k}^{*} D_{k}\right\|^{2} \geqslant \boldsymbol{\varepsilon}\left\|S_{k}^{1 / 2} D_{k}\right\|_{\mathrm{HS}}^{2} /(8 \boldsymbol{M})\right]+e^{-c_{N} n} \\
& \leqslant e^{1-\varepsilon n /(16 \boldsymbol{M})}+e^{-c_{N} n}
\end{aligned}
$$

where we used Lemma 15-1 along with Proposition 10 in the second inequality, and Lemma 15-5 in the third one. We have a similar bound for $\mathbb{P}\left[\left\|y_{k}^{*} P_{k}\right\|^{2} \geqslant \varepsilon\left\|P_{k}\right\|_{\mathrm{HS}}^{2} /\left(2 C^{2}\right)\right]$. Consequently, there exists $c>0$ such that $\mathbb{P}\left[\mathcal{E}^{\prime}(\varepsilon) \cap \mathcal{E}_{\text {op }}(C)\right] \leqslant e^{-c \varepsilon n}$ for all large $N$.

Observing that $\mathcal{E}_{\text {Den }}(\varepsilon) \cap \mathcal{E}^{\mathrm{c}} \subset \mathcal{E}^{\prime}\left(\varepsilon /\left(C_{21}^{-2}+1\right)\right)$, we obtain as a consequence of these inequalities that there exists $c_{22}>0$ such that for all large $N$,

$$
\begin{aligned}
\mathbb{P}\left[\mathcal{E}_{\text {Den }}(\varepsilon) \cap \mathcal{E}_{\text {op }}(C)\right] & \leqslant \mathbb{P}\left[\mathcal{E}_{\text {Den }}(\varepsilon) \cap \mathcal{E}^{\mathrm{c}} \cap \mathcal{E}_{\mathrm{op}}(C)\right]+\mathbb{P}\left[\mathcal{E} \cap \mathcal{E}_{\mathrm{op}}(C)\right] \\
& \leqslant \mathbb{P}\left[\mathcal{E}^{\prime}\left(\varepsilon /\left(C_{21}^{-2}+1\right)\right) \cap \mathcal{E}_{\mathrm{op}}(C)\right]+\mathbb{P}\left[\mathcal{E} \cap \mathcal{E}_{\mathrm{op}}(C)\right] \\
& \leqslant \exp \left(-c_{22} \varepsilon n\right)
\end{aligned}
$$

The proof of the lemma is now complete.
Now we get back to the expression in (21). As in Subsection 3.1, we use the shorthand notation $S_{k}=S\left(e^{2 \pi \pi k / n}\right)$. Given $k \in \mathcal{C}_{\text {good }}\left(\alpha_{19}\right)$, we have

$$
\begin{aligned}
& \mathbb{P}\left[\left[\operatorname{dist}\left(h_{k}, H_{k}\right) \leqslant t / n\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right]=\mathbb{P}\left[\left[\mathrm{Num}_{k} / \operatorname{Den}_{k} \leqslant t / n\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right] \\
& \leqslant \mathbb{P}\left[\left[\frac{\mathrm{Num}_{k}}{\operatorname{Den}_{k}} \leqslant t / n\right] \cap \mathcal{E}_{\text {Den }}(\varepsilon)^{\mathrm{c}} \cap \mathcal{E}_{\mathrm{op}}(C)\right]+\mathbb{P}\left[\mathcal{E}_{\text {Den }}(\varepsilon) \cap \mathcal{E}_{\mathrm{op}}(C)\right] \\
& \leqslant \mathbb{P}\left[\frac{n \mathrm{Num}_{k}}{\left\|D_{k}\right\|_{\mathrm{HS}}} \leqslant t \sqrt{\varepsilon}\right]+\exp \left(-c_{22} \varepsilon n\right) \\
& \leqslant \mathbb{P}\left[\frac{n \mathrm{Num}_{k}}{\left\|S_{k}^{1 / 2} D_{k} S_{k}^{1 / 2}\right\|_{\mathrm{HS}}} \leqslant t \sqrt{\varepsilon} / \alpha_{19}\right]+\exp \left(-c_{22} \varepsilon n\right) .
\end{aligned}
$$

Writing $\operatorname{Num}_{k}=\left|\omega^{-k L}+\left(\tilde{y}_{k}+\check{y}_{k}\right)^{*} D_{k}\left(\tilde{y}_{k}+\check{y}_{k}\right)\right|$, we obtain from Lemma 15-3 and Proposition 10 that for some constant $c>0$,

$$
\begin{aligned}
\mathbb{P}\left[\left[\operatorname{dist}\left(h_{k}, H_{k}\right) \leqslant t / n\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right] & \leqslant \mathbb{E}_{Y_{k}} \mathcal{L}_{\tilde{y}_{k}}\left(\frac{\left(\tilde{y}_{k}+\check{y}_{k}\right)^{*} D_{k}\left(\tilde{y}_{k}+\check{y}_{k}\right)}{\left\|S_{k}^{1 / 2} D_{k} S_{k}^{1 / 2}\right\|_{\mathrm{HS}}}, \frac{t \sqrt{\varepsilon}}{\alpha_{19}}\right)+\exp \left(-c_{22} \varepsilon n\right) \\
& \leqslant c t \sqrt{\varepsilon}+\exp \left(-c_{22} \varepsilon n\right) .
\end{aligned}
$$

Proposition 19 is obtained by combining the last inequality with Inequality (21).

## Theorem 1: End of proof.

The proof is now completed by combining Proposition 18, Proposition 19, and Lemma 16.

### 3.5 Proof of Theorem 2

Throughout the proof, we fix $k \in\left[\left\lfloor N^{\beta}\right\rfloor,\lfloor N / 2\rfloor\right]$. None of the constants that will appear in the course of the proof will depend on $k$.

The following lemma refines Inequality (7).
Lemma 23. For any $k \in[N]$, we have

$$
s_{N+n-k-1}\left(H_{\cdot,[N+n-k]}\right) \leqslant s_{N-k-1}\left(Y \Omega^{L} Y^{*}-z\right) .
$$

Proof. Write any vector $u \in \mathbb{C}^{N+n-k}$ as $u=\left[v^{\top} w^{\top}\right]^{\top}$ with $w \in \mathbb{C}^{N-k}$. Writing

$$
H_{\cdot,[N+n-k]} u=\left[\begin{array}{cc}
\Omega^{-L} & \left(Y_{[N-k], \cdot}\right)^{*} \\
Y & z\left(I_{N}\right)_{\cdot,[N-k]}
\end{array}\right]\left[\begin{array}{c}
v \\
w
\end{array}\right]=\left[\begin{array}{l}
\Omega^{-L} v+\left(Y_{[N-k], .}\right)^{*} w \\
Y v+z\left(I_{N}\right)_{\cdot,[N-k]} w
\end{array}\right],
$$

we have

$$
\begin{aligned}
s_{N+n-k-1}\left(H_{\cdot,[N+n-k]}\right)=\min _{u \in \mathbb{S}^{N+n-k-1}}\left\|H_{\cdot,[N+n-k]} u\right\| & \leqslant \min _{u \in \mathbb{S}^{N+n-k-1}: v=-\Omega^{L}\left(Y_{[N-k], \cdot}\right) * w}\left\|H_{\cdot,[N+n-k]} u\right\| \\
& \leqslant \min _{w \in \mathbb{S}^{N-k-1}}\left\|\left(z I_{N}-Y \Omega^{L} Y^{*}\right)_{\cdot,[N-k]} w\right\| .
\end{aligned}
$$

On the other hand, using the variational characterization of the singular values of a matrix, see [15, Th. 4.3.15], we can write

$$
\min _{w \in \mathbb{S}^{N-k-1}}\left\|\left(Y \Omega^{L} Y^{*}-z I_{N}\right)_{\cdot,[N-k]} w\right\| \leqslant s_{N-k-1}\left(Y \Omega^{L} Y^{*}-z I_{N}\right)
$$

and hence the result follows.
Similar to what we did for controlling the smallest singular value of $H$, we use the characterization $s_{N+n-k-1}\left(H_{\cdot,[N+n-k]}\right)=\min _{u \in \mathbb{S}^{N+n-k-1}}\left\|H_{\cdot,[N+n-k]} u\right\|$, and we partition the set $\mathbb{S}^{N+n-k-1}$ into sets of compressible and incompressible vectors. So, write any vector $u \in \mathbb{S}^{N+n-k-1}$ as $u=\left[v^{\top} w^{\top}\right]^{\top}$ with $w \in \mathbb{C}^{N-k}$. Let the set $\mathcal{S} \subset \mathbb{S}^{N+n-k-1}$ be as in the statement of Proposition 18 above. We of course have

$$
\begin{equation*}
s_{N+n-k-1}\left(H_{\cdot,[N+n-k]}\right)=\inf _{u \in \mathcal{S}}\left\|H_{\cdot,[N+n-k]} u\right\| \wedge_{u \in \mathbb{S}^{N}+n-k-1} \inf _{\mathcal{S}}\left\|H_{\cdot,[N+n-k]} u\right\| . \tag{23}
\end{equation*}
$$

It can be readily checked that Proposition 18 remains true, once we change the values of the constants that appear in its statement as needed. So, finally, the contribution of the term $\inf _{u \in \mathcal{S}}\left\|H_{\cdot,[N+n-k]} u\right\|$ has been estimated, and we are left with the incompressible vectors. Regarding these, we have the following proposition.

Proposition 24 (incompressible vectors for the small singular values). There exist constants $a_{24}>0$ and $c_{24}>0$ such that for all $N$ larger than an integer that is independent of $k \in$ $\left[\left\lfloor N^{\beta}\right\rfloor,\lfloor N / 2\rfloor\right]$,

$$
\mathbb{P}\left[\left[\inf _{u \in \mathbb{S}^{N+n-k-1} \backslash \mathcal{S}}\left\|H_{\cdot,[N+n-k]} u\right\| \leqslant a_{24} \sqrt{k} / N\right] \cap \mathcal{E}_{\text {op }}(C)\right] \leqslant \exp \left(-c_{24} k\right)
$$

Remember the definition of the set $\mathcal{C}_{\text {good }}(\alpha) \subset[n]$ in Lemma 17 above. To prove this proposition, we begin by mimicking the argument that lead to Inequality (21) above. Specifically, there exist positive constants $c, c^{\prime}$ and $\alpha$ such that

$$
\begin{align*}
\mathbb{P}\left[\left[\inf _{u \in \mathbb{S}^{N+n-k-1} \backslash \mathcal{S}}\left\|H_{\cdot,[N+n-k]} u\right\|\right.\right. & \left.\left.\leqslant \frac{c t}{n^{1 / 2}}\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right] \\
& \leqslant \frac{c^{\prime}}{n} \sum_{\ell \in \mathcal{C}_{\mathrm{good}}(\alpha)} \mathbb{P}\left[\left[\operatorname{dist}\left(h_{\ell}, H_{\cdot,[N+n-k] \backslash\{\ell\}}\right) \leqslant t\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right], \tag{24}
\end{align*}
$$

and $\left|\mathcal{C}_{\text {good }}(\alpha)\right|$ is of order $N$.
This inequality calls for the following remark.
Remark 1. By checking the structure of the vector $h_{\ell}$ and the matrix $H_{\cdot,[N+n-k] \backslash\{\ell\}}$ (see below), one can intuitively infer that with high probability, $\operatorname{dist}\left(h_{\ell}, H_{\cdot,[N+n-k] \backslash\{\ell\}}\right)^{2}$ is of order $\operatorname{codim}\left(H_{\cdot,[N+n-k] \backslash\{\ell\}}\right) / N=(k-1) / N$. Taking $t=\sqrt{k / N}$ and recalling the characterization (23), we get that $s_{N+n-k-1} \gtrsim \sqrt{k} / N$, which is what is predicted by Theorem 2. However, one naturally expects that these singular values grow linearly in $k$ (as $k / N$ ) which would make the result of this theorem sub-optimal. Actually, this sub-optimality is due in the first place to Inequality (24), which is too conservative for obtaining optimal bounds on the small singular values that we are dealing with here.

In [27], a tighter inequality is used to control the smallest singular value of a rectangular matrix. Obtaining a corresponding inequality appears to be a quite involved task in the present context.

Proof of Proposition 24. We need to bound the summands at the right hand side of Inequality (24). To this end, assume for notational simplicity that $\ell=0$ (assuming without loss of generality that $\left.0 \in \mathcal{C}_{\text {good }}(\alpha)\right)$. In this case,

$$
h_{0}=\left[\begin{array}{c}
1 \\
0_{n-1} \\
y_{0}
\end{array}\right]
$$

and

$$
H_{\cdot,[N+n-k] \backslash\{0\}}=\left[\begin{array}{cc}
0_{n-1}^{\top} & \left(y_{0}\right)_{[N-k]}^{*} \\
\left(\Omega^{-L}\right)_{[1: n],[1: n]} & \left(Y_{[N-k],[1: n]}\right)^{*} \\
Y_{\cdot,[1: n]} & {\left[\begin{array}{c}
z I_{N-k} \\
0_{k \times(N-k)}
\end{array}\right]}
\end{array}\right] \in \mathbb{C}^{(N+n) \times(N+n-k-1)} .
$$

Write

$$
\underline{h}_{0}=\left[\begin{array}{c}
0_{n-1} \\
y_{0}
\end{array}\right] \in \mathbb{C}^{N+n-1}
$$

and let $\underline{H}_{0}$ be the matrix obtained by taking out the row 0 from $H_{\cdot,[N+n-k] \backslash\{0\}}$, i.e.,

$$
\underline{H}_{0}=\left[\begin{array}{cc}
\left(\Omega^{-L}\right)_{[1: n],[1: n]} & \left(Y_{[N-k],[1: n]}\right)^{*} \\
Y_{\cdot,[1: n]} & {\left[\begin{array}{c}
z I_{N-k} \\
0_{k \times(N-k)}
\end{array}\right]}
\end{array}\right] \in \mathbb{C}^{(N+n-1) \times(N+n-k-1)} .
$$

Note that $\underline{h}_{0}$ and $\underline{H}_{0}$ are "almost" independent. However, this is not so for $h_{0}$ and $H_{\text {,, }[N+n-k] \backslash\{0\}}$. Let $a \in \mathbb{C}^{N+n-k-1}$ be the vector such that $H_{\cdot,[N+n-k] \backslash\{0\}} a=\Pi_{H_{\cdot,[N+n-k] \backslash 0\}}} h_{0}$. Then,

$$
\operatorname{dist}\left(h_{0}, H_{\cdot,[N+n-k] \backslash\{0\}}\right)^{2}=\left\|h_{0}-H_{\cdot,[N+n-k] \backslash\{0\}} a\right\|^{2} \geqslant\left\|\underline{h}_{0}-\underline{H}_{0} a\right\|^{2} \geqslant \operatorname{dist}\left(\underline{h}_{0}, \underline{H}_{0}\right)^{2} .
$$

Consider again the decomposition $y_{0}=\tilde{y}_{0}+\check{y}_{0}$ with $\check{y}_{0}=\mathbb{E}\left[y_{0} \mid Y_{\cdot,[1: n]}\right]$, and write

$$
\tilde{h}_{0}=\left[\begin{array}{c}
0_{n-1} \\
\tilde{y}_{0}
\end{array}\right], \quad \text { and } \quad \check{h}_{0}=\left[\begin{array}{c}
0_{n-1} \\
\check{y}_{0}
\end{array}\right] \quad \in \mathbb{C}^{N+n-1}
$$

Then, writing $G_{0}=\left[\begin{array}{ll}\check{h}_{0} & \underline{H}_{0}\end{array}\right]$, we have

$$
\operatorname{dist}\left(\underline{h}_{0}, \underline{H}_{0}\right)=\operatorname{dist}\left(\tilde{h}_{0}+\check{h}_{0}, \underline{H}_{0}\right) \geqslant \operatorname{dist}\left(\tilde{h}_{0}+\check{h}_{0}, G_{0}\right)=\operatorname{dist}\left(\tilde{h}_{0}, G_{0}\right) .
$$

Observe that $\tilde{h}_{0}$ and $G_{k}$ are independent, and that $N+n-k-1 \leqslant \operatorname{rank}\left(G_{0}\right) \leqslant N+n-k$ with probability one. Let $r=N+n-1-\operatorname{rank}\left(G_{0}\right)$, and let $A=\left[V^{\top} W^{\top}\right]^{\top} \in \mathbb{C}^{(N+n-1) \times r}$ be an isometry matrix such that $A A^{*}=\Pi_{G_{0}}^{\perp}$. Here the partitioning of $A$ is such that $V \in \mathbb{C}^{(n-1) \times r}$, and $W \in \mathbb{C}^{N \times r}$. Note that $r \in\{k-1, k\}$ w.p. 1 . Since $A A^{*}=\Pi_{G_{0}}^{\perp}$, it holds that $\underline{H}_{0}^{*} A=0$, which can be developed as

$$
\begin{align*}
& 0=\left(\Omega^{L}\right)_{[1: n],[1: n]} V+\left(Y_{,[1: n]}\right)^{*} W  \tag{25}\\
& 0=Y_{[N-k],[1: n]} U+\left[\begin{array}{cc}
\bar{z} I_{N-k} & 0_{(N-k) \times k}
\end{array}\right] W
\end{align*}
$$

Assume that $\|W\|_{\text {HS }}^{2}<a r$ with $a=/\left(1+C^{2}\right)$. By Equation (25), on the event $\mathcal{E}_{\text {op }}(C)$, we have $\|V\|_{\mathrm{HS}}^{2} \leqslant C^{2}$ ar, which contradicts the identity $\|V\|_{\mathrm{HS}}^{2}+\|W\|_{\mathrm{HS}}^{2}=r$. Thus,

$$
\mathbb{P}\left[\left[\|W\|_{\mathrm{HS}}^{2}<a r\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right]=0
$$

Using this result, we now provide a control over $\operatorname{dist}\left(\tilde{h}_{0}, G_{0}\right)=\left\|W^{*} \tilde{y}_{0}\right\|$. It is clear that $\tilde{y}_{0}$ and $W$ are independent. On the event $\left[\|W\|_{\mathrm{HS}}^{2} \geqslant a r\right]$, we know by Lemma 14 that

$$
s_{\left\lfloor\frac{a}{2-a} r\right\rfloor-1}(W)^{2} \geqslant a / 2
$$

Moreover, since $0 \in \mathcal{C}_{\text {good }}(\alpha)$, Proposition 10, with the help of Assumptions 1 and 3, shows that $n \mathbb{E} \tilde{y}_{0} \tilde{y}_{0}^{*} \geqslant(\alpha / 2) I_{N}$ for all $N$ large. Consequently, on the event $\left[\|W\|_{\mathrm{HS}}^{2} \geqslant a r\right]$, the conditional distribution of $W^{*} \tilde{y}_{k}$ with respect to the sigma-field $\sigma(W)$ is Gaussian with a covariance matrix $\Sigma \in \mathcal{H}_{+}^{r}$ such that $s_{\left\lfloor\frac{a}{2-a} r\right\rfloor-1}(\Sigma) \geqslant a \alpha /(4 n)$ for all large $N$. Put $t^{2}=a \alpha\left\lfloor\frac{a}{2-a} r\right\rfloor /(8 n) \sim r / n$. By what precedes and by Lemma 15-2, we have

$$
\begin{aligned}
& \mathbb{P}\left[\left[\operatorname{dist}\left(h_{0}, H_{\cdot,[N+n-k] \backslash\{0\}}\right)^{2} \leqslant t^{2}\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right] \\
& \leqslant \mathbb{P}\left[\left[\operatorname{dist}\left(\tilde{h}_{0}, G_{0}\right)^{2} \leqslant t^{2}\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right] \\
& \leqslant \mathbb{P}\left[\left[\operatorname{dist}\left(\tilde{h}_{0}, G_{0}\right)^{2} \leqslant t^{2}\right] \cap\left[\|W\|_{\mathrm{HS}}^{2} \geqslant a r\right]\right]+\mathbb{P}\left[\left[\|W\|_{\mathrm{HS}}^{2}<a r\right] \cap \mathcal{E}_{\mathrm{op}}(C)\right] \\
& \left.\leqslant \exp \left(-c_{15,2} \frac{a}{2-a} r\right\rfloor\right) .
\end{aligned}
$$

The proof of Proposition 24 is completed by using Inequality (24).

## Proof of Theorem 2

To complete the proof, just combine Lemma 23 with the characterization (23), apply Propositions 18 and 24 respectively to the two terms in this characterization, and use Lemma 16 to bound $\mathbb{P}\left[\mathcal{E}_{\text {op }}(C)\right]$.

## 4 Proofs regarding the singular value distribution of $\widehat{R}_{L}^{(N)}-z$

### 4.1 Proof of Theorem 5

Lemma 25. Given an integer $m>0$, let $M \in \mathcal{M}_{+}^{m}$. Then $M$ is invertible, and $\left\|M^{-1}\right\| \leqslant$ $\left\|(\Im M)^{-1}\right\|$. Moreover, $-M^{-1} \in \mathcal{M}_{+}^{m}$.

Proof. For each non-zero vector $u$, we have

$$
\|u\|\|M u\| \geqslant\left|u^{*} M u\right|=\left|u^{*} \Re M u+\imath u^{*} \Im M u\right| \geqslant u^{*} \Im M u>0
$$

and hence $M$ is invertible. For an arbitrary non-zero vector $v$, there is $u \neq 0$ such that $v=M u$, and we have from the former display that

$$
\left\|M^{-1} v\right\|\|v\| \geqslant\left(M^{-1} v\right)^{*}(\Im M) M^{-1} v \geqslant\left\|(\Im M)^{-1}\right\|^{-1}\left\|M^{-1} v\right\|^{2}
$$

and hence the second result follows. We finally have

$$
\Im\left(M^{-1}\right)=M^{-1} \frac{M^{*}-M}{2 \imath} M^{-*}=-M^{-1}(\Im M) M^{-*}<0 .
$$

The following result is well-known. We provide its proof for completeness.
Lemma 26. A probability measure $\check{\nu}$ is symmetric if and only if its Stieltjes transform $g_{\check{\nu}}$, seen as an analytic function on $\mathbb{C} \backslash \mathbb{R}$, satisfies $g_{\check{\nu}}(-\eta)=-g_{\check{\nu}}(\eta)$.

Proof. The necessity is obvious from the definition of the Stieltjes transform and from the fact that $\check{\nu}(d \lambda)=\check{\nu}(-d \lambda)$. To prove the sufficiency, we use the Perron inversion formula, which says that for any function $\varphi \in C_{c}(\mathbb{R})$,

$$
\int_{\mathbb{R}} \varphi(x) \check{\nu}(d x)=\lim _{\varepsilon \downarrow 0} \frac{1}{\pi} \int_{\mathbb{R}} \varphi(x) \Im g_{\check{\nu}}(x+\imath \varepsilon) d x .
$$

By a simple variable change at the right hand side, and by using the equalities $g_{\check{\nu}}(-\eta)=-g_{\check{\nu}}(\eta)$ and $g_{\check{\nu}}(\bar{\eta})=\bar{g}_{\check{\nu}}(\eta)$, we obtain that $\int \varphi(x) \check{\nu}(d x)=\int \varphi(-x) \check{\nu}(d x)$, as desired.

The operator $\mathscr{T}$ introduced before Theorem 5 has the following properties.
Lemma 27. Suppose $M \in \mathbb{C}^{2 N \times 2 N}$ and $S \in \mathcal{H}_{+}^{N}$. Then,

$$
\left\|\mathscr{T}\left(\left(I_{2} \otimes S\right) M\right)\right\| \leqslant\|M\| \operatorname{tr} S / n \leqslant(N / n)\|S\|\|M\|
$$

If $\Im M \geqslant \rho I_{2 N}$ for some $\rho>0$, then

$$
\Im \mathscr{T}\left(\left(I_{2} \otimes S\right) M\right) \geqslant \rho n^{-1}(\operatorname{tr} S) I_{2} .
$$

Proof. To obtain the first result, we write

$$
\mathscr{T}\left(\left(I_{2} \otimes S\right) M\right)=\mathscr{T}\left(\left(I_{2} \otimes S^{1 / 2}\right) M\left(I_{2} \otimes S^{1 / 2}\right)\right)=\frac{1}{n} \sum_{\ell \in[N]}\left[\begin{array}{lll}
e_{\ell}^{*} S^{1 / 2} & \\
& e_{\ell}^{*} S^{1 / 2}
\end{array}\right] M\left[\begin{array}{ll}
S^{1 / 2} e_{\ell} & \\
& S^{1 / 2} e_{\ell}
\end{array}\right]
$$

Hence,

$$
\left\|\mathscr{T}\left(\left(I_{2} \otimes S\right) M\right)\right\| \leqslant n^{-1}\|M\| \sum_{\ell \in[N]}\left\|S^{1 / 2} e_{\ell}\right\|^{2}=\|M\| \operatorname{tr} S / n
$$

To obtain the second inequality, just observe that $\Im \mathscr{T}\left(\left(I_{2} \otimes S\right) M\right)=\mathscr{T}\left(\left(I_{2} \otimes S\right) \Im M\right)$ and follow the same derivation as above.

## Proof of Theorem 5

Given a function $M(\eta) \in \mathcal{S}^{2 N}$, let

$$
A\left(M(\eta), e^{\imath \theta}\right)=\mathscr{T}\left(\left(I_{2} \otimes \Sigma\left(e^{\imath \theta}\right)\right) M(\eta)\right)+U_{L}\left(e^{\imath \theta}\right)
$$

We shall show that $A\left(M(\eta), e^{2 \theta}\right)^{-1}$ exists, is holomorphic as a function of $\eta$, is continuous as a function of $\theta$, and satisfies $\Im A\left(M(\eta), e^{\imath \theta}\right)^{-1} \leqslant 0$.

First let $\rho>0$ be such that $\Im M(\eta) \geqslant \rho I$. By Lemma 27, $\Im A\left(M(\eta), e^{\imath \theta}\right) \geqslant \rho n^{-1} \operatorname{tr} \Sigma\left(e^{\imath \theta}\right) I_{2}$. Thus, if $\Sigma\left(e^{\imath \theta}\right) \neq 0$, then $A\left(M(\eta), e^{\imath \theta}\right)^{-1}$ exists, is holomorphic in $\eta$, and $\Im A\left(M(\eta), e^{2 \theta}\right)^{-1}<0$ by Lemma 25.

Otherwise, $A\left(M(\eta), e^{\imath \theta}\right)=U_{L}\left(e^{\imath \theta}\right)$, which is trivially invertible, and $\Im A\left(M(\eta), e^{\imath \theta}\right)^{-1}=0$. In summary, $A\left(M(\eta), e^{\imath \theta}\right)^{-1}$ is holomorphic in $\eta$ and satisfies $\Im A\left(M(\eta), e^{\imath \theta}\right)^{-1} \leqslant 0$, for each $\theta \in[0,2 \pi)$. Moreover, the continuity of $A\left(M(\eta), e^{\imath \cdot}\right)^{-1}$ follows from the continuity of $A\left(M(\eta), e^{\ell \cdot}\right)$.

From these properties of $A$, it follows that

$$
B(M(\eta))=(2 \pi)^{-1} \int_{0}^{2 \pi} A\left(M(\eta), e^{\imath \theta}\right)^{-1} \otimes \Sigma\left(e^{\imath \theta}\right) d \theta
$$

exists, is holomorphic as a function of $\eta$, and satisfies $\Im B(M(\eta)) \leqslant 0$. Since $\Im\left[\begin{array}{ll}\eta & z \\ \bar{z} & \eta\end{array}\right]=\Im \eta I_{2}$, we get by Lemma 25 that the function $\mathcal{F}_{\Sigma, z}(M(\eta), \eta)$ is holomorphic in $\eta \in \mathbb{C}_{+}$, and takes values in $\mathcal{M}_{+}^{2 N}$. Furthermore, since $M \in \mathcal{S}^{2 N}$, it holds that $\|M(\eta)\| \leqslant 1 /(\Im \eta)$, and it is easy to show that $\lim _{t \rightarrow \infty} \imath t \mathcal{F}_{\Sigma, z}(M, \imath t)=-I_{2 N}$. In summary, $\mathcal{F}_{\Sigma, z}(M(\eta), \eta) \in \mathcal{S}^{2 N}$ as a function of $\eta$ when $M \in \mathcal{S}^{2 N}$.

Let us now establish the uniqueness of the solution of Equation (2) in the class $\mathcal{S}^{2 N}$. Assume that $P(z, \cdot)$ and $P^{\prime}(z, \cdot)$ are two such solutions. Then,

$$
\begin{aligned}
& \mathcal{F}_{\Sigma, z}(P, \eta)-\mathcal{F}_{\Sigma, z}\left(P^{\prime}, \eta\right) \\
& =\mathcal{F}_{\Sigma, z}(P, \eta)\left(B\left(P^{\prime}\right)-B(P)\right) \mathcal{F}_{\Sigma, z}\left(P^{\prime}, \eta\right) \\
& =\mathcal{F}_{\Sigma, z}(P, \eta)\left\{\frac{1}{2 \pi} \int_{0}^{2 \pi}\left(A\left(P^{\prime}, e^{\imath \theta}\right)^{-1} \mathscr{T}\left(\left(I_{2} \otimes \Sigma\left(e^{\imath \theta}\right)\right)\left(P-P^{\prime}\right)\right) A\left(P, e^{\imath \theta}\right)^{-1}\right) \otimes \Sigma\left(e^{\imath \theta}\right) d \theta\right\} \mathcal{F}_{\Sigma, z}\left(P^{\prime}, \eta\right) .
\end{aligned}
$$

Define the domain

$$
\mathcal{D}=\left\{\eta \in \mathbb{C}_{+}: \Im \eta>4\|\Sigma\|_{\infty}\left(\frac{N}{n} \vee \sqrt{\frac{N}{n}}\right)\right.
$$

Let $\eta \in \mathcal{D}$. It can be checked using Inequality $\|P\|,\left\|P^{\prime}\right\| \leqslant 1 / \Im \eta$ along with Lemma 27 that $\left\|A\left(P, e^{\imath \theta}\right)^{-1}\right\|,\left\|A\left(P^{\prime}, e^{\imath \theta}\right)^{-1}\right\| \leqslant 2$ when $\eta \in \mathcal{D}$. Using Lemma 27 again, we have

$$
\begin{equation*}
\left\|P-P^{\prime}\right\|=\left\|\mathcal{F}_{\Sigma, z}(P, \eta)-\mathcal{F}_{\Sigma, z}\left(P^{\prime}, \eta\right)\right\| \leqslant \frac{4}{(\Im \eta)^{2}} \frac{N}{n}\|\Sigma\|_{\infty \infty}^{2}\left\|P-P^{\prime}\right\| \leqslant \frac{1}{2}\left\|P-P^{\prime}\right\| \tag{26}
\end{equation*}
$$

which shows that $P(z, \eta)=P^{\prime}(z, \eta)$ for $\eta \in \mathcal{D}$, and hence, for $\eta \in \mathbb{C}_{+}$.
To show the existence of the solution, set $P_{0}(z, \eta)=-\eta^{-1} I_{2 N} \in \mathcal{S}^{2 N}$, and consider the iterations

$$
P_{k+1}(z, \eta)=\mathcal{F}_{\Sigma, z}\left(P_{k}(z, \eta), \eta\right) .
$$

Then, $P_{k}(z, \cdot) \in \mathcal{S}^{2 N}$ for each $k$, and furthermore, the sequence $P_{k}(z, \eta)$ converges on $\mathcal{D}$ to a function $P_{\infty}(z, \eta)$ which satisfies $P_{\infty}(z, \eta)=\mathcal{F}_{\Sigma, z}\left(P_{\infty}(z, \eta), \eta\right)$ by Banach's fixed point theorem. From standard arguments using the normal family theorem, $P_{\infty}(z, \cdot)$ solves Equation (2) in $\mathcal{S}^{2 N}$.

It remains to establish the last result of Theorem 5. Extending the domain of $\mathcal{F}_{\Sigma, z}(P(z, \eta), \eta)$ in the parameter $\eta$ to $\mathbb{C} \backslash \mathbb{R}$, we show that $P$ is the solution of the equation $P=\mathcal{F}_{\Sigma, z}(P, \eta)$ if and only if the matrix function $P^{\prime}=\left[\begin{array}{cc}-P_{00} & P_{01} \\ P_{10} & -P_{11}\end{array}\right]$ is the solution of the equation $P^{\prime}=\mathcal{F}_{\Sigma}\left(P^{\prime},-\eta\right)$. This can be demonstrated by a direct calculation: Writing

$$
T\left(e^{\imath \theta}, z, \eta\right)=\mathscr{T}\left(\left(I_{2} \otimes \Sigma\left(e^{\imath \theta}\right)\right) P(z, \eta)\right)=\left[\begin{array}{cc}
t_{00} & t_{01} \\
t_{10} & t_{11}
\end{array}\right]
$$

we have

$$
\left(T\left(e^{\imath \theta}, z, \eta\right)+U_{L}\left(e^{\imath \theta}\right)\right)^{-1}=\frac{1}{\Delta}\left[\begin{array}{cc}
t_{11} & -\left(t_{01}+e^{-\imath L \theta}\right) \\
-\left(t_{10}+e^{2 L \theta}\right) & t_{00}
\end{array}\right]
$$

with $\Delta=t_{00} t_{11}-\left(t_{01}+e^{-\imath L \theta}\right)\left(t_{10}+e^{\imath L \theta}\right)$. Thus,

$$
P=\left[\begin{array}{ll}
P_{00} & P_{01} \\
P_{10} & P_{11}
\end{array}\right]=\left[\begin{array}{cc}
\frac{1}{2 \pi} \int \frac{t_{11} \Sigma-\eta}{\Delta}{ }^{2 \pi} & -\frac{1}{2 \pi} \int \frac{\left(t_{01}+e^{-2 L \theta}\right)}{\Delta} \Sigma-z \\
-\frac{1}{2 \pi} \int \frac{\left(t_{10}+e^{2 L \theta}\right)}{\Delta} \Sigma-\bar{z} & \frac{1}{2 \pi} \int \frac{t_{00}}{\Delta} \Sigma-\eta
\end{array}\right]^{-1} .
$$

Recalling the formula for the inverse of a partitioned matrix (see [15, §0.7.3])

$$
\left[\begin{array}{ll}
M_{00} & M_{01} \\
M_{10} & M_{11}
\end{array}\right]^{-1}=\left[\begin{array}{cc}
\left(M_{00}-M_{01} M_{11}^{-1} M_{10}\right)^{-1} & -M_{00}^{-1} M_{01}\left(M_{11}-M_{10} M_{00}^{-1} M_{01}\right)^{-1} \\
-\left(M_{11}-M_{10} M_{00}^{-1} M_{01}\right)^{-1} M_{10} M_{00}^{-1} & \left(M_{11}-M_{10} M_{00}^{-1} M_{01}\right)^{-1}
\end{array}\right]
$$

we get the required result by a direct checking.
Using this result in conjunction with Lemma 26, we obtain that for each deterministic vector $u \in \mathbb{C}^{N}$, the scalar measures $\left(u^{*} \Lambda_{i i} u\right)(d t)$ are symmetric for $i=0,1$. This shows at once that the measures $\Lambda_{i i}$ are symmetric. The proof of Theorem 5 is now complete.

### 4.2 Proof of Theorem 6

Recall that $\boldsymbol{M}$ is the bound provided by Assumption 1-2. From now on we shall write $\gamma_{\text {sup }}=$ $\sup _{N} N / n$.

We first establish the tightness of the sequence $\left(\check{\boldsymbol{\nu}}_{z, N}\right)_{N}$. As is well-known [10], this is equivalent to showing that $-\imath t g_{\check{\nu}_{z, N}}(\imath t) \rightarrow 1$ as $t \rightarrow \infty$ uniformly in $N$.

Given $M \in \mathcal{M}_{+}^{2 N}$, let

$$
A\left(M, e^{\imath \theta}\right) \triangleq \mathscr{T}\left(\left(I_{2} \otimes S\left(e^{\imath \theta}\right)\right) M\right)+U_{L}\left(e^{\imath \theta}\right)
$$

By Lemma 27,

$$
\left\|\mathscr{T}\left(\left(I_{2} \otimes S\left(e^{\imath \theta}\right)\right) G(z, \imath t)\right)\right\| \leqslant \boldsymbol{\gamma}_{\mathrm{sup}} \boldsymbol{M} / t
$$

Thus, for $t \geqslant 2 \gamma_{\text {sup }} \boldsymbol{M},\left\|A\left(G(z, \imath t), e^{\imath \theta}\right)^{-1}\right\| \leqslant 2$. Since

$$
G(z, \imath t)=\mathcal{F}_{S, z}(G(z, \imath t), \imath t)=\left(\frac{1}{2 \pi} \int A\left(G(z, \imath t), e^{\imath \theta}\right)^{-1} \otimes S\left(e^{\imath \theta}\right) d \theta-\left[\begin{array}{cc}
\imath t & z \\
\bar{z} & \imath t
\end{array}\right] \otimes I_{N}\right)^{-1}
$$

the required conclusion can be obtained by a simple argument.
The remainder of the proof is devoted towards establishing the convergence (3). Recalling the expression $\widehat{R}_{L}=X J^{L} X^{*}$ provided at the beginning of Section 3.1, we have

$$
Q(z, \eta)=\left(\boldsymbol{H}\left(\widehat{R}_{L}-z\right)-\eta\right)^{-1}=\left[\begin{array}{cc}
-\eta & X J^{L} X^{*}-z \\
X J^{-L} X^{*}-\bar{z} & -\eta
\end{array}\right]^{-1}=\left[\begin{array}{ll}
Q_{00} & Q_{01} \\
Q_{10} & Q_{11}
\end{array}\right]
$$

We begin by bounding the variance of $(2 N)^{-1} \operatorname{tr} D Q$ at the left hand side of (3).
Proposition 28. Under Assumption 1-2), for each deterministic matrix $B \in \mathbb{C}^{N \times N}$ and each $u, v \in\{0,1\}$,

$$
\operatorname{Var}\left(\operatorname{tr} B Q_{u v}(z, \eta)\right) \leqslant 8 \boldsymbol{\gamma}_{\mathrm{sup}} \boldsymbol{M}^{2}\|B\|^{2} /(\Im \eta)^{4}
$$

The proof of this proposition will be based on the well-known Poincaré-Nash (PN) inequality [6], [21], which is also a particular case of the Brascamp-Lieb inequality. Let $\boldsymbol{v}=\left[v_{0}, \ldots, v_{m-1}\right]^{\top}$ be a complex Gaussian random vector with $\mathbb{E} \boldsymbol{v}=0, \mathbb{E} \boldsymbol{v} \boldsymbol{v}^{\top}=0$, and $\mathbb{E}\left[\boldsymbol{v} \boldsymbol{v}^{*}\right]=\Sigma$. Let $\varphi=$
$\varphi\left(v_{0}, \ldots, v_{m-1}, \bar{v}_{0}, \ldots, \bar{v}_{m-1}\right)$ be a $C^{1}$ complex function which is polynomially bounded together with its derivatives. Then, writing

$$
\nabla_{\boldsymbol{v}} \varphi=\left[\partial \varphi / \partial v_{0}, \ldots, \partial \varphi / \partial v_{m-1}\right]^{\top} \text { and } \nabla_{\overline{\boldsymbol{v}}} \varphi=\left[\partial \varphi / \partial \bar{v}_{0}, \ldots, \partial \varphi / \partial \bar{v}_{m-1}\right]^{\top},
$$

the $P N$ inequality is

$$
\begin{equation*}
\mathbb{V a r}(\varphi(\boldsymbol{v})) \leqslant \mathbb{E}\left[\nabla_{\boldsymbol{v}} \varphi(\boldsymbol{v})^{T} \Sigma \overline{\nabla_{\boldsymbol{v}} \varphi(\boldsymbol{v})}\right]+\mathbb{E}\left[\left(\nabla_{\overline{\boldsymbol{v}}} \varphi(\boldsymbol{v})\right)^{*} \Sigma \nabla_{\overline{\boldsymbol{v}}} \varphi(\boldsymbol{v})\right] \tag{27}
\end{equation*}
$$

If we rewrite $Q(z, \eta)$ as $Q(z, \eta)=Q^{X}$ to emphasize the dependence of the resolvent on the matrix $X$, then, given a matrix $\Delta \in \mathbb{C}^{N \times n}$, the resolvent identity implies that

$$
\begin{aligned}
Q^{X+\Delta}-Q^{X} & =-Q^{X+\Delta}\left(\left(Q^{X+\Delta}\right)^{-1}-\left(Q^{X}\right)^{-1}\right) Q^{X} \\
& =-Q^{X+\Delta}\left[(X+\Delta) J^{-L}(X+\Delta)^{*}-X J^{-L} X^{*} \quad(X+\Delta) J^{L}(X+\Delta)^{*}-X J^{L} X^{*}\right] Q^{X}
\end{aligned}
$$

Using this equation, we can obtain the expression of $\partial a^{*} Q_{u v} b / \partial \bar{x}_{i j}$, where $a, b \in \mathbb{C}^{N}, u, v \in\{0,1\}$, $i \in[N]$, and $j \in[n]$. Indeed, taking $\Delta=e_{N, i} e_{n, j}^{*}$ in the former expression, we get after a simple derivation that

$$
\begin{equation*}
\frac{\partial a^{*} Q_{u v} b}{\partial \bar{x}_{i j}}=-\left[a^{*} Q_{u 1} X J^{-L}\right]_{j}\left[Q_{0 v} b\right]_{i}-\left[a^{*} Q_{u 0} X J^{L}\right]_{j}\left[Q_{1 v} b\right]_{i} \tag{28}
\end{equation*}
$$

Proof of Proposition 28. We apply Inequality (27) by respectively replacing $\boldsymbol{v}$ and $\varphi$ with vec $X$ and $\operatorname{tr} B Q_{u v}$ (seen as a function of $X$ ).

Given $k, i \in[N]$, and $j \in[n]$, we have

$$
\frac{\partial\left[e_{k}^{*} B Q_{u v} e_{k}\right]}{\partial \bar{x}_{i j}}=-\left[B Q_{u 1} X J^{-L}\right]_{k j}\left[Q_{0 v}\right]_{i k}-\left[B Q_{u 0} X J^{L}\right]_{k j}\left[Q_{1 v}\right]_{i k}
$$

and thus,

$$
\frac{\partial \operatorname{tr} B Q_{u v}}{\partial \bar{x}_{i j}}=-\left[Q_{0 v} B Q_{u 1} X J^{-L}\right]_{i j}-\left[Q_{1 v} B Q_{u 0} X J^{L}\right]_{i j}
$$

Let us focus on the second term at the right hand side of Inequality (27). Observing that $\mathbb{E}\left[x_{i_{1} j_{1}} \bar{x}_{i_{2} j_{2}}\right]=n^{-1}\left[R_{j_{1}-j_{2}}\right]_{i_{1}, i_{2}}$, and recalling the expression of the block-Toeplitz matrix $\mathcal{R}$ given by Equation (10), we have

$$
\begin{aligned}
& \quad \sum_{i_{1}, i_{2} \in[N]} \sum_{j_{1}, j_{2} \in[n]} \mathbb{E}\left[\frac{\partial \overline{\operatorname{tr} B Q_{u v}}}{\partial \bar{x}_{i_{1} j_{1}}} \mathbb{E}\left[x_{i_{1} j_{1}} \bar{x}_{i_{2} j_{2}}\right] \frac{\partial \operatorname{tr} B Q_{u v}}{\partial \bar{x}_{i_{2} j_{2}}}\right] \\
& \leqslant \\
& \leqslant \mathbb{E} \operatorname{vec}\left(Q_{0 v} B Q_{u 1} X J^{-L}\right)^{*} \mathcal{R} \operatorname{vec}\left(Q_{0 v} B Q_{u 1} X J^{-L}\right)+2 \mathbb{E} \operatorname{vec}\left(Q_{1 v} B Q_{u 0} X J^{L}\right)^{*} \mathcal{R} \operatorname{vec}\left(Q_{1 v} B Q_{u 0} X J^{L}\right) \\
& \leqslant \\
& \leqslant \frac{2 \boldsymbol{M}}{n}\left(\mathbb{E}\left\|Q_{0 v} B Q_{u 1} X J^{-L}\right\|_{\mathrm{HS}}^{2}+\mathbb{E}\left\|Q_{1 v} B Q_{u 0} X J^{L}\right\|_{\mathrm{HS}}^{2}\right) \\
& \leqslant \\
& \leqslant \frac{4 \boldsymbol{M}\|B\|^{2}}{(\Im \eta)^{4} n} \mathbb{E}\|X\|_{\mathrm{HS}}^{2} \\
& \leqslant \frac{4 \boldsymbol{\gamma}_{\mathrm{sup}} \boldsymbol{M}^{2}\|B\|^{2}}{(\Im \eta)^{4}}
\end{aligned}
$$

The first term at the right hand side of Inequality (27) is treated similarly, leading to the bound given in the statement, and the proof is complete.

In order to establish the convergence (3), using Proposition 28 and the Borel-Cantelli lemma, it will be enough to show that

$$
\begin{equation*}
\forall \eta \in \mathbb{C}_{+}, \quad \frac{1}{2 N} \operatorname{tr} D^{(N)}\left(\mathbb{E} Q^{(N)}(z, \eta)-G^{(N)}(z, \eta)\right) \xrightarrow[N \rightarrow \infty]{ } 0 \tag{29}
\end{equation*}
$$

Following the general canvas of [23], we approximate our process $\boldsymbol{x}^{(N)}$ by a Moving Average process with a finite memory. We shall construct from this MA process a resolvent that will be more easily manageable than $Q^{(N)}$.

With reference to Theorem 5, first consider a discrete analogue of the integral within the expression of $\mathcal{F}_{\Sigma, z}$. A straightforward adaptation of its proof yields the following proposition and we omit the details.

Proposition 29. Let $\Sigma: \mathbb{T} \rightarrow \mathcal{H}_{+}^{N}$ be a continuous function, and let $z \in \mathbb{C}$. Then, the conclusions of Theorem 5 remain true if the function $\mathcal{F}_{\Sigma, z}$ there is replaced with

$$
\left.\left.\left.\begin{array}{l}
\check{\mathcal{F}}_{\Sigma, z}(M(\eta), \eta)= \\
\left(\frac { 1 } { n } \sum _ { \ell \in [ n ] } \left(\mathscr{T}\left(\left(I_{2} \otimes \Sigma\left(e^{2 \imath \pi \ell / n}\right)\right) M(\eta)\right)+\left[e^{-2 \imath \pi \ell L / n}\right.\right.\right.
\end{array} e^{-2 \imath \pi \ell L / n}\right]\right)^{-1} \otimes \Sigma\left(e^{2 \imath \pi \ell / n}\right)-\left[\begin{array}{ll}
\eta & z \\
\bar{z} & \eta
\end{array}\right] \otimes I_{N}\right)^{-1} .
$$

Now, given an integer constant $K>0$, let us define the function $\widehat{S}^{(N, K)}$ on $\mathbb{T}$ as

$$
\widehat{S}^{(N, K)}\left(e^{\imath \theta}\right)=\frac{1}{2 \pi} \int_{0}^{2 \pi} F_{K}\left(e^{\imath(\theta-\psi)}\right) S^{(N)}\left(e^{\imath \psi}\right) d \psi
$$

where $F_{K}$ is the Fejér kernel (see Equation (12)). It can be easily checked that this function has the following properties:

1. $\widehat{S}^{(N, K)}$ is a spectral density.
2. $\widehat{S}^{(N, K)}$ is a Laurent trigonometric polynomial of the form $\widehat{S}^{(N, K)}\left(e^{\imath \theta}\right)=\sum_{\ell=-K}^{K} e^{\imath \ell \theta} \widetilde{R}_{\ell}^{(N, K)}$.
3. With Assumptions 1-1 and 1-2, it holds that

$$
\begin{equation*}
\sup _{N}\left\|\widehat{S}^{(N, K)}-S^{(N)}\right\|_{\infty}^{\mathbb{T}} \xrightarrow[K \rightarrow \infty]{ } 0 \tag{30}
\end{equation*}
$$

This last property can be established by splitting the integral that defines $\widehat{S}^{(N, K)}\left(e^{\imath \theta}\right)$ into two pieces as $\int_{0}^{2 \pi}=\int_{\psi:|\theta-\psi| \leqslant \delta}+\int_{\psi:|\theta-\psi|>\delta}$ for a properly chosen $\delta>0$, and by using the properties of the Fejér kernel provided after Equation (12).

Consider the implicit equation

$$
\widehat{G}^{(N, K)}(z, \eta)=\check{\mathcal{F}}_{\widehat{S}^{(N, K)}, z}\left(\widehat{G}^{(N, K)}(z, \eta), \eta\right) \quad\left(\text { in } \quad \mathcal{S}^{2 N}\right)
$$

From the previous proposition, the solution $\widehat{G}^{(N, K)}(z, \cdot)$ exists and is unique. The following three propositions will be proved in Section 4.3.

Proposition 30. For each $\eta$ such that $\Im \eta \geqslant C_{30}$ with $C_{30}=C_{30}\left(\gamma_{\text {sup }}, \boldsymbol{M}\right)>0$,

$$
\underset{N}{\limsup }\left\|\widehat{G}^{(N, K)}(z, \eta)-G^{(N)}(z, \eta)\right\| \xrightarrow[K \rightarrow \infty]{ } 0
$$

Let $\left(\hat{x}_{k}^{(N, K)}\right)_{k \in \mathbb{Z}}$ be a $\mathbb{C}^{N}$-valued stationary centered Gaussian process with the spectral density $n^{-1} \widehat{S}^{(N, K)}$. Define

$$
\hat{X}^{(N, K)}=\left[\begin{array}{lll}
\hat{x}_{0}^{(N, K)} & \cdots & \hat{x}_{n-1}^{(N, K)}
\end{array}\right] \text { and } \hat{Q}^{(N, K)}(z, \eta)=\left(\boldsymbol{H}\left(\hat{X}^{(N, K)} J^{L}\left(\hat{X}^{(N, K)}\right)^{*}-z\right)-\eta I\right)^{-1} .
$$

We then have the following proposition.

Proposition 31. Fix $K>0$. Then, for an arbitrary deterministic matrix $D^{(N)} \in \mathbb{C}^{2 N \times 2 N}$ with $\left\|D^{(N)}\right\|=1$, we have

$$
\frac{1}{N}\left|\operatorname{tr} D^{(N)}\left(\mathbb{E} \widehat{Q}^{(N, K)}(z, \eta)-\widehat{G}^{(N, K)}(z, \eta)\right)\right| \leqslant \frac{C K}{\sqrt{N}}
$$

where $C>0$ depends only on $\eta, \boldsymbol{M}$, and $\boldsymbol{\gamma}_{\text {sup }}$.
We note here that the bound provided in the statement of this proposition is not optimal but is enough for our purpose.
Proposition 32. With $D^{(N)}$ as in the previous proposition,

$$
\limsup _{N} \frac{1}{N}\left|\operatorname{tr} D^{(N)}\left(\mathbb{E} Q^{(N)}(z, \eta)-\mathbb{E} \widehat{Q}^{(N, K)}(z, \eta)\right)\right| \underset{K \rightarrow \infty}{ } 0
$$

## Theorem 6: end of the proof

We write

$$
\begin{aligned}
\frac{1}{N} \operatorname{tr} D^{(N)}\left(\mathbb{E} Q^{(N)}-G^{(N)}\right)= & \frac{1}{N} \operatorname{tr} D^{(N)}\left(\mathbb{E} Q^{(N)}-\mathbb{E} \widehat{Q}^{N, K}\right)+\frac{1}{N} \operatorname{tr} D^{(N)}\left(\mathbb{E} \widehat{Q}^{N, K}-\widehat{G}^{(N, K)}\right) \\
& +\frac{1}{N} \operatorname{tr} D^{(N)}\left(\widehat{G}^{(N, K)}-G^{(N)}\right) \\
& \triangleq \chi_{1}(N, K)+\chi_{2}(N, K)+\chi_{3}(N, K)
\end{aligned}
$$

Fix an arbitrarily small $\varepsilon>0$. Let $K_{0}, N_{0}>0$ be such that, by Propositions 32 and 30 ,

$$
\left|\chi_{1}\left(N, K_{0}\right)\right|,\left|\chi_{3}\left(N, K_{0}\right)\right| \leqslant \varepsilon \quad \text { for all } N \geqslant N_{0} \text { and } \Im \eta>C_{30}
$$

By Proposition 32, $\chi_{2}\left(N, K_{0}\right) \rightarrow_{N \rightarrow \infty} 0$. Thus, $(2 N)^{-1} \operatorname{tr} D^{(N)}\left(\mathbb{E} Q^{(N)}-G^{(N)}\right) \rightarrow_{N} 0$ for $\Im \eta>C_{30}$, and thus, for each $\eta \in \mathbb{C}_{+}$by analyticity, and (29) is established. This completes the proof of Theorem 6.

### 4.3 Remaining proofs for Section 4.2

## Proof of Proposition 30

Let $\bar{S}^{(N, K)}\left(e^{\imath \theta}\right)=\widehat{S}^{(N, K)}\left(e^{2 \imath \pi k / n}\right)$ and $\bar{U}_{L}\left(e^{\imath \theta}\right)=U_{L}\left(e^{2 \imath \pi k / n}\right)$ for $\theta \in[2 \pi k / n, 2 \pi(k+1) / n)$ be the respective stepwise continuous versions of the functions $\widehat{S}^{(N, K)}$ and $U_{L}$ with step size $2 \pi / n$. Within this proof, we re-denote the function $\mathcal{F}_{\Sigma, z}$ defined in the statement of Theorem 5 as $\mathcal{F}_{\Sigma, z, U_{L}}$ to stress the dependence on $U_{L}$. With this notation, it is obvious that $\check{\mathcal{F}}_{\widehat{S}^{(N, K)}, z}=\mathcal{F}_{\bar{S}^{(N, K)}, z, \bar{U}_{L}}$.

In the rest of the proof, we often drop the superscripts ${ }^{(N)}$ and ${ }^{(N, K)}$ and the subscript ${ }_{L}$ for brevity. Given $M \in \mathcal{S}^{2 N}$, put $A_{\boldsymbol{S}, \boldsymbol{U}}\left(M, e^{\imath \theta}\right) \triangleq \mathscr{T}\left(\left(I_{2} \otimes \boldsymbol{S}\left(e^{\imath \theta}\right)\right) M\right)+\boldsymbol{U}\left(e^{\imath \theta}\right)$, where $(\boldsymbol{S}, \boldsymbol{U})=(S, U)$ or $(\bar{S}, \bar{U})$. Write $B_{\boldsymbol{S}, \boldsymbol{U}}(M)=(2 \pi)^{-1} \int_{0}^{2 \pi} A_{\boldsymbol{S}, \boldsymbol{U}}\left(M, e^{2 \theta}\right)^{-1} \otimes \boldsymbol{S}\left(e^{\imath \theta}\right) d \theta$.

We also assume that $K$ is large enough so that

$$
\sup _{N}\|\bar{S}\|_{\infty}^{\mathbb{T}} \leqslant 2 M
$$

By dropping the unnecessary parameters from the notations, we write

$$
\begin{aligned}
G-\widehat{G}= & \mathcal{F}_{S, U}(G)-\mathcal{F}_{\bar{S}, \bar{U}}(\widehat{G})=\mathcal{F}_{S, U}(G)\left(B_{\bar{S}, \bar{U}}(\widehat{G})-B_{S, U}(G)\right) \mathcal{F}_{\bar{S}, \bar{U}}(\widehat{G}) \\
= & \mathcal{F}_{S, U}(G)\left(\frac{1}{2 \pi} \int_{0}^{2 \pi}\left(A_{\bar{S}, \bar{U}}\left(\widehat{G}, e^{\imath \theta}\right)^{-1} \otimes \bar{S}\left(e^{\imath \theta}\right)-A_{S, U}\left(G, e^{\imath \theta}\right)^{-1} \otimes S\left(e^{\imath \theta}\right)\right) d \theta\right) \mathcal{F}_{\bar{S}, \bar{U}}(\widehat{G}) \\
= & \mathcal{F}_{S, U}(G)\left(\frac { 1 } { 2 \pi } \int \left(A_{\bar{S}, \bar{U}}(\widehat{G})^{-1} \otimes \bar{S}-A_{S, U}(\widehat{G})^{-1} \otimes S\right.\right. \\
& \left.\left.+\left(A_{S, U}(\widehat{G})^{-1}-A_{S, U}(G)^{-1}\right) \otimes S\right) d \theta\right) \mathcal{F}_{\bar{S}, \bar{U}}(\widehat{G}) \\
= & \mathcal{F}_{S, U}(G)\left(\frac { 1 } { 2 \pi } \int \left(A_{\bar{S}, \bar{U}}(\widehat{G})^{-1} \otimes(\bar{S}-S)+\left(A_{\bar{S}, \bar{U}}(\widehat{G})^{-1}-A_{S, U}\left(\widehat{G}^{-1}\right) \otimes S\right.\right.\right. \\
& \left.\left.+\left(A_{S, U}(\widehat{G})^{-1}-A_{S, U}(G)^{-1}\right) \otimes S\right) d \theta\right) \mathcal{F}_{\bar{S}, \bar{U}}(\widehat{G}) \\
\cong & \mathcal{F}_{S, U}(G)\left(\frac{1}{2 \pi} \int\left(\chi_{1}+\chi_{2}+\chi_{3}\right) d \theta\right) \mathcal{F}_{\bar{S}, \bar{U}}(\widehat{G}) .
\end{aligned}
$$

By Lemma 27 , we have that $\left\|\mathscr{T}\left(\left(I_{2} \otimes \boldsymbol{S}\left(e^{\imath \theta}\right)\right) \boldsymbol{G}\right)\right\| \leqslant 2 \boldsymbol{\gamma}_{\text {sup }} \boldsymbol{M} / \Im \eta$ for any of the possibilities for $\boldsymbol{S}$ and for $\boldsymbol{G}=G, \widehat{G}$. Thus, for

$$
2 \gamma_{\mathrm{sup}} M / \Im \eta \leqslant 1 / 2
$$

we have $\left\|A_{\boldsymbol{S}}\left(\boldsymbol{G}, e^{2 \theta}\right)^{-1}\right\| \leqslant 2$. Therefore,

$$
\left\|\chi_{1}\right\| \leqslant 2\|\bar{S}-S\|_{\infty}^{\mathbb{T}}
$$

Moreover,
$\chi_{2}=\left(A_{\bar{S}, \bar{U}}(\widehat{G})^{-1}-A_{S, U}(\widehat{G})^{-1}\right) \otimes S=\left(A_{\bar{S}, \bar{U}}(\widehat{G})^{-1}\left(\mathscr{T}\left(\left(I_{2} \otimes(\bar{S}-S)\right) \widehat{G}\right)+U-\bar{U}\right) A_{S, U}(\widehat{G})^{-1}\right) \otimes S$
satisfies

$$
\left\|\chi_{2}\right\| \leqslant 4 \boldsymbol{M} \boldsymbol{\gamma}_{\text {sup }}(\Im \eta)^{-1}\|\bar{S}-S\|_{\infty}^{\mathbb{T}}+4 \boldsymbol{M}\|U-\bar{U}\|
$$

for the same values of $\eta$. By mimicking the calculation that lead to Inequality (26), we obtain

$$
\left\|\chi_{3}\right\| \leqslant \frac{4}{(\Im \eta)^{2}} \gamma_{\text {sup }} \boldsymbol{M}^{2}\|G-\widehat{G}\|
$$

Using the inequality $\left\|\mathcal{F}_{\boldsymbol{S}}(\boldsymbol{G})\right\| \leqslant 1 /(\Im \eta)$, we obtain altogether
$\left(1-\frac{4 \gamma_{\mathrm{sup}} \boldsymbol{M}^{2}}{(\Im \eta)^{4}}\right)\left\|G^{(N)}-\widehat{G}^{(N, K)}\right\| \leqslant \frac{2}{(\Im \eta)^{2}}\left(1+\frac{4 \boldsymbol{\gamma}_{\mathrm{sup}} \boldsymbol{M}}{\Im \eta}\right)\left\|\bar{S}^{(N, K)}-S^{(N)}\right\|_{\infty}^{\mathbb{T}}+\frac{4 \boldsymbol{M}}{(\Im \eta)^{2}}\left\|\bar{U}_{L}-U_{L}\right\|_{\infty}^{\mathbb{T}}$.
Thus, if $\Im \eta>\left(8 \gamma_{\text {sup }} \boldsymbol{M}^{2}\right)^{1 / 4} \vee\left(4 \gamma_{\text {sup }} \boldsymbol{M}\right)$, then

$$
\left\|G^{(N)}-\widehat{G}^{(N, K)}\right\| \leqslant C\left\|\bar{S}^{(N, K)}-S^{(N)}\right\|_{\infty}^{\mathbb{T}}+C^{\prime}\left\|\bar{U}_{L}-U_{L}\right\|_{\infty}^{\mathbb{T}}
$$

for some constants $C, C^{\prime}>0$, and now the proof can be completed by suing the convergence (30).

## Proof of Proposition 31

From now on, $C$ is a positive constant that depends on $\eta, \boldsymbol{M}$, and $\boldsymbol{\gamma}_{\text {sup }}$ at most, and can change from line to line. Recalling the properties 1 and 2 of the density $\widehat{S}^{(N, K)}$ that were stated in Section 4.2 above, our first step is to apply the well-known operator version of the Fejér-Riesz theorem (see [24, Sec. 6.6]) to $\widehat{S}^{(N, K)}$. This implies that for each $(N, K)$, there exists an $N \times N$ matrix trigonometric polynomial

$$
P^{(N, K)}\left(e^{\imath \theta}\right)=\sum_{\ell=0}^{K} e^{\imath \ell \theta} B_{\ell}^{(N, K)}
$$

such that $\widehat{S}^{(N, K)}\left(e^{\imath \theta}\right)=P^{(N, K)}\left(e^{\imath \theta}\right) P^{(N, K)}\left(e^{\imath \theta}\right)^{*}$.
Letting $\boldsymbol{\xi}^{(N)}=\left(\boldsymbol{\xi}_{k}^{(N)}\right)_{k \in \mathbb{Z}}$ be an i.i.d. process with $\boldsymbol{\xi}_{k}^{(N)} \sim \mathcal{N}_{\mathbb{C}}\left(0, I_{N}\right)$, the process $\left(\hat{x}_{k}^{(N, K)}\right)_{k \in \mathbb{Z}}$ that we used to construct the resolvent $\widehat{Q}^{(N, K)}$ can be defined as

$$
\hat{x}_{k}^{(N, K)}=\frac{1}{\sqrt{n}} \sum_{\ell=0}^{K} B_{\ell}^{(N, K)} \boldsymbol{\xi}_{k-\ell}^{(N)} .
$$

Define the finite sequence of random vectors $\left(\tilde{x}_{k}^{(N, K)}\right)_{k \in[n]}$ as

$$
\tilde{x}_{k}^{(N, K)}=\frac{1}{\sqrt{n}} \sum_{\ell=0}^{K} B_{\ell}^{(N, K)} \boldsymbol{\xi}_{(k-\ell) \bmod n}^{(N)}
$$

(thus, $\tilde{x}_{k}^{(N, K)}$ is the analogue of $\hat{x}_{k}^{(N, K)}$ obtained through a circular convolution). Define
$\tilde{X}^{(N, K)}=\left[\begin{array}{lll}\tilde{x}_{0}^{(N, K)} & \cdots & \tilde{x}_{n-1}^{(N, K)}\end{array}\right] \in \mathbb{C}^{N \times n}$ and $\widetilde{Q}^{(N, K)}(z, \eta)=\left(\boldsymbol{H}\left(\tilde{X}^{(N, K)} J^{L}\left(\tilde{X}^{(N, K)}\right)^{*}-z\right)-\eta I\right)^{-1}$.
Observing that $\operatorname{rank}(\hat{X}-\tilde{X}) \leqslant K$, we get

$$
\operatorname{rank}\left(\boldsymbol{H}\left(\widehat{X} J^{L} \widehat{X}^{*}-z\right)-\boldsymbol{H}\left(\tilde{X} J^{L} \tilde{X}^{*}-z\right)\right) \leqslant 4 K
$$

Thus, for each matrix $D \in \mathbb{C}^{2 N \times 2 N}$, the following inequality

$$
\begin{equation*}
|\operatorname{tr} D(\widehat{Q}(z, \eta)-\widetilde{Q}(z, \eta))| \leqslant \frac{4 K\|D\|}{\Im \eta} \tag{31}
\end{equation*}
$$

holds (see [29, Lemma 2.6]). We can thus work with $\widetilde{Q}$ in place of $\widehat{Q}$ for establishing Proposition 31.
For $k \in[n]$, let

$$
w_{k}^{(N, K)}=\frac{1}{\sqrt{n}} \sum_{\ell=0}^{n-1} e^{2 \imath \pi k \ell / n} \tilde{x}_{\ell}^{(N, K)}
$$

be the discrete Fourier transform of the finite sequence $\left(\tilde{x}_{0}, \ldots, \tilde{x}_{n-1}\right)$, and define the matrix

$$
W^{(N, K)}=\left[\begin{array}{lll}
w_{0}^{(N, K)} & \cdots & w_{n-1}^{(N, K)}
\end{array}\right]=\tilde{X}^{(N, K)} \mathrm{F} \in \mathbb{C}^{N \times n}
$$

where F is the Fourier matrix defined in (6). Since the $\tilde{x}_{\ell}$ are built through a circular convolution, the vectors $w_{k}$ are independent, and $w_{k} \sim \mathcal{N}_{\mathbb{C}}\left(0, n^{-1} \widehat{S}\left(e^{2 \imath \pi k / n}\right)\right)$. Recalling that $J=\mathrm{F} \Omega \mathrm{F}^{*}$ with $\Omega=\operatorname{diag}\left(\omega^{\ell}\right)_{\ell=0}^{n-1}$ and $\omega=\exp (-2 \imath \pi / n)$, we can write

$$
\widetilde{Q}(z, \eta)=\left(\boldsymbol{H}\left(\tilde{X} J^{L} \tilde{X}^{*}-z\right)-\eta I\right)^{-1}=\left(\boldsymbol{H}\left(W \Omega^{L} W^{*}-z\right)-\eta I\right)^{-1}
$$

The remainder of the proof will be devoted towards showing that

$$
\begin{equation*}
\left\|\mathbb{E} \widetilde{Q}^{(N, K)}(z, \eta)-\widehat{G}^{(N, K)}(z, \eta)\right\| \leqslant C N^{-1 / 2} \tag{32}
\end{equation*}
$$

taking advantage of the independence of the columns of $W$. This bound, used in conjunction with the bound (31), immediately leads to the result of Proposition 31.

The proof of Inequality (32) relies on the NP inequality that we used above, as well as on the well-known Integration by Parts (IP) formula for Gaussian functionals [13, 17]. Recalling the definition of the vector $\boldsymbol{v}$ and the functional $\varphi$ after Proposition 28 above, the IP formula reads as

$$
\mathbb{E} v_{k} \varphi(\boldsymbol{v})=\sum_{\ell=0}^{n-1}[\Sigma]_{k \ell} \mathbb{E}\left[\frac{\partial \varphi(\boldsymbol{v})}{\partial \bar{v}_{\ell}}\right]
$$

Write $\widetilde{Q}(z, \eta)=\left[\begin{array}{ll}\widetilde{Q}_{00} & \widetilde{Q}_{01} \\ \widetilde{Q}_{10} & \widetilde{Q}_{11}\end{array}\right]$, and $W=\left[w_{i j}\right]_{i \in[N], j \in[n]}$. By reproducing verbatim the derivation that we made to obtain the Identity (28), we have

$$
\frac{\partial a^{*} \widetilde{Q}_{u v} b}{\partial \bar{w}_{i j}}=-\left[a^{*} \widetilde{Q}_{u 1} W \Omega^{-L}\right]_{j}\left[\widetilde{Q}_{0 v} b\right]_{i}-\left[a^{*} \widetilde{Q}_{u 0} W \Omega^{L}\right]_{j}\left[\widetilde{Q}_{1 v} b\right]_{i}
$$

In the subsequent derivations, we write $\widehat{S}_{k}=\widehat{S}\left(e^{2 \imath \pi k / n}\right)$. Let $a, b$ be two constant vectors in $\mathbb{C}^{N}$. Write $b=\left[b_{0}, \ldots, b_{N-1}\right]^{\top}$, and let $\alpha_{u v}(\ell)=\mathbb{E}\left[\left[a^{*} \widetilde{Q}_{u v} W\right]_{\ell}\left[W^{*} b\right]_{\ell}\right]$. By the IP formula, we have

$$
\begin{aligned}
\alpha_{u v}(\ell)= & \sum_{i \in[N]} \mathbb{E}\left[\left[a^{*} \widetilde{Q}_{u v}\right]_{i} w_{i \ell}\left[W^{*} b\right]_{\ell}\right]=\frac{1}{n} \sum_{i, m \in[N]}\left[\widehat{S}_{\ell}\right]_{i m} \mathbb{E} \frac{\partial\left(\left[a^{*} \widetilde{Q}_{u v}\right]_{i}\left[W^{*} b\right]_{\ell}\right)}{\partial \bar{w}_{m \ell}} \\
= & \frac{1}{n} \sum_{i, m}\left[\widehat{S}_{\ell}\right]_{i m} \mathbb{E}\left[-\left[a^{*} \widetilde{Q}_{u 1} W \Omega^{-L}\right]_{\ell}\left[\widetilde{Q}_{0 v}\right]_{m i}\left[W^{*} b\right]_{\ell}-\left[a^{*} \widetilde{Q}_{u 0} W \Omega^{L}\right]_{\ell}\left[\widetilde{Q}_{1 v}\right]_{m i}\left[W^{*} b\right]_{\ell}\right. \\
& \left.+\left[a^{*} \widetilde{Q}_{u v}\right]_{i} b_{m}\right] \\
= & -\mathbb{E}\left[\left[a^{*} \widetilde{Q}_{u 1} W\right]_{\ell}\left[W^{*} b\right]_{\ell} \omega^{-\ell L} \frac{\operatorname{tr} \widetilde{Q}_{0 v} \widehat{S}_{\ell}}{n}\right]-\mathbb{E}\left[\left[a^{*} \widetilde{Q}_{u 0} W\right]_{\ell}\left[W^{*} b\right]_{\ell} \omega^{\ell L} \frac{\operatorname{tr} \widetilde{Q}_{1 v} \widehat{S}_{\ell}}{n}\right] \\
+ & \frac{\mathbb{E}\left[a^{*} \widetilde{Q}_{u v} \widehat{S}_{\ell} b\right]}{n} .
\end{aligned}
$$

Write $\widetilde{q s}{ }_{u v}(\ell)=\mathbb{E} \operatorname{tr} \widetilde{Q}_{u v} \widehat{S}_{\ell} / n$. We shall isolate the terms $\widetilde{q s}_{0 v}(\ell)$ and $\widetilde{q s}_{1 v}(\ell)$ in the last display by resorting to the following lemma.

Lemma 33. It holds that $\operatorname{Var} \widetilde{q s} u v(\ell) \leqslant C / n^{2}$, and $\operatorname{Var}\left[a^{*} \widetilde{Q}_{u v} W\right]_{\ell}\left[W^{*} b\right]_{\ell} \leqslant C\|a\|^{2}\|b\|^{2} / n$.
Proof. The first bound is obtained by repeating almost word for word the proof of Proposition 28. To obtain the second bound, we also use the NP inequality again. We start by writing

$$
\begin{aligned}
\frac{\partial\left[a^{*} \widetilde{Q}_{u v} W\right]_{\ell}\left[W^{*} b\right]_{\ell}}{\partial \bar{w}_{i j}}= & \sum_{k} \frac{\partial\left[a^{*} \widetilde{Q}_{u v}\right]_{k} w_{k \ell}\left[W^{*} b\right]_{\ell}}{\partial \bar{w}_{i j}} \\
= & {\left[a^{*} \widetilde{Q}_{u v} W\right]_{\ell} b_{i} \mathbb{1}_{j=\ell}-\left[a^{*} \widetilde{Q}_{u 1} W \Omega^{-L}\right]_{j}\left[Q_{0 v} W\right]_{i \ell}\left[W^{*} b\right]_{\ell} } \\
& -\left[a^{*} \widetilde{Q}_{u 0} W \Omega^{L}\right]_{j}\left[Q_{1 v} W\right]_{i \ell}\left[W^{*} b\right]_{\ell} .
\end{aligned}
$$

We focus on the second term at the right hand side of Inequality (27), treating separately the three terms at the right hand side of the last display. Starting with the first term, we get

$$
\frac{1}{n} \sum_{i_{1}, i_{2} \in[N]} \mathbb{E} \overline{\left[a^{*} \widetilde{Q}_{u v} W\right]_{\ell} b_{i_{1}}}\left[\widehat{S}_{\ell}\right]_{i_{1}, i_{2}}\left[a^{*} \widetilde{Q}_{u v} W\right]_{\ell} b_{i_{2}}=\frac{1}{n} b^{*} \widehat{S}_{\ell} b \mathbb{E}\left|\left[a^{*} \widetilde{Q}_{u v} W\right]_{\ell}\right|^{2} \leqslant\|a\|^{2}\|b\|^{2} C / n
$$

Turning to the second of these terms, we have

$$
\begin{aligned}
& \frac{1}{n} \sum_{i_{1}, i_{2} \in[N]} \sum_{j \in[n]} \mathbb{E}\left|\left[W^{*} b\right]_{\ell}\right|^{2}\left|\left[a^{*} \widetilde{Q}_{u 1} W \Omega^{-L}\right]_{j}\right|^{2} \overline{\left.Q_{0 v} W\right]_{i_{1} \ell}}\left[\widehat{S}_{j}\right]_{i_{1}, i_{2}}\left[Q_{0 v} W\right]_{i_{2} \ell} \\
& \leqslant \frac{\boldsymbol{M}}{n} \mathbb{E}\left|\left[W^{*} b\right]_{\ell}\right|^{2}\left\|a^{*} \widetilde{Q}_{u 1} W\right\|^{2}\left\|\left[Q_{0 v} W\right]_{\cdot, \ell}\right\|^{2} \\
& \leqslant\|a\|^{2}\|b\|^{2} C / n^{2},
\end{aligned}
$$

where the last inequality can be obtained by applying, e.g., Lemma 16 along with standard inequalities. The third term can be handled similarly.

Thanks to these bounds and to Cauchy-Schwarz inequality, we obtain the identity

$$
\alpha_{u v}(\ell)=-\alpha_{u 0}(\ell) \omega^{\ell L} \widetilde{q s}_{1 v}(\ell)-\alpha_{u 1}(\ell) \omega^{-\ell L} \widetilde{q s}_{0 v}(\ell)+n^{-1} \mathbb{E}\left[a^{*} \widetilde{Q}_{u v} \widehat{S}_{\ell} b\right]+\varepsilon
$$

with $\|\varepsilon\| \leqslant C\|a\|\|b\| n^{-3 / 2}$. This leads to the system of equations

$$
\left[\begin{array}{cc}
1+\omega^{\ell L} \widetilde{q s}_{10}(\ell) & \omega^{-\ell L} \widetilde{q s}_{00}(\ell) \\
\omega^{\ell L} \widetilde{q s}_{11}(\ell) & 1+\omega^{-\ell L} \widetilde{q s}_{01}(\ell)
\end{array}\right]\left[\begin{array}{c}
\alpha_{u 0}(\ell) \\
\alpha_{u 1}(\ell)
\end{array}\right]=\frac{1}{n}\left[\begin{array}{c}
\mathbb{E}\left[a^{*} \widetilde{Q}_{u 0} \widehat{S}_{\ell} b\right] \\
\mathbb{E}\left[a^{*} \widetilde{Q}_{u 1} \widehat{S}_{\ell} b\right]
\end{array}\right]+\varepsilon,
$$

with $\|\varepsilon\| \leqslant C\|a\|\|b\| n^{-3 / 2}$. The matrix at the left hand side of this expression, that we denote as $T_{\ell}$, is written as

$$
T_{\ell}=I_{2}+\mathscr{T}\left(\left(I_{2} \otimes \widehat{S}_{\ell}\right)\left(\widetilde{Q}^{\top}\left[\begin{array}{ll}
\omega^{\ell L} & \omega^{-\ell L}
\end{array}\right]\right)\right)
$$

Assume that $K$ is large enough so that $\left\|\widehat{S}_{\ell}\right\| \leqslant 2 \boldsymbol{M}$ for all $\ell \in[n]$, and take $\Im \eta \geqslant 4 \gamma_{\text {sup }} \boldsymbol{M}$. Then by Lemma 27 we get that $\left\|T_{\ell}-I_{2}\right\| \leqslant 1 / 2$. Thus, the determinant

$$
d(\ell)=\left(1+\omega^{\ell L} \widetilde{q s}_{10}(\ell)\right)\left(1+\omega^{-\ell L} \widetilde{q s}_{01}(\ell)\right)-\widetilde{q} s_{00}(\ell) \widetilde{q} s_{11}(\ell)
$$

of $T_{\ell}$ is such that $|d(\ell)|$ is bounded away from zero uniformly in $N$ and $\ell$. Solving our system, and reusing henceforth the notations $\varepsilon$ and $\boldsymbol{\varepsilon}$ at will, we get that

$$
\left[\begin{array}{l}
\alpha_{u 0}(\ell)  \tag{33}\\
\alpha_{u 1}(\ell)
\end{array}\right]=\frac{1}{n} \frac{1}{d(\ell)}\left[\begin{array}{cc}
\left(1+\omega^{-\ell L} \widetilde{q s}_{01}(\ell)\right) & -\omega^{-\ell L} \widetilde{q s}_{00}(\ell) \\
-\omega^{\ell L} \widetilde{q s_{11}}(\ell) & \left(1+\omega^{\ell L} \widetilde{q s_{10}}(\ell)\right)
\end{array}\right]\left[\begin{array}{c}
\mathbb{E}\left[a^{*} \widetilde{Q}_{u 0} \widehat{S}_{\ell} b\right] \\
\mathbb{E}\left[a^{*} \widetilde{Q}_{u 1} \widehat{S}_{\ell} b\right]
\end{array}\right]+\varepsilon
$$

with $\|\varepsilon\| \leqslant C\|a\|\|b\| n^{-3 / 2}$.
Now, keeping in mind the identity $\widetilde{Q}\left(\boldsymbol{H}\left(W \Omega^{L} W^{*}\right)-\left[\begin{array}{ll}\eta & z \\ \bar{z} & \eta\end{array}\right] \otimes I_{N}\right)=I$, our purpose is to find an approximant of the matrix

$$
\widetilde{Q} \boldsymbol{H}\left(W \Omega^{L} W^{*}\right)=\left[\begin{array}{ll}
\widetilde{Q}_{01} W \Omega^{-L} W^{*} & \widetilde{Q}_{00} W \Omega^{L} W^{*} \\
\widetilde{Q}_{11} W \Omega^{-L} W^{*} & \widetilde{Q}_{10} W \Omega^{L} W^{*}
\end{array}\right]
$$

To that end, we write

$$
\mathbb{E}\left[a^{*} \widetilde{Q}_{u 0} W \Omega^{L} W^{*} b\right]=\sum_{\ell \in[n]} \omega^{\ell L} \alpha_{u 0}(\ell), \quad \text { and } \quad \mathbb{E}\left[a^{*} \widetilde{Q}_{u 1} W \Omega^{-L} W^{*} b\right]=\sum_{\ell \in[n]} \omega^{-\ell L} \alpha_{u 1}(\ell),
$$

and we use Equation (33) to obtain

$$
\begin{align*}
& {\left[\mathbb{E} a^{*} \widetilde{Q}_{u 1} W \Omega^{-L} W^{*} b \quad \mathbb{E} a^{*} \widetilde{Q}_{u 0} W \Omega^{L} W^{*} b\right]} \\
& =a^{*} \mathbb{E}\left[\begin{array}{ll}
\widetilde{Q}_{u 0} & \widetilde{Q}_{u 1}
\end{array}\right] \times \frac{1}{n} \sum_{\ell \in[n]} \frac{1}{d(\ell)}\left[\begin{array}{cc}
-\widetilde{q s_{11}}(\ell) & \omega^{\ell L}+\widetilde{q s_{01}}(\ell) \\
\omega^{-\ell L}+\widetilde{q s_{10}}(\ell) & -\widetilde{q s_{00}}(\ell)
\end{array}\right] \otimes\left(\widehat{S}_{\ell} b\right)+\varepsilon \tag{34}
\end{align*}
$$

with $\|\varepsilon\| \leqslant C\|a\|\|b\| n^{-1 / 2}$ (note that we lost a factor of $n^{-1}$ because of the summation $\sum_{\ell \in[n]}$ ). Let $U_{L, \ell} \triangleq U_{L}\left(e^{2 \imath \pi \ell / n}\right)=\left[\begin{array}{ll}\omega^{-\ell L} & \omega^{\ell L}\end{array}\right]$, and define the matrix

$$
\begin{aligned}
C(z, \eta) & \triangleq \frac{1}{n} \sum_{\ell \in[n]} \frac{1}{d(\ell)}\left[\begin{array}{cc}
-\widetilde{q s}_{11}(\ell) & \omega^{\ell L}+\widetilde{q s} s_{01}(\ell) \\
\omega^{-\ell L}+\widetilde{q s_{10}}(\ell) & -\widetilde{q s} s_{00}(\ell)
\end{array}\right] \otimes \widehat{S}_{\ell} \\
& =\frac{1}{n} \sum_{\ell \in[n]}\left(\left[\begin{array}{cc}
\widetilde{\widetilde{q s}_{00}}(\ell) & \widetilde{q} s_{01}(\ell) \\
\widetilde{q} s_{10}(\ell) & \widetilde{q} s_{11}(\ell)
\end{array}\right]+U_{L, \ell}\right)^{-1} \otimes \widehat{S}_{\ell} \\
& =\frac{1}{n} \sum_{\ell \in[n]}\left(\mathscr{T}\left(\left(I_{2} \otimes \widehat{S}_{\ell}\right) \mathbb{E} \widetilde{Q}\right)+U_{L, \ell}\right)^{-1} \otimes \widehat{S}_{\ell} .
\end{aligned}
$$

Then, Equation (34) can be rewritten as
$\mathbb{E} a^{*} \widetilde{Q}_{u 1} W \Omega^{-L} W^{*} b=a^{*} \mathbb{E}\left[\begin{array}{cc}\widetilde{Q}_{u 0} & \widetilde{Q}_{u 1}\end{array}\right] C\left[\begin{array}{l}b \\ 0\end{array}\right]+\varepsilon, \quad \mathbb{E} a^{*} \widetilde{Q}_{u 0} W \Omega^{L} W^{*} b=a^{*} \mathbb{E}\left[\begin{array}{cc}\widetilde{Q}_{u 0} & \widetilde{Q}_{u 1}\end{array}\right] C\left[\begin{array}{l}0 \\ b\end{array}\right]+\varepsilon^{\prime}$,
with $|\varepsilon|,\left|\varepsilon^{\prime}\right| \leqslant C\|a\|\|b\| n^{-1 / 2}$. Let $\boldsymbol{a}$ and $\boldsymbol{b}$ be two constant vectors in $\mathbb{C}^{2 N}$. Recalling the expression of $\widetilde{Q} \boldsymbol{H}\left(W \Omega^{L} W^{*}\right)$ above, the last display can be written compactly as

$$
\left|\boldsymbol{a}\left(\mathbb{E} \widetilde{Q} \boldsymbol{H}\left(W \Omega^{L} W^{*}\right)-\mathbb{E} \widetilde{Q} C\right) \boldsymbol{b}\right| \leqslant C\|\boldsymbol{a}\|\|\boldsymbol{b}\| n^{-1 / 2}
$$

equivalently,

$$
\left\|\mathbb{E} \widetilde{Q} \boldsymbol{H}\left(W \Omega^{L} W^{*}\right)-\mathbb{E} \widetilde{Q} C\right\| \leqslant C n^{-1 / 2}
$$

Since $\mathbb{E} \widetilde{Q} \in \mathcal{S}^{2 N}$, it is easy to prove, mostly by mimicking the first part of the proof of Theorem 5 , that the matrix function

$$
R(z, \eta)=\left(C(z, \eta)-\left[\begin{array}{ll}
\eta & z \\
\bar{z} & \eta
\end{array}\right] \otimes I_{N}\right)^{-1}
$$

is well defined for $\eta \in \mathbb{C} \backslash \mathbb{R}$, and $R(z, \cdot) \in \mathcal{S}^{2 N}$. In particular, $\|R(z, \eta)\| \leqslant 1 / \Im \eta$ for $\eta \in \mathbb{C}_{+}$. We therefore have

$$
\begin{aligned}
\|\mathbb{E} \widetilde{Q}-R\| & =\left\|\mathbb{E} \widetilde{Q}\left(R^{-1}-\widetilde{Q}^{-1}\right) R\right\|=\left\|\left(\mathbb{E} \widetilde{Q} C-\mathbb{E} \widetilde{Q} \boldsymbol{H}\left(W \Omega^{L} W^{*}\right)\right) R\right\| \\
& \leqslant \|\left(\mathbb{E} \widetilde{Q} C-\mathbb{E} \widetilde{Q} \boldsymbol{H}\left(W \Omega^{L} W^{*}\right)\| \| R \|\right. \\
& \leqslant C n^{-1 / 2}
\end{aligned}
$$

To complete the proof of Proposition 31, it remains to control the norm $\|R-\widehat{G}\|$. Remembering that $\widehat{G}$ is defined through the implicit equation in Proposition 30, we use the contraction property of $\breve{\mathcal{F}}_{\widehat{S}, z}(\cdot, \eta)$ to gain this control. By mimicking the derivation that led to Inequality (26), we obtain that if $K$ is large enough so that $\|\widehat{S}\|_{\infty}^{\mathbb{T}} \leqslant 2 \boldsymbol{M}$, and if $\Im \eta$ is large enough, then

$$
\left\|\check{\mathcal{F}}_{\widehat{S}, z}(M, \eta)-\check{\mathcal{F}}_{\widehat{S}, z}\left(M^{\prime}, \eta\right)\right\| \leqslant \frac{1}{2}\left\|M-M^{\prime}\right\| .
$$

Notice that $R=\check{\mathcal{F}}_{\widehat{S}, z}(\mathbb{E} \widetilde{Q}, \eta)$. Therefore, if $\Im \eta$ is large enough, we have

$$
\|R-\widehat{G}\|=\left\|\check{\mathcal{F}}_{\widehat{S}, z}(\mathbb{E} \widetilde{Q}, \eta)-\check{\mathcal{F}}_{\widehat{S}, z}(R, \eta)+\check{\mathcal{F}}_{\widehat{S}, z}(R, \eta)-\check{\mathcal{F}}_{\widehat{S}, z}(\widehat{G}, \eta)\right\| \leqslant \frac{1}{2}\|\mathbb{E} \widetilde{Q}-R\|+\frac{1}{2}\|R-\widehat{G}\|
$$

leading to

$$
\|R-\widehat{G}\| \leqslant\|\mathbb{E} \widetilde{Q}-R\| \leqslant C n^{-1 / 2}
$$

Finally $\|\mathbb{E} \widetilde{Q}-\widehat{G}\| \leqslant\|\mathbb{E} \widetilde{Q}-R\|+\|R-\widehat{G}\| \leqslant C n^{-1 / 2}$, and the proof of Proposition 31 is completed by combining this bound with the inequality (31).

## Proof of Proposition 32

We can assume that the processes $\left(x_{k}^{(N)}\right)_{k}$ and $\left(\hat{x}_{k}^{(N, K)}\right)_{k}$ that constitute the columns of the matrices $X^{(N)}$ and $\hat{X}^{(N, K)}$ respectively, are generated by applying the filters with Fourier transforms $n^{-1 / 2} S^{(N)}\left(e^{2 \theta}\right)^{1 / 2}$ and $n^{-1 / 2} \widehat{S}^{(N, K)}\left(e^{\imath \theta}\right)^{1 / 2}$ respectively to the same i.i.d. process $\boldsymbol{\xi}^{(N)}=\left(\boldsymbol{\xi}_{k}^{(N)}\right)_{k \in \mathbb{Z}}$ with $\boldsymbol{\xi}_{k}^{(N)} \sim \mathcal{N}_{\mathbb{C}}\left(0, I_{N}\right)$. This being the case, we have

$$
\mathbb{E}\left\|x_{k}^{(N)}-\hat{x}_{k}^{(N, K)}\right\|^{2}=\frac{1}{n} \frac{1}{2 \pi} \int_{0}^{2 \pi}\left\|S^{(N)}\left(e^{\imath \theta}\right)^{1 / 2}-\widehat{S}^{(N, K)}\left(e^{\imath \theta}\right)^{1 / 2}\right\|_{\mathrm{HS}}^{2} d \theta
$$

Hence, by (30), we get that

$$
\sup _{N} \mathbb{E}\left\|x_{k}^{(N)}-\hat{x}_{k}^{(N, K)}\right\|^{2} \xrightarrow[K \rightarrow \infty]{ } 0
$$

For any two square matrices $M_{1}$ and $M_{2}$ of same order, it is well-known that $\left|\operatorname{tr} M_{1} M_{2}\right| \leqslant$ $\left\|M_{1}\right\|_{\text {HS }}\left\|M_{2}\right\|_{\text {HS }}$. Thus, by the resolvent identity,

$$
\begin{aligned}
\frac{1}{N^{2}}|\operatorname{tr} D(\mathbb{E} Q(z, \eta)-\mathbb{E} \hat{Q}(z, \eta))|^{2} & \leqslant \frac{1}{N^{2}} \mathbb{E}\left|\operatorname{tr} D Q(z, \eta)\left(\boldsymbol{H}\left(\hat{X} J^{L} \hat{X}^{*}\right)-\boldsymbol{H}\left(X J^{L} X^{*}\right)\right) \widehat{Q}(z, \eta)\right|^{2} \\
& \leqslant \frac{1}{(\Im \eta)^{4}} \frac{1}{N} \mathbb{E}\left\|\boldsymbol{H}\left(\widehat{X} J^{L} \widehat{X}^{*}\right)-\boldsymbol{H}\left(X J^{L} X^{*}\right)\right\|_{\mathrm{HS}}^{2}
\end{aligned}
$$

Writing $\mathbb{E}\left\|X J^{L} X^{*}-\widehat{X} J^{L} \hat{X}^{*}\right\|_{\mathrm{HS}}^{2}=\mathbb{E}\left\|(X-\hat{X}) J^{L} X^{*}+\hat{X} J^{L}(X-\hat{X})^{*}\right\|_{\mathrm{HS}}^{2}$, it is enough to bound $\mathbb{E}\left\|(X-\hat{X}) J^{L} X^{*}\right\|_{\text {HS }}^{2}$. Given a constant $\kappa>0$, we have

$$
\begin{aligned}
\frac{1}{n} \mathbb{E}\left\|(X-\hat{X}) J^{L} X^{*}\right\|_{\mathrm{HS}}^{2} & \leqslant \frac{1}{n} \mathbb{E}\|(X-\hat{X})\|_{\mathrm{HS}}^{2}\|X\|^{2} \\
& \leqslant \frac{\kappa^{2}}{n} \mathbb{E}\|(X-\hat{X})\|_{\mathrm{HS}}^{2}+\frac{1}{n} \mathbb{E}\|(X-\hat{X})\|_{\mathrm{HS}}^{2}\|X\|^{2} \mathbb{1}_{\|X\|>\kappa} \\
& \leqslant \kappa^{2} \mathbb{E}\left\|x_{k}-\hat{x}_{k}\right\|^{2}+\frac{1}{n}\left(\mathbb{E}\left(\|X-\hat{X}\|_{\mathrm{HS}}^{4}\right)^{1 / 2}\left(\mathbb{E}\|X\|^{4} \mathbb{1}_{\|X\|>\kappa}\right)^{1 / 2}\right.
\end{aligned}
$$

With the help of Lemma 16, the second term in the last expression can be made as small as desired, independently of $N$, when $N$ is large enough, by choosing $\kappa$ large enough. The first term converges to zero as $K \rightarrow \infty$ as shown above.

## A Proof of Lemma 15

1. Using Markov's inequality, an obvious integration with respect to the exponential distribution, and using the inequality $-\log (1-x) \leqslant 2 x$ for $x \in[0,1 / 2]$, we have

$$
\begin{aligned}
\mathbb{P}\left[\left\|\Sigma^{1 / 2} \xi\right\| \geqslant \sqrt{2 N t}\right] & =\mathbb{P}\left[\sum_{\ell=0}^{N-1} s_{\ell}(\Sigma)\left|\xi_{\ell}\right|^{2} \geqslant 2 N t\right]=\mathbb{P}\left[\exp \frac{\sum_{\ell=0}^{N-1} s_{\ell}(\Sigma)\left|\xi_{\ell}\right|^{2}}{2\|\Sigma\|} \geqslant \exp \frac{N t}{\|\Sigma\|}\right] \\
& \leqslant e^{-N t /\|\Sigma\|} \mathbb{E}\left[\exp \frac{\sum_{\ell=0}^{N-1} s_{\ell}(\Sigma)\left|\xi_{\ell}\right|^{2}}{2\|\Sigma\|}\right]=e^{-\frac{N t}{\|\Sigma\|}} \prod_{\ell=0}^{N-1}\left(1-s_{\ell}(\Sigma) /(2\|\Sigma\|)\right)^{-1} \\
& \leqslant \exp \left(-\frac{N t}{\|\Sigma\|}+\sum_{\ell=0}^{N-1} \frac{s_{\ell}(\Sigma)}{\|\Sigma\|}\right) \\
& \leqslant e^{-(t /\|\Sigma\|-1) N} .
\end{aligned}
$$

2. We have $\left\|\Sigma^{1 / 2} \xi\right\|^{2} \stackrel{\mathcal{L}}{=} \sum_{\ell=0}^{N-1} s_{\ell}(\Sigma)\left|\xi_{\ell}\right|^{2} \geqslant \alpha \sum_{\ell=0}^{m-1}\left|\xi_{\ell}\right|^{2}$. Thus, by a calculation similar to above,

$$
\mathbb{P}\left[\left\|\Sigma^{1 / 2} \xi\right\| \leqslant \sqrt{\alpha m / 2}\right] \leqslant \mathbb{P}\left[\sum_{\ell=0}^{m-1}\left|\xi_{\ell}\right|^{2} \leqslant m / 2\right] \leqslant \exp \left(-c_{15,2} m\right)
$$

3. We obviously have that

$$
(\xi+a)^{*} M(\xi+a)=(\xi+a)^{*} \Re M(\xi+a)+\imath(\xi+a)^{*} \Im M(\xi+a),
$$

and both $(\xi+a)^{*} \Re M(\xi+a)$ and $(\xi+a)^{*} \Im M(\xi+a)$ are real. Furthermore,

$$
\|M\|_{\mathrm{HS}}^{2}=\|\Re M\|_{\mathrm{HS}}^{2}+\|\Im M\|_{\mathrm{HS}}^{2}
$$

Let us assume that $\|\Re M\|_{\text {HS }} \geqslant\|M\|_{\text {HS }} / \sqrt{2}$, otherwise, we replace $\Re M$ with $\Im M$. From these facts, we deduce that

$$
\mathcal{L}\left((\xi+a)^{*} M(\xi+a) /\|M\|_{\mathrm{HS}}, t\right) \leqslant \mathcal{L}\left((\xi+a)^{*} \Re M(\xi+a) /\|\Re M\|_{\mathrm{HS}}, \sqrt{2} t\right) .
$$

By a spectral factorization of the Hermitian matrix $\Re M$,

$$
(\xi+a)^{*} \Re M(\xi+a) /\|\Re M\|_{\mathrm{HS}} \stackrel{\mathcal{L}}{=} \sum_{\ell \in[N]} \beta_{\ell}\left|\xi_{\ell}+u_{\ell}\right|^{2}
$$

where the $u_{\ell}$ are deterministic complex numbers, and the $\beta_{\ell}$ are deterministic reals that satisfy $\sum_{\ell} \beta_{\ell}^{2}=1$, since they are the eigenvalues of $\|\Re M\|_{\mathrm{HS}}^{-1} \Re M$. The random variable $\left|\xi_{\ell}+u_{\ell}\right|^{2}$ is non-central chi-squared with two degrees of freedom and has the density

$$
f_{\ell}(x)=\exp \left(-\left(x+\left|u_{\ell}\right|^{2}\right)\right) \boldsymbol{I}_{0}\left(2\left|u_{\ell}\right| \sqrt{x}\right)
$$

on $\mathbb{R}_{+}$, where $\boldsymbol{I}_{0}$ is the modified Bessel function. Since these densities are bounded by one [1], we can get the result from [28, Th. 1.2].
4. By the restriction property of Lévy's anti-concentration function,

$$
\mathcal{L}\left(\Sigma^{1 / 2} \xi, \sqrt{m} t\right) \leqslant \sup _{\left(d_{0}, \ldots, d_{m-1}\right) \in \mathbb{C}^{m}} \mathbb{P}\left[\sum_{\ell=0}^{m-1}\left|s_{\ell}(\Sigma)^{1 / 2} \xi_{\ell}-d_{\ell}\right|^{2} \leqslant m t^{2}\right]
$$

We furthermore use the following well-known tensorization result (see [26, Lemma 2.2]): Suppose $\left\{w_{0}, \ldots, w_{m-1}\right\}$ is a collection of independent non-negative random variables such that there is a constant $c>0$ such that for each $t \geqslant 0$,

$$
\mathbb{P}\left[w_{\ell} \leqslant t\right] \leqslant c t
$$

Then there exists a constant $c^{\prime}>0$ so that

$$
\mathbb{P}\left[\sum_{\ell=0}^{m-1} w_{\ell}^{2} \leqslant m t^{2}\right] \leqslant\left(c^{\prime} t\right)^{m}
$$

In the present case, for each $\ell \in[m]$ and each $d_{\ell} \in \mathbb{C}, w_{\ell}^{2}=\left[\left|s_{\ell}(\Sigma)^{1 / 2} \xi_{\ell}-d_{\ell}\right|\right]^{2}$ is a noncentral chi-squared random variable with two degrees of freedom and has a density bounded by a constant that depends only on $\alpha$, as can be checked from the previous item. Thus, the tensorization argument applies and the result follows.
5. By a singular value decomposition of $M$, we obtain that

$$
\mathbb{P}\left[\|M \xi\|^{2} \geqslant t\|M\|_{\mathrm{HS}}^{2}\right]=\mathbb{P}\left[\sum_{\ell=0}^{N-1} \sigma_{\ell}^{2}\left|\xi_{\ell}\right|^{2} \geqslant t\right]
$$

where $\sum \sigma_{\ell}^{2}=1$. Writing

$$
\mathbb{P}\left[\sum_{\ell} \sigma_{\ell}^{2}\left|\xi_{\ell}\right|^{2} \geqslant t\right]=\mathbb{P}\left[\exp \left(\sum_{\ell} \sigma_{\ell}^{2}\left|\xi_{\ell}\right|^{2} / 2\right) \geqslant \exp (t / 2)\right]
$$

and following the arguments given for Item 1, we easily get the result.
Lemma 15 is proven.

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