Spectral fluctuations for the multi-dimensional Anderson model

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Abstract. In this paper, we examine fluctuations of polynomial linear statistics for the Anderson model on \mathbb{Z}^d for any potential with finite moments. We prove that if normalized by the square root of the size of the truncated operator, these fluctuations converge to a Gaussian limit. For a vast majority of potentials and polynomials, we show that the variance of the limiting distribution is strictly positive, and we classify in full the rare cases in which this does not happen.

1. Introduction

The purpose of this paper is to study fluctuations of the eigenvalue counting measure for the Anderson model on \mathbb{Z}^d . We denote $|n| = \sum_{v=1}^d n_v$ for any $n \in \mathbb{Z}^d$, and write $n \sim m$ for $n, m \in \mathbb{Z}^d$ if and only if |n - m| = 1. Define the operator $H: \ell^2(\mathbb{Z}^d) \to \ell^2(\mathbb{Z}^d)$ by

$$(Hu)_n = (\Delta u)_n + (Xu)_n = \sum_{m \sim n} u_m + X_n \cdot u_n$$

where $\{X_n\}_{n \in \mathbb{Z}^d}$ is an array of independent, identically distributed (iid) random variables with finite moments, satisfying $\mathbb{E}[X_n] = 0$. We denote the distribution of each variable X_n by $d\rho$, which will henceforth be referred to as *the underlying distribution*.

In this paper, we aim to study the fluctuations of the counting measure for the eigenvalues of finite volume approximations. Explicitly, we study fluctuations of polynomial linear statistics of finite volume truncations of H: for any $L \in \mathbb{N}$, denote

$$\Lambda_L = [-L, L] \cap \mathbb{Z},$$

and let H_L be the truncation of H to the cube $\Lambda_L^d \subset \mathbb{Z}^d$. That is, $H_L = 1_{\Lambda_L^d} H 1_{\Lambda_L^d}$, where

$$(1_{\Lambda_L^d}(u))_n = \begin{cases} u_n & n \in \Lambda_L^d, \\ 0 & n \notin \Lambda_L^d. \end{cases}$$

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We denote by $N(0, \sigma^2)$ the normal distribution on \mathbb{R} with mean 0 and variance σ^2 , and denote by \xrightarrow{d} convergence in distribution. We agree that the zero random variable is also normal, by allowing $\sigma^2 = 0$ (in this case we say the distribution is *degenerate*).

The *empirical measure* of H_L is the measure

$$\mathrm{d}\nu_L = \frac{1}{(2L+1)^d} \sum_{i=0}^{|\Lambda_L^d|} \delta_{\lambda_i^{(\Lambda_L^d)}}$$

where $\{\lambda_1^{(\Lambda_L^d)}, \lambda_2^{(\Lambda_L^d)}, \dots, \lambda_{|\Lambda_L^d|}^{(\Lambda_L^d)}\} = \sigma(H_L)$ are the eigenvalues of H_L (counting multiplicity), and δ_{λ} is the Dirac measure at λ . When the empirical measure has a limit as $L \to \infty$, this limit is known as the *density of states* of H. In our case, it is known that the random measure $d\nu_L$ converges weakly almost surely to a deterministic measure $d\nu$ (see e.g., [1] and references within).

We want to focus on asymptotics of the fluctuations of $d\nu_L$. A natural way to study this is using *linear statistics* for polynomials, i.e., random variables of the form $\int f d\nu_L = \frac{1}{(2L+1)^d} \operatorname{Tr}(f(H_L))$ for some polynomial $f(x) \in \mathbb{R}[x]$.

Fluctuations of the truncated eigenvalues $\lambda_i^{(\Lambda_L^d)}$ are assumed to be associated to continuity properties of the spectral measures. There are several results indicating this is indeed true. Minami [6] studied the microscopic scale of the eigenvalues of the Anderson model in \mathbb{Z}^d , after Molchanov [7] did the same for the continuous case in one dimension. Minami proved that, under certain conditions that ensure localization with exponentially decaying eigenfunctions, the eigenvalues of the Anderson model have Poisson behavior on the microscopic scale. For d = 1, it is well known that localization holds for any ergodic non-deterministic potential [1]. However, for $d \ge 3$ and for sufficiently low energies, it is conjectured that H has extended states, i.e., the spectrum of H has an absolutely continuous component.

We now state our main theorem:

Theorem 1.1. Let $f(x) \in \mathbb{R}[x]$ be a non-constant polynomial. Then

$$\frac{\operatorname{Tr}(f(H_L)) - \mathbb{E}[\operatorname{Tr}(f(H_L))]}{(2L+1)^{d/2}} \xrightarrow{d} N(0, \sigma(f)^2)$$

as $L \to \infty$, where

- 1. *if the underlying distribution* $(d\rho)$ *is supported by more than three points, then* $\sigma(f)^2 > 0$;
- 2. *if the underlying distribution is supported by exactly two points, there exist polynomials* $g_2, g_3, g_5 \in \mathbb{R}[x]$, *of degrees* 2, 3, 5 *respectively, such that* $\sigma(f)^2 = 0$ *if and only if* $f \in \text{span}_{\mathbb{R}}\{g_5, g_3, g_2, 1\}$;

3. *if the underlying distribution is supported by exactly three points, there exists* a polynomial $\tilde{g}_3 \in \mathbb{R}[x]$ of degree 3, such that $\sigma(f)^2 = 0$ if and only if $f \in \text{span}_{\mathbb{R}}{\{\tilde{g}_3, 1\}}$.

The polynomials $g_2, g_3, g_5, \tilde{g}_3$ depend on $d\rho$ as well as on the dimension d, and are given explicitly in Propositions 4.3 and 4.4 below.

The study of fluctuations of finite truncations of the Anderson Model has received a considerable amount of attention, although most results focus on the one-dimensional case. Reznikova [10] proved a central limit theorem (CLT) for the eigenvalue counting function of the truncated Anderson model in 1-dimension. Kirsch and Pastur [5] proved a CLT for the trace of truncations of the Green function of the Anderson model in one dimension. Recently, Pastur and Shcherbina [9] extended this result to other functions of H.

In our proof we shall compute the trace of powers of H_L by counting paths on the associated lattice. Path counting and weighted path counting is commonly used in the study of random Schrodinger operators and in the study of random matrices (see, e.g., [1,2] and references therein).

This paper can be viewed as a second paper in a series, continuing the work of Breuer with the authors [3]. In the previous paper, path counting was used to prove a CLT for a decaying model over \mathbb{N} . In this paper, the methods have been modified to apply to the Anderson model over \mathbb{Z}^d for general $d \in \mathbb{N}$. Each of the papers is self contained, but there are many parallels in the overall structure of the paper and propositions.

The rest of the paper is organized as follows. In Section 2 we set up our definitions, and prove that "typical" diagonal elements in the matrix representation of H_L^k have a combinatorial description (using path counting). In Sections 3 and 4 we prove our main theorem – in Section 3 we show that fluctuations of $\text{Tr}(f(H_L))$ converge to a normally distributed random variable, and in Section 4 we classify all cases in which the limit distribution is non-degenerate. The final proof of Theorem 1.1 appears at the end of Section 4. We conclude with Section 5 (which is independent from the rest of the paper) in which we state and prove a CLT for *m*-dependent random variables indexed by \mathbb{Z}^d , which implies the CLT we use in Section 3.

2. Definitions and preliminaries

Fix $d \in \mathbb{N}$. As stated in the introduction, we explore the random operator

$$H: \ell^2(\mathbb{Z}^d) \to \ell^2(\mathbb{Z}^d).$$

It is useful to decompose H as

$$H = V + \sum_{\nu=1}^{d} U_{\nu} + \sum_{\nu=1}^{d} D_{\nu}, \qquad (2.1)$$

where V is the random potential operator, and each U_v (respectively D_v) is the operator shifting forward (respectively backward) in direction v. In other words, let e_1, e_2, \ldots, e_d denote the standard generators of \mathbb{Z}^d as a free abelian group. Then, for every $n \in \mathbb{Z}^d$ and $u \in \ell^2(\mathbb{Z}^d)$, we have $(Vu)_n = X_n u_n$, and for every integer $1 \le v \le d$ we have $(U_v u)_n = u_{n+e_v}$ and $(D_v u)_n = u_{n-e_v}$. A corresponding decomposition is also given for every finite volume truncation, H_L .

Our theorem deals with the asymptotic behavior (as $L \to \infty$) of $\text{Tr}(f(H_L))$, for polynomials $f \in \mathbb{R}[x]$. We consider $\text{Tr}(f(H_L))$ as a polynomial in the variables $\{X_n \mid n \in \mathbb{Z}^d\}$. To slightly ease notation, we denote our variables by a lowercase Latin letter (such as x, z) when referring to a single variable in a polynomial ring, and by uppercase letters (such as X_n, Z_n, Z) when referring to variables in polynomial rings which can also be understood as random variables with some distribution.

To work with such multivariate monomials, we introduce the following definitions:

Definition 2.1. A finitely supported function $\beta : \mathbb{Z}^d \to \mathbb{N} \cup \{0\}$ will be called a *multi-index*. Let β_n denote the value $\beta(n)$ for every $n \in \mathbb{Z}^d$. Let X^β denote the monomial $\prod_{n \in \mathbb{Z}^d} X_n^{\beta_n}$.

Fix a multi-index δ , by

$$\delta_n = \begin{cases} 1 & n = 0, \\ 0 & n \neq 0. \end{cases}$$

Definition 2.2. For every multi-index β and $i \in \mathbb{Z}^d$, define β^i (β shifted by i), by $\beta_n^i = \beta_{n-i}$ for every $n \in \mathbb{Z}^d$.

Note that using these definitions, for $n, i \in \mathbb{Z}^d$, δ_n^i is 1 if n = i and 0 otherwise. Additionally, $\beta = \sum_{i \in \mathbb{Z}^d} \beta_i \delta^i$ for every multi-index β (this is a finite sum as β is finitely supported).

Next, we fix $k \in \mathbb{N}$ and begin exploring the asymptotic behavior (as $L \to \infty$) of $\text{Tr}(H_L^k)$. As we shall see, the coefficient of any monomial X^β in $\text{Tr}(H_L^k)$ is fixed for sufficiently large L, and has a concrete combinatorial description. Furthermore, these coefficients are invariant under translations of the monomials in \mathbb{Z}^d . The precise statement is given in Proposition 2.7 below, which requires some more definitions.

Definition 2.3. Let $S = \{V, U_1, U_2, \dots, U_d, D_1, D_2, \dots, D_d\}$ be considered as formal symbols. Then S^k denotes the set of all ordered *k*-tuples with elements from S, or all strings of length *k* from the alphabet S.

Definition 2.4. For every $s \in S^k$, we define a finite sequence $y_0(s), y_1(s), \ldots, y_k(s) \in \mathbb{Z}^d$ as follows:

• $y_0(s) = (0, 0, \dots, 0);$ $\begin{cases} y_{j-1}(s) + e_v & s_j = U_v, \end{cases}$

•
$$y_j(s) = \begin{cases} y_{j-1}(s) - e_v & s_j = D_v, \\ y_{j-1}(s) & s_j = V. \end{cases}$$

We say that *s* is *balanced* if $y_k(s) = y_0(s)$.

Note that $s \in S^k$ is balanced if and only if for every v = 1, 2, ..., d, the symbols U_v and D_v appear in s the same number of times.

Definition 2.5. For every $s \in S^k$, define a multi-index $\varphi(s)$ by

$$\varphi(s)_n = \#\{1 \le j \le k \mid y_j(s) = y_{j-1}(s) = n\}$$

for every $n \in \mathbb{Z}^d$.

In other words, given a string $s \in S^k$, the following process describes the multiindex $\varphi(s)$ in terms of the symbols of s: we initialize $\varphi(s) = 0$ at the position $0 \in \mathbb{Z}^d$, and then read the symbols of s consecutively. When we encounter a symbol U_v , we update our position in \mathbb{Z}^d by adding e_v , when we encounter D_v we subtract e_v (not changing the multi-index), and when we encounter V we increase the value of $\varphi(s)_n$ by 1, where $n \in \mathbb{Z}^d$ is the current position (which is not changed by the symbol V).

Thus, a symbol V which appears in s at the j-th position, contributes 1 to $\varphi(s)_n$ for a unique $n \in \mathbb{Z}^d$, where n_v counts the difference between the number of times U_v and D_v appear in s in positions < j. For example, if d = 2 and

$$s = U_1 V U_1 D_2 V D_1 U_2 V D_1,$$

then $\beta = \varphi(s)$ has $\beta_{(1,0)} = 2$, $\beta_{(2,-1)} = 1$, and $\beta_n = 0$ for all other $n \in \mathbb{Z}^2$.

Definition 2.6. For every multi-index β , let $p^k(\beta)$ be the number of balanced strings $s \in S^k$ satisfying $\varphi(s)^i = \beta$, for some $i \in \mathbb{Z}^d$.

Note that for every $s \in S^k$ and multi-index β , there is at most one $i \in \mathbb{Z}^d$ for which $\varphi(s)^i = \beta$.

Proposition 2.7. For every non-zero multi-index β , and $k, L \in \mathbb{N}$, let $a_L^k(\beta)$ denote the coefficient of X^{β} in the polynomial $\text{Tr}(H_L^k)$. Then

- 1. $0 \le a_L^k(\beta) \le p^k(\beta);$
- 2. if $\beta_n > 0$ for some $n \in \Lambda^d_{L-k}$, we have $a^k_L(\beta) = p^k(\beta)$;
- 3. if $\beta_n > 0$ for some $n \notin \Lambda_L^d$, we have $a_L^k(\beta) = 0$.

Proof. Use (2.1) to expand H^k . This gives us a bijection between operators in the expansion of H^k and strings in S^k . Furthermore, let M_L be any matrix in the expansion of H_L^k corresponding to a string $s \in S^k$. It is straightforward to verify that if $s \in S^k$ is balanced, and $i \in \Lambda_L^d$, and $y^j(s) + i \in \Lambda_L^d$ for every j = 1, 2, ..., k, we have $(M_L)_{i,i} = X^{\varphi(s)^i}$. Otherwise, we have $(M_L)_{i,i} = 0$.

Therefore, fixing a multi-index β , the coefficient $a_L^k(\beta)$ equals the number of strings $s \in S^k$, for which $\varphi(s)^i = \beta$ and the additional conditions $y_j(s) + i \in \Lambda_L^d$ are fulfilled (we simply compute the trace as the sum over all diagonal entries from all matrices in the expansion). The number of such strings is at least 0 and at most $p^k(\beta)$ (which is the number of such strings without the additional conditions), proving (1).

Note that for any balanced $s \in S^k$, we have $|y_j(s)| \le \frac{k}{2}$ for every j = 0, 1, ..., k. We deduce that whenever β takes a non-zero value in Λ_{L-k}^d , if $\beta = \varphi(s)^i$ we must have $y_j(s) + i \in \Lambda_{L-k}^d$ for *some* j, therefore $y_j(s) + i \in \Lambda_L^d$ for *every* j = 0, 1, ..., k. For such β , any $i \in \mathbb{Z}^d$ and $s \in S^k$ satisfying $\varphi(s)^i = \beta$ automatically fulfill the additional conditions, proving (2).

Similarly, if β obtains a non-zero value outside of Λ_L^d , satisfying $\varphi(s)^i = \beta$ guarantees that $y_j(s) + i \notin \Lambda_L^d$ for some j, therefore X^β does not appear anywhere on the diagonal of M_L , proving (3).

Note that, from Definition 2.6, it is clear that $p^k(\beta^i) = p^k(\beta)$, for any multiindex β , any $i \in \mathbb{Z}^d$, and any $k \in \mathbb{N}$. Therefore, when considering the integers $p^k(\beta)$ which appear as coefficients in the polynomials $\text{Tr}(H_L^k)$, we may restrict our attention to a set of non-zero multi-indices which contains some shifting of every multi-index exactly once. We denote this set by B:

Definition 2.8. Two multi-indices β and γ are said *equivalent* if $\gamma = \beta^i$ for some $i \in \mathbb{Z}^d$. From each equivalence class other than zero, choose a unique representative β , satisfying $\beta_0 > 0$ (one way to make such choices, is to require the lexicographic minimum of the support of β to be 0). Let *B* be the set of all chosen representatives.

In other words, *B* is any set of multi-indices with the properties:

- 1. for any non-zero multi-index γ , we have $\gamma^i \in B$ for a unique $i \in \mathbb{Z}^d$;
- 2. $\beta_0 > 0$ for every $\beta \in B$.

3. A central limit theorem for polynomial linear statistics

In this section, we prove that for every polynomial $f(x) \in \mathbb{R}[x]$,

$$\frac{\operatorname{Tr}(f(H_L)) - \mathbb{E}[\operatorname{Tr}(f(H_L))]}{(2L+1)^{d/2}}$$

converges in distribution (as $L \to \infty$) to a normal distribution with variance $\sigma(f)^2 \in [0, \infty)$ (see Proposition 3.8 below). We start by proving this CLT in the case where $f(x) = x^k$ is a monomial, which is easier to prove for an approximated version of the random variable $Tr(H_L^k)$:

Definition 3.1. For every $k, L \in \mathbb{N}$, let

$$T_L^k = \sum_{\beta \in B} p^k(\beta) \sum_{i \in \Lambda_L^d} X^{\beta^i}, \qquad (3.1)$$

which we consider both as a random variable, and as a polynomial in the variables $\{X_n \mid n \in \mathbb{Z}^d\}$.

Note that the above sum is finite, since $p^k(\beta) = 0$ for all but finitely many $\beta \in B$.

We start by proving that T_L^k can indeed approximate $Tr(H_L^k)$, in the following sense:

Proposition 3.2. For every $k \in \mathbb{N}$, the random variables

$$\frac{\text{Tr}(H_L^k) - \mathbb{E}[\text{Tr}(H_L^k)]}{(2L+1)^{d/2}} - \frac{T_L^k - \mathbb{E}[T_L^k]}{(2L+1)^{d/2}}$$

converge in probability (as $L \to \infty$) to 0.

Proof. It is sufficient to show that $Var(T_L^k - Tr(H_L^k)) = o(L^d)$. From Proposition 2.7 and (3.1), we have

$$T_L^k - \operatorname{Tr}(H_L^k) = \sum_{\beta \in B} \sum_{i \in \Lambda_L^d} (p^k(\beta) - a_L^k(\beta^i)) X^{\beta^i},$$

where $p^k(\beta) - a_L^k(\beta^i) = 0$ whenever $i \in \Lambda_{L-k}^d$. Therefore, the number of non-zero terms in the above sum is at most

$$|B_k|(|\Lambda_L^d| - |\Lambda_{L-k}^d|) = O(L^{d-1}),$$

where $B_k = \{\beta \in B \mid p^k(\beta) \neq 0\}$ is finite. Next, consider the sum

$$\operatorname{Var}(T_{L}^{k} - \operatorname{Tr}(H_{L}^{k})) = \sum_{\beta, \gamma \in B_{k}} \sum_{i, j \in \Lambda_{L}^{d}} (p^{k}(\beta) - a_{L}^{k}(\beta^{i}))(p^{k}(\gamma) - a_{L}^{k}(\gamma^{j})) \operatorname{Cov}(X^{\beta^{i}}, X^{\gamma^{j}}).$$
(3.2)

Fixing $\beta, \gamma \in B_k$ and $i \in \Lambda_L^d \setminus \Lambda_{L-k}^d$, we see that whenever $j - i \notin \Lambda_k^d$, the supports of β^i and γ^j are disjoint, therefore X^{β^i} and X^{γ^j} are independent. This tells us that there are at most

$$|B_k|^2 \cdot (|\Lambda_L^d| - |\Lambda_{L-k}^d|) \cdot |\Lambda_k^d| = O(L^{d-1})$$

non-zero terms in (3.2). We know from Proposition 2.7 (1) that

$$0 \le (p^k(\beta) - a_L^k(\beta^i))(p^k(\gamma) - a_L^k(\gamma^j)) \le p^k(\beta)p^k(\gamma).$$
(3.3)

Since $\{X_n \mid n \in \mathbb{Z}^d\}$ are identically distributed, we have

$$\operatorname{Cov}(X^{\beta^{i}}, X^{\gamma^{j}}) = \operatorname{Cov}(X^{\beta}, X^{\gamma^{j-i}})$$

From here we deduce that for fixed $\beta, \gamma \in B_k$, the term $\text{Cov}(X^{\beta^i}, X^{\gamma^j})$ only obtains a finite number of values: it is either 0 or uniquely determined by $\beta, \gamma \in B_k$, the value of $j - i \in \Lambda_k^d$, and some of the (finite) moments of the underlying distribution $d\rho$. Together with (3.3), this gives us a uniform bound on all terms in (3.2), showing that indeed

$$\operatorname{Var}(T_L^k - \operatorname{Tr}(H_L^k)) = O(L^{d-1}).$$

Our next step is a central limit theorem for the random variables T_L^k . Although our initial random variables $\{X_n \mid n \in \mathbb{Z}^d\}$ were iid, for a fixed multi-index β the random variables $\{X^{\beta^i} \mid i \in \mathbb{Z}^d\}$ are generally not independent. However, for any *i* and *j* sufficiently far apart $(j - i \notin \Lambda_k^d)$ is sufficient), the variables X^{β^i} and X^{β^j} are independent. We use CLTs for weakly dependent random variables, by Hoeffding and Robbins [4] and Neumann [8], to prove:

Theorem 3.3. Let $\{X_n\}_{n \in \mathbb{Z}^d}$ be an array of iid random variables with finite moments, satisfying $\mathbb{E}[X_n] = 0$. Let B be any set of multi-indices, and let $\{a_\beta\}_{\beta \in B}$ be a set of coefficients such that $a_\beta = 0$ for all but finitely many $\beta \in B$. Then

$$\frac{1}{(2L+1)^{d/2}} \sum_{\beta \in \mathcal{B}} a_{\beta} \sum_{i \in \Lambda_L^d} (X^{\beta^i} - \mathbb{E}[X^{\beta^i}]) \stackrel{d}{\to} N(0, \sigma^2),$$

as $L \to \infty$, for some $\sigma^2 \ge 0$.

The proof is postponed to Section 5.

Corollary 3.4. For every $k \in \mathbb{N}$,

$$\frac{T_L^k - \mathbb{E}[T_L^k]}{(2L+1)^{d/2}} \stackrel{d}{\to} N(0, \sigma_k^2)$$

as $L \to \infty$, for some $\sigma_k^2 \ge 0$.

Now, that we have a central limit theorem for our approximating random variables, we would like to compute the limit variances, and more generally, the limit covariances. We do this first for individual multi-indices: **Lemma 3.5.** For every two multi-indices β and γ , we have

$$\lim_{L \to \infty} \frac{1}{(2L+1)^d} \operatorname{Cov}\left(\sum_{i \in \Lambda_L^d} X^{\beta^i}, \sum_{i \in \Lambda_L^d} X^{\gamma^i}\right) = \sum_{j \in \mathbb{Z}^d} \operatorname{Cov}(X^\beta, X^{\gamma^j})$$
(3.4)

Note that the sum on the right-hand side of (3.4) is uniquely determined by β , γ , and the moments of the underlying distribution $d\rho$, and it is in fact a finite sum: since β and γ are finitely supported, the supports of β and γ^j are disjoint (and therefore X^{β} and X^{γ^j} are independent) for all but finitely many $j \in \mathbb{Z}^d$.

Proof. Since $\{X_n \mid n \in \mathbb{Z}^d\}$ are identically distributed, the covariances are invariant to translations, and we may write

$$\operatorname{Cov}\left(\sum_{i \in \Lambda_{L}^{d}} X^{\beta^{i}}, \sum_{i \in \Lambda_{L}^{d}} X^{\gamma^{i}}\right) = \sum_{i,i' \in \Lambda_{L}^{d}} \operatorname{Cov}(X^{\beta^{i}}, X^{\gamma^{i'}})$$
$$= \sum_{i,i' \in \Lambda_{L}^{d}} \operatorname{Cov}(X^{\beta}, X^{\gamma^{i'-i}})$$
$$= \sum_{j \in \mathbb{Z}^{d}} z_{j}(L) \cdot \operatorname{Cov}(X^{\beta}, X^{\gamma^{j}}).$$

where $z_j(L) = #\{i, i' \in \Lambda_L^d \mid j = i - i'\}$. Clearly

$$\lim_{L \to \infty} \frac{z_j(L)}{(2L+1)^d} = 1$$

for any $j \in \mathbb{Z}^d$, and since $\text{Cov}(X^{\beta}, X^{\gamma^j}) \neq 0$ only for finitely many $j \in \mathbb{Z}^d$, the claim follows.

Corollary 3.6. For every $k, \ell \in \mathbb{N}$,

$$\lim_{L \to \infty} \operatorname{Cov} \left(\frac{T_L^k}{(2L+1)^{d/2}}, \frac{T_L^\ell}{(2L+1)^{d/2}} \right)$$
$$= \lim_{L \to \infty} \frac{1}{(2L+1)^d} \operatorname{Cov}(T_L^k, T_L^\ell)$$
$$= \sum_{\beta, \gamma \in B} p^k(\beta) p^\ell(\gamma) \sum_{j \in \mathbb{Z}^d} \operatorname{Cov}(X^\beta, X^{\gamma^j})$$
(3.5)

This result allows us to deduce results for the asymptotic behavior of the trace of monomials:

Corollary 3.7. *For any* $k \in \mathbb{N}$ *,*

$$\frac{\operatorname{Tr}(H_L^k) - \mathbb{E}[\operatorname{Tr}(H_L^k)]}{(2L+1)^{d/2}} \xrightarrow{d} N(0, \sigma_k^2)$$

as $L \to \infty$, where

$$\sigma_k^2 = \sum_{\beta,\gamma \in B} p^k(\beta) p^k(\gamma) \sum_{j \in \mathbb{Z}^d} \operatorname{Cov}(X^\beta, X^{\gamma^j}).$$

Furthermore, for every $k, \ell \in \mathbb{N}$ *,*

$$\lim_{L \to \infty} \operatorname{Cov} \left(\frac{\operatorname{Tr}(H_L^k) - \mathbb{E}[\operatorname{Tr}(H_L^k)]}{(2L+1)^{d/2}}, \frac{\operatorname{Tr}(H_L^\ell) - \mathbb{E}[\operatorname{Tr}(H_L^\ell)]}{(2L+1)^{d/2}} \right)$$
$$= \sum_{\beta, \gamma \in \mathcal{B}} p^k(\beta) p^\ell(\gamma) \sum_{j \in \mathbb{Z}^d} \operatorname{Cov}(X^\beta, X^{\gamma^j}).$$

Proof. Follows directly from Proposition 3.2 and Corollaries 3.4 and 3.6.

And now we can prove the CLT for any polynomial:

Proposition 3.8. Let $f(x) = \sum_{k=0}^{m} a_k x^k \in \mathbb{R}[x]$ be a polynomial. Then $\frac{\operatorname{Tr}(f(H_L)) - \mathbb{E}[\operatorname{Tr}(f(H_L))]}{(2L+1)^{d/2}} \xrightarrow{d} N(0, \sigma(f)^2)$

as $L \rightarrow \infty$, where

$$\sigma(f)^2 = \sum_{k,\ell=1}^m a_k a_\ell \sum_{\beta,\gamma \in B} p^k(\beta) p^\ell(\gamma) \sum_{j \in \mathbb{Z}^d} \operatorname{Cov}(X^\beta, X^{\gamma^j}).$$

Proof. Since $\text{Tr}(f(H_L)) - \mathbb{E}[\text{Tr}(f(H_L))]$ does not depend on a_0 , we may assume without loss of generality that $a_0 = 0$ and deg(f) = m > 0. Then

$$\frac{\operatorname{Tr}(f(H_L)) - \mathbb{E}[\operatorname{Tr}(f(H_L))]}{(2L+1)^{d/2}} = \sum_{k=1}^m a_k \frac{\operatorname{Tr}(H_L^k) - \mathbb{E}[\operatorname{Tr}(H_L^k)]}{(2L+1)^{d/2}},$$

and from Corollary 3.7 we obtain the value of the variance $\sigma(f)^2$. Using Proposition 3.2 and (3.1), we now rewrite

$$\lim_{L \to \infty} \frac{\operatorname{Tr}(f(H_L)) - \mathbb{E}[\operatorname{Tr}(f(H_L))]}{(2L+1)^{d/2}} = \lim_{L \to \infty} \frac{1}{(2L+1)^{d/2}} \sum_{k=1}^{m} a_k (T_L^k - \mathbb{E}[T_L^k]) \\
= \lim_{L \to \infty} \frac{1}{(2L+1)^{d/2}} \sum_{k=1}^{m} a_k \sum_{\beta \in B} p^k(\beta) \sum_{i \in \Lambda_L^d} (X^{\beta^i} - \mathbb{E}[X^{\beta^i}]) \\
= \lim_{L \to \infty} \frac{1}{(2L+1)^{d/2}} \sum_{\beta \in B} \left(\sum_{k=1}^{m} a_k p^k(\beta)\right) \sum_{i \in \Lambda_L^d} (X^{\beta^i} - \mathbb{E}[X^{\beta^i}]). \quad (3.6)$$

Note that the first equality holds in the sense that both limit random variables have the same distribution. Theorem 3.3 now applies, proving that the limit has a normal distribution.

4. Degenerate and non-degenerate cases

Now that we proved the convergence in Theorem 1.1, it remains to determine under which conditions the limit distribution is non-degenerate, that is when $\sigma(f)^2 > 0$ for a non-constant polynomial $f \in \mathbb{R}[x]$. It turns out that $\sigma(f)^2$ is always positive if $\deg(f) \neq 2, 3, 5$, but for some polynomials of degree 2, 3, 5 and some specific underlying distributions, the variance may vanish. We first demonstrate positive variance in degrees $\neq 2, 3, 5$:

Proposition 4.1. Let $f(x) = \sum_{k=0}^{m} a_k x^k \in \mathbb{R}[x]$ be a non-constant polynomial of degree $m \neq 2, 3, 5$. Then $\sigma(f)^2 > 0$.

Proof. Using (3.6), we write

$$\sigma(f)^{2} = \operatorname{Var}\left(\lim_{L \to \infty} \frac{\operatorname{Tr}(f(H_{L})) - \mathbb{E}[\operatorname{Tr}(f(H_{L}))]}{(2L+1)^{d/2}}\right)$$
$$= \lim_{L \to \infty} \operatorname{Var}\left(\frac{1}{(2L+1)^{d/2}} \sum_{\beta \in B} \left(\sum_{k=1}^{m} a_{k} p^{k}(\beta)\right) \sum_{i \in \Lambda_{L}^{d}} (X^{\beta^{i}} - \mathbb{E}[X^{\beta^{i}}])\right). \quad (4.1)$$

We follow the same general method used in [3] – it is sufficient to find a multi-index $\gamma \in B$, with the following properties:

- 1. $p^m(\gamma) \neq 0;$
- 2. $p^k(\gamma) = 0$ for every k < m;
- 3. $\sum_{i \in \mathbb{Z}^d} \operatorname{Cov}(X^{\gamma}, X^{\gamma^j}) > 0;$
- 4. $\operatorname{Cov}(X^{\gamma}, X^{\beta^{j}}) = 0$ for every $j \in \mathbb{Z}^{d}$ and every $\gamma \neq \beta \in B$ satisfying $p^{k}(\beta) \neq 0$ for some $1 \leq k \leq m$.

If we find such γ , we deduce from property (4) that the random variables

$$Y_L^1 \equiv a_m p^m(\gamma) \sum_{i \in \Lambda_L^d} (X^{\gamma^i} - \mathbb{E}[X^{\gamma^i}])$$

and

$$Y_L^2 \equiv \sum_{\beta \in B \setminus \{\gamma\}} \left(\sum_{k=1}^m a_k p^k(\beta) \right) \sum_{i \in \Lambda_L^d} (X^{\beta^i} - \mathbb{E}[X^{\beta^i}])$$

are uncorrelated (for any $L \in \mathbb{N}$), and (4.1) becomes

$$\sigma(f)^{2} = \lim_{L \to \infty} \operatorname{Var}\left(\frac{Y_{L}^{1} + Y_{L}^{2}}{(2L+1)^{d/2}}\right)$$
$$= \lim_{L \to \infty} \operatorname{Var}\left(\frac{Y_{L}^{1}}{(2L+1)^{d/2}}\right) + \lim_{L \to \infty} \operatorname{Var}\left(\frac{Y_{L}^{2}}{(2L+1)^{d/2}}\right)$$
$$\geq \lim_{L \to \infty} \operatorname{Var}\left(\frac{Y_{L}^{1}}{(2L+1)^{d/2}}\right) = a_{m} p^{m}(\gamma) \sum_{j \in \mathbb{Z}^{d}} \operatorname{Cov}(X^{\gamma}, X^{\gamma^{j}}) > 0,$$

where the final equality is due to Lemma 3.5. We make the following choices for γ :

- 1. if m = 1, choose $\gamma = \delta$;
- 2. if m > 4 is even, choose $\gamma = \delta + \delta^{(\frac{m}{2}-1)e_1}$;
- 3. if $m \ge 7$ is odd, choose $\gamma = \delta + \delta^{e_1} + \delta^{(\frac{m-3}{2})e_1}$.

The proof that these γ satisfy properties (1) and (2) is straightforward path counting. For (3) and (4), recall that

$$Cov(X^{\gamma}, X^{\beta^{j}}) = \mathbb{E}[X^{\gamma}X^{\beta^{j}}] - \mathbb{E}[X^{\gamma}]\mathbb{E}[X^{\beta^{j}}]$$
$$= \mathbb{E}\Big[\prod_{n \in \mathbb{Z}^{d}} X_{n}^{\gamma_{n}} X_{n}^{\beta_{n}^{j}}\Big] - \mathbb{E}\Big[\prod_{n \in \mathbb{Z}^{d}} X_{n}^{\gamma_{n}}\Big]\mathbb{E}\Big[\prod_{n \in \mathbb{Z}^{d}} X_{n}^{\beta_{n}^{j}}\Big]$$
$$= \prod_{n \in \mathbb{Z}^{d}} \mathbb{E}[X_{n}^{\gamma_{n}} + \beta_{n}^{j}] - \prod_{n \in \mathbb{Z}^{d}} \mathbb{E}[X_{n}^{\gamma_{n}}]\mathbb{E}[X_{n}^{\beta_{n}^{j}}].$$

If there exists any $n \in \mathbb{Z}^d$ such that $\gamma_n = 1$ and $\beta_n^j = 0$, the term $\mathbb{E}[X_n] = 0$ appears in both products, thus $\operatorname{Cov}(X^{\gamma}, X^{\beta^j}) = 0$ and any β^j for which $\operatorname{Cov}(X^{\gamma}, X^{\beta^j}) \neq 0$ must have $\beta_n^j \ge \gamma_n$ for every $n \in \mathbb{Z}^d$. If $\beta^j = \gamma$, since $\beta, \gamma \in B$ we must also have j = 0and $\beta = \gamma$. Otherwise, we have $\beta_n^j > \gamma_n$ for some $n \in \mathbb{Z}^d$, and it is straightforward to verify that every string *s* with $\varphi(s)^i = \beta^j$ must have length > *m*, therefore $p^k(\beta) = 0$ for every $1 \le k \le m$.

Note that there is some freedom in the choice of the representative set *B*, but one may choose *B* such that $\gamma \in B$ in all of the above cases, or alternatively replace the above choice of γ with some $\gamma^i \in B$.

For polynomials f of degree 2, 3, or 5, we must carefully analyze all cases. Since there are specific underlying distributions and polynomials f for which $\sigma(f)^2 = 0$, and we want an explicit description of all such cases, we need to explicitly compute all non-zero values of $p^k(\beta)$, for $1 \le k \le 5$.

Lemma 4.2. If $k \in \{1, 2, 3, 4, 5\}$ and γ is a multi-index with $p^k(\gamma) > 0$, then

1. γ is either equivalent to a unique β which equals $m \cdot \delta$ (for some $m \in \{1, 2, 3, 4, 5\}$), or to one of $\delta + \delta^e$, $2\delta + \delta^e$, or $2\delta + \delta^{-e}$ (for some $e \in \{e_1, e_2, \dots, e_d\}$).

2. The value of $p^{k}(\gamma) = p^{k}(\beta)$ is given in the table below (empty entries correspond to $p^{k}(\beta) = 0$):

To prove the lemma, we found no alternative to enumerating the relevant strings in S^k (for k = 1, 2, 3, 4, 5). We omit this technical proof.

Our method of verifying which polynomials $f(x) = \sum_{k=0}^{5} a_k x^k$ satisfy the condition $\sigma(f)^2 > 0$, is to describe random variables W_1, W_2, W_3, W_4, W_5 such that $\operatorname{Var}(\sum_{k=1}^{5} a_k W_k) = \sigma(f)^2$ (for any choice of coefficients a_0, a_1, \ldots, a_5). We then explore the random variable $\sum_{k=1}^{5} a_k W_k$ and determine under which conditions it is almost surely constant. If deg $(f) \le 3$, we may replace $\{W_i\}$ with a simpler set of random variables, T_1, T_2, T_3 . We verify this case before approaching polynomials of degree 5:

Proposition 4.3. Let $f(x) = \sum_{k=0}^{m} a_k x^k \in \mathbb{R}[x]$ be a polynomial of degree $1 \le m \le 3$. *Then*

- 1. *if the underlying distribution* $(d\rho)$ *is supported by more than three values, then* $\sigma(f)^2 > 0$;
- 2. *if the underlying distribution is supported by exactly three values, denoted by* $a, b, c \in \mathbb{R}$, then $\sigma(f)^2 = 0$ if and only if $f = a_3\tilde{g}_3 + a_0$, where

$$\tilde{g}_3(x) = x^3 - (a+b+c)x^2 + (ab+ac+bc-6d)x;$$

3. *if the underlying distribution is supported by exactly two values, denoted by* $a, b \in \mathbb{R}$, then $\sigma(f)^2 = 0$ if and only if

$$f = a_3g_3 + a_2g_2 + a_0,$$

where

$$g_3(x) = x^3 - (a^2 + ab + b^2 + 6d)x,$$

$$g_2(x) = x^2 - (a + b)x.$$

Proof. Let Z denote both a random variable distributed by $d\rho$, and the variable in polynomial ring $\mathbb{R}[Z]$. Define

$$T_1 = Z$$
, $T_2 = Z^2$, $T_3 = Z^3 + 6dZ$.

Using Lemma 4.2, we see that for every k = 1, 2, 3, we have $T_k = \sum_{n=1}^{3} p^k(n\delta)Z^n$, thus for every $k, \ell = 1, 2, 3$,

$$\operatorname{Cov}(T_k, T_\ell) = \sum_{n,m=1,2,3} p^k(n\delta) p^\ell(m\delta) \operatorname{Cov}(Z^n, Z^m)$$
$$\sum_{\substack{\beta, \gamma = \delta, 2\delta, 3\delta}} p^k(\beta) p^\ell(\gamma) \operatorname{Cov}(X^\beta, X^\gamma)$$
$$\sum_{\substack{\beta, \gamma \in B}} p^k(\beta) p^\ell(\gamma) \sum_{j \in \mathbb{Z}^d} \operatorname{Cov}(X^\beta, X^{\gamma^j})$$

(we may assume without loss of generality that δ , 2δ , $3\delta \in B$). We now deduce from Proposition 3.8 that

$$\sigma(f)^2 = \operatorname{Var}(a_3T_3 + a_2T_2 + a_1T_1),$$

which is zero if and only if $F = a_3T_3 + a_2T_2 + a_1T_1$ is almost surely constant, as a random variable. As a polynomial, $F \in \mathbb{R}[Z]$ has at most 3 distinct roots, so if Z is supported by more than 3 points, F is non-constant as a random variable, thus Var(F) > 0, proving (1).

Observe that any assignment of a value to the random variable Z corresponds to a ring homomorphism $\mathbb{R}[Z] \to \mathbb{R}$. Furthermore, if we only assign values from $\{a, b, c\}$, all three assignment homomorphisms factor through the quotient ring $\mathbb{R}[Z]/((Z-a)(Z-b)(Z-c))$. Write $P \equiv Q$ for two polynomials P, Q, if they have the same projection in the quotient. Note that $P \equiv Q$ if and only if as random variables, P = Q almost surely. Clearly, Var(F) = 0 as a random variable if and only if $F \equiv \text{const in } \mathbb{R}[Z]$. Now, write

$$(Z-a)(Z-b)(Z-c) = Z^{3} - (a+b+c)Z^{2} + (ab+ac+bc)Z - abc$$

= $T_{3} - (a+b+c)T_{2} + (ab+ac+bc-6d)T_{1} - abc$,

and deduce (2): the polynomial $\tilde{g}_3(x)$ has $\sigma(\tilde{g}_3) = 0$ from the above, therefore

$$\sigma(a_3\tilde{g}_3+a_0)=0.$$

If $f \neq a_3 \tilde{g}_3 + a_0$, we see that $F = a_3 T_3 + a_2 T_2 + a_1 T_2$ is equivalent to a polynomial of degree 1 or 2 in $\mathbb{R}[Z]$, and therefore is not fixed under assignments from $\{a, b, c\}$.

Finally, if $d\rho$ is supported on $\{a, b\}$, the same arguments hold with a different quotient ring, $\mathbb{R}[Z]/((Z-a)(Z-b))$. Now, note that

$$0 \equiv (Z - a)(Z - b) = Z^{2} - (a + b)Z + ab,$$

therefore

$$T_2 - (a+b)T_1 = Z^2 - (a+b)Z \equiv \text{const}$$

proving $\sigma(g_2)^2 = 0$. We also have

$$Z^{3} \equiv (a+b)Z^{2} - abZ \equiv (a^{2} + ab + b^{2})Z - ab(a+b),$$

therefore

$$T_3 - (a^2 + ab + b^2 + 6d)T_1 = Z_3 - (a^2 + ab + b^2)Z \equiv \text{const},$$

proving $\sigma(g_3)^2 = 0$. If $f \neq a_3g_3 + a_2g_2 + a_0$, then the above computations show that *F* is equivalent to a polynomial of degree 1, which is not equivalent to any constant, therefore Var(*F*) > 0.

Proposition 4.4. Let $f(x) = \sum_{k=0}^{5} a_k x^k \in \mathbb{R}[x]$ be a polynomial of degree 5. Then

- 1. *if the underlying distribution* $(d\rho)$ *is supported by more than two values, then* $\sigma(f)^2 > 0$,
- 2. *if the underlying distribution is supported by exactly two values, denoted by* $a, b \in \mathbb{R}$, then $\sigma(f)^2 = \sigma(f a_5g_5)^2$, where

$$2g_5(x) = 2x^5 - 5(a+b)x^4 + [3(a^4 + b^4) + 8(a^3b + a^2b^2 + ab^3) + 20d(a^2 + b^2) + 100dab - 120d^2 + 60d]x.$$

In particular, $\sigma(g_5)^2 = 0$.

Proof. For every $n \in \Lambda_1^d$, let $Z_n = X_n$. We regard the variables $\{Z_n\}$ both as 3^d independent random variables distributed by $d\rho$, and as the variables in the polynomial ring $R = \mathbb{R}[Z_n \mid n \in \Lambda_1^d]$. Define

$$W_1 = 3^{-d/2} \sum_{n \in \Lambda_1^d} Z_n, \quad W_2 = 3^{-d/2} \sum_{n \in \Lambda_1^d} Z_n^2.$$

and

$$W_3 = 3^{-d/2} \sum_{n \in \Lambda_1^d} (Z_n^3 + 6dZ_n).$$

Let *E* consist of all unordered pairs $\{n, m\}$, such that $n, m \in \Lambda_1^d$ differ in exactly one coordinate, that is

$$E = \{\{n, m\} \mid n, m \in \Lambda_1^d, \#\{1 \le v \le d \mid n_v \ne m_v\} = 1\}.$$

Now, define

$$W_4 = 3^{-d/2} \sum_{n \in \Lambda_1^d} (Z_n^4 + 8dZ_n^2) + 3^{-d/2} \sum_{\{n,m\} \in E} 4Z_n Z_m$$

and

$$W_{5} = 3^{-d/2} \sum_{n \in \Lambda_{1}^{d}} (Z_{n}^{5} + 10dZ_{n}^{3} + (60d^{2} - 30d)Z_{n}) + 3^{-d/2} \sum_{\{n,m\} \in E} 5(Z_{n}^{2}Z_{m} + Z_{n}Z_{m}^{2}).$$

Following Lemma 4.2 and a straightforward computation that we omit, we verify that

$$\operatorname{Cov}(W_k, W_\ell) = \sum_{\beta \in B} p^k(\beta) p^\ell(\gamma) \sum_{j \in \mathbb{Z}^d} \operatorname{Cov}(X^\beta, X^{\gamma^j})$$

for every $k, \ell \in \{1, 2, 3, 4, 5\}$ then deduce from Proposition 3.8 that

$$\sigma(f)^2 = \operatorname{Var}\Big(\sum_{k=1}^5 a_k W_k\Big).$$

Denote

$$F = \sum_{k=1}^{5} a_k W_k.$$
 (4.2)

As in the proof of Proposition 4.3, we note that Var(F) = 0 if and only if F is almost surely constant as a random variable. This shows that Var(F) > 0 if $d\rho$ is not finitely supported: generally if $f \in \mathbb{R}[x_1, \ldots, x_m]$ is a non-constant multivariate polynomial, and S is a set such that $f(s_1, \ldots, s_m) = 0$ for every $s_1, \ldots, s_m \in S$, then straightforward induction on m shows that $|S| \leq \deg(f)$.

So, let us assume henceforth that the variables Z_n are supported by a finite set $\operatorname{supp}(d\rho) \subset \mathbb{R}$. Denote $q(x) = \prod_{a \in \operatorname{supp}(d\rho)} (x-a)$, let $Q_n = q(Z_n) \in R$, and let $I \subset R$ be the ideal generated by the polynomials $\{Q_n\}_{n \in \Lambda_1^d}$. Every possible assignment of values to $\{Z_n\}$ corresponds to a ring homomorphism $R \to \mathbb{R}$. If we only assign values from $\operatorname{supp}(d\rho)$, the homomorphism factors through the quotient ring R/I. Write $P \equiv Q$ for two polynomials P, Q, if they have the same projection in the quotient. Note that $P \equiv Q$ in R if and only if, as random variables, P = Q almost surely. Clearly, $\operatorname{Var}(F) = 0$ as a random variable if and only if $F \equiv \operatorname{const}$ in R.

Next, we denote $\omega_1 = \omega_2 = \omega_3 = 0$,

$$\omega_4 = 3^{-d/2} \sum_{\{n,m\} \in E} 4Z_n Z_m,$$

and

$$\omega_5 = 3^{-d/2} \sum_{\{n,m\}\in E} 5(Z_n^2 Z_m + Z_n Z_m^2)$$
(4.3)

(so each ω_k is the part of W_k which is a sum of products involving more than one variable). Now, rewrite (4.2) as

$$F = a_5\omega_5 + a_4\omega_4 + \sum_{k=1}^5 a_k(W_k - \omega_k),$$

and note that if $|\operatorname{supp}(d\rho)| \le k$, then $W_k - \omega_k$ is equivalent to a polynomial of degree lowith respect to $|\operatorname{supp}(d\rho)|$: every term of the form Z_n^k is equivalent to

$$Z_n^k - Z_n^{k-|\operatorname{supp}(\mathrm{d}\rho)|} q(Z_n),$$

with degree strictly lowith respect to k. Thus

$$\sum_{n \in \Lambda_1} Z_n^k \equiv \sum_{n \in \Lambda_1^d} Z_n^k - Z_n^{k-|\operatorname{supp}(d\rho)|} q(Z_n),$$

and summing over $n \in \Lambda_1^d$ allows us to reduce $W_k - \omega_k$ to an equivalent combination of $W_1 - \omega_1, \ldots, W_{k-1} - \omega_{k-1}$, to eventually obtain

$$F \equiv \tilde{F} = a_5\omega_5 + a_4\omega_4 + \sum_{k=1}^{|\operatorname{supp}(d\rho)| - 1} \tilde{a}_k(W_k - \omega_k)$$
(4.4)

for some $\tilde{a}_1, \ldots, \tilde{a}_{|\operatorname{supp}(d\rho)|-1} \in \mathbb{R}$. We are now ready to prove that $\operatorname{Var}(\tilde{F}) > 0$, whenever $|\operatorname{supp}(d\rho)| \ge 3$. Otherwise, $\operatorname{Var}(\tilde{F}) = 0$ implies that $\tilde{F} - c \in I$ for some constant *c*, so we can find polynomials $H_n \in R$, such that

$$\widetilde{F} - c = \sum_{n \in \Lambda_1^d} H_n \cdot Q_n \tag{4.5}$$

in R. Fix some $a \in \text{supp}(d\rho)$, and let $\psi_a \colon R \to \mathbb{R}[x]$ be the ring homomorphism, defined by

$$\psi_a(Z_n) = \begin{cases} x, & n = 0, \\ a, & n \neq 0. \end{cases}$$

We have $\psi_a(Q_n) = q(a) = 0$ for every $n \neq 0$; so, when we apply ψ_a to (4.5), we obtain the equality

$$\psi_a(\tilde{F}) - c = h(x)q(x) \tag{4.6}$$

in $\mathbb{R}[x]$, where $h(x) = \psi_a(H_0)$. Note that $W_k - \omega_k$ has degree k in R, therefore $\psi_a(W_k - \omega_k)$ has degree at most k. Clearly, $\psi_a(a_5\omega_5 + a_4\omega_4)$ has degree 2, so from (4.4) the polynomial in the left-hand side of (4.6) has degree strictly less than $|\operatorname{supp}(d\rho)|$. But q(x) has degree $|\operatorname{supp}(d\rho)|$, so we must have h(x) = 0 (otherwise

the right-hand side of (4.6) would have degree $|\operatorname{supp}(d\rho)|$ or higher). We deduce that $\psi_a(\tilde{F}) - c = 0$ as a polynomial in $\mathbb{R}[x]$.

Since for every $n \in \Lambda_1^d$ there are $\#\{m \mid \{n, m\} \in E\} = 2d$ values of *m* for which $5 \cdot Z_n^2 Z_m$ appears in the sum (4.3), the coefficient of x^2 in $\psi_a(\omega_5)$ is $2d \cdot 3^{-d/2} \cdot 5a$. We deduce that the coefficient of x^2 in $\psi_a(\tilde{F}) - c = 0$ is

$$a_5 \cdot 10d \cdot 3^{-d/2} \cdot a + c' = 0,$$

where c' does not depend on our choice of $a \in \text{supp}(d\rho)$. Since $a_5 \neq 0$, there is at most one $a \in \mathbb{R}$ satisfying the above equation. However, for any $b \in \text{supp}(d\rho)$, applying ψ_b to (4.5) allows us to obtain $a_5 \cdot 10d \cdot 3^{-d/2} \cdot b + c' = 0$, which is a contradiction. This concludes the proof of (1).

If supp $(d\rho) = \{a, b\}$ then $q(Z_n) = (Z_n - a)(Z_n - b) \in I$, therefore

$$Z_n^2 \equiv (a+b)Z_n - ab \tag{4.7}$$

for every $n \in \Lambda_1^d$, thus (4.3) becomes $\omega_5 \equiv \frac{5}{2}(a+b)\omega_4 - 20dabW_1$, which allows us to deduce

$$a_5\omega_5 + a_4\omega_4 \equiv -20a_5dabW_1 \tag{4.8}$$

whenever $a_4 = -\frac{5}{2}(a+b)a_5$.

Finally, from (4.7) we verify

$$Z_n^2 \equiv (a+b)Z_n - ab, \tag{4.9a}$$

$$Z_n^3 \equiv (a^2 + ab + b^2)Z_n - ab(a+b),$$
(4.9b)

$$Z_n^4 \equiv (a^3 + a^2b + ab^2 + b^3)Z_n - ab(a^2 + ab + b^2), \qquad (4.9c)$$

$$Z_n^5 \equiv (a^4 + a^3b + a^2b^2 + ab^3 + b^4)Z_n - \text{const.}$$
(4.9d)

Summing over $n \in \Lambda_1^d$ allows us to reduce $3^{-d/2} \sum_n Z_n^k$ (for k = 2, 3, 4, 5) to equivalent expressions involving W_1 and constants, and along with (4.8) and the definitions of W_1, W_4, W_5 we deduce

$$2W_5 - 5(a+b)W_4 + \text{const} \equiv (-3a^4 - 8a^3b - 8a^2b^2 - 8ab^3 - 3b^4 - 20da^2 - 100dab - 20db^2 + 120d^2 - 60d)W_1.$$

From here it follows that $\sigma(g_5)^2 = 0$ and that $\sigma(f)^2 = \sigma(f - cg_5)^2$ for any polynomial *F* and constant *c*, concluding our proof.

Proof of Theorem 1.1. Given a polynomial $f(x) = \sum_{k=0}^{m} a_k x^k \in \mathbb{R}[x]$, we have

$$\frac{\operatorname{Tr}(f(H_L)) - \mathbb{E}[\operatorname{Tr}(f(H_L))]}{(2L+1)^{d/2}} \xrightarrow{d} N(0, \sigma(f)^2)$$

for $\sigma(f)^2 \in [0, \infty)$ as $L \to \infty$, from Proposition 3.8. From Propositions 4.1 and 4.3 we determine the cases in which $\sigma(f)^2 > 0$ whenever deg $(f) \neq 5$. If deg(f) = 5, we know from Proposition 4.4 that $\sigma(f)^2 = \sigma(f - a_5g_5)^2$. If $f - a_5g_5$ is non-constant and deg $(f - a_5g_5)$ is 1 or 4, we determine that $\sigma(f)^2 = \sigma(f - a_5g_5)^2 > 0$ from Proposition 4.1, otherwise we use Proposition 4.3 to determine the positivity.

5. Appendix – Proof of Theorem 3.3

In the setting of Theorem 3.3, we consider a d-dimensional array of weakly dependent random variables. Explicitly, we prove a central limit theorem which is valid in the setting of m-dependent random variables, which we now define:

Definition 5.1. Let $\{Y_i\}_{i \in \mathbb{Z}^d}$ be a sequence of random variables. We say that the sequence is *m*-dependent if for any two finite sets of indices $I, J \subset \mathbb{Z}^d$ which satisfy |i - j| > m for every $i \in I$ and $j \in J$, the corresponding sets of random variables

$$\{Y_i\}_{i\in I}, \quad \{Y_j\}_{j\in J}$$

are independent.

Note that this definition extends a notion of *m*-dependence from [4] defined for sequences of variables indexed by \mathbb{N} (the definition of *m*-dependence in [4] is equivalent to *m*-dependence as defined above, when we take d = 1 and $Y_i = 0$ for every $i \notin \mathbb{N}$). In [4], Hoeffding and Robbins proved the following central limit theorem:

Theorem 5.2 (Hoeffding–Robbins). Let $\{X_i\}_{i \in \mathbb{N}}$ be an *m*-dependent sequence of random variables satisfying $\mathbb{E}[X_i] = 0$ and $\mathbb{E}[|X_i|^3] \leq R^3 < \infty$ for every $i \in \mathbb{N}$, and

$$\lim_{p \to \infty} p^{-1} \sum_{h=1}^{p} A_{i+h} = A$$

uniformly for all $i \in \mathbb{N}$, where

$$A_{i} = \mathbb{E}[X_{i+m}^{2}] + 2\sum_{j=1}^{m} \mathbb{E}[X_{i+m-j}X_{i+m}].$$

Then

$$\frac{X_1 + \dots + X_n}{n^{\frac{1}{2}}} \stackrel{d}{\to} N(0, A).$$

Theorem 5.2 allows us to deduce a central limit theorem for d = 1, and the following theorem by Neumann [8] will allow us to prove an induction argument on d:

Theorem 5.3 (Neumann). Suppose that $\{X_{n,k} \mid n \in \mathbb{N}, k = 1, 2, ..., n\}$ is a triangular scheme of random variables with $\mathbb{E}[X_{n,k}] = 0$ and

$$\sum_{k=1}^{n} \mathbb{E}[X_{n,k}^2] \le C$$

for all n, k and some $C < \infty$. We assume that

$$\sigma_n^2 = \operatorname{Var}(X_{n,1} + \dots + X_{n,n}) \xrightarrow[n \to \infty]{} \sigma^2 \in [0, \infty),$$

and that

$$\sum_{k=1}^{n} \mathbb{E}[X_{n,k}^2 \mathbb{1}(|X_{n,k}| > \varepsilon)] \xrightarrow[n \to \infty]{} 0$$

holds for all $\varepsilon > 0$. Furthermore, we assume that there exists a summable sequence $(\theta_r)_{r \in \mathbb{N}}$ such that for all $u \in \mathbb{N}$ and all indices

$$1 \le s_1 < s_2 < \dots < s_u < s_u + r = t_1 \le t_2 \le n_1$$

the following upper bounds for covariances hold true: for all measurable functions $g: \mathbb{R}^u \to \mathbb{R}$ with $||g||_{\infty} = \sup_{x \in \mathbb{R}^u} |g(x)| \le 1$, we have

$$|\operatorname{Cov}(g(X_{n,s_1},\ldots,X_{n,s_u})X_{n,s_u},X_{n,t_1})| \le \left(\mathbb{E}[X_{n,s_u}^2] + \mathbb{E}[X_{n,t_1}^2] + \frac{1}{n}\right)\theta_r \quad (5.1)$$

and

$$|\operatorname{Cov}(g(X_{n,s_1},\ldots,X_{n,s_u}), X_{n,t_1}X_{n,t_2})| \le \left(\mathbb{E}[X_{n,t_1}^2] + \mathbb{E}[X_{n,t_2}^2] + \frac{1}{n}\right)\theta_r.$$
 (5.2)

Then

$$X_{n,1} + \dots + X_{n,n} \xrightarrow{d} N(0,\sigma^2)$$

as $n \to \infty$.

Our central limit theorem for *m*-dependent random variables follows:

Proposition 5.4. Let $\{Y_i\}_{i \in \mathbb{Z}^d}$ be an identically distributed *d*-dimensional *m*-dependent array of random variables such that $\mathbb{E}[Y_i] = 0$, and $\mathbb{E}[|Y_i|^3] < \infty$.

Then

$$\frac{1}{(2L+1)^{d/2}} \sum_{i \in \Lambda_L^d} Y_i \stackrel{d}{\to} N(0,\sigma^2),$$

where

$$\sigma^2 = \lim_{L \to \infty} \frac{1}{(2L+1)^d} \operatorname{Var}\Big(\sum_{i \in \Lambda_L^d} Y_i\Big).$$

Proof. By induction on d. For d = 1, this is a straightforward application of Theorem 5.2 to the random variables $\{X_i\}_{i \in \mathbb{N}}$, defined by $X_i = Y_{i+m} + Y_{-i-m}$ (noting that for i > m, $\{X_i\}_{i \in \mathbb{N}}$ are identically distributed and *m*-dependent, and the exclusion of a finite set of random variables $\{Y_i : |i| \le m\}$ from the sum has no effect on the limit distribution).

We now assume by induction that the proposition holds for some $d \in \mathbb{N}$, and prove it in dimension d + 1. For every $L \in \mathbb{N}$ we denote n = 2L + 1, rewrite

$$\frac{1}{(2L+1)^{(d+1)/2}} \sum_{i \in \Lambda_L^{d+1}} Y_i = \sum_{j=-L}^L Z_{n,j},$$

where

$$Z_{n,j} = \frac{1}{n^{1/2}} \cdot \frac{1}{n^{d/2}} \sum_{i \in I_{n,j}} Y_i$$

and

$$I_{n,j} = \Lambda_L^d \times \{j\} = \{(i_1, \dots, i_{d+1}) \in \Lambda_L^{d+1} \mid i_{d+1} = j\}$$

are defined for every $j \in \Lambda_L$. Our proof will be completed by applying Theorem 5.3 to the random variables

$$X_{n,k} = \begin{cases} Z_{n,k-L-1}, & n = 2L+1, \\ Z_{n+1,k-L-1}, & n = 2L, \end{cases}$$

which are defined for every $n \in \mathbb{N}$ and k = 1, 2, ..., n. We will apply the requirements of the theorem to the corresponding variables $Z_{n,j}$ (we henceforth ignore even values of n).

Fixing any $j \in \mathbb{Z}$, we may identify $I_{n,j}$ with Λ_L^d , and note that the *d*-dimensional array $\{Y_i \mid i \in \mathbb{Z}^{d+1}, i_{d+1} = j\}$ is identically distributed and *m*-dependent (the distribution of the array is independent of $j \in \mathbb{Z}$ as well). The induction hypothesis now applies, and we deduce

$$\sqrt{n}Z_{n,j} = \frac{1}{n^{d/2}} \sum_{i \in I_{n,j}} Y_i \stackrel{d}{\to} N(0, \sigma_d^2)$$
(5.3)

as $n \to \infty$, uniformly in j, for some $\sigma_d^2 \ge 0$. The variables $Z_{n,j}$ are "well behaved," in the sense that for any sufficiently large n,

$$\mathbb{E}[Z_{n,j}^2] = \operatorname{Var}(Z_{n,j}) \le \frac{1}{n}(\sigma_d^2 + 1)$$

(thus there exists C > 0 such that $\mathbb{E}[Z_{n,j}^2] \leq \frac{C}{n}$ for all $n \in \mathbb{N}$ and $j \in \Lambda_L$). We deduce that

$$\mathbb{E}[Z_{n,j}] = 0, \quad \sum_{j=-L}^{L} \mathbb{E}[Z_{n,j}^2] \le C.$$

Additionally, since the finite sequence $\{Z_{n,j}\}_{j \in \Lambda_L}$ is both identically distributed and *m*-dependent (for every $n = 2L + 1 \in \mathbb{N}$), one can verify that

$$\operatorname{Var}\left(\sum_{j=-L}^{L} Z_{n,j}\right) \xrightarrow[n \to \infty]{} \sigma^{2} < \infty.$$

Next, we prove that

$$\sum_{j=-L}^{L} \mathbb{E}[Z_{n,j}^2 \mathbb{1}(|Z_{n,j}| > \varepsilon)] \to [n \to \infty]0$$

for every $\varepsilon > 0$. Note that

$$\sum_{j=-L}^{L} \mathbb{E}[Z_{n,j}^2 \mathbb{1}(|Z_{n,j}| > \varepsilon)] = n \mathbb{E}[Z_{n,j}^2 \mathbb{1}(|Z_{n,j}| > \varepsilon)]$$
$$= \mathbb{E}[n(Z_{n,j})^2 \mathbb{1}(|\sqrt{n}Z_{n,j}| > \varepsilon\sqrt{n})].$$
(5.4)

From the induction hypothesis, we know that $\sqrt{n}Z_{n,j} \xrightarrow{d} N(0, \sigma_d^2)$. We deduce that for every M > 0 we have

$$\sqrt{n}Z_{n,j}1(|\sqrt{n}Z_{n,j}| > M) \xrightarrow{d} \Phi_M,$$
(5.5)

where Φ_M is a random variable satisfying $\mathbb{E}[\Phi_M] = 0$, and

$$\operatorname{Var}(\Phi_M) = \begin{cases} 2\int\limits_{M}^{\infty} \frac{t^2}{\sigma_d \sqrt{2\pi}} \exp\left(-\frac{t^2}{2\sigma_d^2}\right) dt, & \sigma_d^2 > 0, \\ 0, & \sigma_d^2 = 0. \end{cases}$$

Choose some M > 0 so that $Var(\Phi_M)$ is arbitrarily close to 0. For every $\varepsilon > 0$, any sufficiently large $n \in \mathbb{N}$ satisfies $\varepsilon \sqrt{n} > M$, so

$$1(|\sqrt{n}Z_{n,j}| > \varepsilon\sqrt{n}) \le 1(|\sqrt{n}Z_{n,j}| > M),$$

and (5.4) now becomes

$$\sum_{j=-L}^{L} \mathbb{E}[Z_{n,j}^2 \mathbb{1}(|Z_{n,j}| > \varepsilon)] = \mathbb{E}[n(Z_{n,j})^2 \mathbb{1}(|\sqrt{n}Z_{n,j}| > \varepsilon\sqrt{n})]$$

$$\leq \mathbb{E}[n(Z_{n,j})^2 \mathbb{1}(|\sqrt{n}Z_{n,j}| > M)]$$

$$= \operatorname{Var}(\sqrt{n}Z_{n,j} \mathbb{1}(|\sqrt{n}Z_{n,j}| > M)) \xrightarrow[n \to \infty]{} \operatorname{Var}(\Phi_M)$$

(due to (5.5)).

It remains to show that there exists a summable sequence $(\theta_r)_{r \in \mathbb{N}}$ so that the upper bounds for covariances required in Neumann's theorem hold (equations (5.1) and (5.2), for all relevant cases). From the *m*-dependence of the finite sequence $\{Z_{n,j}\}_{j \in \Lambda_L}$, we deduce that the left-hand sides of (5.1) and (5.2) equal 0 whenever r > m, so we conclude by finding some $\theta_1, \ldots, \theta_m < \infty$. A straightforward computation shows that (5.1) holds as long as $\theta_r \ge 1$. To prove (5.2), we use

$$\operatorname{Var}(g(Z_{n,s_1},\ldots,Z_{n,s_u})) \leq \mathbb{E}[g(Z_{n,s_1},\ldots,Z_{n,s_u})^2] \leq 1$$

(as $||g||_{\infty} \leq 1$) to obtain

$$|\operatorname{Cov}(g(Z_{n,s_1},\ldots,Z_{n,s_u}), Z_{n,t_1}Z_{n,t_2})| \le \sqrt{\operatorname{Var}(g(Z_{n,s_1},\ldots,Z_{n,s_u}))\operatorname{Var}(Z_{n,t_1}Z_{n,t_2})} \le \sqrt{\operatorname{Var}(Z_{n,t_1}Z_{n,t_2})},$$

and we conclude by showing that for some $\theta < \infty$,

$$\sqrt{\operatorname{Var}(Z_{n,t_1}Z_{n,t_2})} \le \frac{1}{n}\theta$$

holds for every n = 2L + 1 and $t_1, t_2 \in \Lambda_L$. Equivalently, we will show that

$$\sup_{n,t_1,t_2} \operatorname{Var}(\sqrt{n}Z_{n,t_1} \cdot \sqrt{n}Z_{n,t_2}) < \infty.$$

From (5.3), we deduce that

$$\sup_{n} \operatorname{Var}(\sqrt{n} Z_{n,t_1} \cdot \sqrt{n} Z_{n,t_2}) < \infty$$

for every $t_1, t_2 \in \mathbb{Z}$. Furthermore, since our initial variables $\{Y_i\}_{i \in \mathbb{Z}^d}$ are identically distributed, the value of $\operatorname{Var}(\sqrt{nZ_{n,t_1}} \cdot \sqrt{nZ_{n,t_2}})$ depends only on *n* and $t_2 - t_1$, and since our variables are *m*-dependent, it is enough to consider $|t_2 - t_1| \in \{0, 1, \dots, m+1\}$. This concludes our proof.

Proof of Theorem 3.3. Theorem 3.3 will follow from Proposition 5.4, applied to the variables

$$Y_i = \sum_{\beta \in B} a_\beta (X^{\beta^i} - \mathbb{E}[X^{\beta^i}]).$$

Clearly, the variables $\{Y_i\}_{i \in \mathbb{Z}^d}$ are identically distributed (since $\{X_n\}_{n \in \mathbb{Z}^d}$ are), and $\mathbb{E}[Y_i] = 0$. Since every X_n has finite moments, so do Y_i (as a finite sum of products of the variables $\{X_n\}_{n \in \mathbb{Z}^d}$). In particular, $\mathbb{E}[|Y_i|^3] < \infty$.

Since $a_{\beta} \neq 0$ only for finitely many $\beta \in B$, one can find sufficiently large *m*, such that whenever |j - i| > m and $a_{\beta}, a_{\gamma} \neq 0$, the supports of β^i and γ^j are disjoint. From here it follows that $\{Y_i\}_{i \in \mathbb{Z}^d}$ is *m*-dependent.

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