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A phase transition for the limiting spectral density of random matrices^{*}

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Abstract

We analyze the spectral distribution of symmetric random matrices with correlated entries. While we assume that the diagonals of these random matrices are stochastically independent, the elements of the diagonals are taken to be correlated. Depending on the strength of correlation, the limiting spectral distribution is either the famous semicircle distribution, the distribution derived for Toeplitz matrices by Bryc, Dembo and Jiang (2006), or the free convolution of the two distributions.

Keywords: random matrices; dependent random variables; Toeplitz matrices; semicircle law; Curie-Weiss model.

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1 Introduction

Historically, the theory of random matrices is fed by two sources. They were introduced in mathematical statistics by the seminal work of Wishart [20]. On the other hand, Wigner used random matrices as a toy model for the energy levels and excitation spectra of heavy nuclei [19]. From these two roots random matrix theory has grown into an independent mathematical theory with applications in many areas of science.

A central role in the study of random matrices with growing dimension is played by their eigenvalues. To introduce them let, for any $n \in \mathbb{N}$, $\{a_n(p,q), 1 \le p \le q \le n\}$ be a real valued random field. Define the symmetric random $n \times n$ matrix \mathbf{X}_n by

$$\mathbf{X}_n(q,p) = \mathbf{X}_n(p,q) = \frac{1}{\sqrt{n}}a_n(p,q), \qquad 1 \le p \le q \le n.$$

We will denote the (real) eigenvalues of \mathbf{X}_n by $\lambda_1^{(n)} \leq \lambda_2^{(n)} \leq \ldots \leq \lambda_n^{(n)}$. Let μ_n be the empirical eigenvalue distribution, i.e.

$$\mu_n = \frac{1}{n} \sum_{k=1}^n \delta_{\lambda_k^{(n)}}.$$

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Wigner proved in his fundamental work [19] that, if the entries $a_n(p,q)$ are independent Bernoulli variables, the expected empirical eigenvalue distribution converges weakly to the so called semicircle distribution (or law), i.e. the probability distribution ν on \mathbb{R} with density

$$\nu(dx) = \frac{1}{2\pi} \sqrt{4 - x^2} \mathbb{1}_{|x| \le 2}.$$

Quite some effort has been spent in investigating the universality of this result. Arnold [2] showed that the convergence to the semicircle law is also true if one replaces the Gaussian distributed random variables by independent and identically distributed (i.i.d.) random variables with a finite fourth moment. Also the identical distribution may be replaced by some other assumptions (see e.g. [9]). Recently, it was observed by Erdös et al. ([10]) that the convergence of the spectral measure towards the semicircle law holds in a local sense. More precisely, it can be proved that on intervals with width going to zero sufficiently slowly, the empirical eigenvalue distribution still converges to the semicircle distribution.

This result therefore interpolates between the global and the local behavior of the eigenvalues in the bulk of the spectrum, which was rather recently proved to be universal as well in the so-called "four-moment-theorem" ([18]).

Other generalizations of Wigner's semicircle law concern matrix ensembles with entries drawn according to weighted Haar measures on classical (e.g., orthogonal, unitary, symplectic) groups. Such results are particularly interesting since such random matrices also play a major role in non-commutative probability (see e.g. [13], or the very recommendable book Anderson, Guionnet, and Zeitouni [1]).

A slightly different approach to universality was taken in [14], [12], [16] and [11]. Here, matrices with correlated entries are studied. In [11] it is shown that, if the diagonals of \mathbf{X}_n are independent and the correlation between elements along a diagonal decays sufficiently quickly, again the limiting spectral distribution is the semicircle law.

Universality, however, does have its limitations. As was shown by Bryc et al. [5] the limiting spectral distribution of large random Toeplitz or Hankel matrices is not the semicircle law. In fact, not much is known about the limiting measures, apart from their moments (which are the result of the proof by a moment method, a technique, that will also be employed by the present paper).

The present note tries to explore the borderline between the weak correlations studied in [11] and the strong correlations that lead to a limiting spectral distribution that is not of Wigner type. We will again assume that \mathbf{X}_n has independent diagonals and we will see, which quantity determines whether the limiting measure of the empirical eigenvalue distribution is a semicircle law or not. A particularly nice example is borrowed from statistical mechanics. There the Curie-Weiss model is the easiest model of a ferromagnet. Here a magnetic substance has little atoms that carry a magnetic spin, that is either +1 or -1. These spins interact in cooperative way, the strength of the interaction being triggered by a parameter, the so-called inverse temperature. The model exhibits phase transition from paramagnetic to magnetic behavior (the standard reference for the Curie-Weiss model is [8]). We will see that this phase transition can be recovered on the level of the limiting spectral distribution of random matrices, if we fill their diagonals independently with the spins of Curie-Weiss models. For small interaction parameter, this limiting spectral distribution is the semicircle law, while for a large interaction parameter we obtain a distribution similar to the Toeplitz case.

The rest of this paper is organized as follows. Section 2 contains the technical assumptions we have to make together with the statement of our main results. Section 3 characterizes the various limiting distributions we obtain. Section 4 contains some interesting examples, while Sections 5 and 6 are devoted to the proofs of the two main

theorems.

2 Main Result

This section contains the general theorem that describes the various limiting spectral distributions for the matrices \mathbf{X}_n introduced above. In order to be able to state the theorem we will have to impose the following conditions on \mathbf{X}_n :

(C1)
$$\mathbb{E}[a_n(p,q)] = 0$$
, $\mathbb{E}[a_n(p,q)^2] = 1$ and

$$m_k := \sup_{n \in \mathbb{N}} \max_{1 \le p \le q \le n} \mathbb{E}\left[|a_n(p,q)|^k \right] < \infty, \quad k \in \mathbb{N}.$$
(2.1)

- (C2) the diagonals of \mathbf{X}_n , i.e. the families $\{a_n(p, p+r), 1 \le p \le n-r\}, 0 \le r \le n-1$, are independent,
- (C3) the covariance of two entries on the same diagonal depends only on n, i.e. for any $0 \le r \le n-1$ and $1 \le p, q \le n-r$, $p \ne q$, we can define

$$\operatorname{Cov}(a_n(p, p+r), a_n(q, q+r)) =: c_n,$$

(C4) the limit $c := \lim_{n \to \infty} c_n$ exists.

Remark 2.1. Note that the assumptions above imply that $0 \le c \le 1$. Indeed, take the process $\{a_n(p,p), 1 \le p \le n\}$ on the main diagonal, and calculate

$$0 \le \mathbb{V}\left(\sum_{p=1}^{n} a_n(p, p)\right) = \sum_{p=1}^{n} \mathbb{V}(a_n(p, p)) + \sum_{\substack{p,q=1, \\ p \ne q}}^{n} \operatorname{Cov}(a_n(p, p), a_n(q, q))$$
$$= n + n(n-1)c_n,$$

implying that $c_n \ge -(1/(n-1))$. Since the right hand side tends to zero, we can conclude that $c = \lim_{n\to\infty} c_n \ge 0$. On the other hand, Hölder's inequality yields $c_n \le 1$ since $\mathbb{E}\left[a_n(p,p)^2\right] = 1$ by (C1). Thus, we have $c \le 1$.

With these notations and conditions we are able to formulate the central result of this note.

Theorem 2.2. Assume that the symmetric random matrix X_n as defined above satisfies the conditions (C1), (C2), (C3) and (C4). Then, with probability 1, the empirical spectral distribution μ_n of X_n converges weakly to a nonrandom probability distribution ν_c which does not depend on the distribution of the entries of X_n .

Since the proof of Theorem 2.2 relies on the so-called moment-method, we will describe ν_c in terms of its moments in Section 3. However, to give an idea of the kind of measure we deal with, we first want to recall the notion of the *free convolution*. Therefore, let μ_1 and μ_2 be two probability measures on \mathbb{R} which are uniquely determined by their moments. Let \mathcal{A} be a unital C^* -algebra over \mathbb{C} and $\varphi : \mathcal{A} \to \mathbb{C}$ a unital linear functional satisfying $\varphi(a^*a) \geq 0$ for any $a \in \mathcal{A}$. Then, (\mathcal{A}, φ) is a C^* -probability space. We say that two elements $x_1, x_2 \in \mathcal{A}$ are *freely independent* if for any $k \in \mathbb{N}$, polynomials P_1, \ldots, P_k , and $i(1), \ldots, i(k) \in \{1, 2\}$ with $i(1) \neq i(2) \neq \ldots \neq i(k)$, we have

$$\varphi(P_j(x_{i(j)})) = 0$$
 for any $j = 1, \dots, k \implies \varphi(P_1(x_{i(1)}) \cdots P_k(x_{i(k)})) = 0.$

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Assume that $x_1, x_2 \in A$ are selfadjoint and freely independent with distributions μ_1 and μ_2 , respectively, i.e.

$$\varphi\left(x_{i}^{k}\right) = \int_{\mathbb{R}} t^{k} d\mu_{i}(t), \quad i = 1, 2, \ k \in \mathbb{N}.$$

Then the distribution of the sum $x_1 + x_2$ is called the *free convolution* of μ_1 and μ_2 and is denoted by $\mu_1 \boxplus \mu_2$. For more details, we refer to [15]. Returning to the measure ν_c , we now have the following statement.

Theorem 2.3. For any $0 \le c \le 1$, we have $\nu_c = \nu_{0,1-c} \boxplus \nu_{1,c}$ with $\nu_{0,1-c}$ denoting the rescaled semicircle law with variance 1 - c, and $\nu_{1,c}$ the rescaled Toeplitz law with variance c. In particular, ν_c is a symmetric measure with a bounded density. If c > 0, ν_c has an unbounded support, and if 0 < c < 1, the density is smooth.

3 The Limiting Distribution ν_c

It is not surprising that ν_c is some combination of the semicircle distribution and the limiting distribution of Toeplitz matrices as described in [5]. Indeed, c = 0 covers the case of independent entries implying that ν_0 is the semicircle law. On the other hand, considering symmetric Toeplitz matrices, we have c = 1, and thus ν_1 is the corresponding limiting distribution we want to introduce in the following (cf. [5]). Therefore, we have to start with some notation. For any even $k \in \mathbb{N}$, let $\mathcal{PP}(k)$ denote the set of all pair partitions π of $\{1, \ldots, k\}$. If i and j are in the same block of π , we also write $i \sim_{\pi} j$. The measure ν_1 can be defined with the help of *Toeplitz volumes*. Thus, we associate to any partition $\pi \in \mathcal{PP}(k)$ the following system of equations in unknowns x_0, \ldots, x_k :

$$\begin{array}{l}
x_{1} - x_{0} + x_{l_{1}} - x_{l_{1}-1} = 0, \quad \text{if } 1 \sim_{\pi} l_{1}, \\
x_{2} - x_{1} + x_{l_{2}} - x_{l_{2}-1} = 0, \quad \text{if } 2 \sim_{\pi} l_{2}, \\
\vdots \\
x_{i} - x_{i-1} + x_{l_{i}} - x_{l_{i}-1} = 0, \quad \text{if } i \sim_{\pi} l_{i}, \\
\vdots \\
x_{k} - x_{k-1} + x_{l_{k}} - x_{l_{k}-1} = 0, \quad \text{if } k \sim_{\pi} l_{k}.
\end{array}$$
(3.1)

Since π is a pair partition, we in fact have only k/2 equations although we have listed k. However, we have k + 1 variables. If $\pi = \{\{i_1, j_1\}, \ldots, \{i_{k/2}, j_{k/2}\}\}$ with $i_l < j_l$ for any $l = 1, \ldots, k/2$, we solve (3.1) for $x_{j_1}, \ldots, x_{j_{k/2}}$, and leave the remaining variables undetermined. We further impose the condition that all variables x_0, \ldots, x_k lie in the interval I = [0, 1]. Solving the equations above in this way determines a cross section of the cube $I^{k/2+1}$. The volume of this will be denoted by $p_T(\pi)$.

Returning to the measure ν_1 , we can use the results in [5] to see that all odd moments of ν_1 are zero, and for any even $k \in \mathbb{N}$, the k-th moment is given by

$$\int x^k d\nu_1(x) = \sum_{\pi \in \mathcal{PP}(k)} p_T(\pi).$$

The expression above is bounded by (k-1)!!. Hence, Carleman's condition is satisfied implying that the distribution ν_1 is uniquely determined by its moments. Moreover, it has an unbounded support as verified in [5]. To describe ν_c for general $c \in [0, 1]$, we need a further definition which was introduced in [5] to analyze Markov matrices.

Definition 3.1. Let $k \in \mathbb{N}$ be even, and fix $\pi \in \mathcal{PP}(k)$. The height $h(\pi)$ of π is the number of elements $i \sim_{\pi} j$, i < j, such that either j = i + 1 or the restriction of π to $\{i + 1, \ldots, j - 1\}$ is a pair partition.

Note that the property that the restriction of π to $\{i + 1, \ldots, j - 1\}$ is a pair partition in particular requires that the distance $j - i - 1 \ge 1$ is even. To give an example how to calculate the height of a partition, take $\pi = \{\{1, 6\}, \{2, 4\}, \{3, 5\}\}$. Considering the block $\{1, 6\}$, we see that the restriction of π to $\{2, 3, 4, 5\}$ is a pair partition, namely $\{\{2, 4\}, \{3, 5\}\}$. However, this is not true for both remaining blocks. Hence, $h(\pi) = 1$.

In the following, we say that a pair partition π is crossing if there are indices i < j < l < m with $i \sim_{\pi} l$ and $j \sim_{\pi} m$. Otherwise, we call π non-crossing. We will denote the set of all crossing pair partitions of $\{1, \ldots, k\}$ by CPP(k), and the set of non-crossing pair partitions of $\{1, \ldots, k\}$ by NPP(k). Note that for $\pi \in NPP(k)$, we have the height $h(\pi) = k/2$ and the Toeplitz volume $p_T(\pi) = 1$.

In Section 5, we will see that all odd moments of ν_c vanish, implying that ν_c is symmetric. The even moments are given by

$$\int x^k d\nu_c(x) = C_{\frac{k}{2}} + \sum_{\pi \in \mathcal{CPP}(k)} p_T(\pi) c^{\frac{k}{2} - h(\pi)} = \sum_{\pi \in \mathcal{PP}(k)} p_T(\pi) c^{\frac{k}{2} - h(\pi)},$$
(3.2)

where $C_k = \frac{(2k)!}{k!(k+1)!}$ denotes the k-th Catalan number. Note that the number of elements in $\mathcal{NPP}(k)$ coincides with the Catalan number $C_{k/2}$. The latter is exactly the k-th moment of the semicircle distribution. As for the limiting distribution in the Toeplitz case, we can verify the Carleman condition to see that ν_c is uniquely determined by its moments.

4 Examples

In this section, we want to give some examples of processes satisfying the assumptions of Theorem 2.2.

4.1 Toeplitz Matrices

Consider a symmetric Toeplitz matrix. The limiting spectral distribution calculated in [5] can be deduced from Theorem 2.2 as well. Indeed, assuming that the entries are centered with unit variance and have existing moments of any order, we see that all conditions are satisfied with $c = c_n = 1$. Thus, we get

$$\int x^k d\nu_1(x) = \begin{cases} C_{\frac{k}{2}} + \sum_{\pi \in \mathcal{CPP}(k)} p_T(\pi) = \sum_{\pi \in \mathcal{PP}(k)} p_T(\pi), & \text{if } k \text{ is even,} \\ 0, & \text{if } k \text{ is odd.} \end{cases}$$

4.2 Exchangeable Random Variables

In [6], it was shown that symmetric matrices with exchangeable entries above the main diagonal, and an appropriate scaling, still obey the semicircle law. In our situation, we suppose that for any $n \in \mathbb{N}$, we have a family $\{x_n(p), 1 \le p \le n\}$ of exchangeable random variables, i.e. the distribution of the vector $(x_n(1), \ldots, x_n(n))$ is the same as that of $(x_n(\sigma(1)), \ldots, x_n(\sigma(n)))$ for any permutation σ of $\{1, \ldots, n\}$. In this case, we can conclude that for any $1 \le p < q \le n$, we have

$$Cov(x_n(p), x_n(q)) = Cov(x_n(1), x_n(2)) =: c_n.$$

Now assume that $c_n \to c \in \mathbb{R}$ as $n \to \infty$. Define for any $n \in \mathbb{N}$, $r \in \{0, \ldots, n-1\}$, the process $\{a_n(p, p+r), 1 \le p \le n-r\}$ to be an independent copy of $\{x_n(p), 1 \le p \le n-r\}$.

Then, all conditions of Theorem 2.2 are satisfied if we ensure that the moment condition (C1) holds. The resulting limiting distribution for different choices of c is depicted in Figure 1.



Figure 1: Histograms of the empirical spectral distribution of 100 realizations of 1000×1000 matrices \mathbf{X}_{1000} with standard Gaussian entries.

An example for a process with exchangeable variables is the Curie-Weiss model with inverse temperature $\beta > 0$. Here, the vector $x_n = (x_n(1), \ldots, x_n(n))$ takes values in $\{-1, 1\}^n$, and for any $\omega = (\omega(1), \ldots, \omega(n)) \in \{-1, 1\}^n$, we have

$$\mathbb{P}(x_n = \omega) = \frac{1}{Z_{n,\beta}} \exp\left(\frac{\beta}{2n} \left(\sum_{i=1}^n \omega(i)\right)^2\right),\,$$

where $Z_{n,\beta}$ is the normalizing constant. Since $\mathbb{P}(x_n(1) = -1) = \mathbb{P}(x_n(1) = 1) = \frac{1}{2}$, we obtain $\mathbb{E}[x_n(1)] = 0$. Further, we clearly have $\mathbb{E}[x_n(1)^2] = 1$. It remains to determine $c = \lim_{n \to \infty} c_n$. Therefore, we want to make use of the identity

$$c_n = \operatorname{Cov}(x_n(1), x_n(2)) = \mathbb{E}[x_n(1)x_n(2)] = \frac{n}{n-1}\mathbb{E}[m_n^2] - \frac{1}{n-1},$$

where $m_n := \frac{1}{n} \sum_{i=1}^n x_n(i)$ is the so-called magnetization of the system. Since $|m_n| \leq 1$, we see that $(m_n^2)_{n \in \mathbb{N}}$ is uniformly integrable. Thus, m_n converges in \mathcal{L}^2 to some random variable m if and only if $m_n \to m$ in probability. In [7], it was verified that $m_n \to 0$ in probability if $\beta \leq 1$, and $m_n \to m$ in probability with $m \sim \frac{1}{2} \delta_{m(\beta)} + \frac{1}{2} \delta_{-m(\beta)}$ for some $m(\beta) > 0$ if $\beta > 1$. The mapping $\beta \mapsto m(\beta)$ is monotonically increasing on $(1, \infty)$, and satisfies $m(\beta) \to 0$ as $\beta \searrow 1$ and $m(\beta) \to 1$ as $\beta \to \infty$. We now obtain

$$c = \lim_{n \to \infty} c_n = \begin{cases} 0, & \text{if } \beta \le 1, \\ m(\beta)^2, & \text{if } \beta > 1. \end{cases}$$

Thus, the limiting spectral distribution of \mathbf{X}_n is the semicircle law if $\beta \leq 1$, and approximately the Toeplitz limit if β is large. This is insofar not surprising as the different

sites in the Curie-Weiss model show little interaction, i.e. behave almost independently, if the temperature is high, or, in other words, β is small. However, if the temperature is low, i.e. β is large, the magnetization of the sites strongly depends on each other. The phase transition at the critical inverse temperature $\beta = 1$ in the Curie-Weiss model is thus reflected in the limiting spectral distribution of \mathbf{X}_n as well.

Proof of Theorem 2.2 5

The main technique we want to apply is the method of moments. The idea is to first determine the weak limit of the expected empirical spectral distribution. Therefore, the similar structure of the matrices under consideration allows us to repeat some concepts presented in [11]. However, we need to develop new ideas when calculating the expectations of the entries.

5.1 The expected empirical spectral distribution

To determine the limit of the k-th moment of the expected empirical spectral distribution μ_n of \mathbf{X}_n , we write

$$\mathbb{E}\left[\int x^k d\mu_n(x)\right] = \frac{1}{n} \mathbb{E}\left[\operatorname{tr}\left(\mathbf{X}_n^k\right)\right]$$
$$= \frac{1}{n^{\frac{k}{2}+1}} \sum_{p_1,\dots,p_k=1}^n \mathbb{E}\left[a_n(p_1,p_2)a_n(p_2,p_3)\cdots a_n(p_{k-1},p_k)a_n(p_k,p_1)\right].$$

The main task is now to compute the expectations on the right hand side. However, we have to face the problem that some of the entries involved are independent and some are not. To be more precise, $a_n(p_1,q_1),\ldots,a_n(p_j,q_j)$ are independent whenever they can be found on different diagonals of \mathbf{X}_n , i.e. the distances $|p_1 - p_1|$ $q_1|,\ldots,|p_j-q_j|$ are distinct. Hence, a first step in our proof is to consider the expectation $\mathbb{E}[a_n(p_1,p_2)a_n(p_2,p_3)\cdots a_n(p_{k-1},p_k)a_n(p_k,p_1)]$, and to identify entries with the same distance of their indices. Therefore, we want to adapt some concepts of [16] and [5] to our situation.

To start with, fix $k \in \mathbb{N}$, and define $\mathcal{T}_n(k)$ to be the set of k-tuples of *consistent pairs*, that is multi-indices (P_1, \ldots, P_k) satisfying for any $j = 1, \ldots, k$,

- (i) $P_j = (p_j, q_j) \in \{1, \dots, n\}^2$, (ii) $q_j = p_{j+1}$, where k + 1 is cyclically identified with 1.

With this notation, we find that

$$\frac{1}{n}\mathbb{E}\left[\operatorname{tr}\left(\mathbf{X}_{n}^{k}\right)\right] = \frac{1}{n^{\frac{k}{2}+1}}\sum_{\left(P_{1},\ldots,P_{k}\right)\in\mathcal{T}_{n}(k)}\mathbb{E}\left[a_{n}(P_{1})\cdots a_{n}(P_{k})\right].$$

To reflect the dependency structure among the entries $a_n(P_1) \dots a_n(P_k)$, we want to make use of the set $\mathcal{P}(k)$ of partitions of $\{1, \ldots, k\}$. Thus, take $\pi \in \mathcal{P}(k)$. We say that an element $(P_1, \ldots, P_k) \in \mathcal{T}_n(k)$ is a π -consistent sequence if

$$|p_i - q_i| = |p_j - q_j| \quad \Longleftrightarrow \quad i \sim_{\pi} j.$$

According to condition (C2), this implies that $a_n(P_{i_1}), \ldots, a_n(P_{i_l})$ are stochastically independent if i_1, \ldots, i_l belong to l different blocks of π . The set of all π -consistent sequences $(P_1, \ldots, P_k) \in \mathcal{T}_n(k)$ is denoted by $S_n(\pi)$. Note that the sets $S_n(\pi)$, $\pi \in \mathcal{P}(k)$, are pairwise disjoint, and $\bigcup_{\pi \in \mathcal{P}(k)} S_n(\pi) = \mathcal{T}_n(k)$. Consequently, we can write

$$\frac{1}{n}\mathbb{E}\left[\operatorname{tr}\left(\mathbf{X}_{n}^{k}\right)\right] = \frac{1}{n^{\frac{k}{2}+1}}\sum_{\pi\in\mathcal{P}(k)}\sum_{(P_{1},\ldots,P_{k})\in S_{n}(\pi)}\mathbb{E}\left[a_{n}(P_{1})\cdots a_{n}(P_{k})\right].$$
(5.1)

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In a next step, we want to exclude partitions that do not contribute to (5.1) as $n \to \infty$. These are those partitions satisfying either $\#\pi > \frac{k}{2}$ or $\#\pi < \frac{k}{2}$, where $\#\pi$ denotes the number of blocks of π . We want to treat the two cases separately.

First case: $\#\pi > \frac{k}{2}$. Since π is a partition of $\{1, \ldots, k\}$, there is at least one singleton, i.e. a block containing only one element *i*. Consequently, $a_n(P_i)$ is independent of $\{a_n(P_j), j \neq i\}$ if $(P_1, \ldots, P_k) \in S_n(\pi)$. Since we assumed the entries to be centered, we obtain

$$\mathbb{E}\left[a_n(P_1)\cdots a_n(P_k)\right] = \mathbb{E}\left[\prod_{j\neq i}a_n(P_j)\right]\mathbb{E}\left[a_n(P_i)\right] = 0.$$

This yields

$$\frac{1}{n^{\frac{k}{2}+1}} \sum_{(P_1,\dots,P_k)\in S_n(\pi)} \mathbb{E}\left[a_n(P_1)\cdots a_n(P_k)\right] = 0$$

Second case: $r := \#\pi < \frac{k}{2}$. Here, we want to argue that π gives vanishing contribution to (5.1) as $n \to \infty$ by calculating $\#S_n(\pi)$. To fix an element $(P_1, \ldots, P_k) \in S_n(\pi)$, we first choose the pair $P_1 = (p_1, q_1)$. There are at most n possibilities to assign a value to p_1 , and another n possibilities for q_1 . To fix $P_2 = (p_2, q_2)$, note that the consistency of the pairs implies $p_2 = q_1$. If now $1 \sim_{\pi} 2$, the condition $|p_1 - q_1| = |p_2 - q_2|$ allows at most two choices for q_2 . Otherwise, if $1 \not\sim_{\pi} 2$, we have at most n possibilities. We now proceed sequentially to determine the remaining pairs. When arriving at some index i, we check whether i is in the same block as some preceding index $1, \ldots, i - 1$. If this is the case, then we have at most two choices for P_i and otherwise, we have n. Since there are exactly $r = \#\pi$ different blocks, we can conclude that

$$\#S_n(\pi) \le n^2 n^{r-1} 2^{k-r} \le C \ n^{r+1}$$
(5.2)

with a constant C = C(r, k) depending on r and k.

Now the uniform boundedness of the moments (2.1) and the Hölder inequality together imply that for any sequence (P_1, \ldots, P_k) ,

$$\left|\mathbb{E}\left[a_{n}(P_{1})\cdots a_{n}(P_{k})\right]\right| \leq \left[\mathbb{E}\left|a_{n}(P_{1})\right|^{k}\right]^{\frac{1}{k}}\cdots \left[\mathbb{E}\left|a_{n}(P_{k})\right|^{k}\right]^{\frac{1}{k}} \leq m_{k}.$$
(5.3)

Consequently, taking account of the relation $r < \frac{k}{2}$, we get

$$\frac{1}{n^{\frac{k}{2}+1}} \sum_{(P_1,\dots,P_k)\in S_n(\pi)} \left| \mathbb{E}\left[a_n(P_1)\cdots a_n(P_k) \right] \right| \le C \ \frac{\#S_n(\pi)}{n^{\frac{k}{2}+1}} \le C \ \frac{1}{n^{\frac{k}{2}-r}} = o(1).$$

Combining the calculations in the first and the second case, we can conclude that

$$\frac{1}{n}\mathbb{E}\left[\operatorname{tr}\left(\mathbf{X}_{n}^{k}\right)\right] = \frac{1}{n^{\frac{k}{2}+1}}\sum_{\substack{\pi\in\mathcal{P}(k),\ (P_{1},\ldots,P_{k})\in S_{n}(\pi)\\ \#\pi=\frac{k}{2}}}\mathbb{E}\left[a_{n}(P_{1})\cdots a_{n}(P_{k})\right] + o(1).$$

Now assume that k is odd. Then the condition $\#\pi = \frac{k}{2}$ cannot be satisfied, and the considerations above immediately yield

$$\lim_{n \to \infty} \frac{1}{n} \mathbb{E}\left[\operatorname{tr} \left(\mathbf{X}_{n}^{k} \right) \right] = 0.$$

It remains to determine the even moments. Thus, let $k \in \mathbb{N}$ be even. Recall that we denoted by $\mathcal{PP}(k) \subset \mathcal{P}(k)$ the set of all pair partitions of $\{1, \ldots, k\}$. In particular, $\#\pi = \frac{k}{2}$ for any $\pi \in \mathcal{PP}(k)$. On the other hand, if $\#\pi = \frac{k}{2}$ but $\pi \notin \mathcal{PP}(k)$, we can conclude

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that π has at least one singleton and hence, as in the first case above, the expectation corresponding to the π -consistent sequences will become zero. Consequently,

$$\frac{1}{n}\mathbb{E}\left[\operatorname{tr}\left(\mathbf{X}_{n}^{k}\right)\right] = \frac{1}{n^{\frac{k}{2}+1}}\sum_{\pi\in\mathcal{PP}(k)}\sum_{(P_{1},\dots,P_{k})\in S_{n}(\pi)}\mathbb{E}\left[a_{n}(P_{1})\cdots a_{n}(P_{k})\right] + o(1).$$
(5.4)

We have now reduced the original set $\mathcal{P}(k)$ to the subset $\mathcal{PP}(k)$. Next we want to fix a $\pi \in \mathcal{PP}(k)$ and concentrate on the set $S_n(\pi)$. The following lemma will help us to calculate that part of (5.4) which involves non-crossing partitions.

Lemma 5.1 (cf. [5], Proposition 4.4.). Let $S_n^*(\pi) \subseteq S_n(\pi)$ denote the set of π -consistent sequences (P_1, \ldots, P_k) satisfying

$$i \sim_{\pi} j \implies q_i - p_i = p_j - q_j$$

for all $i \neq j$. Then, we have

$$\#\left(S_n(\pi)\backslash S_n^*(\pi)\right) = o\left(n^{\frac{k}{2}+1}\right).$$

Proof. If $(P_1, \ldots, P_k) \in S_n(\pi) \setminus S_n^*(\pi)$, we can find some $i \sim_{\pi} j$, $i \neq j$, such that $q_i - p_i \neq p_j - q_j$. However, $i \sim_{\pi} j$ implies $|p_i - q_i| = |p_j - q_j|$. We can thus conclude that $q_i - p_i = q_j - p_j$.

To fix $(P_1, \ldots, P_k) \in S_n(\pi) \setminus S_n^*(\pi)$, we first choose a π -block $\{i, j\}$ satisfying $q_i - p_i = q_j - p_j$, and then fix the signs of the differences $q_l - p_l$, $l = 1, \ldots, k$. The number of possibilities to accomplish this depends only on k and not on n. Now we choose one of n possible values for p_i , and continue with assigning values to the distances $|q_l - p_l|$ for all $l \in \{1, \ldots, k\} \setminus \{i, j\}$. The fact that π is a pair partition ensures that we have at most $n^{k/2-1}$ possibilities for the latter. Since $\sum_{l=1}^k q_l - p_l = 0$ by consistency, we find that

$$2(q_i - p_i) = q_i - p_i + q_j - p_j = \sum_{l \in \{1, \dots, k\} \setminus \{i, j\}} p_l - q_l.$$

Since we have already chosen the signs of the differences $q_l - p_l$, $l \neq i, j$, as well as their absolute values, we know the value of the sum on the right hand side. Hence, the difference $q_i - p_i = q_j - p_j$ is fixed. We thus made $C n^{k/2}$ choices to obtain the index p_i and all differences $q_l - p_l, l \in \{1, \ldots, k\}$. Starting at P_i , we can use the consistency property and go systematically through the whole sequence (P_1, \ldots, P_k) to see that it is indeed uniquely determined. Consequently, our considerations lead to

$$\#\left(S_n(\pi)\backslash S_n^*(\pi)\right) \le C \ n^{\frac{k}{2}} = o\left(n^{\frac{k}{2}+1}\right).$$

A consequence of Lemma 5.1 and relation (5.3) is the identity

$$\frac{1}{n}\mathbb{E}\left[\operatorname{tr}\left(\mathbf{X}_{n}^{k}\right)\right] = \frac{1}{n^{\frac{k}{2}+1}}\sum_{\pi\in\mathcal{PP}(k)}\sum_{(P_{1},\dots,P_{k})\in S_{n}^{*}(\pi)}\mathbb{E}\left[a_{n}(P_{1})\cdots a_{n}(P_{k})\right] + o(1).$$
(5.5)

As already mentioned, the sets $S_n^*(\pi)$ help us to deal with the set $\mathcal{NPP}(k)$ of noncrossing pair partitions.

Lemma 5.2. Let $\pi \in \mathcal{NPP}(k)$. For any $(P_1, \ldots, P_k) \in S_n^*(\pi)$, we have

$$\mathbb{E}\left[a_n(P_1)\cdots a_n(P_k)\right] = 1.$$

Proof. Let l < m with $l \sim_{\pi} m$. Since π is non-crossing, the number l - m - 1 of elements between l and m must be even. In particular, there is $l \leq i < j \leq m$ with $i \sim_{\pi} j$ and j = i + 1. By the properties of $S_n^*(\pi)$, we have $a_n(P_i) = a_n(P_j)$, and the sequence $(P_1, \ldots, P_l, \ldots, P_{i-1}, P_{i+2}, \ldots, P_m, \ldots, P_k)$ is still consistent. Applying this argument successively, all pairs between l and m vanish and we see that the sequence $(P_1, \ldots, P_l, P_m, \ldots, P_k)$ is consistent, that is $q_l = p_m$. Then, the identity $p_l = q_m$ also holds. In particular, $a_n(P_l) = a_n(P_m)$. Since this argument applies for arbitrary $l \sim_{\pi} m$, we obtain

$$\mathbb{E}\left[a_n(P_1)\cdots a_n(P_k)\right] = \prod_{\substack{l < m, \\ l \sim_{\pi} m}} \mathbb{E}\left[a_n(P_l)a_n(P_m)\right] = 1.$$

By Lemma 5.2, we can conclude that

$$\frac{1}{n^{\frac{k}{2}+1}} \sum_{\pi \in \mathcal{NPP}(k)} \sum_{(P_1, \dots, P_k) \in S_n^*(\pi)} \mathbb{E}\left[a_n(P_1) \cdots a_n(P_k)\right] = \frac{1}{n^{\frac{k}{2}+1}} \sum_{\pi \in \mathcal{NPP}(k)} \#S_n^*(\pi).$$

The following lemma allows us to finally calculate the term on the right hand side. Lemma 5.3. For any $\pi \in \mathcal{NPP}(k)$, we have

$$\lim_{n\to\infty}\frac{\#S_n^*(\pi)}{n^{\frac{k}{2}+1}}=1.$$

Proof. Since π is non-crossing, we can find a nearest neighbor pair $i \sim_{\pi} i + 1$. Now fix $(P_1, \ldots, P_k) \in S_n^*(\pi)$, and write $P_l = (p_l, p_{l+1})$, $l = 1, \ldots, k$, where k + 1 is identified with 1. Then the properties of $S_n^*(\pi)$ ensure that $(p_i, p_{i+1}) = (p_{i+2}, p_{i+1})$. Hence, we can eliminate P_i, P_{i+1} to obtain a sequence $(P_1^{(1)}, \ldots, P_{k-2}^{(1)}) := (P_1, \ldots, P_{i-1}, P_{i+2}, \ldots, P_k)$ which is still consistent. Denote by π' the partition obtained from π by deleting the block $\{i, i + 1\}$, and relabeling any $l \ge i + 2$ to l - 2. Since π is non-crossing, we have $\pi' \in \mathcal{NPP}(k-2)$. Moreover, $(P_1^{(1)}, \ldots, P_{k-2}^{(1)}) \in S_n^*(\pi')$. Thus we see that any $(P_1, \ldots, P_k) \in S_n^*(\pi)$ can be reconstructed from a tuple $(P_1^{(1)}, \ldots, P_{k-2}^{(1)}) \in S_n^*(\pi')$ and a choice of p_{i+1} . The latter admits $n - \frac{k-2}{2}$ possibilities since $\{i, i+1\}$ forms a block on its own in π . Consequently,

$$\frac{\#S_n^*(\pi)}{n^{\frac{k}{2}+1}} = \frac{\#S_n^*(\pi')}{n^{\frac{k}{2}}} + o(1).$$
(5.6)

Now if k = 2, we get $S_n^*(\pi) = \{((p,q), (q,p)) : p, q \in \{1, ..., n\}\}$, implying $\frac{\#S_n^*(\pi)}{n^2} = 1$. For arbitrary even $k \in \mathbb{N}$, the statement of Lemma 5.3 follows then by induction using the identity in (5.6).

Taking account of the relation $\#\mathcal{NPP}(k) = C_{\frac{k}{\alpha}}$, we now arrive at

$$\frac{1}{n} \mathbb{E}\left[\operatorname{tr} \left(\mathbf{X}_{n}^{k} \right) \right] = C_{\frac{k}{2}} + \frac{1}{n^{\frac{k}{2}+1}} \sum_{\pi \in \mathcal{CPP}(k)} \sum_{(P_{1}, \dots, P_{k}) \in S_{n}^{*}(\pi)} \mathbb{E}\left[a_{n}(P_{1}) \cdots a_{n}(P_{k}) \right] + o(1), \quad (5.7)$$

with CPP(k) being the set of all crossing pair partitions of $\{1, \ldots, k\}$. Since we consider only pair partitions, we know that the expectation on the right hand side is of the form

$$\mathbb{E}\left[a_{n}(p_{1},q_{1})a_{n}(p_{1}+\tau_{1},q_{1}+\tau_{1})\right]\cdots\mathbb{E}\left[a_{n}(p_{r},q_{r})a_{n}(p_{r}+\tau_{r},q_{r}+\tau_{r})\right]$$

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for $r := \frac{k}{2}$ and some choices of $p_1, q_1, \tau_1, \ldots, p_r, q_r, \tau_r \in \mathbb{N}$. In order to calculate this expectation, assumption (C3) indicates that we only need to distinguish for any $i = 1, \ldots, k$, whether we have $\tau_i = 0$ or not. In the first case, we get the identity $\mathbb{E}[a_n(p_i, q_i)a_n(p_i + \tau_i, q_i + \tau_i)] = 1$, and in the second case, we can conclude that $\mathbb{E}[a_n(p_i, q_i)a_n(p_i + \tau_i, q_i + \tau_i)] = c_n$. Now fix some pair partition $\pi \in \mathcal{PP}(k)$, and take $(P_1, \ldots, P_k) \in S_n^*(\pi)$. Motivated by these considerations, we put $P_i = (p_i, q_i)$, and define

$$m(P_1, \ldots, P_k) := \#\{1 \le i < j \le k : (p_i, q_i) = (q_j, p_j)\}.$$

Note that for any $(P_1, \ldots, P_k) \in S_n^*(\pi)$, we have $(p_i, q_i) = (q_j, p_j)$ if and only if the random variables $a_n(P_i)$ and $a_n(P_j)$ are equal. Obviously, we have $0 \le m(P_1, \ldots, P_k) \le \frac{k}{2}$. With this notation, we find that

$$\frac{1}{n^{\frac{k}{2}+1}} \sum_{(P_1,\dots,P_k)\in S_n^*(\pi)} \mathbb{E}\left[a_n(P_1)\cdots a_n(P_k)\right] = \frac{1}{n^{\frac{k}{2}+1}} \sum_{l=0}^{k/2} c_n^{\frac{k}{2}-l} \# A_n^{(l)}(\pi) , \qquad (5.8)$$

where

$$A_n^{(l)}(\pi) := \{ (P_1, \dots, P_k) \in S_n^*(\pi) : m (P_1, \dots, P_k) = l \}.$$

The following lemma states that if a pair P_i, P_j contributes to $m(P_1, \ldots, P_k)$, then we can assume that the block $\{i, j\}$ in π is not crossed by any other block.

Lemma 5.4. Let $\pi \in \mathcal{PP}(k)$ and fix $i \sim_{\pi} j$, i < j. Define

$$S_n^*(\pi; i, j) := \{ (P_1, \dots, P_k) \in S_n^*(\pi) : P_i = (p_i, q_i), P_j = (p_j, q_j), p_i = q_j, q_i = p_j \}.$$

Assume that there is some $i' \sim_{\pi} j'$ such that i < i' < j, and either j' < i or j < j'. Then,

$$#S_n^*(\pi; i, j) = o\left(n^{\frac{k}{2}+1}\right).$$

To illustrate Lemma 5.4, we want to give an example. Therefore, take k = 4 and $\pi = \{\{1,3\},\{2,4\}\}$. Let i = 1 and j = 3. Here, the set $S_n^*(\pi; i, j)$ consists of all multiindices $((p_1, p_2), (p_2, p_2), (p_2, p_1), (p_1, p_1))$ with $p_1, p_2 \in \{1, \ldots, n\}, p_1 \neq p_2$. In particular, we have $\#S_n^*(\pi; i, j) = \mathcal{O}(n^2)$ implying the statement of Lemma 5.4 in this case.

Proof. To fix some $(P_1, \ldots, P_k) \in S_n^*(\pi; i, j)$, we first choose a value for $p_i = q_j$ and $q_i = p_j$. This allows for at most n^2 possibilities. Hence, P_i and P_j are fixed. Now consider the pairs $P_{i+1}, \ldots, P_{i'-1}$. p_{i+1} is uniquely determined by consistency. For q_{i+1} , there are at most n choices. Then, $p_{i+2} = q_{i+1}$. If $i + 2 \sim_{\pi} i + 1$, we have one choice for q_{i+2} . Otherwise, there are at most n. Proceeding in the same way, we see that we have n possibilities whenever we start a new equivalence class. Similarly, we can assign values to the pairs $P_{j+1}, \ldots, P_{i'+1}$ in this order. Now $P_{i'}$ is determined by consistency. When fixing $P_{i-1}, \ldots, P_1, P_k, \ldots, P_{j+1}$, we again have n choices for any new equivalence class. To sum up, we are left with at most

$$n^2 n^{\frac{k}{2}-2} = n^{\frac{k}{2}}$$

possible values for an element in $S_n^*(\pi; i, j)$.

Recall Definition 3.1 where we introduced the notion of the height $h(\pi)$ of a pair partition π . Lemma 5.4 in particular implies that only those $(P_1, \ldots, P_k) \in S_n^*(\pi)$ with

$$0 \le m \left(P_1, \ldots, P_k \right) \le h(\pi)$$

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contribute to the limit of (5.8). Indeed, if $m(P_1, \ldots, P_k) > h(\pi)$, we can find some $i \sim_{\pi} j, i < j$, such that $(P_1, \ldots, P_k) \in S_n^*(\pi; i, j)$ and neither j = i+1 nor is the restriction of π to $\{i+1, \ldots, j-1\}$ a pair partition. Hence, the crossing property in Lemma 5.4 is satisfied, and (P_1, \ldots, P_k) is contained in a set that is negligible in the limit. The identity in (5.8) thus becomes

$$\frac{1}{n^{\frac{k}{2}+1}} \sum_{(P_1,\ldots,P_k)\in S_n^*(\pi)} \mathbb{E}\left[a_n(P_1)\cdots a_n(P_k)\right] = \frac{1}{n^{\frac{k}{2}+1}} \sum_{l=0}^{h(\pi)} c_n^{\frac{k}{2}-l} \# B_n^{(l)}(\pi) + o(1),$$

where

$$\begin{split} B_n^{(l)}(\pi) &:= \{ (P_1, \dots, P_k) \in S_n^*(\pi) : m \, (P_1, \dots, P_k) = l; \\ (p_i, q_i) &= (q_j, p_j), i < j \ \Rightarrow \ j = i+1 \text{ or } \pi|_{\{i+1, \dots, j-1\}} \text{ is a pair partition} \} \,. \end{split}$$

In the next step, we want to simplify the expression above further by showing that $B_n^{(l)}(\pi) = \emptyset$ whenever $0 \le l < h(\pi)$. This is ensured by

Lemma 5.5. Let $\pi \in \mathcal{PP}(k)$. For any $(P_1, \ldots, P_k) \in S_n^*(\pi)$, we have

$$m(P_1,\ldots,P_k) \ge h(\pi).$$

To give a simple example, consider k = 4 and $\pi = \{\{1,2\},\{3,4\}\}$. Thus, π is a non-crossing partition with $h(\pi) = 2$. Further, the set $S_n^*(\pi)$ contains all multi-indices $(P_1, P_2, P_3, P_4) = ((p_1, p_2), (p_2, p_1), (p_1, p_3), (p_3, p_1))$ with $p_1, p_2, p_3 \in \{1, \ldots, n\}$ and $p_2 \neq p_3$. In particular, we have $m(P_1, P_2, P_3, P_4) = 2 = h(\pi)$.

Proof. If $h(\pi) = 0$, there is nothing to prove. Thus, suppose that $h(\pi) \ge 1$ and take some $i \sim_{\pi} j, i < j$, such that either j = i + 1 or $j - i - 1 \ge 2$ is even and the restriction of π to $\{i + 1, \ldots, j - 1\}$ is a pair partition. Fix $(P_1, \ldots, P_k) \in S_n^*(\pi)$, and write $P_l = (p_l, p_{l+1})$ for any $l = 1, \ldots, k$. We need to verify that $p_{i+1} = p_j$. If we achieve this, the definition of $S_n^*(\pi)$ will also ensure that $p_i = p_{j+1}$. As a consequence, the π -block $\{i, j\}$ will contribute to $m(P_1, \ldots, P_k)$. Since there are $h(\pi)$ such blocks, we will obtain $m(P_1, \ldots, P_k) \ge h(\pi)$ for any choice of $(P_1, \ldots, P_k) \in S_n^*(\pi)$.

If j = i + 1, we immediately obtain $p_{i+1} = p_j$. To show this property in the second case, note that the sequence $(P_{i+1}, \ldots, P_{j-1})$ solves the following system of equations:

$$p_{i+2} - p_{i+1} + p_{l_1+1} - p_{l_1} = 0, \quad \text{if } i + 1 \sim_{\pi} l_1,$$

$$p_{i+3} - p_{i+2} + p_{l_2+1} - p_{l_2} = 0, \quad \text{if } i + 2 \sim_{\pi} l_2,$$

$$\vdots$$

$$p_{i+m+1} - p_{i+m} + p_{l_m+1} - p_{l_m} = 0, \quad \text{if } i + m \sim_{\pi} l_m,$$

$$\vdots$$

$$p_j - p_{j-1} + p_{l_{j-i-1}+1} - p_{l_{j-i-1}} = 0, \quad \text{if } j - 1 \sim_{\pi} l_{j-i-1}$$

Start with solving the first equation for p_{i+2} which yields

$$p_{i+2} = p_{i+1} - p_{l_1+1} + p_{l_1}$$

Then, insert this in the second equation, and solve it for p_{i+3} to obtain

$$p_{i+3} = p_{i+1} - p_{l_1+1} + p_{l_1} - p_{l_2+1} + p_{l_2}.$$

In the j - i - 1-th step, we substitute $p_{j-1} = p_{i+(j-i-1)}$ in the j - i - 1-th equation, and solve it for $p_j = p_{i+(j-i-1)+1}$. We then have

$$p_j = p_{i+1} - \sum_{m=1}^{j-i-1} (p_{l_m+1} - p_{l_m}).$$

Since the restriction of π to $\{i+1,\ldots,j-1\}$ is a pair partition, we can conclude that the sets $\{l_1,\ldots,l_{j-i-1}\}$ and $\{i+1,\ldots,j-1\}$ are equal. Hence, we obtain $\sum_{m=1}^{j-i-1} (p_{l_m+1}-p_{l_m}) = p_j - p_{i+1}$, implying $p_j = p_{i+1}$.

With the help of Lemma 5.5, we thus arrive at

$$\frac{1}{n^{\frac{k}{2}+1}} \sum_{(P_1,\dots,P_k)\in S_n^*(\pi)} \mathbb{E}\left[a_n(P_1)\cdots a_n(P_k)\right] = \frac{\#B_n^{(h(\pi))}(\pi)}{n^{\frac{k}{2}+1}} c_n^{\frac{k}{2}-h(\pi)} + o(1).$$

Note that any element $(P_1, \ldots, P_k) \in S_n^*(\pi)$ satisfying the condition

$$(p_i, q_i) = (q_j, p_j), \ i < j \quad \Rightarrow \quad j = i+1 \text{ or } \pi|_{\{i+1,\dots,j-1\}} \text{ is a pair partition},$$
 (5.9)

fulfills the condition $m(P_1, \ldots, P_k) = h(\pi)$ as well. Indeed, (5.9) guarantees that $m(P_1, \ldots, P_k) \le h(\pi)$, and Lemma 5.5 ensures that $m(P_1, \ldots, P_k) \ge h(\pi)$. Thus, we can write

$$\begin{split} B_n^{(h(\pi))}(\pi) &= \{(P_1, \dots, P_k) \in S_n^*(\pi) : \\ (p_i, q_i) &= (q_j, p_j), i < j \ \Rightarrow \ j = i+1 \text{ or } \pi|_{\{i+1, \dots, j-1\}} \text{ is a pair partition} \} \,. \end{split}$$

Now any element in the complement of $B_n^{(h(\pi))}(\pi)$ satisfies for some $i \sim_{\pi} j$ the crossing assumption in Lemma 5.4. This yields

$$\frac{\#\left(B_n^{(h(\pi))}(\pi)\right)^c}{n^{\frac{k}{2}+1}} = o(1)$$

Since $B_n^{(h(\pi))}(\pi) \cup \left(B_n^{(h(\pi))}(\pi)\right)^c = S_n^*(\pi)$, we obtain that

$$\frac{1}{n^{\frac{k}{2}+1}} \sum_{(P_1,\dots,P_k)\in S_n^*(\pi)} \mathbb{E}\left[a_n(P_1)\cdots a_n(P_k)\right] = \frac{\#S_n^*(\pi)}{n^{\frac{k}{2}+1}} c_n^{\frac{k}{2}-h(\pi)} + o(1).$$
(5.10)

To calculate the limit on the right-hand side, we have

Lemma 5.6 (cf. [5], Lemma 4.6). For any $\pi \in \mathcal{PP}(k)$, it holds that

$$\lim_{n \to \infty} \frac{\# S_n^*(\pi)}{n^{\frac{k}{2}+1}} = p_T(\pi),$$

where $p_T(\pi)$ is the Toeplitz volume defined by solving the system of equations (3.1).

Proof. Fix $\pi \in \mathcal{PP}(k)$. Note that if $P = \{(p_i, p_{i+1}), i = 1, \ldots, k\} \in S_n^*(\pi)$, then we have x_0, x_1, \ldots, x_k with $x_i = p_{i+1}/n$ is a solution of the system of equations (3.1). On the other hand, if $x_0, x_1, \ldots, x_k \in \{1/n, 2/n, \ldots, 1\}$ is a solution of (3.1) and $p_{i+1} = nx_i$, then either $\{(p_i, p_{i+1}), i = 1, \ldots, k\} \in S_n^*(\pi)$ or $\{(p_i, p_{i+1}), i = 1, \ldots, k\} \in S_n(\eta)$ for some partition $\eta \in \mathcal{P}(k)$ such that $i \sim_{\pi} j \Rightarrow i \sim_{\eta} j$, but $\#\eta < \#\pi$.

In (3.1), we have k + 1 variables and only k/2 equations. Denote the k/2 + 1 undetermined variables by $y_1, \ldots, y_{k/2+1}$. We thus need to assign values from the set $\{1/n, 2/n, \ldots, 1\}$ to $y_1, \ldots, y_{k/2+1}$, and then to calculate the remaining k/2 variables from the equations. Since the latter are also supposed to be in the range $\{1/n, 2/n, \ldots, 1\}$, it might happen that not all values for the undetermined variables are admissible. Let $p_n(\pi)$ denote the admissible fraction of the $n^{k/2+1}$ choices for $y_1, \ldots, y_{k/2+1}$. By our remark at the beginning of the proof and estimate (5.2), we have that

$$\lim_{n \to \infty} \frac{\# S_n^*(\pi)}{n^{\frac{k}{2}+1}} = \lim_{n \to \infty} p_n(\pi),$$

if the limits exist. Now we can interpret $y_1, \ldots, y_{k/2+1}$ as independent random variables with a uniform distribution on $\{1/n, 2/n, \ldots, 1\}$. Then, $p_n(\pi)$ is the probability that the computed values stay within the interval (0, 1]. As $n \to \infty$, $y_1, \ldots, y_{k/2+1}$ converge in law to independent random variables uniformly distributed on [0, 1]. Hence, $p_n(\pi) \to p_T(\pi)$.

Applying Lemma 5.6 and assumption (C4) to equation (5.10), we arrive at

$$\lim_{n \to \infty} \frac{1}{n^{\frac{k}{2}+1}} \sum_{(P_1, \dots, P_k) \in S_n^*(\pi)} \mathbb{E} \left[a_n(P_1) \cdots a_n(P_k) \right] = p_T(\pi) c^{\frac{k}{2} - h(\pi)}.$$

Substituting this result in (5.7), we find that for any even $k \in \mathbb{N}$,

$$\lim_{n \to \infty} \frac{1}{n} \mathbb{E} \left[\operatorname{tr} \left(\mathbf{X}_{n}^{k} \right) \right] = C_{\frac{k}{2}} + \sum_{\pi \in \mathcal{CPP}(k)} p_{T}(\pi) c^{\frac{k}{2} - h(\pi)}.$$

To obtain the alternative expression in (3.2) for the even moments of the limiting measure ν_c , note that the considerations above were not restricted to crossing partitions. In particular, we can start from identity (5.5) instead of (5.7) to see that

$$\lim_{n \to \infty} \frac{1}{n} \mathbb{E} \left[\operatorname{tr} \left(\mathbf{X}_n^k \right) \right] = \lim_{n \to \infty} \sum_{\pi \in \mathcal{PP}(k)} \frac{\# S_n^*(\pi)}{n^{\frac{k}{2}+1}} c_n^{\frac{k}{2}-h(\pi)} = \sum_{\pi \in \mathcal{PP}(k)} p_T(\pi) c^{\frac{k}{2}-h(\pi)}.$$

5.2 Almost Sure Convergence

The almost sure convergence of the empirical distribution is a consequence of the following concentration inequality proven in [5] and [11].

Lemma 5.7. Suppose that conditions (C1) and (C2) hold. Then, for any $k, n \in \mathbb{N}$,

$$\mathbb{E}\left[\left(\operatorname{tr}\left(\boldsymbol{X}_{n}^{k}\right)-\mathbb{E}\left[\operatorname{tr}\left(\boldsymbol{X}_{n}^{k}\right)\right]\right)^{4}\right]\leq C\ n^{2}.$$

From Lemma 5.7 and Chebyshev's inequality, we can now conclude that for any $\varepsilon > 0$ and any $k, n \in \mathbb{N}$,

$$\mathbb{P}\left(\left|\frac{1}{n}\operatorname{tr}\left(\mathbf{X}_{n}^{k}\right)-\mathbb{E}\left[\frac{1}{n}\operatorname{tr}\left(\mathbf{X}_{n}^{k}\right)\right]\right|>\varepsilon\right)\leq\frac{C}{\varepsilon^{4}n^{2}}$$

Applying the Borel-Cantelli lemma, we see that

$$\frac{1}{n}\operatorname{tr}\left(\mathbf{X}_{n}^{k}\right) - \mathbb{E}\left[\frac{1}{n}\operatorname{tr}\left(\mathbf{X}_{n}^{k}\right)\right] \to 0, \quad \text{a.s.}.$$
(5.11)

Let Y be a random variable distributed according to ν_c . The convergence of the moments of the expected empirical distributions and relation (5.11) yield

$$\frac{1}{n} \operatorname{tr} \left(\mathbf{X}_{n}^{k} \right) \to \mathbb{E}[Y^{k}], \quad \text{a.s.}$$

Since the distribution of *Y* is uniquely determined by its moments, we obtain almost sure weak convergence of the empirical spectral distribution of \mathbf{X}_n to ν_c .

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6 Proof of Theorem 2.3

We want to give a proof of Theorem 2.3. Therefore, we start with showing that the free cumulants of the free convolution of rescaled versions of ν_0 and ν_1 coincide with the free cumulants of ν_c . Since the involved distributions are uniquely determined by their moments, and hence by their cumulants, we conclude that ν_c is the free convolution of rescaled versions of ν_0 and ν_1 . Therefore, we want to adapt some concepts of Bożejko and Speicher [4] which were picked up by Bryc, Dembo and Jiang [5]. Hence, let $\pi \in \mathcal{PP}(2k)$. We say that $\eta \neq \pi$ is a *sub-partition* of π if for some $i, j \in \{1, \ldots, k\}, \eta$ is a pair partition of $\{i, i + 1, \ldots, j\}$, and any block of η is also a block of π . Further, we denote by $\tilde{\eta}$ the pair partition which consists of all blocks of π not contained in η , i.e. π is the disjoint union of η and $\tilde{\eta}$.

Definition 6.1. We say that $p : \mathcal{PP}(2k) \to \mathbb{R}$ is pyramidally multiplicative, if for every $\pi \in \mathcal{PP}(2k)$ and any sub-partition η of π , we have $p(\pi) = p(\eta)p(\tilde{\eta})$.

In the following, we denote by $\mathcal{PP}_0(2k) \subset \mathcal{PP}(2k)$ the set of all pair partitions without sub-partitions.

Lemma 6.2 ([4], page 152, [5], Lemma A.4). Suppose that the moments of some ditribution are given by

$$m_k = \begin{cases} \sum_{\pi \in \mathcal{PP}(k)} p(\pi), & \text{if } k \text{ is even}, \\ 0, & \text{if } k \text{ is odd}. \end{cases}$$

If $p(\pi)$ is pyramidally multiplicative, then the free cumulants satisfy

$$\kappa_k = \begin{cases} \sum_{\pi \in \mathcal{PP}_0(k)} p(\pi), & \text{if } k \text{ is even}, \\ 0, & \text{if } k \text{ is odd}. \end{cases}$$

Note that the weights $c^{k-h(\pi)}$, $\pi \in \mathcal{PP}(2k)$, are pyramidally multiplicative since the height $h(\pi)$ satisfies the relation $h(\pi) = h(\eta) + h(\tilde{\eta})$ for any sub-partition η of π . Moreover, p_T is pyramidally multiplicative as well. Indeed, $p_T(\pi)$ is the volume of the cross section of the cube $[0,1]^{k+1}$ defined by the system of equations (3.1). If η is a sub-partition, we can decompose the system of equations into two parts corresponding to η and $\tilde{\eta}$, respectively, and calculate the volumes $p_T(\eta)$ and $p_T(\tilde{\eta})$. Since $\eta \cup \tilde{\eta} = \pi$, we conclude that $p_T(\pi) = p_T(\eta)p_T(\tilde{\eta})$. As a consequence of Lemma 6.2, we now have that the even free cumulants of ν_c are given by

$$\kappa_{2k}(\nu_c) = \sum_{\pi \in \mathcal{PP}_0(2k)} p_T(\pi) c^{k-h(\pi)}.$$

For k = 1, the set $\mathcal{PP}_0(2k)$ contains exactly one partition, namely $\pi = \{\{1, 2\}\}$. Here, we have $h(\pi) = 1$ and $p_T(\pi) = 1$, implying that $\kappa_2(\nu_c) = 1 = \kappa_2(\nu_1)$. If $k \ge 2$, any partition $\pi \in \mathcal{PP}_0(2k)$ has no sub-partition so that $h(\pi) = 0$. Thus,

$$\kappa_{2k}(\nu_c) = c^k \sum_{\pi \in \mathcal{PP}_0(2k)} p_T(\pi) = c^k \kappa_{2k}(\nu_1), \quad k \ge 2.$$

In particular, we obtain for the semicircle law ν_0 that $\kappa_{2k}(\nu_0) = \delta_1(k)$. Consequently,

$$(1-c)^k \kappa_{2k}(\nu_0) + c^k \kappa_{2k}(\nu_1) = \kappa_{2k}(\nu_c).$$

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Recall that according to Remark 2.1, we have $c \ge 0$. Assuming that $X \sim \nu_0$, $Y \sim \nu_1$ and $Z \sim \nu_c$, we thus see that Z is the free convolution of $\sqrt{1-c} X$ and $\sqrt{c} Y$.

In [17], it is shown that ν_1 has a bounded density. By [3], Corollary 2, the free convolution of any measure with the semicircle distribution ν_0 has a density, in particular ν_c for c < 1. Moreover, for c < 1, the density is smooth and bounded by Corollary 4 and Proposition 5 in [3]. To see that ν_c has an unbounded support for 0 < c < 1, recall that ν_1 has an unbounded support, and the moments satisfy

$$c^{\frac{k}{2}}\int x^k d\nu_1(x) \leq \int x^k d\nu_c(x), \quad k \in \mathbb{N}.$$

Finally, ν_c is symmetric for any $0 \le c \le 1$ since all odd moments vanish. This proves Theorem 2.3.

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