## Breakdown of universality in multi-cut matrix models

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# Breakdown of universality in multi-cut matrix models 

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#### Abstract

We solve the puzzle of the disagreement between orthogonal polynomials methods and mean-field calculations for random $N \times N$ matrices with a disconnected eigenvalue support. We show that the difference does not stem from a $\mathbb{Z}_{2}$ symmetry breaking, but from the discreteness of the number of eigenvalues. This leads to additional terms (quasiperiodic in $N$ ) which must be added to the naive mean-field expressions. Our result invalidates the existence of a smooth topological large- $N$ expansion and some postulated universality properties of correlators. We derive the large$N$ expansion of the free energy for the general two-cut case. From it we rederive by a direct and easy mean-field-like method the two-point correlators and the asymptotic orthogonal polynomials. We extend our results to any number of cuts and to non-real potentials.


## 1. Introduction

Random matrix models have been introduced in order to give an approximate statistical description for quantum systems involving disorder, chaos or complexity. In particular, they allow a study of the level distribution and response to external fields. Those models are described by a matrix (Hamiltonian, transfer matrix or scattering matrix) of large size $N$, which is too complicated to be diagonalized exactly, and for which only statistical observations of the spectrum are available (see [1,2] for a review on random matrix theory (RMT)).

Most of the quantities of interest and observables are related to the short-range (in energy scale) behaviour of the spectrum (indeed small energies correspond to long-time evolution, i.e. to equilibrium thermodynamical properties). This is why the short-range correlation functions are the most studied.

At first, the simplest models assumed a Gaussian weight for the random matrix and gave good agreement with observations, provided that the ensemble of matrices (Hermitian, orthogonal, quaternionic, etc) has the required symmetries (time reversibility, etc) [2, 3].

It has been conjectured [4], and observed experimentally [5] and numerically [6], that the correlation functions of the spectrum possess universal properties at short range, which do not depend on the probability weight, Gaussian or not. This universality has been proved for a wide range of models by several approaches [3, 7, 8], though a very general proof and the exact hypothesis which lead to it are still under investigations.

Moreover, the long-range correlation functions appear to also share some universal properties, which depend on the probability weight, only through a few parameters [7,9]. The most striking example is the two-point correlation function that we shall discuss below.

In the following we will restrict our attention to the so-called Hermitian one-matrix-model $\dagger$ $[2,10]$ (Hermiticity corresponds to a system with broken time reversibility, for instance in the presence of a magnetic field).

We consider a Hermitian matrix $M$ of size $N \times N$ with a probability law of the form

$$
\mathcal{P}(M)=\mathrm{e}^{-N \operatorname{tr} V(M)}
$$

where $V$ is a polynomial potential bounded from below.
We wish to study the statistical properties of the eigenvalues $\left(\lambda_{1}, \ldots, \lambda_{N}\right)$ of $M$ in the large$N$ limit, in particular, the density of eigenvalues $\rho(\lambda)$ and the correlation function $R(\lambda, \mu)$, which measures the probability that two of the eigenvalues take the values $\lambda$ and $\mu$.

Roughly speaking, the eigenvalues tend to occupy a finite interval centred around the bottom of the potential well, and in the large- $N$ limit, the density of eigenvalues $\rho(\lambda)$ is a continuous function with a compact support.

The simplest case, where the support is connected, known as the 'one-cut case', has been studied extensively [2]. It is found that the density $\rho(\lambda)$ is not universal, but depends on the details of the potential $V$, while the connected correlation function $R_{c}(\lambda, \mu)$ is universal in the short-range regime $(|\lambda-\mu| \sim \mathrm{O}(1 / N)$ ), but also in the long-range regime $(|\lambda-\mu| \sim \mathrm{O}(1))$ once the short-range oscillations (of period $\sim \mathrm{O}(1 / N)$ ) have been smoothed out.

What happens when the potential $V$ possesses several wells, of approximately the same depths? Then the density $\rho(\lambda)$ has a disconnected support, $\left[a_{1}, b_{1}\right] \cup\left[a_{2}, b_{2}\right] \cup \cdots \cup\left[a_{s}, b_{s}\right]$, each interval $\left[a_{i}, b_{i}\right]$ being centred around one well of the potential $V$. This case is known as the multicut (or multiband) case (here $s$ cuts).

In the multicut case the density is still not universal, whereas the two-point correlation function is universal in the short-range regime and seems to have some universal properties in the long-range regime after smoothing: in [11] an explicit form of the two-point connected correlation function was given, and is claimed to be universal: indeed, according to the authors of [11] 'it depends only on the number of connected components of the support and on the position of the endpoints, but not on the potential'. However, more recently several authors [12-15] have studied the two-cut case $s=2$ : they concentrated on the case of an even potential $V$ (the two cuts are thus symmetric $[a, b]$ and $[-b,-a]$ ). Using an ansatz for the asymptotic expression of orthogonal polynomials in the large- $N$ limit, and rederiving the twopoint function from this ansatz, they observed that the connected correlation function is still universal in the short-distance regime (which was expected), but more surprisingly, that the smoothed connected correlation function in the long-range regime depends on the parity of $N$ ( $N$ being the size of the matrix). This seems to contradict the former result of [11]!

In this paper, we will solve this paradox.
We will show that the semiclassical method of [11] gives the two-point connected correlation function only up to an additional non-universal term, which is already present in the free energy, but subdominant at large $N$ in this case. We correct the semiclassical argument of [11] and give a simple (and physically appealing) derivation of the origin of the additional term, that we compute explicitly. This allows us to recover the results of $[13,15]$ for the symmetric case, and to generalize them to non-symmetric potentials, without using orthogonal polynomials.

Using the same semiclassical argument, we are able to derive large- $N$ asymptotics for the orthogonal polynomials, recovering the results of $[13,15]$ as well as the general $s$ cuts asymptotics which appeared recently in the mathematical literature [17], and to extend these results to the case of complex potentials.
$\dagger$ The other ensembles are of course worth considering, but this one is the simplest.

The effect leading to the new term in the semiclassical calculation is simple enough to be explained briefly in this introductory section.

In [11], the free energy $F$ of the matrix model is derived by a saddle-point approximation. In particular, one has to extremize the action with respect to variations of the number $n_{i}$ $(i=1, \ldots, s)$ of eigenvalues in each connected part of the support, or, in other words, with respect to the occupation ratio $x_{i}=n_{i} / N$ :

$$
\begin{equation*}
F=F\left(x_{c}\right) \quad \text { where }\left.\quad \frac{\partial F}{\partial x}\right|_{x=x_{c}}=0 \tag{1.1}
\end{equation*}
$$

However, here one has missed the crucial fact that $n_{i}=N x_{i}$ are not real numbers but integers. When $N x_{c}$ is not an integer, the extremum of $F(x)$ is never reached, and the saddle-point approximation has to be slightly modified. Roughly speaking the discrete sum cannot be approximated by an integral:

$$
\begin{equation*}
\sum_{n} \mathrm{e}^{-g\left(n-N x_{c}\right)^{2}} \neq \int \mathrm{d} x \mathrm{e}^{-N^{2} g\left(x-x_{c}\right)^{2}} \tag{1.2}
\end{equation*}
$$

The discrete sum actually depends on how far from an integer $N x_{c}$ is. For instance, in the symmetric case, we have $x_{c}=\frac{1}{2}$, and the result depends on the parity of $N$.

We will show that (as expected from equation (1.2)), in general, the result involves elliptic theta functions depending on $N x_{c}$, thus leading to a quasi-periodic dependence on $N$. This effect is of order $N^{-2}$ for the free energy, but is of order one for the computation of the orthogonal polynomials and for the correlation functions. It implies, in particular, that there is no regular large- $N$ topological expansion (involving only power series in $N^{-2}$ ) for the two-cut matrix model. This should not come as a surprise, indeed 't Hooft's argument [10, 16] is not valid in this case, since it relies on a perturbative expansion around a quadratic potential, which can never have several wells.

We will find that the short-range correlation function is universal, while the long-range smoothed correlation depends on $N$ quasi-periodically.

The paper is divided as follows. In section 2 we introduce the method and notation for the two-cut model, and we compute the free energy. In section 3 we derive the two-point correlation function, and we recover the expression of $[13,15,17]$ in the symmetric case. In section 4, we give an asymptotic expression for the orthogonal polynomials, which we use to rederive the universal short-range properties of the spectrum, as well as the smoothed long-range two-point correlation function.

The generalizations to a complex potential or to an arbitrary number of cuts are presented in appendices B and C. Appendix A contains some technical details and appendix B is a summary of some relationships between elliptical functions in case the reader is not familiar with them.

## 2. The free energy

### 2.1. Basics

We start from the standard Hermitian matrix model defined by the partition function

$$
\begin{equation*}
Z[V ; N]=\int \mathrm{d}_{N}[M] \mathrm{e}^{-N \operatorname{tr} V(M)} \tag{2.1}
\end{equation*}
$$

where $N$ is the dimension of the matrix $M, V$ is an analytic-in general, polynomial-and for the moment real-function, and $\mathrm{d}_{N}[M]$ is the standard $U(N)$ invariant measure over Hermitian
matrices

$$
\begin{equation*}
\mathrm{d}_{N}[M]=\prod_{i=1}^{N} \mathrm{~d} M_{i i} \prod_{1 \leqslant i<j \leqslant N} 2 \mathrm{~d} \operatorname{Re}\left(M_{i j}\right) \mathrm{d} \operatorname{Im}\left(M_{i j}\right) \tag{2.2}
\end{equation*}
$$

Integrating out the 'angular part' of $M, Z$ can be rewritten as an integral over the $N$ eigenvalues $\lambda_{1}, \ldots, \lambda_{N}$ of $M$ [2]
$Z[V ; N]=C_{N} \tilde{Z}[V ; N]$
$\tilde{Z}[V ; N]=\int \prod_{k=1}^{N} \mathrm{~d} \lambda_{k} \mathrm{e}^{-N \sum_{k} V\left(\lambda_{k}\right)} \prod_{k<l}\left(\lambda_{k}-\lambda_{l}\right)^{2}=\int \prod_{k=1}^{N} \mathrm{~d} \lambda_{k} \mathrm{e}^{-N^{2} S\left(\lambda_{1}, \ldots, \lambda_{N}\right)}$
with the measure factor

$$
\begin{equation*}
C_{N}=\operatorname{Vol}\left[\frac{\mathrm{U}(N)}{\mathrm{U}(1)^{N} \times \mathfrak{S}_{N}}\right]=\frac{1}{N!} \prod_{K=1}^{N} \frac{(2 \pi)^{K-1}}{\Gamma(K)} \tag{2.5}
\end{equation*}
$$

and with the action $S\left(\lambda_{k}\right)$

$$
\begin{equation*}
S\left(\lambda_{1}, \ldots, \lambda_{N}\right)=\frac{1}{N} \sum_{k=1}^{N} V\left(\lambda_{k}\right)-\frac{1}{2 N^{2}} \sum_{1 \leqslant k \neq l \leqslant N} \ln \left(\lambda_{k}-\lambda_{l}\right)^{2} . \tag{2.6}
\end{equation*}
$$

In the simplest 'one-cut' case, corresponding in particular to a concave potential, it is known [10] that the free energy $F$ defined as $\dagger$

$$
\begin{equation*}
Z[V ; N]=\left(\frac{2 \pi}{N}\right)^{N^{2} / 2} \mathrm{e}^{-F[V ; N]} \tag{2.7}
\end{equation*}
$$

has a topological large- $N$ expansion

$$
\begin{equation*}
F=N^{2} F_{0}+F_{1}+N^{-2} F_{2}+\cdots \tag{2.8}
\end{equation*}
$$

obtained for instance by reorganizing the perturbative expansion according to the topology of the Feynmann diagrams. The large- $N$ limit (planar limit) can be described by a 'master field' configuration where the eigenvalues are described by a continuous density $\rho(\lambda)$ with a connected compact support $\mathcal{C}=[a, b]$, with the constraints

$$
\begin{equation*}
\int_{\mathcal{C}} \mathrm{d} \lambda \rho(\lambda)=1 \quad \text { and } \quad \rho(\lambda) \geqslant 0 \quad \text { if } \lambda \in \mathcal{C} \tag{2.9}
\end{equation*}
$$

and the action (2.6) becomes

$$
\begin{equation*}
S[\rho]=\int_{\mathcal{C}} \mathrm{d} \lambda V(\lambda) \rho(\lambda)-\int_{\mathcal{C} \times \mathcal{C}} \mathrm{d} \lambda \mathrm{~d} \mu \rho(\lambda) \rho(\mu) \ln |\lambda-\mu| . \tag{2.10}
\end{equation*}
$$

The leading term of the free energy $F_{0}$ can be obtained by the saddle-point method: the effective action $S[\rho]$ is extremized for a continuous distribution $\rho_{c}$ and we have simply

$$
\begin{equation*}
F_{0}=S\left[\rho_{c}\right] \tag{2.11}
\end{equation*}
$$

(up to an additive (potential-independent) constant). To compute $\rho_{c}$ we include the constraint (2.9) in the effective action by a Lagrange multiplier $\Gamma$

$$
\begin{equation*}
\bar{S}[\rho]=S[\rho]+\Gamma\left(1-\int_{\mathcal{C}} \rho\right) \tag{2.12}
\end{equation*}
$$

$\dagger Z$ is normalized here so that $F$ is zero for the Gaussian model $V=\frac{1}{2} M^{2}$

The saddle-point equation for $\rho$ reads as

$$
\begin{equation*}
\frac{\partial \bar{S}}{\partial \rho(\lambda)}=V(\lambda)-2 \int_{\mathcal{C}} \mathrm{d} \mu \rho(\mu) \ln (|\lambda-\mu|)-\Gamma=0 \quad \forall \lambda \in \mathcal{C} \tag{2.13}
\end{equation*}
$$

which simply means that the real part of the effective potential

$$
\begin{equation*}
V_{\text {eff }}(\lambda)=V(\lambda)-2 \int_{\mathcal{C}} \mathrm{d} \mu \rho(\mu) \ln (\lambda-\mu) \tag{2.14}
\end{equation*}
$$

is constant on the EV support $\mathcal{C}$, and equal to $\Gamma$. The derivative of (2.13) with respect to $\lambda$ gives the well known equation

$$
\begin{equation*}
\operatorname{Re}\left(\omega_{0}(\lambda)\right)=f_{\mathcal{C}} \mathrm{d} \mu \rho(\mu) \frac{1}{\lambda-\mu}=V^{\prime}(\lambda) / 2 \quad \forall \lambda \in \mathcal{C} \tag{2.15}
\end{equation*}
$$

where $\omega_{0}$ is the large- $N$ resolvent

$$
\begin{equation*}
\omega_{0}(\lambda)=\lim _{N \rightarrow \infty}\left\langle\frac{1}{N} \operatorname{Tr}\left[\frac{1}{\lambda-M}\right]\right\rangle=\int_{\mathcal{C}} \mathrm{d} \mu \frac{\rho(\mu)}{\lambda-\mu} \tag{2.16}
\end{equation*}
$$

Finally, let us recall that in the one-cut case $\mathcal{C}=[a, b]$, if the potential $V$ is a polynomial of degree $P, \omega$ is of the form
$\omega_{0}(\lambda)=\frac{1}{2} V^{\prime}(\lambda)-\frac{1}{2} M(\lambda) \sqrt{\sigma(\lambda)} \quad$ with $\quad \sigma(\lambda)=(\lambda-a)(\lambda-b)$
where $M(\lambda)$ is a polynomial with degree $P-2 . a, b$ and $M$ are entirely determined by the constraint that

$$
\begin{equation*}
\omega_{0}(\lambda)=\lambda^{-1}+\mathcal{O}\left(\lambda^{-2}\right) \quad \text { for } \quad \lambda \rightarrow \infty . \tag{2.18}
\end{equation*}
$$

The EV density is given by the discontinuity of $\omega$

$$
\begin{equation*}
\rho(\lambda)=\frac{\mathrm{i}}{2 \pi}\left[\omega\left(\lambda+\mathrm{i} 0_{+}\right)-\omega\left(\lambda-\mathrm{i} 0_{+}\right)\right]=\frac{M(\lambda) \sqrt{|\sigma(\lambda)|}}{2 \pi} . \tag{2.19}
\end{equation*}
$$

### 2.2. The two-cut case

2.2.1. Mean field. If the potential $V$ is real but has more than one minimum, the large- $N$ limit may be described by an EV distribution on several disconnected intervals. For simplicity we shall first consider the case where there are two intervals

$$
\begin{equation*}
\mathcal{C}=\mathcal{C}_{1} \cup \mathcal{C}_{2} \quad \mathcal{C}_{1}=[a, b] \quad \mathcal{C}_{2}=[c, d] \quad a<b<c<d \tag{2.20}
\end{equation*}
$$

In this case, as we shall see, there is no topological large- $N$ expansion, even for the free energy $F$. As shown in $[18,19]$, to describe the large- $N$ limit, we have to consider as an additional variable the 'average' proportion of eigenvalues $x_{1}=n_{1} / N$ and $x_{2}=n_{2} / N$ in each interval $\mathcal{C}_{1}$ and $\mathcal{C}_{2}$, and introduce the associated Lagrange multipliers $\Gamma_{1}$ and $\Gamma_{2}$ for the constraints

$$
\begin{equation*}
x_{\alpha}=\int_{\mathcal{C}_{\alpha}} \rho(\lambda) \mathrm{d} \lambda \quad \alpha=1,2 \tag{2.21}
\end{equation*}
$$

The effective action (2.12) now reads, with

$$
\begin{align*}
& x=x_{1}  \tag{2.22}\\
& \bar{S}[\rho ; x]=S[\rho]+\sum_{\alpha=1}^{2} \Gamma_{\alpha}\left(x_{\alpha}-\int_{\mathcal{C}_{\alpha}} \rho(\lambda) \mathrm{d} \lambda\right) \quad x_{1}+x_{2}=1 \tag{2.23}
\end{align*}
$$

with $S[\rho]$ given by (2.10) as before. The saddle-point equation with respect to $\rho(\lambda)$ gives as before the equation (2.13), which implies that the effective potential defined by (2.14) is constant on each interval

$$
\begin{equation*}
V_{\text {eff }}(\lambda)=\Gamma_{\alpha} \quad \text { when } \quad \lambda \in \mathcal{C}_{\alpha} \tag{2.24}
\end{equation*}
$$

but the corresponding EV density $\rho_{c}(\lambda)$ and the effective action $\bar{S}_{c}$ still depend explicitly on the EV proportion $x$, since we have

$$
\begin{equation*}
\bar{S}_{c}[x]=\frac{1}{2}\left(\int_{\mathcal{C}} \rho_{c}(\lambda) V(\lambda) \mathrm{d} \lambda+\sum_{\alpha} \Gamma_{\alpha} x_{\alpha}\right) . \tag{2.25}
\end{equation*}
$$

The saddle-point equation with respect to $x$ implies the equality of the effective potentials for each interval

$$
\begin{equation*}
\frac{\partial \bar{S}}{\partial x}=\Gamma_{1}-\Gamma_{2}=0 \tag{2.26}
\end{equation*}
$$

This fixes the value of $x$, and it is known that with this last equation the EV density $\rho_{c}$ is uniquely determined in the two-cut case. The large- $N$ free energy is then given simply by

$$
\begin{equation*}
F_{0}=\bar{S}\left[\rho_{c} ; x_{c}\right]=\bar{S}_{c}\left[x_{c}\right] . \tag{2.27}
\end{equation*}
$$

For an explicit polynomial potential $V$ of degree $P$, and for fixed $x$, the two-cut mean-field solution for the resolvent is
$\omega_{0}(\lambda, x)=\frac{1}{2} V^{\prime}(\lambda)-\frac{1}{2} M(\lambda) \sqrt{\sigma(\lambda)} \quad$ with $\quad \sigma(\lambda)=(\lambda-a)(\lambda-b)(\lambda-c)(\lambda-d)$
and $M(\lambda)$ a polynomial with degree $P-3$. The EV density $\rho(\lambda, x)$ is still given by the discontinuity of $\omega_{0}$. The coefficients of $M$ and the four end-points $a, b, c, d$ are entirely determined by the constraint that $\omega_{0}(\lambda) \simeq \lambda^{-1}$ when $\lambda \rightarrow \infty$ and by the fact that $x$ must be given by

$$
\begin{equation*}
x=\int_{a}^{b} \rho(\lambda, x) \mathrm{d} \lambda=\frac{1}{2 \pi} \int_{a}^{b}|M(\lambda, x)| \sqrt{|\sigma(\lambda)|} \mathrm{d} \lambda . \tag{2.29}
\end{equation*}
$$

Finally, the equation (2.26) which fixes $x=x_{c}$ reads
$0=V_{\text {eff }}(b)-V_{\text {eff }}(c)=\int_{b}^{c} \mathrm{~d} \lambda\left(2 \omega_{0}(\lambda, x)-V^{\prime}(\lambda)\right)=-\int_{b}^{c} \mathrm{~d} \lambda M(\lambda, x) \sqrt{|\sigma(\lambda)|}$.
2.2.2. Discreteness of number of EVs. This is sufficient if one is interested in the leading term in the large- $N$ limit (planar approximation). However, in order to understand the structure of the subdominant terms of the large- $N$ expansion, it turns out that we cannot neglect the fact that the number of $\mathrm{EV} n_{\alpha}=N x_{\alpha}$ in each interval $\mathcal{C}_{\alpha}$ must be an integer.

We consider the simple case where the potential $V$ is such that the support of the eigenvalues consists of two intervals. For simplicity we assume that the potential $V(z)$ has only two separate minima $z_{1}$ and $z_{2}$, though our argument is independent of the number of minima of $V$, but depends only on the number of intervals.

What has to be done is first to fix the number of eigenvalues $n_{1}=n$ (respectively, $n_{2}=N-n$ ) in the vicinity of $z_{1}$ (respectively, $z_{2}$ ) in the partition function (2.5) by writing

$$
\begin{equation*}
\tilde{Z}[V ; N]=\sum_{n=0}^{N} \frac{N!}{n!(N-n)!} \tilde{Z}[V ; n, N-n] \tag{2.31}
\end{equation*}
$$

where
$\tilde{Z}[V ; n, N-n]=\int_{-\infty}^{E} \prod_{i \leqslant n} \mathrm{~d} \lambda_{i} \int_{E}^{+\infty} \prod_{j>n} \mathrm{~d} \lambda_{j} \mathrm{e}^{-N \sum_{k} V\left(\lambda_{k}\right)} \prod_{k<l}\left(\lambda_{k}-\lambda_{l}\right)^{2}$
with $E$ a 'frontier', $b<E<c$, between the two semiclassical cuts $[a, b]$ and $[c, d]$. We now claim that each term of this discrete sum has a well defined large- $N$ topological expansion. Indeed, we can rewrite (2.32) as a matrix integral over two separate matrices: an $n_{1} \times n_{1}$ matrix $M_{1}$ with the $n_{1}=n \mathrm{EV}<E$ and an $n_{2} \times n_{2}$ matrix $M_{2}$ with the $n_{2}=N-n \mathrm{EV}>E$, as

$$
\begin{equation*}
\tilde{Z}[V ; n]=\frac{1}{C_{n} C_{N-n}} \int \mathrm{~d}_{n_{1}}\left[M_{1}\right] \int \mathrm{d}_{n_{2}}\left[M_{2}\right] \mathrm{e}^{-N \operatorname{Tr}\left(V\left(M_{1}\right)\right)-N \operatorname{Tr}\left(V\left(M_{2}\right)\right)+2 \operatorname{Tr}\left(\ln \left(M_{1} \otimes \mathrm{Id}-\operatorname{Id} \otimes M_{2}\right)\right)} . \tag{2.33}
\end{equation*}
$$

This last matrix integral has a topological large- $N$ expansion of the form (2.8) in the 't Hooft limit $N \rightarrow \infty, x=n / N$ fixed. This expansion is obtained by doing a classical perturbative expansion around the smallest minimum $z_{1}$ of $V$ for $M_{1}$ and around the largest minimum $z_{2}$ of $V$ for $M_{2}$, and by reorganizing the perturbative expansion according to the topology of the Feynmann diagrams, as usual (see appendix A for details). Taking into account carefully the measure factors $C_{n}$ and $C_{N-n}$, and using their large- $N$ asymptotics
$C_{N}=\frac{1}{N!}\left(\frac{2 \pi}{N}\right)^{N^{2} / 2} \mathrm{e}^{\frac{3}{4} N^{2}}(2 \pi)^{-N} N^{1 / 12} \operatorname{cst}\left(1+\mathcal{O}\left(N^{-1}\right)\right) \quad$ when $\quad N \rightarrow \infty$
(easily derived from the Stirling formula), we obtain that

$$
\begin{equation*}
Z[V ; N]=\left(\frac{2 \pi}{N}\right)^{N^{2} / 2} N^{-1 / 12} \sum_{n=0}^{N} \mathrm{e}^{-F[V ; N, x]} \tag{2.35}
\end{equation*}
$$

where each $F[V ; N, x]$ has a regular large- $N$ asymptotic expansion of the form

$$
\begin{equation*}
F[V ; N, x]=\sum_{h=0}^{\infty} N^{2-2 h} F_{h}[V, x] \quad \text { where } \quad x=n / N \tag{2.36}
\end{equation*}
$$

with each $F_{h}[V, x]$ a regular function of $x=n / N$. In particular, the leading large- $N$ term is given (up to an additive, $V$ and $x$ independent, constant) by the classical effective action (2.27)

$$
\begin{equation*}
F_{0}=\bar{S}\left[\rho_{c} ; x_{c}\right]=\bar{S}_{c}\left[x_{c}\right] . \tag{2.37}
\end{equation*}
$$

Finally, let us stress that although this decomposition depends on the arbitrary parameter $E$, since $E$ is in the interval ] $b, c$ [ where the density of eigenvalues is exponentially small with $N$, the integral (2.32) depends on $E$ only through exponentially small terms of order $\mathrm{e}^{-\mathrm{cst} \cdot N}$, which are 'non-perturbative' in the topological expansion equation (2.36).
2.2.3. Beyond mean field. We can now easily calculate the subleading terms of order $\mathcal{O}\left(N^{-2}\right)$ for the full partition function. In the large- $N$ limit we can approximate the sum (2.36) by

$$
\begin{equation*}
Z[V ; N] \propto \sum_{n=0}^{N} \mathrm{e}^{-N^{2} F_{0}[V ; x]-F_{1}[V ; x]+\cdots} \tag{2.38}
\end{equation*}
$$

If $x_{c}$ denotes the saddle point of $F_{0}[x]$ given by (2.26), the sum is dominated by the $n$ s such that

$$
\begin{equation*}
\left|n-N x_{c}\right|=\mathcal{O}(1) \tag{2.39}
\end{equation*}
$$

Thus we can still use a quadratic approximation for $F_{0}[x]$

$$
\begin{equation*}
Z[V ; N] \propto \mathrm{e}^{-\left(N^{2} F_{0}\left[V ; x_{c}\right]+F_{1}\left[V ; x_{c}\right]+\cdots\right)} \sum_{n} \mathrm{e}^{-\left(n-N x_{c}\right)^{2} F_{0}^{\prime \prime}\left[V ; x_{c}\right] / 2} \tag{2.40}
\end{equation*}
$$

where $F_{0}^{\prime \prime}=\partial^{2} F_{0} / \partial x^{2}$ and where the $\cdots$ represent terms of order $\mathcal{O}\left(N^{-2}\right)$. The last sum over $n$ gives simply an elliptic Jacobi theta function $\theta_{3}$

$$
\begin{equation*}
\sum_{n} \mathrm{e}^{-\left(n-N x_{c}\right)^{2} F_{0}^{\prime \prime}\left[V ; x_{c}\right] / 2}=\left(2 \pi F_{0}^{\prime \prime}\left[V ; x_{c}\right]\right)^{-1 / 2} \theta_{3}\left(N x_{c} \mid \tau\right) \tag{2.41}
\end{equation*}
$$

with modular parameter $\tau$ given by

$$
\begin{equation*}
\tau=\frac{2 \mathrm{i} \pi}{F_{0}^{\prime \prime}\left[V ; x_{c}\right]} \tag{2.42}
\end{equation*}
$$

and where the theta function is defined as

$$
\begin{equation*}
\theta_{3}(z \mid \tau)=\theta_{3}(z)=\sum_{n \in \mathbb{Z}} q^{n^{2}} \mathrm{e}^{2 \mathrm{i} \pi n z} \quad \text { with } \quad q=\mathrm{e}^{\mathrm{i} \pi \tau} \tag{2.43}
\end{equation*}
$$

It obeys the periodicity relations

$$
\begin{equation*}
\theta_{3}(z+1)=\theta_{3}(z) \quad \theta_{3}(z+\tau)=\mathrm{e}^{-\mathrm{i} \pi(2 z+\tau)} \theta_{3}(z) \tag{2.44}
\end{equation*}
$$

(For details on elliptic functions see, e.g., [26-28].) Eventually, we have for the free energy

$$
\begin{equation*}
F[V ; N]=N^{2} F_{0}\left[V, x_{c}\right]-\ln \left(\theta_{3}\left(N x_{c}\right)\right)+F_{1}\left(V ; x_{c}\right)+\frac{1}{2} \ln \left(2 \pi F_{0}^{\prime \prime}\left[V ; x_{c}\right]\right)+\mathcal{O}\left(N^{-2}\right) \tag{2.45}
\end{equation*}
$$

where $F_{1}$ is the torus contribution in the topological expansion of (2.33). The next terms of this expansion can be calculated along the same line.

Let us stress that this is not a topological expansion, since the second term $\ln \left(\theta_{3}\left(N x_{c}\right)\right)$, seemingly $\mathcal{O}(1)$ and contributing at the torus order, is not regular in $N$. Indeed, it is periodic in $x_{c}$ with period $1 / N$. When computing some observables or quantities of the matrix model, one must take derivatives of $F$ with respect to some parameters of the potential $V$. Since the saddle point $x_{c}$ depends implicitly on $V$, every derivative will give a factor $N$, and this term may become of the same order as the first term $N^{2} F_{0}\left[x_{c}\right]$ given by the planar limit. Note that the last two terms depend on $x$ and not on $N x$, and they will remain subdominant once we take derivatives of $F$.

### 2.3. The modular parameter

Finally, we can express the modular parameter $\tau$ defined by equation (2.42) simply in terms of the end-points $a, b, c, d$ of the support of EV For this purpose, we introduce the function $\sigma$

$$
\begin{equation*}
\sigma(\lambda)=(\lambda-a)(\lambda-b)(\lambda-c)(\lambda-d) \tag{2.46}
\end{equation*}
$$

and the function $u$

$$
\begin{equation*}
u(\lambda)=\frac{1}{2 K} \int_{d}^{\lambda} \frac{\mathrm{d} z}{\sqrt{\sigma(z)}} \tag{2.47}
\end{equation*}
$$

where $K$ is
$K=\int_{b}^{c} \frac{\mathrm{~d} z}{\sqrt{|\sigma(z)|}}=\frac{2}{\sqrt{(c-a)(d-b)}} K[m] \quad$ with $\quad m=\frac{(d-a)(c-b)}{(d-b)(c-a)}$
and $K[m]$ is the complete elliptic integral of the first kind. Similarly, we define
$K^{\prime}=\int_{a}^{b} \frac{\mathrm{~d} z}{\sqrt{|\sigma(z)|}}=\frac{2}{\sqrt{(c-a)(d-b)}} K\left[m^{\prime}\right] \quad$ with $\quad m^{\prime}=1-m$.
We shall show that the modular parameter $\tau$ of equation (2.42) coincides with the standard modular parameter of the torus associated with the mapping $u$, i.e. of the elliptic curve $y^{2}=\sigma(z)$. Indeed, $\tau$ is simply given by

$$
\begin{equation*}
\tau=\mathrm{i} \frac{K^{\prime}}{K}=\mathrm{i} \frac{K[1-m]}{K[m]} . \tag{2.50}
\end{equation*}
$$

So we have
$u(d)=0 \quad u(a)=\frac{1}{2} \quad u(b)=\frac{1}{2}(1+\tau) \quad u(c)=\frac{1}{2} \tau \quad u(\infty)=u_{\infty}$
and (see figure 1) $u$ maps the upper half $\lambda$-plane onto the half-periods rectangle $(1 / 2, \tau / 2)$ and the double-sheeted complex $\lambda$-plane onto the period rectangle $(1, \tau)$.


Figure 1. The upper half-plane is mapped onto a rectangle $\left(\frac{1}{2}, \frac{1}{2} \tau\right)$.
To show equation (2.50), we use the fact that in the two-cut case, if we fix $x$ (the EV ratio in the first cut) the semiclassical EV density (extrema of the effective action $\bar{S}$ ) is now a function $\rho(\lambda, x)$ of $\lambda$ and $x$, and the end-points $a, b, c, d$ depend on $x$. Therefore, the large- $N$ resolvent $\omega_{0}$ is of the form

$$
\begin{equation*}
\omega_{0}(\lambda, x)=\frac{1}{2} V^{\prime}(\lambda)-\frac{1}{2} M(\lambda, x) \sqrt{\sigma(\lambda)} \tag{2.52}
\end{equation*}
$$

with $M(\lambda, x)$ a polynomial with degree $P-3$ in $\lambda$ ( $P$ being the degree of $V$ ), entirely fixed by the constraints (2.18) and (2.29). Therefore, the partial derivative of $\omega_{0}(\lambda, x)$ with respect to $x$ is necessarily of the form

$$
\begin{equation*}
\frac{\partial \omega_{0}(\lambda, x)}{\partial x}=\frac{C}{\sqrt{\sigma(\lambda)}} \tag{2.53}
\end{equation*}
$$

with $C=C(\lambda, x)$ a priori a polynomial in $\lambda$. Since (2.18) still holds independently of $x$ we must have

$$
\begin{equation*}
\frac{\partial \omega_{0}(\lambda, x)}{\partial x}=\mathcal{O}\left(\lambda^{-2}\right) \quad \text { for } \quad \lambda \rightarrow \infty \tag{2.54}
\end{equation*}
$$

which implies that $C(\lambda, x)$ is of degree zero in $\lambda$, i.e. is a constant (depending only on $x$ )

$$
\begin{equation*}
C=C(x) \tag{2.55}
\end{equation*}
$$

This constant can be easily determined by using that

$$
\begin{equation*}
x=\int_{a}^{b} \rho(\lambda, x) \mathrm{d} \lambda=\int_{\mathcal{C}^{\prime}} \frac{\mathrm{d} \lambda}{2 \mathrm{i} \pi} \omega_{0}(\lambda, x) \tag{2.56}
\end{equation*}
$$

with $\mathcal{C}^{\prime}$ a clockwise contour encircling the interval $[a, b]$. Therefore, we have

$$
\begin{equation*}
\frac{\partial x}{\partial x}=1=\int_{\mathcal{C}^{\prime}} \frac{\mathrm{d} \lambda}{2 \mathrm{i} \pi} \frac{C}{\sqrt{\sigma(\lambda)}}=-\frac{C K^{\prime}}{\pi} \Rightarrow C=-\frac{\pi}{K^{\prime}} \tag{2.57}
\end{equation*}
$$

with $K^{\prime}$ the half-period defined in equation (2.49). Now we use equations (2.26), (2.27) and the definition of the effective potential $V_{\text {eff }}$ of equation (2.14) to write the derivative of the free energy with respect to $x$ as

$$
\begin{equation*}
\frac{\partial F_{0}}{\partial x}=V_{\text {eff }}(b)-V_{\text {eff }}(c)=\int_{b}^{c} \mathrm{~d} \lambda\left(2 \omega_{0}(\lambda)-V^{\prime}(\lambda)\right) \tag{2.58}
\end{equation*}
$$

Now we take the derivative with respect to $x$ of this equation and obtain

$$
\begin{equation*}
F_{0}^{\prime \prime}=\frac{\partial^{2} F_{0}}{\partial x^{2}}=2 \int_{b}^{c} \mathrm{~d} \lambda \frac{\partial \omega_{0}(\lambda)}{\partial x}=2 \int_{b}^{c} \mathrm{~d} \lambda \frac{C}{\sqrt{\sigma(\lambda)}}=-2 C K=\frac{2 \pi K}{K^{\prime}} . \tag{2.59}
\end{equation*}
$$

Using equation (2.42) we thus obtain the result (2.50).

## 3. Two-point correlation function

### 3.1. The basic formula

As a first application we compute the large- $N$ smoothed connected two-point correlation function (first obtained by [13, 15, 17]), defined as
$\omega^{c}(\lambda, \mu)=\left\langle\operatorname{Tr}\left[\frac{1}{\lambda-M}\right] \operatorname{Tr}\left[\frac{1}{\mu-M}\right]\right\rangle-\left\langle\operatorname{Tr}\left[\frac{1}{\lambda-M}\right]\right\rangle\left\langle\operatorname{Tr}\left[\frac{1}{\mu-M}\right]\right\rangle$.
Adding source terms to the potential of the form

$$
\begin{equation*}
V_{\epsilon_{\lambda}}=V(z)-\epsilon_{\lambda} \frac{1}{\lambda-z} \quad V_{\epsilon_{\lambda}, \epsilon_{\mu}}(z)=V(z)-\epsilon_{\lambda} \frac{1}{\lambda-z}-\epsilon_{\mu} \frac{1}{\mu-z} \tag{3.2}
\end{equation*}
$$

we have

$$
\begin{equation*}
\omega^{c}(\lambda, \mu)=-\left.\frac{1}{N^{2}} \frac{\partial}{\partial \epsilon_{\lambda}} \frac{\partial}{\partial \epsilon_{\mu}} F\left[V_{\epsilon_{\lambda}, \epsilon_{\mu}}, N\right]\right|_{\epsilon_{\lambda}=\epsilon_{\mu}=0} \tag{3.3}
\end{equation*}
$$

and for the resolvent (one-point function)

$$
\begin{equation*}
\omega(\lambda)=\frac{1}{N}\left\langle\operatorname{Tr}\left[\frac{1}{\lambda-M}\right]\right\rangle=-\left.\frac{1}{N^{2}} \frac{\partial}{\partial \epsilon_{\lambda}} F\left[V_{\epsilon_{\lambda}}, N\right]\right|_{\epsilon_{\lambda}=0} . \tag{3.4}
\end{equation*}
$$

If the $\epsilon_{\lambda} \mathrm{s}$ are small and the $\lambda \mathrm{s}$ are not too close to the cuts, the mean-field solution is still a two-cut EV distribution, with $x_{c}=x_{c}(\epsilon)$ an explicit function of the $\lambda \mathrm{s}$. So from equation (2.45) for the free energy we have for the two-point function
$\omega^{c}(\lambda, \mu)=\left[-\frac{\partial}{\partial \epsilon_{\lambda}} \frac{\partial}{\partial \epsilon_{\mu}} F_{0}\left[V_{\epsilon_{\lambda}, \epsilon_{\mu}}\right]+\frac{\partial x_{c}}{\partial \epsilon_{\lambda}} \frac{\partial x_{c}}{\partial \epsilon_{\mu}}\left[\ln \left(\theta_{3}\left(N x_{c}\right)\right)\right]^{\prime \prime}\right]_{\epsilon_{\lambda}=\epsilon_{\mu}=0}+\mathcal{O}\left(N^{-1}\right)$.
The first term in the right-hand side of (3.5) is the mean-field contribution already calculated in $[11,20]$, the second term involving a second derivative of an elliptic function, characterizes the multi-cut solution.

### 3.2. The mean-field contribution

For completeness let us first rederive the mean-field contribution of [11,20]. Taking the derivative with respect to $\epsilon_{\lambda}$ we obtain the mean-field resolvent for the potential $V_{\epsilon_{\mu}}$

$$
\begin{equation*}
-\left.\frac{\partial}{\partial \epsilon_{\lambda}} F_{0}\left[V_{\epsilon_{\lambda}, \epsilon_{\mu}}\right]\right|_{\epsilon_{\lambda}=0}=\omega_{0}\left(\lambda ; V_{\epsilon_{\mu}}\right) \tag{3.6}
\end{equation*}
$$

which must be of the form

$$
\begin{equation*}
\omega_{0}\left(z ; V_{\epsilon_{\mu}}\right)=\frac{1}{2}\left[V^{\prime}(z)-\frac{\epsilon_{\mu}}{(z-\mu)^{2}}+\frac{M(z) \sqrt{\sigma(z)}}{(z-\mu)^{2}}\right] \tag{3.7}
\end{equation*}
$$

with $M(z)$ a polynomial of degree $P-1$ (here both the coefficients of $M$ and of $\sigma$ depend on $\mu$ and $\epsilon_{\mu}$ ). In addition to the $P-2$ constraints (coming from (2.18) and (2.30)) $\omega_{0}\left(z ; V_{\epsilon_{\mu}}\right)$ must be regular at $z=\mu$. This determines $M$ entirely. Taking the derivative with respect to $\epsilon_{\mu}$ and using the symmetry $\lambda \leftrightarrow \mu$ we obtain

$$
\begin{equation*}
\omega_{0}^{c}(\lambda, \mu)=-\frac{\partial}{\partial \epsilon_{\lambda}} \frac{\partial}{\partial \epsilon_{\mu}} F_{0}=-\frac{1}{2} \frac{1}{(\lambda-\mu)^{2}}\left[1+\frac{Q(\lambda, \mu)}{\sqrt{\sigma(\lambda)} \sqrt{\sigma(\mu)}}\right] \tag{3.8}
\end{equation*}
$$

with $Q(\lambda, \mu)$ a symmetric polynomial in $\lambda$ and in $\mu$. The constraints on $\omega_{0}^{c}$ are:
(a) $\omega_{0}^{c}=\mathcal{O}\left(\lambda^{-2}\right)$ as $\lambda \rightarrow \infty$ which implies that $Q$ is of degree at most two;
(b) $\omega_{0}^{c}$ is regular at $\lambda=\mu$ which implies that $Q(\lambda, \mu)=-\sigma((\lambda+\mu) / 2)+\mathcal{O}\left((\lambda-\mu)^{2}\right)$;
(c) finally, the equality of the effective potential on the two cuts implies that

$$
\begin{equation*}
\int_{b}^{c} \omega_{0}^{c}(\lambda, \mu) \mathrm{d} \mu=0 \tag{3.9}
\end{equation*}
$$

Conditions (a) and (b) fix uniquely $Q$

$$
Q(\lambda, \mu)=-\frac{1}{2}\left[\begin{array}{c}
(\lambda-a)(\mu-b)(\mu-c)(\lambda-d)  \tag{3.10}\\
+(\mu-a)(\lambda-b)(\lambda-c)(\mu-d)
\end{array}\right]+S(\lambda-\mu)^{2}
$$

up to a constant $S$ fixed by condition (c), which is found to be

$$
\begin{equation*}
S=-\frac{1}{2}(c-a)(d-b) \frac{E[m]}{K[m]} \quad \text { with } \quad m=\frac{(c-b)(d-a)}{(c-a)(d-b)} \tag{3.11}
\end{equation*}
$$

and where $K[m]$ and $E[m]$ are the standard elliptic integrals of the first and second kind.

### 3.3. The non-regular contribution

In order to compute the second contribution to (3.5), we simply need $\frac{\partial x_{c}}{\partial \epsilon_{\lambda}}$. Since $x_{c}$ is fixed by the constraint $\frac{\partial F_{0}}{\partial x}=0$ we can write

$$
\begin{equation*}
\frac{\partial x_{c}}{\partial \epsilon_{\lambda}}=-\frac{\partial^{2} F_{0}}{\partial x \partial \epsilon_{\lambda}} / \frac{\partial^{2} F_{0}}{\partial x^{2}} \tag{3.12}
\end{equation*}
$$

Using the results of subsection 2.3 we have

$$
\begin{equation*}
\frac{\partial^{2} F_{0}}{\partial x^{2}}=\frac{2 \pi K}{K^{\prime}} \quad \text { and } \quad \frac{\partial^{2} F_{0}}{\partial x \partial \epsilon_{\lambda}}=-\frac{\partial}{\partial x} \omega_{0}(\lambda, x)=\frac{\pi}{K^{\prime}} \frac{1}{\sqrt{\sigma(\lambda)}} \tag{3.13}
\end{equation*}
$$

So we have eventually

$$
\begin{equation*}
\frac{\partial x_{c}}{\partial \epsilon_{\lambda}}=-\frac{1}{2 K \sqrt{\sigma(\lambda)}} \tag{3.14}
\end{equation*}
$$

and the second non-regular term is

$$
\begin{equation*}
\frac{\partial x_{c}}{\partial \epsilon_{\lambda}} \frac{\partial x_{c}}{\partial \epsilon_{\mu}}\left[\ln \left(\theta_{3}\left(N x_{c} \mid \tau\right)\right)\right]^{\prime \prime}=\frac{1}{2 K \sqrt{\sigma(\lambda)}} \frac{1}{2 K \sqrt{\sigma(\mu)}}\left[\ln \left(\theta_{3}\left(N x_{c} \mid \tau\right)\right)\right]^{\prime \prime} \tag{3.15}
\end{equation*}
$$

with $K$ defined by equation (2.48). Using standard relations on elliptic functions, this can be rewritten as

$$
\begin{equation*}
\frac{(c-a)(d-b)}{4 \sqrt{\sigma(\lambda)} \sqrt{\sigma(\mu)}}\left[-\frac{E[m]}{K[m]}+\operatorname{dn}^{2}\left(N x_{c}+\frac{1}{2}\right)\right] \tag{3.16}
\end{equation*}
$$

with

$$
\begin{equation*}
\operatorname{dn}(u)=\operatorname{dn}(2 K[m] u \mid m) \tag{3.17}
\end{equation*}
$$

where $\operatorname{dn}(u \mid m)$ is the Jacobi elliptic function dn . Its periods are $2 K[m]$ and $4 \mathrm{i} K\left[m^{\prime}\right]$, and $\operatorname{dn}^{2}(z)$ has periods 1 and $\tau$.

### 3.4. The final result

Combining (3.8), (3.10), (3.11) and (3.16) we obtain the final result for the two-point correlation function

$$
\begin{align*}
\omega^{c}(\lambda, \mu)=- & \frac{1}{4(\lambda-\mu)^{2}}\left[\left(1-\sqrt{\frac{(\lambda-a)(\lambda-b)(\mu-c)(\mu-d)}{(\mu-a)(\mu-b)(\lambda-c)(\lambda-d)}}\right)+(\lambda \leftrightarrow \mu)\right] \\
& -\frac{(c-a)(d-b)}{4 \sqrt{\sigma(\lambda)} \sqrt{\sigma(\mu)}} \operatorname{sn}^{2}\left(N x_{c}+\frac{1}{2}\right) . \tag{3.18}
\end{align*}
$$

We have used the relation $\mathrm{dn}^{2}(u)=1-m \operatorname{sn}^{2}(u)$, where similarly to (3.17) we note

$$
\begin{equation*}
\operatorname{sn}(u)=\operatorname{sn}(2 K[m] u \mid m) . \tag{3.19}
\end{equation*}
$$

Surprisingly, the ratio $E[m] / K[m]$ characteristic of the mean-field solution of $[11,20]$ has disappeared.

The smoothed two-point connected density correlator $\rho^{c}(\lambda, \mu)$, defined as
$\rho^{c}(\lambda, \mu)=\langle\operatorname{Tr}[\delta(\lambda-M)] \operatorname{Tr}[\delta(\mu-M)]\rangle-\langle\operatorname{Tr}[\delta(\lambda-M)]\rangle\langle\operatorname{Tr}[\delta(\mu-M)]\rangle$
can be obtained easily from the discontinuity of $\omega^{c}(\lambda, \mu)$. One obtains in the large- $N$ limit, if $\lambda$ and $\mu$ are on the support of EV

$$
\begin{align*}
& \rho^{c}(\lambda, \mu)=- \frac{1}{4 \pi^{2}}\left[\frac{1}{(\lambda-\mu)^{2}}\left(\sqrt{\left|\frac{(\lambda-a)(\lambda-b)(\mu-c)(\mu-d)}{(\mu-a)(\mu-b)(\lambda-c)(\lambda-d)}\right|}+\lambda \leftrightarrow \mu\right)\right. \\
&\left.+\varepsilon_{\lambda} \varepsilon_{\mu} \frac{(c-a)(d-b)}{\sqrt{|\sigma(\lambda)|} \sqrt{|\sigma(\mu)|}} \operatorname{sn}^{2}\left(N x_{c}+\frac{1}{2}\right)\right]  \tag{3.21}\\
& \varepsilon_{\lambda}=1 \quad \text { if } \quad \lambda \in[a, b] \quad-1 \quad \text { if } \quad \lambda \in[c, d] \tag{3.22}
\end{align*}
$$

and zero otherwise.
The new non-regular term $\operatorname{sn}^{2}\left(N x_{c}+\frac{1}{2}\right)$ is an even periodic function of $N x_{c}$ with period 1 which varies between 0 and 1 . Therefore, as $N$ varies, depending on the rationality or the irrationality of $x_{c}$, the two-point function will be varying with $N$ in a periodic or quasiperiodic way.

### 3.5. The symmetric case

It is now very easy to recover the results of $[13,15]$ for a symmetric potential. Indeed, if the potential $V$ is symmetric, the two cuts are also symmetric

$$
\begin{equation*}
a=-d \quad b=-c \tag{3.23}
\end{equation*}
$$

and we have automatically

$$
\begin{equation*}
x_{c}=\frac{1}{2} \tag{3.24}
\end{equation*}
$$

so that

$$
\operatorname{sn}^{2}\left(N x_{c}+\frac{1}{2}\right)= \begin{cases}\operatorname{sn}^{2}\left(\frac{1}{2}\right)=1 & \text { if } N \text { is even }  \tag{3.25}\\ \operatorname{sn}^{2}(0)=0 & \text { if } N \text { is odd }\end{cases}
$$

Equations (3.18) and (3.22) become
$\omega^{c}(\lambda, \mu)=-\frac{1}{2(\lambda-\mu)^{2}}\left[1-\frac{\left(a^{2}-\lambda \mu\right)\left(b^{2}-\lambda \mu\right)}{\sqrt{\sigma(\lambda)} \sqrt{\sigma(\mu)}}\right]-\frac{(-1)^{N}}{2} \frac{a b}{\sqrt{\sigma(\lambda)} \sqrt{\sigma(\mu)}}$
$\rho^{c}(\lambda, \mu)=\frac{1}{2 \pi^{2}} \frac{\varepsilon_{\lambda} \varepsilon_{\mu}}{\sqrt{|\sigma(\lambda)|} \sqrt{|\sigma(\mu)|}}\left(\frac{\left(a^{2}-\lambda \mu\right)\left(b^{2}-\lambda \mu\right)}{(\lambda-\mu)^{2}}-(-1)^{N} a b\right)$
with $\sigma(\lambda)=\left(\lambda^{2}-a^{2}\right)\left(\lambda^{2}-b^{2}\right)$.

### 3.6. The two-point function as an elliptic function

It is interesting to consider the two-point correlator in terms of the elliptic coordinates defined by equation (2.47)

$$
\begin{equation*}
u=u(\lambda) \quad v=u(\mu) \tag{3.28}
\end{equation*}
$$

Let us thus consider

$$
\begin{equation*}
\bar{\omega}^{c}(u, v)=\frac{\partial \lambda}{\partial u} \frac{\partial \mu}{\partial v} \omega^{c}(\lambda, \mu)=2 K \sqrt{\sigma(\lambda)} 2 K \sqrt{\sigma(\mu)} \omega^{c}(\lambda, \mu) . \tag{3.29}
\end{equation*}
$$

It is easy to see (from the properties of $\omega^{c}$ ) that $\bar{\omega}^{c}(u, v)$ satisfies:
(a) $\bar{\omega}^{c}(u, v)$ is a doubly periodic function of $u$ (and of $v$ ) with periods 1 and $\tau$;
(b) $\bar{\omega}^{c}(u, v)$ is regular at $u$ and $v=u(a), u(b), u(c), u(d)$ and $u_{\infty}$;
(c) $\bar{\omega}^{c}(u, v)$ is regular when $u=v$, but has a double pole at $u=-v$ (corresponding to the double pole of $\omega^{c}(\lambda, \mu)$ when $\lambda=\mu$ but with $\lambda$ in the first sheet and $\mu$ in the second sheet), with residue 1 .

This implies that $\bar{\omega}^{c}(u, v)$ is a Weirstrass elliptic function

$$
\begin{equation*}
\bar{\omega}^{c}(u, v)=\wp(u+v \mid \tau)+\text { constant } \tag{3.30}
\end{equation*}
$$

where the constant depends on $N x_{c}(\wp$ has periods 1 and $\tau)$. Using classical identities between the Weirstrass $\wp$ function and the Jacobi elliptic functions, it can be easily calculated. We find the remarkably simple result

$$
\begin{equation*}
\bar{\omega}^{c}(u, v)=\wp(u+v \mid \tau)-\wp\left(\left.N x_{c}+\frac{1}{2} \tau \right\rvert\, \tau\right) \tag{3.31}
\end{equation*}
$$

or equivalently

$$
\begin{equation*}
\bar{\omega}^{c}(u, v)=-\left[\ln \left(\theta_{1}(u+v \mid \tau)\right)\right]^{\prime \prime}+\left[\ln \left(\theta_{3}\left(N x_{c} \mid \tau\right)\right)\right]^{\prime \prime} . \tag{3.32}
\end{equation*}
$$

## 4. The orthogonal polynomials

Let us briefly recall some basic facts about the well known method of orthogonal polynomials [21], which is a powerful tool for studying the spectral properties of random matrices [2]. Asymptotic expressions for the orthogonal polynomials have been obtained recently [17] in the mathematical literature, by solving a Rieman-Hilbert problem. Here we will derive them from the free energy directly.

Consider the partition function (2.4):

$$
\begin{equation*}
\tilde{Z}=\int \mathrm{d} \lambda_{1} \ldots \mathrm{~d} \lambda_{N} \mathrm{e}^{-N \sum_{i} V\left(\lambda_{i}\right)} \prod_{i<j}\left(\lambda_{i}-\lambda_{j}\right)^{2} . \tag{4.1}
\end{equation*}
$$

The last term is a Vandermonde determinant [2]:

$$
\begin{equation*}
\prod_{i<j}\left(\lambda_{i}-\lambda_{j}\right)=\operatorname{det}_{i, j}\left(\left(\lambda_{i}\right)^{j-1}\right)=\operatorname{det}_{i, j}\left(\mathcal{P}_{j-1}\left(\lambda_{i}\right)\right) \tag{4.2}
\end{equation*}
$$

where the last equality is obtained by linearly mixing columns of the determinant, and holds for arbitrary monic polynomials $\mathcal{P}_{n}(\lambda)$ with leading coefficient $\mathcal{P}_{n}(\lambda)=\lambda^{n}+\cdots$.

The method of orthogonal polynomials consists of choosing a family of polynomials suitable for the computation of (4.1), namely, the family of polynomials orthogonal with respect to the weight $\exp [-N V(\lambda)]$ :

$$
\begin{equation*}
\int \mathrm{d} \lambda \mathcal{P}_{n}(\lambda) \mathcal{P}_{m}(\lambda) \mathrm{e}^{-N V(\lambda)}=h_{n} \delta_{n m} \tag{4.3}
\end{equation*}
$$

With this particular choice of polynomials, the integral (4.1) is merely

$$
\begin{equation*}
\tilde{Z}=N!\prod_{n=0}^{N-1} h_{n} \tag{4.4}
\end{equation*}
$$

and the joint probability density of all the eigenvalues takes the form of a Slater determinant:

$$
\begin{equation*}
R_{N}\left(\lambda_{1}, \ldots, \lambda_{N}\right)=\frac{1}{N!}\left(\operatorname{iet}_{\substack{0 \leqslant n<N \\ 1 \leqslant i \leqslant N}}\left[\psi_{n-1}\left(\lambda_{i}\right)\right]\right)^{2} \tag{4.5}
\end{equation*}
$$

where the wavefunctions $\psi_{n}(\lambda)=\frac{1}{\sqrt{h_{n}}} \mathcal{P}_{n}(\lambda) \mathrm{e}^{-\frac{1}{2} N V(\lambda)}$ are orthonormal.

### 4.1. The Kernel $K(\lambda, \mu)$

The square of a determinant can be rewritten as the determinant of a product:

$$
\left(\operatorname{det}_{n, i}\left(\psi_{n-1}\left(\lambda_{i}\right)\right)\right)^{2}=\operatorname{det}_{1 \leqslant i, j \leqslant N}\left[\sum_{n=0}^{N-1} \psi_{n}\left(\lambda_{i}\right) \psi_{n}\left(\lambda_{j}\right)\right]
$$

we are thus led to introduce the kernel $K(\lambda, \mu)$ [22]:

$$
\begin{equation*}
K(\lambda, \mu)=\frac{1}{N} \sum_{n=0}^{N-1} \psi_{n}(\lambda) \psi_{n}(\mu) \tag{4.6}
\end{equation*}
$$

In terms of which the joint density of eigenvalues is now a determinant

$$
\begin{equation*}
R_{N}\left(\lambda_{1}, \ldots, \lambda_{N}\right)=\frac{N^{N}}{N!} \operatorname{det}\left[K\left(\lambda_{i}, \lambda_{j}\right)\right] . \tag{4.7}
\end{equation*}
$$

The orthonormality properties of the polynomials imply the projection relations

$$
\begin{equation*}
\int \mathrm{d} \lambda K(\lambda, \lambda)=1 \quad \text { and } \quad \int \mathrm{d} \lambda K(\mu, \lambda) K(\lambda, \nu)=\frac{1}{N} K(\mu, \nu) \tag{4.8}
\end{equation*}
$$

which make any partial integration of (4.7) easy to perform (theorem of Dyson [23]).
In particular, the integration over $N-1$ eigenvalues gives the density of eigenvalues

$$
\rho\left(\lambda_{1}\right)=\int \mathrm{d} \lambda_{2} \ldots \mathrm{~d} \lambda_{N} R_{N}\left(\lambda_{1}, \ldots, \lambda_{N}\right)=K\left(\lambda_{1}, \lambda_{1}\right)
$$

and the integration over $N-2$ eigenvalues gives the correlation function

$$
\begin{aligned}
R_{2}\left(\lambda_{1}, \lambda_{2}\right) & =\int \mathrm{d} \lambda_{3} \ldots \mathrm{~d} \lambda_{N} R_{N}\left(\lambda_{1}, \ldots, \lambda_{N}\right) \\
& =\frac{N}{N-1}\left(K\left(\lambda_{1}, \lambda_{1}\right) K\left(\lambda_{2}, \lambda_{2}\right)-K\left(\lambda_{1}, \lambda_{2}\right) K\left(\lambda_{2}, \lambda_{1}\right)\right)
\end{aligned}
$$

In short,

$$
\begin{equation*}
\rho(\lambda)=K(\lambda, \lambda) \quad \rho(\lambda, \mu)=\left(K(\lambda, \lambda) K(\mu, \mu)-K(\lambda, \mu)^{2}\right) \tag{4.9}
\end{equation*}
$$

In addition, the Darboux-Christoffel theorem [24, 27], asserts that

$$
\begin{equation*}
K(\lambda, \mu)=\frac{1}{N h_{N-1}} \frac{\mathcal{P}_{N}(\lambda) \mathcal{P}_{N-1}(\mu)-\mathcal{P}_{N}(\mu) \mathcal{P}_{N-1}(\lambda)}{\lambda-\mu} \mathrm{e}^{-\frac{1}{2} N(V(\lambda)+V(\mu))} \tag{4.10}
\end{equation*}
$$

which means that we need to evaluate $\mathcal{P}_{n}$ only for $n=N$ and $n=N-1$.
Thus, we shall now aim at finding asymptotic expressions for the orthogonal polynomials $\mathcal{P}_{n}(\lambda)$, and the kernel $K(\lambda, \mu)$ in the large- $N$ limit, and $n$ close to $N$. This has been done in the one-cut case [7] and in the symmetric two-cut case [13, 14, 17]. Here we will generalize it to the non-symmetric case, with the method used in [24].

### 4.2. WKB approximation for the orthogonal polynomials $\mathcal{P}_{n}(\lambda)$

The orthogonal polynomials have the following integral representation (see appendix 1 of [24] or [21]):

$$
\begin{equation*}
\mathcal{P}_{n}(\lambda)=\frac{\int \mathrm{d} M_{n \times n} \operatorname{det}(\lambda-M) \mathrm{e}^{-N \operatorname{tr} V(M)}}{\int \mathrm{d} M_{n \times n} \mathrm{e}^{-N \operatorname{tr} V(M)}} \tag{4.11}
\end{equation*}
$$

where the integral is restricted to Hermitian matrices of size $n \times n$.
Thus the orthogonal polynomial is given by the ratio of two matrix integrals of the same type as the partition function (2.1):

$$
\begin{equation*}
\mathcal{P}_{n}(\lambda)=\frac{Z\left[V+\delta V_{1}+\delta V_{2} ; n\right]}{Z\left[V+\delta V_{1} ; n\right]}=\frac{\mathrm{e}^{-F\left[V+\delta V_{1}+\delta V_{2} ; n\right]}}{\mathrm{e}^{-F\left[V+\delta V_{1} ; n\right]}} \tag{4.12}
\end{equation*}
$$

where

$$
\begin{equation*}
\delta V_{1}(z)=\frac{N-n}{n} V(z) \quad \text { and } \quad \delta V_{2}=-\frac{1}{n} \ln (z-\lambda) . \tag{4.13}
\end{equation*}
$$

We have seen in the previous section (equation (2.45)) that

$$
\begin{equation*}
F[V ; n]=n^{2} F_{0}\left[V ; x_{c}\right]-\ln \theta_{3}\left(n x_{c}[V]\right)+\cdots \tag{4.14}
\end{equation*}
$$

We will use the fact that under a variation $\delta V$ of the potential, the variation of $F_{0}$ is [9,24]:

$$
\begin{equation*}
\delta F_{0}=\frac{1}{2 \mathrm{i} \pi} \oint \omega(z) \delta V(z) \mathrm{d} z \tag{4.15}
\end{equation*}
$$

where the anticlockwise contour encloses the support of the density of eigenvalues and $\omega(z)$ is the resolvent (equation (3.4) and (2.28)):

$$
\omega(z)=\frac{1}{2}\left(V^{\prime}(z)-M(z) \sqrt{\sigma(z)}\right)
$$

It is convenient to introduce two sources $t_{1}$ and $t_{2}$ for the variations $\delta V_{1}$ and $\delta V_{2}$ of the potential, and consider a generalized potential $\mathcal{V}(z)$ :

$$
\mathcal{V}(z)=V(z)+t_{1} \delta V_{1}(z)+t_{2} \delta V_{2}(z)=V(z)+t_{1} V(z)+t_{2} \ln |\lambda-z|
$$

Since $t_{1}=\frac{N-n}{n}$ and $t_{2}=-\frac{1}{n}$ are both small of order $\mathrm{O}\left(N^{-1}\right)$, we will expand $F$ in Taylor's series:
$F[\mathcal{V}(z) ; n]=F[V(z) ; n]+t_{1} \partial_{1} F+t_{2} \partial_{2} F+\frac{1}{2} t_{1}^{2} \partial_{11} F+t_{1} t_{2} \partial_{12} F+\frac{1}{2} t_{2}^{2} \partial_{22} F+\cdots$
all the derivatives being taken at the point $t_{1}=t_{2}=0$.
This will give
$\mathcal{P}_{n}(\lambda) \sim \mathrm{e}^{n \partial_{2} F_{0}} \mathrm{e}^{(N-n) \partial_{12} F_{0}} \mathrm{e}^{-\frac{1}{2} \partial_{22} F_{0}} \frac{\theta_{3}\left(n x+(N-n) \partial_{1} x-\partial_{2} x\right)}{\theta_{3}\left(n x+(N-n) \partial_{1} x\right)}\left(1+\mathrm{O}\left(N^{-1}\right)\right)$.
Now, let us compute the derivatives of $F_{0}$ and $x=x_{c}$ with respect to $t_{1}$ and $t_{2}$. The method proceeds similarly to section 3.1.
4.2.1. Derivatives of $F_{0}$ with respect to $t_{1}$ and $t_{2}$. Using (4.15) with (4.13):

$$
\frac{\partial F_{0}}{\partial t_{2}}=\frac{1}{2 \mathrm{i} \pi} \oint \omega(z) \ln (z-\lambda) \mathrm{d} z
$$

After integration by part, the pole in $(z-\lambda)$ picks a residue, and the result is a primitive of $\omega(\lambda)$ :

$$
\begin{equation*}
\frac{\partial F_{0}}{\partial t_{2}}=\int_{\lambda_{0}}^{\lambda} \omega(z) \mathrm{d} z . \tag{4.17}
\end{equation*}
$$

The lower bound of integration $\lambda_{0}$ is to be chosen such that $\mathrm{e}^{n \partial_{2} F_{0}} \sim_{\lambda \rightarrow \infty} \lambda^{n}$; i.e.

$$
\begin{equation*}
\ln \lambda_{0}=\int_{\lambda_{0}}^{\infty}\left(\omega(z)-\frac{1}{z}\right) \mathrm{d} z \tag{4.18}
\end{equation*}
$$

In order to compute the second derivatives $\partial_{12} F_{0}$ and $\partial_{22} F_{0}$, we will need to differentiate $\omega(z)$ with respect to $t_{1}$ and $t_{2}$.
4.2.2. Derivatives of $\omega(z)$ with respect to $t_{1}$ and $t_{2}$. The resolvent $\omega(z)$ computed for the potential $\mathcal{V}(z)$ takes the form

$$
\begin{equation*}
\omega(z)=\frac{1}{2}\left(\mathcal{V}^{\prime}(z)-M(z) \sqrt{\sigma(z)}\right) \tag{4.19}
\end{equation*}
$$

where $M(z)$ is analytic. Note that when $\mathcal{V}^{\prime}(z)$ has a pole in $z=\lambda, M(z)$ may have a pole too. $\omega(z)$ obeys a linear equation:

$$
\begin{equation*}
\omega(z+\mathrm{i} 0)+\omega(z-\mathrm{i} 0)=\mathcal{V}^{\prime}(z) \quad \text { for } \quad z \in[a, b] \cup[c, d] \tag{4.20}
\end{equation*}
$$

Thus its derivatives obey linear equations as well:

$$
\begin{align*}
& \partial_{1} \omega(z+\mathrm{i} 0)+\partial_{1} \omega(z-\mathrm{i} 0)=\delta V_{1}^{\prime}(z)=V^{\prime}(z)  \tag{4.21}\\
& \partial_{2} \omega(z+\mathrm{i} 0)+\partial_{2} \omega(z-\mathrm{i} 0)=\delta V_{2}^{\prime}(z)=\frac{1}{z-\lambda} . \tag{4.22}
\end{align*}
$$

- $\partial \omega_{1}$. The solution of (4.21) is

$$
\begin{equation*}
\partial_{1} \omega(z)=\omega(z)-\frac{f(z)}{\sqrt{\sigma(z)}} \tag{4.23}
\end{equation*}
$$

where $f(z)$ is analytic in $z$. The boundary conditions (2.18) imply that $f(z) \sim z$ when $z \rightarrow \infty$ and $f$ has no pole, thus $f(z)$ is a polynomial of degree one:

$$
\begin{equation*}
\frac{\partial \omega(z)}{\partial t_{1}}=\omega(z)-\frac{z-z_{0}}{\sqrt{\sigma(z)}} \tag{4.24}
\end{equation*}
$$

$z_{0}$ is determined as a function of $a, b, c, d$ by the derivative of (2.30) with respect to $t_{1}$ :

$$
\begin{equation*}
\int_{b}^{c} \mathrm{~d} z \frac{z-z_{0}}{\sqrt{\sigma(z)}}=0 \tag{4.25}
\end{equation*}
$$

It can be checked that in term of elliptic theta functions we have (see appendix A, or [26])

$$
\begin{equation*}
\frac{z-z_{0}}{\sqrt{\sigma(z)}}=\frac{\mathrm{d}}{\mathrm{~d} z} \ln \frac{\theta_{1}\left(u(z)+u_{\infty}\right)}{\theta_{1}\left(u(z)-u_{\infty}\right)} \tag{4.26}
\end{equation*}
$$

and thus

$$
\begin{equation*}
\frac{\partial \omega(z)}{\partial t_{1}}=\omega(z)-\frac{\mathrm{d}}{\mathrm{~d} z} \ln \frac{\theta_{1}\left(u(z)+u_{\infty}\right)}{\theta_{1}\left(u(z)-u_{\infty}\right)} . \tag{4.27}
\end{equation*}
$$

- $\partial \omega_{2}$. Note that the $t_{2}$ source term is the primitive of the $\epsilon_{\lambda}$ source term of (3.2), and that

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} z} \frac{\partial \omega(z)}{\partial t_{2}}=\frac{\partial \omega(z)}{\partial \epsilon_{\lambda}}=-\frac{1}{n^{2}} \frac{\partial^{2} F}{\partial \epsilon_{z} \partial \epsilon_{\lambda}}=\omega_{c}(z, \lambda) \tag{4.28}
\end{equation*}
$$

so, $\partial \omega / \partial t_{2}$ has already been computed in (3.8). The second derivative $\partial_{22} F$ corresponds to $z=\lambda$,

$$
\begin{equation*}
\partial_{22} F_{0}=\ln \sqrt{\sigma(\lambda)}+2 \ln \left(\theta_{1}\left(u(\lambda)-u_{\infty}\right)\right) . \tag{4.29}
\end{equation*}
$$

### 4.2.3. Derivatives of $x$. Recall that

$$
\begin{equation*}
x=\int_{a}^{b} \rho(\lambda) \mathrm{d} \lambda=\frac{1}{2 \mathrm{i} \pi} \int_{a}^{b} M(z) \sqrt{\sigma(z)} \mathrm{d} z \tag{4.30}
\end{equation*}
$$

using (4.27) we obtain

$$
\begin{equation*}
\frac{\partial x}{\partial t_{1}}=x+\frac{1}{2 \mathrm{i} \pi} \int_{a}^{b} \frac{z-z_{0}}{\sqrt{\sigma(z)}} \mathrm{d} z=x+2 u_{\infty} . \tag{4.31}
\end{equation*}
$$

Similarly, from (4.28) and (3.18) or (3.31) (or less tediously, taking the primitive of (3.14)), we obtain

$$
\begin{equation*}
\frac{\partial x}{\partial t_{2}}=-u(\lambda)+u_{\infty} \tag{4.32}
\end{equation*}
$$

### 4.3. Final result

4.3.1. Case $\lambda \notin[a, b] \cup[c, d]$. Eventually, inserting (4.17), (4.27), (4.29), (4.31), (4.32) into (4.16) we obtain

$$
\begin{equation*}
\mathcal{P}_{n}(\lambda)_{\lambda \notin[a, b] \cup[c, d]} \sqrt{u^{\prime}(\lambda)} p_{n}(u(\lambda)) \mathrm{e}^{N \int_{\lambda_{0}}^{\lambda} \omega} \tag{4.33}
\end{equation*}
$$

where

$$
\begin{equation*}
p_{n}(u)=C_{n} \frac{\theta_{3}\left(N x+2(N-n) u_{\infty}+u-u_{\infty}\right) \theta_{1}\left(2 u_{\infty}\right)}{\theta_{3}\left(N x+2(N-n) u_{\infty}\right) \theta_{1}\left(u-u_{\infty}\right)}\left(\frac{\theta_{1}\left(u+u_{\infty}\right)}{\theta_{1}\left(u-u_{\infty}\right)}\right)^{n-N} \tag{4.34}
\end{equation*}
$$

and $C_{n}$ is a normalization such that $\mathcal{P}_{n} \sim \lambda^{n}$ for $\lambda \rightarrow \infty$.
$C_{n}=\sqrt{2 K} A^{n-N+1} \quad$ with $\quad A=-\frac{1}{2 K} \frac{\theta_{1}^{\prime}(0)}{\theta_{1}\left(2 u_{\infty}\right)} \quad$ and $\quad K=\int_{c}^{b} \frac{\mathrm{~d} z}{\sqrt{\sigma(z)}}$
$A=-\frac{1}{4}|d-a-c+b| \frac{\theta_{3}(0)}{\theta_{3}\left(2 u_{\infty}\right)}$.
Note that (4.33) is unchanged under $u \rightarrow u+1$ and $u \rightarrow u+\tau$. Indeed, a shift $u \rightarrow u+\tau$ amounts to a non-trivial circle around the cut $[c, d]$. Thus $\int \omega$ is shifted by $-2 \mathrm{i} \pi \int_{c}^{d} \rho=-2 \mathrm{i} \pi(1-x)$, and $\mathrm{e}^{N \int \omega}$ receives a phase $\mathrm{e}^{2 \mathrm{i} \pi N x}$. In the same time, the $\theta$ functions receive phase factors: $\theta(v+\tau)=\theta(v) \mathrm{e}^{-2 i \pi(v+\tau / 2)}$. One can easily check that the total phase shift is 0 .
4.3.2. Case $\lambda \in[a, b] \cup[c, d]$. Expression (4.34) has been derived by a saddle-point approximation of (4.12) when $\lambda$ does not belong to $[a, b] \cup[c, d]$. When $\lambda$ lies on the cut $[a, b] \cup[c, d]$, equation (4.12) actually has two saddle points, contributing to the same order. They correspond to the two determinations of the square root $\pm \sqrt{\sigma(\lambda)}$. The asymptotic expression for the orthogonal polynomial is then given by a sum of two terms:

$$
\begin{equation*}
\mathcal{P}_{n}(\lambda)_{\lambda \in[a, b] \cup[c, d]} C \sqrt{u^{\prime}}\left[p_{n}(u) \mathrm{e}^{-\mathrm{i} N \pi \zeta(\lambda)}+\mathrm{i} p_{n}(-u) \mathrm{e}^{\mathrm{i} N \pi \zeta(\lambda)}\right] \mathrm{e}^{\frac{1}{2} N V(\lambda)} \tag{4.36}
\end{equation*}
$$

where $\zeta(\lambda)=\int_{d}^{\lambda} \rho(z) \mathrm{d} z$ and

$$
\begin{equation*}
C=\mathrm{e}^{-\frac{1}{2} N\left(V\left(\lambda_{0}\right)+\int_{\lambda_{0}}^{d} M(z) \sqrt{\sigma(z)} \mathrm{d} z\right)}=\mathrm{d}^{N} \mathrm{e}^{-\frac{1}{2} N V(d)} \mathrm{e}^{-N \int_{d}^{\infty}(\omega(z)-1 / z) \mathrm{d} z} . \tag{4.37}
\end{equation*}
$$



Figure 2. Typical behaviour of the wavefunction.

(a)


Figure 3. Deformation of the contour integral.

To summarize, when $\lambda \notin[a, b] \cup[c, d]$, the wavefunction $\psi_{n}(\lambda)=\mathcal{P}_{n}(\lambda) \mathrm{e}^{-N V / 2}$ decays exponentially, and within the support $[a, b] \cup[c, d]$, it oscillates at a frequency of order $N$.
4.3.3. Check of the orthogonality. For completeness, let us check that the functions (4.33) are indeed orthogonal (at leading order in $N^{-1}$ ). Let us compute the integral

$$
\int_{-\infty}^{\infty} \mathrm{d} \lambda \mathcal{P}_{n}(\lambda) \mathcal{P}_{m}(\lambda) \mathrm{e}^{-N V(\lambda)}
$$

The contributions of the integral along $]-\infty, a] \cup[b, c] \cup[d, \infty[$ are exponentially small and do not contribute at leading order.

Along $[a, b] \cup[c, d]$, we use expression (4.36) and get a sum of four terms:
$C^{2} \int \mathrm{~d} \lambda u^{\prime}(\lambda)\left(\begin{array}{cc}\mathrm{i} p_{n}(u) p_{m}(-u) & +\mathrm{i} p_{n}(-u) p_{m}(u) \\ +p_{n}(u) p_{m}(u) \mathrm{e}^{-2 \mathrm{i} N \pi \zeta(\lambda)} & -p_{n}(-u) p_{m}(-u) \mathrm{e}^{2 \mathrm{i} N \pi \zeta(\lambda)}\end{array}\right)$
Since the two last terms have fast oscillations of frequency $N$, they are suppressed as $\mathrm{O}(1 / N)$.
The leading contribution is thus given by the two first terms of (4.38), which can be rewritten as integrals in the $u$ plane along the contour depicted in figure $3(a)$ :

$$
\begin{align*}
\int \mathcal{P}_{n} \mathcal{P}_{m} \mathrm{e}^{-N V} & =\mathrm{i} c_{n m} \int \mathrm{~d} u \frac{\theta_{3}\left(x_{n}+u-u_{\infty}\right) \theta_{3}\left(x_{m}-u-u_{\infty}\right)}{\theta_{3}\left(x_{n}\right) \theta_{3}\left(x_{m}\right) \theta_{1}\left(u-u_{\infty}\right) \theta_{1}\left(u+u_{\infty}\right)} \\
& \times\left(\frac{\theta_{1}\left(u+u_{\infty}\right)}{\theta_{1}\left(u-u_{\infty}\right)}\right)^{n-m}+(u \rightarrow-u) \tag{4.39}
\end{align*}
$$

where $x_{n}$ and $c_{n m}$ are short notation for

$$
\begin{equation*}
x_{n}=N x+2(N-n) u_{\infty} \quad \text { and } \quad c_{n m}=C^{2} C_{n} C_{m} \cdot \theta_{1}^{2}\left(2 u_{\infty}\right) \tag{4.40}
\end{equation*}
$$

If $n>m$ we may deform the contour to a circle around the point $-u_{\infty}$ (figure $3(b)$ ), and the integral vanishes since there is no pole, while if $m>n$ we deform the contour to a circle around $+u_{\infty}$ (figure $3(c)$ ). Therefore, the integral vanishes for $n \neq m$.

When $n=m$, the integral picks a residue:

$$
\begin{equation*}
\int \mathrm{d} \lambda \mathcal{P}_{n} \mathcal{P}_{m} \mathrm{e}^{-N V(\lambda)}=h_{n} \delta_{n m} \tag{4.41}
\end{equation*}
$$

with ( $C, c_{n m}, A, x_{n}$ are defined in (4.35), (4.40), (4.37))

$$
\begin{equation*}
h_{n}=c_{n n} \frac{4 \pi}{\theta_{1}\left(2 u_{\infty}\right) \theta_{1}^{\prime}(0)} \frac{\theta_{3}\left(x_{n+1}\right)}{\theta_{3}\left(x_{n}\right)}=-4 \pi C^{2} \frac{\theta_{3}\left(x_{n+1}\right)}{\theta_{3}\left(x_{n}\right)} A^{2(n-N+1 / 2)} . \tag{4.42}
\end{equation*}
$$

4.3.4. Recurrence equation. It is well known that the orthogonal polynomials satisfy a recurrence equation of the form $[2,21]$

$$
\begin{equation*}
\lambda \mathcal{P}_{n}(\lambda)=\mathcal{P}_{n+1}(\lambda)+\beta_{n} \mathcal{P}_{n}(\lambda)+\alpha_{n} \mathcal{P}_{n-1}(\lambda) . \tag{4.43}
\end{equation*}
$$

Here, we find that (divide (4.43) by $\mathcal{P}_{n}$, and match the poles on both sides)

$$
\begin{equation*}
\alpha_{n}=\frac{h_{n}}{h_{n-1}}=A^{2} \frac{\theta_{3}\left(x_{n+1}\right) \theta_{3}\left(x_{n-1}\right)}{\theta_{3}^{2}\left(x_{n}\right)} \tag{4.44}
\end{equation*}
$$

which can be rewritten more compactly as

$$
\begin{equation*}
\alpha_{n}=\frac{1}{16}\left(((d-a)-(c-b))^{2}+4(d-a)(c-b) \mathrm{cn}^{2}\left(x_{n}+\frac{1}{2}\right)\right) . \tag{4.45}
\end{equation*}
$$

By taking $u=0$ in (4.43) we obtain $\beta_{n}$ :

$$
\begin{equation*}
\beta_{n}-d=A\left[\frac{\theta_{3}\left(x_{n+3 / 2}\right) \theta_{3}\left(x_{n}\right)}{\theta_{3}\left(x_{n+1}\right) \theta_{3}\left(x_{n+1 / 2}\right)}+\frac{\theta_{3}\left(x_{n+1}\right) \theta_{3}\left(x_{n-1 / 2}\right)}{\theta_{3}\left(x_{n}\right) \theta_{3}\left(x_{n+1 / 2}\right)}\right] \tag{4.46}
\end{equation*}
$$

which can be rewritten more compactly as

$$
\begin{equation*}
\beta_{n}=\frac{a+d+(c-b)}{2}-(c-b) \frac{d-b}{\left.c-b+(d-c) /\left[\mathrm{cn}^{2} x_{n}-u_{\infty}+\frac{1}{2}\right)\right]} . \tag{4.47}
\end{equation*}
$$

The sequences $\alpha_{n}$ and $\beta_{n}$ are thus quasi-periodic in $n$. It is interesting to recall that the behaviour of these coefficients has been extensively studied (mainly by numerical methods) by several authors in the early 1990s [25]. The general conclusion was that in the multi-cut case the general behaviour of the recursion coefficients was 'chaotic' in $n$ (and regular or quasi-periodic only in some special cases). It is clear from our expressions that in the two-cut case the behaviour is always periodic or quasi-periodic and never chaotic (in the mathematical sense). This is, in fact, true even if the number of cuts is larger than two (see appendix C).

In the symmetric case, $x=\frac{1}{2}$ and $u_{\infty}=\frac{1}{4}$, we have $x_{n}=n / 2 \bmod 1$, so that we recover $\beta_{n}=0$ and $\alpha_{n}=\frac{1}{4}\left(a-(-1)^{n} b\right)^{2}$.

In the general case, $\alpha_{n}$ and $\beta_{n}$ vary along a periodic curve, between two extrema, given by

$$
\begin{aligned}
& \frac{1}{16}(d-a-(c-b))^{2} \leqslant \alpha_{n} \leqslant \frac{1}{16}(d-a+(c-b))^{2} \\
& \frac{1}{2}(d+a)-\frac{1}{2}(c-b) \leqslant \beta_{n} \leqslant \frac{1}{2}(d+a)+\frac{1}{2}(c-b) .
\end{aligned}
$$

Similarly to the one-cut case, one may relate $\alpha_{n}$ to square width of the distribution of eigenvalues, and $\beta_{n}$ to the centre of the distribution.

### 4.4. The kernel $K(\lambda, \mu)$

We can now evaluate the kernel $K(\lambda, \mu)$ according to (4.10). Let us note $u=u(\lambda)$ and $v=u(\mu)$ and we assume $\lambda, \mu \in[a, b] \cup[c, d]$
$K(\lambda, \mu) \sim \frac{C^{2} \sqrt{u^{\prime} v^{\prime}}}{N h_{N-1}} \sum_{\epsilon, \eta= \pm 1} \frac{\sqrt{\epsilon \eta} p_{N}(\epsilon u) p_{N-1}(\eta v) \mathrm{e}^{-\epsilon N \mathrm{i} \pi \zeta(\lambda)} \mathrm{e}^{-\eta N \mathrm{i} \pi \zeta(\mu)}-(u \rightarrow v)}{(\lambda-\mu)}$
which can be rewritten as a sum of eight terms

$$
\begin{align*}
K(\lambda, \mu)= & \frac{c_{N, N-1}}{h_{N-1} \theta_{3}\left(x_{N}\right) \theta_{3}\left(x_{N-1}\right)} \frac{\sqrt{u^{\prime} v^{\prime}}}{N(\lambda-\mu)} \\
& \times \sum_{\epsilon, \eta, \kappa= \pm 1} \kappa \sqrt{\epsilon \eta} \frac{\theta_{3}\left(N x+\epsilon u-\kappa u_{\infty}\right) \theta_{3}\left(N x+\eta v+\kappa u_{\infty}\right)}{\theta_{1}\left(\epsilon u-\kappa u_{\infty}\right) \theta_{1}\left(\eta v+\kappa u_{\infty}\right)} \mathrm{e}^{-N \mathrm{i} \pi(\epsilon \zeta(\lambda)+\eta \zeta(\mu))} \tag{4.49}
\end{align*}
$$

We will see below that not all the terms contribute to the same order.
4.4.1. Regime $|\lambda-\mu| \sim O(1 / N)$. The eight terms of (4.49) can be rewritten as four combinations of the type

$$
\sin (N \pi(\zeta(\lambda) \pm \zeta(\mu))) \frac{f(u, v) \mp f(v, u)}{N(\lambda-\mu)}
$$

and

$$
\cos (N \pi(\zeta(\lambda) \pm \zeta(\mu))) \frac{g(u, v)-g(v, u)}{N(\lambda-\mu)}
$$

In the limit $|\lambda-\mu|$ small, i.e. $|u-v|$ small, the terms with a cosine will be proportional to derivatives of $g(u, v)$, and there will be an overall $1 / N$ factor. Similarly, the term with a sine and $\mathrm{a}+\operatorname{sign}$ will be proportional to a derivative of $f(u, v)$ and will be of order $1 / N$. Only the term proportional to $\sin N \pi \int_{\lambda}^{\mu} \rho(z) \mathrm{d} z$ can balance the $1 / N$ factor, and is dominant in the short range regime. After calculation we obtain

$$
\begin{equation*}
K(\lambda, \mu)_{|\lambda-\mu| \sim \mathrm{O}(1 / N)}^{\sim} \frac{\sin N \pi \int_{\lambda}^{\mu} \rho(z) \mathrm{d} z}{N \pi(\lambda-\mu)} . \tag{4.50}
\end{equation*}
$$

As expected we have

$$
\begin{equation*}
K(\lambda, \lambda)=\rho(\lambda) \tag{4.51}
\end{equation*}
$$

and we recover the universal short-range correlation function

$$
\begin{equation*}
\rho(\lambda, \mu) \sim \rho(\lambda) \rho(\mu)\left(1-\left(\frac{\sin N \pi \rho(\lambda)(\lambda-\mu)}{N \pi \rho(\lambda)(\lambda-\mu)}\right)^{2}\right) . \tag{4.52}
\end{equation*}
$$

4.4.2. Long-range regime, smoothed oscillations. When $|\lambda-\mu| \sim \mathrm{O}(1), K(\lambda, \mu)$ has highfrequency oscillations, and only a smoothed correlation function obtained by averaging the oscillations can be observed.

Recall that the connected two-point correlation function is related to $K^{2}$ by (4.9):

$$
\begin{equation*}
\rho_{2 c}(\lambda, \mu)=-K(\lambda, \mu)^{2} \tag{4.53}
\end{equation*}
$$

with $K(\lambda, \mu)$ given by (4.49).

Smoothing out the oscillations amounts to kill all terms containing some $\mathrm{e}^{\mathrm{i} N \pi \zeta}$ in the square of equation (4.49), we thus have

$$
\begin{align*}
\overline{K(\lambda, \mu)^{2}}= & \frac{-2 u^{\prime} v^{\prime} c_{N N-1}^{2}}{h_{N-1}^{2} \theta_{3}^{2}\left(x_{N}\right) \theta_{3}^{2}\left(x_{N-1}\right) N^{2}(\lambda-\mu)^{2}} \sum_{\epsilon, \eta, \kappa_{1}, \kappa_{2}= \pm 1} \kappa_{1} \kappa_{2}\left\{\theta_{3}\left(N x+\epsilon u-\kappa_{1} u_{\infty}\right)\right. \\
& \left.\times \theta_{3}\left(N x-\epsilon u-\kappa_{2} u_{\infty}\right) \theta_{3}\left(N x+\eta v+\kappa_{1} u_{\infty}\right) \theta_{3}\left(N x-\eta v+\kappa_{2} u_{\infty}\right)\right\} \\
& \times\left\{\theta_{1}\left(\epsilon u-\kappa_{1} u_{\infty}\right) \theta_{1}\left(-\epsilon u-\kappa_{2} u_{\infty}\right) \theta_{1}\left(\eta v+\kappa_{1} u_{\infty}\right) \theta_{1}\left(-\eta v+\kappa_{2} u_{\infty}\right)\right\}^{-1} . \tag{4.54}
\end{align*}
$$

Using that (see appendix A and [26-28])

$$
\begin{equation*}
\lambda-\mu=-2 A \frac{\theta_{1}(u-v) \theta_{1}(u+v) \theta_{1}^{2}\left(2 u_{\infty}\right)}{\theta_{1}\left(u-u_{\infty}\right) \theta_{1}\left(u+u_{\infty}\right) \theta_{1}\left(v-u_{\infty}\right) \theta_{1}\left(v+u_{\infty}\right)} \tag{4.55}
\end{equation*}
$$

we obtain $\left(\epsilon_{i j}=-\epsilon_{j i}= \pm 1\right)$

$$
\begin{align*}
\overline{K(\lambda, \mu)^{2}}= & \frac{2 u^{\prime} v^{\prime} \theta_{1}^{\prime 2}(0)}{N^{2} 4 \pi^{2} \theta_{3}^{4}(N x) \theta_{1}^{2}\left(2 u_{\infty}\right)}\left(\frac{\theta_{1}\left(u-u_{\infty}\right) \theta_{1}\left(u+u_{\infty}\right) \theta_{1}\left(v-u_{\infty}\right) \theta_{1}\left(v+u_{\infty}\right)}{\theta_{1}(u-v) \theta_{1}(u+v)}\right)^{2} \\
& \times \sum_{i, j, k, l= \pm 1}\left(\epsilon_{i j} \epsilon_{k l}+\epsilon_{i l} \epsilon_{k j}\right) \\
& \times \frac{\theta_{3}\left(N x+u-\mathrm{i} u_{\infty}\right) \theta_{3}\left(N x-u-k u_{\infty}\right) \theta_{3}\left(N x+v-j u_{\infty}\right) \theta_{3}\left(N x-v-l u_{\infty}\right)}{\theta_{1}\left(u-\mathrm{i} u_{\infty}\right) \theta_{1}\left(-u-k u_{\infty}\right) \theta_{1}\left(v-j u_{\infty}\right) \theta_{1}\left(-v-l u_{\infty}\right)} . \tag{4.56}
\end{align*}
$$

We see that (4.56) has no pole when $u= \pm u_{\infty}$, it can have (double) poles only when $u= \pm v$. Thus, equation (4.56) can be rewritten in terms of Weirstrass functions of $u+v$ and $u-v$ :

$$
\overline{K(\lambda, \mu)^{2}}=-\frac{1}{2 N^{2} \pi^{2}} u^{\prime} v^{\prime}\left(C_{1} \wp(u-v)+C_{2} \wp(u+v)-2 S\right) .
$$

Taking $u=v$ and $u=-v$ in (4.56), we find that the residues are $C_{1}=C_{2}=1$, and taking a particular value of $u$ and $v$, we find the constant $S$, equal to what we had in (3.31)

$$
\begin{equation*}
\overline{K(\lambda, \mu)^{2}}=-\frac{1}{2 N^{2} \pi^{2}} u^{\prime} v^{\prime}\left(\wp(u-v)+\wp(u+v)-2 \wp\left(N x+\frac{1}{2} \tau\right)\right) \tag{4.57}
\end{equation*}
$$

and we recover the result (3.31) found in section 3.6.

## 5. Conclusions

In this paper, we have solved the puzzle raised by $[13,15]$ and understood why the naive mean-field method [11] and the orthogonal polynomial ansatz [13,14] approach used in the symmetric case disagree.

We have proven here that this effect has nothing to do with a $\mathbb{Z}_{2}$ symmetry breaking, as was sometimes assumed [13], it is general as soon as the support of the density is not connected.

The apparent paradox comes from the fact that when the support of eigenvalues is not connected, the free energy admits no large- $N$ expansion in powers of $1 / N^{2}$ (topological expansion [10]). This means that the free energy in the multi-cut case is not given by a topological expansion, i.e. the sum of diagrams with a weight $N^{\chi}(\chi$ is the Euler characteristic of the diagram).

The explanation lies in the discreteness of the number of eigenvalues. For instance, in the symmetric two-cut case, the classical approach assumes that the minimum of the free energy is reached when one half $\left(x=\frac{1}{2}\right)$ of the eigenvalues are in each cut. Obviously, this minimum is never reached when the total number of eigenvalues is odd, and, in general, the result depends on the fractional part of $N x$.

At leading order in $N$ only, the free energy is correctly given by the classical saddle-point limit [11, 19], but the first order in $N$ is not sufficient to determine the two-point (or higher) correlation function.

Here we have computed explicitly the two-point connected correlation function. It contains a universal part depending only on the number of cuts, which was obtained by [11], and contains in addition, a non-universal term quasiperiodic in $N[17,25]$.

Let us stress that our calculation holds for any potential, not necessarily symmetric, and it can also be generalized to a potential with complex coefficients (appendix B), and to an arbitrary number of cuts (appendix C).

We have also reobtained directly the asymptotic expressions for the orthogonal polynomials [17], which allows in principle through the Darboux-Christoffel theorem (equation (4.10)) to compute any correlation function of any number of eigenvalues in the short- or long-range domain (and one can smooth it afterwards).

The orthogonal polynomial approach may in turn be used for other random matrix ensembles, and it would be interesting to apply our results to orthogonal or symplectic ensembles [2].

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## Appendix A. Topological expansion at fixed $\boldsymbol{x}$

In this appendix we explain in detail how the perturbative expansion is organized in a topological expansion for a two-cut model, provided that the proportion of eigenvalues in each well is fixed. For simplicity let us consider the (properly normalized) symmetric quartic potential

$$
\begin{equation*}
V(z)=\frac{1}{8}\left(z^{2}-z_{0}^{2}\right)^{2} \tag{A.1}
\end{equation*}
$$

(but the argument is of course general). From subsection 2.2.2 we consider the two-matrix integral

$$
\begin{gather*}
\int \mathrm{d}_{n_{1}}\left[M_{1}\right] \int \mathrm{d}_{n_{2}}\left[M_{2}\right] \exp \left\{-N \operatorname{Tr}\left(V\left(M_{1}\right)\right)-N \operatorname{Tr}\left(V\left(M_{2}\right)\right)\right. \\
\left.+2 \operatorname{Tr}\left(\ln \left(M_{1} \otimes \mathrm{Id}-\mathrm{Id} \otimes M_{2}\right)\right)\right\} \tag{A.2}
\end{gather*}
$$

where the size of the Hermitian matrices $M_{1}$ and $M_{2}$ are, respectively, $n_{1}=x N$ and $n_{2}=(1-x) N$, and where their EVs are concentrated around the two minima $-z_{0}$ and $z_{0}$, respectively. To construct the perturbative expansion we introduce the fiducial expansion


Figure A1. Schematic topology of diagrams generated by perturbation theory at fixed ratio $x$. For this example $h=4$ and $\chi=-6$.
parameter $g$ by shifting $V \rightarrow V / g$ ( $g$ will be set to one at the end). Then we can expand $M_{1}$ and $M_{2}$ around their classical minima

$$
\begin{equation*}
M_{1}=-z_{0} \mathrm{Id}-z_{0}^{-1} \sqrt{g / N} A \quad M_{2}=z_{0} \mathrm{Id}+z_{0}^{-1} \sqrt{g / N} B \tag{A.3}
\end{equation*}
$$

The action reads then

$$
\begin{align*}
& \operatorname{Tr}\left[\frac{1}{2} A^{2}+\frac{1}{2 z_{0}^{2}} \sqrt{\frac{g}{N}} A^{3}+\frac{1}{8 z_{0}^{4}} \frac{g}{N} A^{4}\right]+\operatorname{Tr}\left[\frac{1}{2} B^{2}+\frac{1}{2 z_{0}^{2}} \sqrt{\frac{g}{N}} B^{3}+\frac{1}{8 z_{0}^{4}} \frac{g}{N} B^{4}\right] \\
&-2 n_{1} n_{2} \log \left(-2 Z_{0}\right)-\operatorname{Tr}\left[\log \left(\operatorname{Id} \otimes \operatorname{Id}+\frac{1}{2 z_{0}^{2}} \sqrt{\frac{g}{N}}(A \otimes \operatorname{Id}+\operatorname{Id} \otimes B)\right)\right] . \tag{A.4}
\end{align*}
$$

The first two terms generates standard planar graphs for the two matrices $A$ and $B$ with threeand four-point vertices. Expanding the $\log$ of the last term we generates interaction terms between the two matrices of the form

$$
\begin{equation*}
\left(\frac{g}{N}\right)^{\left(k_{1}+k_{2}\right) / 2} \operatorname{tr}\left(A^{k_{1}}\right) \operatorname{tr}\left(B^{k_{2}}\right) \tag{A.5}
\end{equation*}
$$

If $k_{2}=0\left(\right.$ respectively, $\left.k_{1}=0\right)$ this gives $\mathcal{O}(g)$ counterterms to the pure $A$ (respectively, $B$ ) action of the form

$$
\begin{equation*}
\frac{g^{k / 2}}{N^{k / 2-1}} \operatorname{tr}\left(A^{k}\right) \tag{A.6}
\end{equation*}
$$

If both $k$ are non-zero, this gives $\mathcal{O}(g)$ contact terms between the $A$-graphs and the $B$ graphs. Indeed, in terms of Feynmann graphs, the perturbative expansion can be reorganized in (orientable) fat $A$-graphs and fat $B$-graphs, which are allowed to be in contact $A-B$ at one or several points ( $A-B$ hetero-contacts only being allowed). It is easy to check that the power of $N$ of a given graph is still its Euler characteristics $\chi=2(1-h)$, where $h$ is the genus of the surface (see figure A1).

Hence, for fixed ratio $x=n_{1} / N$, the perturbative expansion in $g$ can still be organized in a topological expansion, provided of course than at a fixed order in $N$ (i.e. for graphs with a fixed Euler characteristics), the perturbative series in $g$ is convergent up to $g=1$. This is ensured as long as we do not reach a critical point as we increase $g$. In our double-well potential case this would occur only if at fixed $x$ the EVs in one of the wells flow over the wall into the second well. This would contradicts the initial condition that we have chosen, which is that $x$ and the potential $V$ are such that the large- $N$ mean-field solution is a two-intervals solution.

## Appendix B. A few useful identities on elliptic functions

Here we collect a few useful identities on elliptic functions used through the paper. For details see [26-28]. We start from

$$
\begin{equation*}
\sigma(\lambda)=(\lambda-a)(\lambda-b)(\lambda-c)(\lambda-d) \quad a<b<c<d \tag{B.1}
\end{equation*}
$$

and the map from the complex plane to the torus

$$
\begin{equation*}
u(\lambda)=\frac{1}{2 K} \int_{d}^{\lambda} \frac{\mathrm{d} z}{\sqrt{\sigma(z)}} \tag{B.2}
\end{equation*}
$$

where the half-period $K$ is

$$
\begin{equation*}
K=\int_{b}^{c} \frac{\mathrm{~d} z}{\sqrt{|\sigma(z)|}}=\frac{2}{\sqrt{(c-a)(d-b)}} K[m]=\frac{\pi \theta_{3}^{2}(0 \mid \tau)}{\sqrt{(c-a)(d-b)}} \tag{B.3}
\end{equation*}
$$

and $K[m]$ is the standard complete elliptic integral [26-28], with the modulus $m$ equal to the biratio of the four points $a, b, c, d$ :

$$
\begin{equation*}
m=\frac{(d-a)(c-b)}{(d-b)(c-a)} \tag{B.4}
\end{equation*}
$$

$m$ is related to the modular parameter $\tau$ of the torus by

$$
\begin{equation*}
m=\mathrm{e}^{\mathrm{i} \pi \tau} \frac{\theta_{3}^{4}\left(\left.\frac{1}{2} \tau \right\rvert\, \tau\right)}{\theta_{3}^{4}(0 \mid \tau)} \quad \text { and conversely } \quad \tau=\mathrm{i} \frac{K[1-m]}{K[m]} \tag{B.5}
\end{equation*}
$$

where we have used the Jacobi theta functions

$$
\begin{equation*}
\theta_{1}(z \mid \tau)=\theta_{1}(z)=-\mathrm{i} \sum_{r \in \mathbb{Z}+1 / 2}(-1)^{r} q^{r^{2}} \mathrm{e}^{2 \mathrm{i} \pi r z} \quad \text { with } \quad q=\mathrm{e}^{\mathrm{i} \pi \tau} \tag{B.6}
\end{equation*}
$$

and

$$
\begin{equation*}
\theta_{3}(z \mid \tau)=\theta_{3}(z)=q^{1 / 4} \mathrm{e}^{\mathrm{i} \pi z} \theta_{1}\left(\left.z+\frac{1}{2}+\frac{1}{2} \tau \right\rvert\, \tau\right) \tag{B.7}
\end{equation*}
$$

With this mapping $u(\lambda)$ between the $\lambda$ complex plane and the periodic rectangle of sides $(1, \tau)$, we have

$$
\begin{equation*}
u(d)=0 \quad u(a)=\frac{1}{2} \quad u(b)=\frac{1}{2}(1+\tau) \quad u(c)=\frac{1}{2} \tau \quad u(\infty)=u_{\infty} \tag{B.8}
\end{equation*}
$$

The inverse mapping can be written in terms of theta functions

$$
\begin{align*}
& \lambda-d=-\frac{\theta_{1}^{\prime}(0)}{2 K} \frac{\theta_{1}^{2}(u) \theta_{1}\left(2 u_{\infty}\right)}{\theta_{1}\left(u+u_{\infty}\right) \theta_{1}\left(u-u_{\infty}\right) \theta_{1}^{2}\left(u_{\infty}\right)}  \tag{B.9}\\
& \sqrt{\sigma(\lambda)}=\frac{\theta_{1}^{\prime 2}(0)}{4 K^{2}} \frac{\theta_{1}(2 u) \theta_{1}\left(2 u_{\infty}\right)}{\theta_{1}^{2}\left(u-u_{\infty}\right) \theta_{1}^{2}\left(u+u_{\infty}\right)} \tag{B.10}
\end{align*}
$$

and in terms of the usual trigonometric elliptic functions $\mathrm{sn}, \mathrm{cn}, \mathrm{dn}$ [26-28] that we normalize to have periods 1 and $\tau$, i.e.

$$
\begin{equation*}
\operatorname{sn}(u)=\operatorname{sn}(2 K[m] u \mid m) \quad \operatorname{dn}(u)=\operatorname{dn}(2 K[m] u \mid m) \ldots \tag{B.11}
\end{equation*}
$$

one has
$\frac{\operatorname{sn}^{2}(u)}{\operatorname{sn}^{2}\left(u_{\infty}\right)}=\frac{\lambda-d}{\lambda-c} \quad \frac{\mathrm{cn}^{2}(u)}{\operatorname{cn}^{2}\left(u_{\infty}\right)}=\frac{\lambda-a}{\lambda-c} \quad \frac{\operatorname{dn}^{2}(u)}{\operatorname{dn}^{2}\left(u_{\infty}\right)}=\frac{\lambda-b}{\lambda-c}$
$\lambda-c=(d-c) \frac{1}{1-\operatorname{sn}^{2}(u) / \operatorname{sn}^{2}\left(u_{\infty}\right)} \quad \lambda-b=(d-b) \frac{\mathrm{dn}^{2}(u)}{1-\operatorname{sn}^{2}(u) / \operatorname{sn}^{2}\left(u_{\infty}\right)}$
$\lambda-a=(d-a) \frac{\mathrm{cn}^{2}(u)}{1-\mathrm{sn}^{2}(u) / \mathrm{sn}^{2}\left(u_{\infty}\right)} \quad \lambda-d=(d-c) \frac{d-a}{c-a} \frac{\mathrm{sn}^{2}(u)}{1-\mathrm{sn}^{2}(u) / \mathrm{sn}^{2}\left(u_{\infty}\right)}$
$\sqrt{\sigma(\lambda)}=(d-c)(d-a) \sqrt{\frac{d-b}{c-a}} \frac{\operatorname{sn}(u) \mathrm{cn}(u) \mathrm{dn}(u)}{\left(1-\operatorname{sn}^{2}(u) / \operatorname{sn}^{2}\left(u_{\infty}\right)\right)^{2}}$
and $u_{\infty}$ is related to $a, b, c, d$ by any of the following relations:

$$
\begin{equation*}
\operatorname{sn}^{2}\left(u_{\infty}\right)=\frac{c-a}{d-a} \quad \operatorname{cn}^{2}\left(u_{\infty}\right)=\frac{d-c}{d-a} \quad \operatorname{dn}^{2}\left(u_{\infty}\right)=\frac{d-c}{d-b} . \tag{B.16}
\end{equation*}
$$

## Appendix C. Complex potentials

The case of complex potentials, that is to say of a polynomial potential $V(\lambda)$ with complex coefficients, is interesting for some applications of the matrix models to two-dimensional gravity and when studying their connections with integrable hierarchies. In this case, the meanfield large- $N$ solution is known to be given by a continuous distribution of the eigenvalues along arcs in the complex plane [29].

In this appendix we show that our results are only slightly modified in this case.
In the two-cut case, we can repeat the analysis of section 2 . We fix $x=n_{1} / N$ (the proportion of EV in the first cut). The resolvent is still of the form (2.28), with the polynomial $M$ and the end points $a, b, c, d$ determined by the constraints (2.18) and (2.29), but they are no more real, in general, as well as the resulting mean-field free energy $F_{0}(x)$.

If we now repeat the calculation of section 2.2 .3 we cannot use a saddle-point approximation for the sum over $n_{1}$ by expanding $F_{0}(x)$ around the saddle point $x_{0}$ which is the true extremum of $F_{0}$.

$$
\begin{equation*}
\frac{\partial F_{0}}{\partial x}\left(x_{0}\right)=0 \tag{C.1}
\end{equation*}
$$

Indeed, this extremum is at a finite non-zero distance of the real axis, i.e. $\operatorname{Im}\left(x_{0}\right)=\mathcal{O}(1)$, while the method of section 2.2 .3 is valid only if $\operatorname{Im}\left(x_{0}\right)=\mathcal{O}(1 / N)$. However, since $N$ is integer, we can expand $F_{0}$ around any $x_{k}$, provided that

$$
\begin{equation*}
F_{0}^{\prime}\left(x_{k}\right)=2 \mathrm{i} \pi \frac{k}{N} \quad k \in \mathbb{Z} \tag{C.2}
\end{equation*}
$$

since the dangerous oscillating term $\mathrm{e}^{-N^{2}\left(x-x_{k}\right) F_{0}^{\prime}\left(x_{c}\right)}$ is then a constant for $x=n / N, n \in \mathbb{Z}$.
Therefore, as in [29] we have to consider the real pseudo-saddle point $x_{c}$ such that

$$
\begin{equation*}
\operatorname{Re}\left(F_{0}^{\prime}\left(x_{c}\right)\right)=0 \quad \text { with } \quad \operatorname{Im}\left(x_{c}\right)=0 \tag{C.3}
\end{equation*}
$$

and denote

$$
\begin{equation*}
\Delta_{c}=\frac{1}{2 \mathrm{i} \pi} F_{0}^{\prime}\left(x_{c}\right)=\frac{1}{2 \pi} \operatorname{Im}\left(F_{0}^{\prime}\left(x_{c}\right)\right) . \tag{C.4}
\end{equation*}
$$

We expand $F_{0}$ around some $x_{k}$ defined by equation (C.2) and such that

$$
\begin{equation*}
x_{k}-x_{c}=\mathcal{O}\left(N^{-1}\right) \tag{C.5}
\end{equation*}
$$

and we obtain for the total free energy (by exactly the same calculation as in section 2.2.3)

$$
\begin{equation*}
F=N^{2} F_{0}\left(x_{k}\right)-\ln \left(\theta_{3}\left(N x_{k}\right)\right)+F_{1}\left(x_{c}\right)+\frac{1}{2} \ln \left(2 \pi F_{0}^{\prime \prime}\left(x_{c}\right)\right)+\mathcal{O}\left(N^{-2}\right) \tag{C.6}
\end{equation*}
$$

where $\theta_{3}$ is the theta function with modular parameter

$$
\begin{equation*}
\tau=\frac{2 \mathrm{i} \pi}{F_{0}^{\prime \prime}\left(x_{c}\right)} . \tag{C.7}
\end{equation*}
$$

Only the first two terms are important for calculating the two-point functions and the orthogonal polynomials in the large- $N$ limit. This leading term does not depend on $k$. Indeed, we have

$$
\begin{equation*}
F_{0}^{\prime}\left(x_{k}\right)-F_{0}^{\prime}\left(x_{c}\right)=\left(x_{k}-x_{c}\right) F_{0}^{\prime \prime}\left(x_{c}\right)+\mathcal{O}\left(N^{-1}\right) \tag{C.8}
\end{equation*}
$$

hence

$$
\begin{equation*}
\left(x_{k}-x_{c}\right)=\frac{1}{N} \tau\left(k-N \Delta_{c}\right)+\mathcal{O}\left(N^{-2}\right) \tag{C.9}
\end{equation*}
$$

and using the periodicity relations of $\theta_{3}$ we can rewrite the leading term for the free energy as

$$
\begin{equation*}
N^{2} F_{0}\left(x_{k}\right)-\ln \left(\theta_{3}\left(N x_{k}\right)\right)=N^{2} F_{0}\left(n_{c} / N\right)-\frac{\mathrm{i} \pi}{\tau} u_{c}^{2}-\ln \left(\theta_{3}\left(u_{c}\right)\right) \tag{C.10}
\end{equation*}
$$

with

$$
\begin{equation*}
n_{c}=\mathrm{E}\left[N x_{c}\right] \quad u_{c}=\left[N x_{c}\right]-\tau\left[N \Delta_{c}\right] \tag{C.11}
\end{equation*}
$$

where $\mathrm{E}[u]$ is the integer part of $u$ (the largest integer smaller than $u$ ) and $[u]=u-\mathrm{E}[u]$ is the fractional part of $u$. This does not depend on $k$ up to negligible terms of order $\mathcal{O}\left(N^{-1}\right)$ (provided that condition (C.5) for $k$ holds).

One can now repeat the calculation of section 3 for the two-point function. Nothing is changed but we simply have to replace $x_{c}$ by $x_{k}$ in the intermediate steps and to use equation (C.9) at the end of the calculation. This amounts to replace the $x_{c}$ in the elliptic function $\mathrm{sn}^{2}$ by $x_{c}-\tau \Delta_{c}$. The final result for the two-point resolvent is

$$
\begin{align*}
\omega^{c}(\lambda, \mu)=- & \frac{1}{4(\lambda-\mu)^{2}}\left[\left(1-\sqrt{\frac{(\lambda-a)(\lambda-b)(\mu-c)(\mu-d)}{(\mu-a)(\mu-b)(\lambda-c)(\lambda-d)}}\right)+(\lambda \leftrightarrow \mu)\right] \\
& -\frac{(c-a)(d-b)}{4 \sqrt{\sigma(\lambda)} \sqrt{\sigma(\mu)}} \operatorname{sn}^{2}\left(N\left(x_{c}-\tau \Delta_{c}\right)+\frac{1}{2}\right) . \tag{C.12}
\end{align*}
$$

Similar results holds for the orthogonal polynomials. We simply have to consider the end-points $a, b, c, d$ for the mean-field real parameter $x_{c}$ and to make the replacement

$$
\begin{equation*}
N x_{c} \rightarrow N\left(x_{c}+\tau \Delta_{c}\right) \tag{C.13}
\end{equation*}
$$

in the elliptic functions involving $N x_{c}$. In any case these terms depend only on the fractional parts of $N x_{c}$ and of $N \Delta_{c}$.

A final interesting remark on the periodicity properties of the non-universal term $\operatorname{sn}^{2}\left(N\left(x_{c}+\tau \Delta_{c}\right)\right)$ can be made. From the definition (C.4) $\Delta_{c}$ corresponds to a 'phase shift' between the two arcs where the density of EV is non-zero.

$$
\begin{equation*}
\Delta_{c}=\frac{1}{2 \mathrm{i} \pi} F_{0}^{\prime}\left(x_{c}\right)=\frac{\Gamma_{1}-\Gamma_{2}}{2 \mathrm{i} \pi} \tag{C.14}
\end{equation*}
$$

where $N \Gamma_{\alpha}$ is the (constant) effective potential on the arc $\alpha$. The two periods 1 and $\tau$ of the $\mathrm{sn}^{2}$ function correspond, respectively, in term of eigenvalues to (a) transfer a single EV from the first arc to the second one $(\delta N x= \pm 1)$, (b) or to shift the phase between the two arcs by $2 \pi\left(\delta N \Delta_{c}= \pm 1\right)$.

## Appendix D. Multicut case

Consider now a support of eigenvalues split into $s$ intervals:

$$
\begin{equation*}
\mathcal{C}=\mathcal{C}_{1} \cup \cdots \cup \mathcal{C}_{s} . \tag{D.1}
\end{equation*}
$$

Let $n_{i}$ be the number of eigenvalues in each $\mathcal{C}_{i}$, and $x_{i}=n_{i} / N$ the occupation ratio, which we denote collectively as a vector

$$
\begin{equation*}
x_{i}=\int_{\mathcal{C}_{i}} \mathrm{~d} \lambda \rho(\lambda) \quad \vec{x}=\left(x_{1}, \ldots, x_{s-1}\right) \tag{D.2}
\end{equation*}
$$

Note that only $s-1$ of them are independent since $x_{1}+\cdots+x_{s}=1$.
As in the two-cut case (equation (2.36)), the free energy at fixed $\vec{n}$ admits a topological large- $N$ expansion

$$
\begin{equation*}
F[V ; \vec{n}]=N^{2} F_{0}[V ; \vec{x}]+N^{0} F_{1}[V ; \vec{x}]+\mathcal{O}\left(1 / N^{2}\right) \tag{D.3}
\end{equation*}
$$

and as in (2.38), the partition function can be written as a sum over $\vec{n}$

$$
\begin{equation*}
Z=\mathrm{e}^{-F}=\sum_{\vec{n}} \mathrm{e}^{-F[V, \vec{n}]} \tag{D.4}
\end{equation*}
$$

The sum is dominated by the vicinity of the extremum $\vec{x}_{c}$ of $F_{0}[V ; \vec{x}]$ :
$Z \sim \sum_{\vec{n}} \mathrm{e}^{-N^{2}\left(F_{0}\left[V, \vec{x}_{c}\right]+\mathrm{i} \pi\left(\vec{n} / N-\vec{x}_{c}\right) \cdot \tau^{-1}\left(\vec{n} / N-\vec{x}_{c}\right)\right)} \quad$ where $\left.\quad \frac{\partial}{\partial \vec{x}} F_{0}(\vec{x})\right|_{\vec{x}=\vec{x}_{c}}=\overrightarrow{0}$
$\tau$ is the $s-1 \times s-1$ matrix defined by

$$
\begin{equation*}
\tau_{i j}^{-1}=\left.\frac{1}{2 \mathrm{i} \pi} \frac{\partial^{2} F_{0}}{\partial x_{i} \partial x_{j}}\right|_{\vec{x}=\bar{x}_{c}} . \tag{D.6}
\end{equation*}
$$

Then, the summation over $\vec{n}$ yields

$$
\begin{equation*}
Z \sim \mathrm{e}^{-N^{2} F_{0}\left[V, \vec{x}_{c}\right]} \theta\left(N \vec{x}_{c} \mid \tau\right) \tag{D.7}
\end{equation*}
$$

where $\theta(\vec{u} \mid \tau)$ is Riemann's theta function [30] in genus $s-1$ :

$$
\begin{equation*}
\theta(\vec{u} \mid \tau)=\theta(\vec{u})=\sum_{\vec{n}} \mathrm{e}^{-\mathrm{i} \pi(\vec{n}-\vec{u}) \cdot \tau^{-1}(\vec{n}-\vec{u})}=\sum_{\vec{n}} \mathrm{e}^{\mathrm{i} \pi \vec{n} \cdot \tau \vec{n}} \mathrm{e}^{-2 \mathrm{i} \pi \vec{n} \cdot \vec{u}} \tag{D.8}
\end{equation*}
$$

where $\tau$ is an $s-1 \times s-1$ matrix, $\vec{u}$ is an $s-1$ component vector, and $\vec{n}$ is a vector with integer coordinates.

The $\theta$ function obeys the relations ( $\vec{k}$ being an arbitrary integer vector):
$\theta(\vec{u}+\vec{k})=\theta(\vec{u}) \quad \theta(\vec{u}+\tau \vec{k})=\mathrm{e}^{-\mathrm{i} \pi(2 \vec{u} \cdot \vec{k}+\vec{k} \cdot \tau \vec{k})} \theta(\vec{u}) \quad \theta(-\vec{u})=\theta(\vec{u})$.
Eventually, the free energy at leading orders in $N$ is

$$
\begin{equation*}
F \sim N^{2}\left[F_{0}\left[V, \vec{x}_{c}\right]-\frac{1}{N^{2}} \ln \theta\left(N \vec{x}_{c} \mid \tau\right)+\frac{1}{N^{2}} F_{1}\left(\vec{x}_{c}\right)+\cdots\right] \tag{D.10}
\end{equation*}
$$

It is now straightforward but lengthy to rederive the two-point correlation function and the orthogonal polynomials from (D.10). One needs to differentiate (D.10) with respect to variations of the potential as in (3.3) or (4.16), and express the hyperelliptical functions involved in the calculation through prime forms (hyperelliptical generalization of the $\theta_{1}$ function) [30]. One should thus obtain expressions similar to those of [17].

## References

[1] Guhr T, Müller-Groeling A, Weidenmüller H A 1998 Phys. Rep. 299 189-421
(Guhr T, Müller-Groeling A, Weidenmüller H A 1997 Preprint cond-mat/9707301)
Mello P A 1994 Theory of Random Matrices (Les Houches) Session LXI ed E Akkermans, G Montambaux, J L Pichard and J Zinn-Justin (Amsterdam: North-Holland)
[2] Mehta M L 1991 Random Matrices 2nd edn (New York: Academic)
[3] Wigner E P 1951 Proc. Camb. Phil. Soc. 47790
Porter C E 1991 Statistical Theories of Spectra: Fluctuations (New York: Academic)
[4] Leff H S 1964 J. Math. Phys. 5763
[5] Hacken G, Werbin R and Rainwater J 1978 Phys. Rev. C 1743
Wilson W M, Bilpuch E G and Mitchell G E 1975 Nucl. Phys. A 245
Haq R U, Pandey A and Bohigas O 1982 Phys. Rev. Lett. 481086
Camarda S and Georgopulos P D 1983 Phys. Rev. Lett. 50492
[6] Camarda H S 1992 Phys. Rev. A 45579
Kobayakawa T S, Hatsugai Y, Kohmoto M and Zee A 1995 Phys. Rev. E 51 5365-70
[7] Brézin E and Zee A 1993 Nucl. Phys. B402
Brézin E and Zee A 1995 Nucl. Phys. B 453531
Brezin E and Zee A 1994 Phys. Rev. E 452588
[8] Beenakker C W J 1994 Nucl. Phys. B 422 [FS] 515
Forrester P J 1995 Nucl. Phys. B 435 [FS] 421
[9] Ambjørn J, Jurkiewicz J and Makeenko Y 1990 Phys. Lett. B 251517
[10] Brézin E, Itzykson C, Parisi G and Zuber J B 1978 Commun. Math. Phys. 5935
Bessis D, Itzykson C and Zuber J B 1980 Adv. Appl. Math. 1109
[11] Ambjørn J and Akemann G 1996 J. Phys. A: Math. Gen. 29 L555
[12] Bleher P and Its A 1999 Semiclassical asymptotics of orthogonal polynomials, Riemann-Hilbert problem and universality in the matrix model Ann. Math. 150 185-266
[13] Kanzieper E and Freilikher V 1998 Phys. Rev. E 576604
(Kanzieper E and Freilikher V 1997 Preprint cond-mat/9709309)
Kanzieper E and Freilikher V 1999 Spectra of large random matrices: a method of study Diffuse Waves in Complex Media (NATO ASI Series vol 531) (Dordrecht: Kluwer) pp 165-211
(Kanzieper E and Freilikher V 1998 Preprint cond-mat/9809365)
[14] Deo N 1997 Nucl. Phys. B 504 609-20 (Deo N 1997 Preprint cond-mat/9703136)
[15] Brézin E and Deo N 1999 Phys. Rev. E 59 3901-10 (Brézin E and Deo N 1998 Preprint cond-mat/98050096)
[16] 't Hooft G 1974 Nucl. Phys. B 72461
[17] Deift P et al 1999 Commun. Pure Appl. Math. 52 1491-552 Deift P et al 1999 Commun. Pure Appl. Math. 52 1335-425
[18] Jurkiewicz J 1990 Phys. Lett. 245 B 178
[19] David F 1991 Nucl. Phys. B 348507
[20] Akemann G 1996 Nucl. Phys. B 482 403-30 (Akemann G 1996 Preprint hep-th/9606004) Akemann G 1995 Preprint hep-th/9503185
[21] Szegö G 1967 Orthogonal Polynomials (Providence, RI: American Mathematical Society)
[22] See Mehta M L 1991 Random Matrices 2nd edn (New York: Academic) ch 19 Dyson F J 1970 Commun. Math. Phys. 19235 More recently, Fox D F and Kahn P B 1964 Phys. Rev. 134B 1151 Leff H S 1964 J. Math. Phys. 5763 Bronk B V 1965 J. Math. Phys. 6228 Zinn-Justin P 1997 Nucl. Phys. B 497725 (Zinn-Justin P 1997 Preprint cond-mat/9703033)
[23] Dyson F J 1970 Commun. Math. Phys. 19 235-50
[24] Eynard B 1997 Nucl. Phys. B 506 633-64 (Eynard B 1997 Preprint cond-mat/9707005)
[25] Jurkiewicz J 1991 Phys. Lett. B 261 260-8 Senechal D 1992 Int. J. Mod. Phys. A 71491 Sasaki M and Suzuki H 1991 Phys. Rev. D 434015

Lechtenfeld O, Ray R and Ray A 1991 Int. J. Mod. Phys. A 64491
Lechtenfeld O 1992 Int. J. Mod. Phys. A 72335
Lechtenfeld O 1992 Int. J. Mod. Phys. A 77097
(Lechtenfeld O 1991 Preprint hep-th/9112045)
[26] Wang Z X, Guo D R 1989 Special Functions (Singapore: World Scientific)
[27] Erdélyi A (ed) 1953 Higher Transcendental Functions (Bateman Manuscript Project Vol II) (New York: McGraw-Hill)
[28] Abramowitz M and Stegun E A (ed) 1955 Handbook of Mathematical Functions (NBS Appl. Math. Ser. vol 55)
[29] David F 1993 Phys. Lett. B 302 403-10
(David F 1992 Preprint hep-th/9212106)
[30] Fay J D 1973 Theta Functions on Riemann Surfaces (Berlin: Springer)

