

Localization for Random Schrödinger Operators with Correlated Potentials

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Abstract. We prove localization at high disorder or low energy for lattice Schrödinger operators with random potentials whose values at different lattice sites are correlated over large distances. The class of admissible random potentials for our multiscale analysis includes potentials with a stationary Gaussian distribution whose covariance function $C(x, y)$ decays as $|x - y|^{-\theta}$, where $\theta > 0$ can be arbitrarily small, and potentials whose probability distribution is a completely analytical Gibbs measure. The result for Gaussian potentials depends on a multivariable form of Nelson's best possible hypercontractive estimate.

1. Introduction

We consider the random Schrödinger operator $H = -\Delta + V$ on $l^2(\mathbf{Z}^d)$, where Δ is the centered finite difference Laplacian, i.e., $\Delta(x, y) = 1$ if $|x - y| = 1$ and zero otherwise, and V is an ergodic potential, i.e., $\{V(x); x \in \mathbf{Z}^d\}$ is an ergodic stochastic process. The motivation for studying this class of operators comes from Solid State Physics, where one is interested in the behavior of an electron in a random background. This model was first introduced by Anderson [1] and is known as the Anderson tight-binding model.

It is well known that the spectrum of the random operator H is independent of the choice of potential with probability one [2, 3, 4]. The same is true of the decomposition of the spectrum into pure point, absolutely continuous and singular continuous spectrum [3, 4].

The random operator H exhibits localization in an energy interval I if it has only pure point spectrum in I with probability one. In this case, if the eigenfunc-

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tions corresponding to eigenvalues in I decay exponentially (polynomially), we say we have exponential (polynomial) localization.

Exponential localization for the Anderson tight-binding model is well understood in one dimension (e.g., [4]), where it was first established in the continuum by Gold'sheid, Molchanov and Pastur [5]. When the $\{V(x), x \in \mathbf{Z}^d\}$ are independent identically distributed random variables, one always has exponential localization [3, 6–9]. The same is true on the strip [10, 11]. For arbitrary ergodic potentials localization also holds if the joint conditional probability of the values of the potential at two neighboring sites with respect to the values of the potential in the remaining sites has an absolutely continuous component [12, 13, 4]. In particular, Simon [12] noticed that the result holds for nondeterministic ergodic Gaussian potentials and for potentials whose probability distribution is a Gibbs measure with an “a priori” measure with an absolutely continuous component and finite range interactions. The result for Gaussian potentials was also noticed by Pastur [14] (see also [4]). In one dimension exponential localization has also been proved for quasi-periodic potentials with large coupling constant [15–17].

In more than one dimension localization results had only been proven for independent potentials. For such potentials exponential localization was proven at either high disorder or low energy [18, 6, 19, 20, 21, 9].

In this article we prove localization at high disorder or low energy for certain correlated potentials, in any dimension. The class of admissible potentials include potentials whose values at different lattice sites are strongly correlated over large distances. Our proof yields only polynomial localization.

Our proof involves a multiscale analysis based on the ideas used by von Dreifus and Spencer [21, 22] and by von Dreifus and Klein [9] to give a simpler proof of localization in the independent case.

The class of potentials for which we prove localization includes stationary Gaussian potentials whose covariance function $C(x, y)$ decays as $|x - y|^{-\theta}$, where $\theta > 0$ can be arbitrarily small, and potentials whose probability distribution is a completely analytical Gibbs measure. The result for Gaussian potentials depends on a multivariable form of Nelson's best possible hypercontractive estimate; for Gibbs potentials we use the work of Dobrushin and Shlosman [23, 24].

2. Statement of Results

We start with some notations and definitions. If $\Lambda \subset \mathbf{Z}^d$, we denote by H_Λ the operator H restricted to $l^2(\Lambda)$ with zero boundary condition outside Λ . The corresponding Green's function is $G_\Lambda(z) = (H_\Lambda - z)^{-1}$, defined for $z \notin \sigma(H_\Lambda)$. We will write $G_\Lambda(z; x, y) = (H_\Lambda - z)^{-1}(x, y)$ for $x, y \in \Lambda$. If $\Lambda = \mathbf{Z}^d$, we simply write $G(z; x, y)$. Notice that we omit the dependence of H_Λ and G_Λ on the potential V .

If $x \in \mathbf{Z}^d$, $x = (x_1, \dots, x_d)$, we let $\|x\| = \|x\|_\infty = \max\{|x_1|, \dots, |x_d|\}$. It will be convenient to use this norm; distances in \mathbf{Z}^d will always be taken with respect to this norm. Occasionally we may also need the usual Euclidean norm $|x| = \|x\|_2 = (x_1^2 + \dots + x_d^2)^{1/2}$.

We will denote by \mathbf{E} and \mathbf{P} the expectation and probability measure on the underlying probability space for the stochastic process $\{V(x), x \in \mathbf{Z}^d\}$.

Definition. The ergodic potential V is a Wegner potential if the conditional probability distribution of $V(0)$ given $V(0)^\perp = \{V(x), x \in \mathbf{Z}^d \setminus \{0\}\}$ is absolutely continuous with respect to Lebesgue measure for \mathbf{P} -a.e. $V(0)^\perp$, and

$$\delta \equiv \left\| \frac{\mathbf{P}(dV(0)|V(0)^\perp)}{dV(0)} \right\|_\infty < \infty,$$

where the L^∞ -norm is taken with respect to \mathbf{P} .

If V is a Wegner potential, Wegner's estimate on the density of states [25, 18, 8, 4] applies and we have

$$\mathbf{P}\{d(E, \sigma(H_\Lambda)) \leq \varepsilon\} \leq \frac{2}{\delta} \varepsilon |\Lambda| \tag{2.1}$$

for all $E \in \mathbf{R}$, $\varepsilon > 0$, $\Lambda \subset \mathbf{Z}^d$. Furthermore, Kotani's trick can be used so the analysis of Simon and Wolff [20] and Delyon, Levi and Souillard [19] can be applied to Wegner potentials.

Let $x \in \mathbf{Z}^d$, $L > 0$. We will denote by $\Lambda_L(x)$ the cube centered at x with sides of length L , i.e.,

$$\Lambda_L(x) = \left\{ y \in \mathbf{Z}^d; \|x - y\| \leq \frac{L}{2} \right\}.$$

Notice that $|\Lambda_L(x)| \leq (L + 1)^d$.

Given $\Lambda \subset \mathbf{Z}^d$ an event A on Λ is an event that depends only on $\{V(y); y \in \Lambda\}$. If A is an event on the cube $\Lambda(0)$, we will use $A(x)$ to denote the same event shifted to the cube $\Lambda_L(x)$.

Definition. Let $1 < \alpha < \rho$, $K \in \{2, 3, 4, \dots\}$. We will say that the ergodic potential V is of type (α, ρ, K) if for all L sufficiently large, given $x_1, \dots, x_K \in \mathbf{Z}^d$ with $\|x_i - x_j\| \geq \frac{1}{2}L^\alpha$ for $i \neq j$, and any event A on $\Lambda_L(0)$, we have

$$\mathbf{P}\left\{ \bigcap_{i=1}^K A(x_i) \right\} \leq c \mathbf{P}\{A\}^\rho \tag{2.2}$$

for some constant $c < \infty$.

If V is an independent potential, i.e., the $V(x)$, $x \in \mathbf{Z}^d$, are i.i.d.r.v.'s, then V is of type (α, K, K) for any $1 < \alpha < K$, $K = 2, 3, \dots$. In this case we have equality in (2.2) with $c = 1$.

Our main result is

Theorem 2.1. Let V be a Wegner potential of type (α, ρ, K) for some $1 < \alpha < \rho$, $K \in \{2, 3, 4, \dots\}$. Then the random Schrödinger operator $H = -\Delta + V$ exhibits polynomial localization at high disorder or low energy. More precisely; there exist $\eta_0 > 0$ such that:

- (i) Let $H_\lambda = -\Delta + \lambda V$, $\lambda \in \mathbf{R}$. Then for $\eta > \eta_0$ we can find $\lambda(\eta) > 0$ such that for any λ with $|\lambda| > \lambda(\eta)$ H_λ has pure point spectrum with probability one and the corresponding eigenfunctions decay at least as fast as $\|x\|^{-\eta}$.
- (ii) For any $\eta > \eta_0$ we can find $E(\eta) > 0$ such that H has pure point spectrum in $(-\infty, -E(\eta)) \cup (E(\eta), \infty)$ with probability one and the corresponding eigenfunctions decay at least as fast as $\|x\|^{-\eta}$.

Theorem 2.1 will be proven by a multiscale analysis in Sect. 3. Condition (2.2) will replace independence in our proof, the price we pay is that the proof only yields polynomial localization. But condition (2.2) is satisfied by potentials that can be strongly correlated at large distances.

We give two examples of ergodic potentials satisfying the hypotheses of Theorem 2.1: Gaussian and Gibbs potentials.

We start with Gaussian potentials, which we can take to have mean zero without loss of generality.

Theorem 2.2. *Let V be a stationary Gaussian process with mean zero and covariance $C(x, y) = \mathbf{E}(V(x)V(y))$ for $x, y \in \mathbf{Z}^d$. Suppose:*

(i) C_Λ is invertible for all $\Lambda \subset \mathbf{Z}^d$ finite and

$$a \equiv \sup\{C_\Lambda^{-1}(0, 0); 0 \in \Lambda \subset \mathbf{Z}^d \text{ finite}\} < \infty, \tag{2.3}$$

where C_Λ is the matrix $\{C(x, y)\}_{x, y \in \Lambda}$.

(ii)

$$|C(x, y)| \leq \frac{b}{1 + \|x - y\|^\theta} \tag{2.4}$$

for some $\theta > 0$, $b < \infty$ and all $x, y \in \mathbf{Z}^d$.

Then there exists τ , $0 \leq \tau \leq d$, such that if $\max\left\{1, \frac{d + \tau}{\theta}\right\} < \alpha < \rho < 2^k$, k a positive integer, we have that V is a Wegner potential of type $(\alpha, \rho, 2^k)$.

In particular, $H = -\Delta + V$ exhibits polynomial localization at high disorder or low energy, in the precise sense of Theorem 2.1.

Remarks. Ergodicity follows from (2.4). Notice that

$$C_\Lambda^{-1}(0, 0) = \frac{\det C_{\Lambda \setminus \{0\}}}{\det C_\Lambda}.$$

If $d = 1$, (2.3) is equivalent to the Gaussian process being nondeterministic (e.g., [26, 27]). If the operator C on $l^2(\mathbf{Z}^d)$ with kernel given by the covariance function $C(x, y)$ is strictly positive, i.e., $C \geq wI > 0$, then (2.3) holds with $a \leq 1/w$. In this case we will see that we can take $\tau = 0$. Notice that the exponent θ in (2.4) can be arbitrarily small so the values of the Gaussian potential can be strongly correlated at large distances.

Examples of such Gaussian potentials can be given by specifying the covariance operator C . For instance let $\bar{\Delta} = \Delta - 2d$, so $\bar{\Delta}$ is the usual finite difference Laplacian on \mathbf{Z}^d . Then $C = (-\bar{\Delta} + m^2)^{-1}$ satisfies the desired hypotheses in any dimension if $m^2 > 0$. The same is true for $C = (-\bar{\Delta})^{-1}$ if $d \geq 3$. Notice that for $C = (-\bar{\Delta})^{-1}$ and $d = 3$ we have $\theta = 1$.

Theorem 2.2 will be proven in Sect. 4. The fact that V is a Wegner potential is a consequence of (2.3). To show V is of type $(\alpha, \rho, 2^k)$, we will derive a multivariable version of Nelson’s best possible hypercontractive estimate from which the result will follow.

Our second example concerns Gibbs potentials, i.e., potentials whose probability distribution is given by a Gibbs measure, which we are going to require to be completely analytical in the sense of Dobrushin and Shlosman [23, 24]. We will call such potentials completely analytical Gibbs potentials.

Theorem 2.3. *Let V be a completely analytical Gibbs potential with the “a priori” single site spin distribution absolutely continuous with compact support and bounded Radon–Nikodym derivative. Then V is a Wegner potential of type (α, K, K) for any $1 < \alpha < K$. In particular, $H = -\Delta + V$ exhibits polynomial localization at high disorder or low energy, in the precise sense of Theorem 2.1.*

The proof of Theorem 2.3 is given in Sect. 5. That V is a Wegner potential follows from the DLR equations and our assumptions on the “a priori” spin distribution, since completely analytical Gibbs fields have finite range interactions. The proof that V is of type (α, K, K) is due to Senya Shlosman [28].

Examples of completely analytical Gibbs fields are given by high-temperature Gibbs fields, or ferromagnetic Gibbs fields with a large magnetic field (arbitrary temperature) [23].

3. The Multiscale Analysis

We start by recalling a Simon-Lieb type inequality for Green’s functions that follows from the resolvent identity [18, 6, 7, 21, 9]. Let

$$\Lambda \subset \Omega \subset \mathbf{Z}^d, E \in \mathbf{R}, \varepsilon \neq 0, x \in \Lambda, y \in \Omega \setminus \Lambda.$$

We have

$$G_{\Omega}(E + i\varepsilon; x, y) = \sum_{\langle u, u' \rangle \in \partial(\Lambda, \Omega)} G_{\Lambda}(E + i\varepsilon; x, u) G_{\Omega}(E + i\varepsilon; u', y),$$

where

$$\partial(\Lambda, \Omega) = \{ \langle u, u' \rangle; u \in \Lambda, u' \in \Omega \setminus \Lambda, |u - u'| = 1 \}.$$

We write

$$\begin{aligned} \partial\Lambda &= \partial(\Lambda, \mathbf{Z}^d), \\ \partial\Lambda^+ &= \{ u' \in \mathbf{Z}^d \setminus \Lambda; \langle u, u' \rangle \in \partial\Lambda \text{ for some } u \in \Lambda \}, \\ \partial\Lambda^- &= \{ u \in \Lambda; \langle u, u' \rangle \in \partial\Lambda \text{ for some } u' \in \mathbf{Z}^d \setminus \Lambda \}, \end{aligned}$$

and

$$G_{\Lambda}(E + i\varepsilon; x, \partial) = \sum_{\langle u, u' \rangle \in \partial\Lambda} |G_{\Lambda}(E + i\varepsilon; x, u)|.$$

Thus

$$|G_{\Omega}(E + i\varepsilon; x, y)| \leq G_{\Lambda}(E + i\varepsilon; x, \partial) |G_{\Omega}(E + i\varepsilon; u'', y)| \tag{3.1}$$

for some $u'' \in \partial\Lambda^+ \cap \Omega$.

Definition. Let $\eta > 0, E \in \mathbf{R}, L > 0$. A site $x \in \mathbf{Z}^d$ is (η, E, L) -regular (for a fixed V) if

$$\tilde{G}_{\Lambda_L(x)}(E; x, \partial) \equiv \sup_{\varepsilon \neq 0} |G_{\Lambda_L(x)}(E + i\varepsilon; x, \partial)| \leq \frac{1}{L^\eta}.$$

A set $\Lambda \subset \mathbf{Z}^d$ is (η, E, L) regular if every $x \in \Lambda$ is (η, E, L) -regular.

Let $x \in \Omega \subset \mathbf{Z}^d$ with $\Lambda_L(x) \cup \partial\Lambda_L(x)^+ \subset \Omega$. If x is (η, E, L) -regular, it follows from (3.1) that for all $y \in \Omega \setminus \Lambda_L(x)$ and $\varepsilon \neq 0$,

$$|G_\Omega(E + i\varepsilon; x, y)| \leq \frac{1}{L^\eta} |G_\Omega(E + i\varepsilon; v, y)| \tag{3.2}$$

for some $v \in \partial\Lambda_L(x)^+$.

Theorem 3.1. *Let V be a Wegner potential of type $(\alpha, \rho K)$ for some $1 < \alpha < \rho$, $K \in \{2, 3, 4, \dots\}$. Suppose*

$$\mathbf{P}\{0 \text{ is } (\eta, E, L_0)\text{-regular}\} \geq 1 - \frac{1}{L_0^r} \tag{3.3}$$

for some

$$E \in \mathbf{R}, L_0 > 1, r > \frac{\alpha(d-1)K}{\rho - \alpha}, \eta(A - \alpha) > (\beta + d - 1)(K - 1 + \alpha(A + K - 1)),$$

where $A > \alpha$ is an integer and $\beta > \alpha r + (A + K - 1)(d - 1) + d$.

Let $L_{k+1} = (A + K - 1)L_k^\alpha$, $k = 0, 1, 2, \dots$. Then, there exists $\bar{L} = \bar{L}(\alpha, d, K, r, A, \beta, \eta) > 0$, such that if $L_0 > \bar{L}$, we have

$$\mathbf{P}\{0 \text{ is } (\eta, E, L_k)\text{-regular}\} \geq 1 - \frac{1}{L_k^r}$$

for all $k = 0, 1, 2, \dots$.

Theorem 3.2. *Let V be a Wegner potential, such that for some $E \in \mathbf{R}$, $\alpha > 1$, $D \geq 1$, $L_0 > 1$, $r > \alpha d$, $\eta > \alpha d$, we have*

$$\mathbf{P}\{0 \text{ is } (\eta, E, L_k)\text{-regular}\} \geq 1 - \frac{1}{L_k^r}$$

for all $k = 0, 1, 2, \dots$, where $L_{k+1} = DL_k^\alpha$. Let $0 < \bar{\eta} < \frac{\eta - \alpha d}{\alpha}$. We have, with probability one,

$$\sup_{\varepsilon \neq 0} |G(E + i\varepsilon; 0, x)| \leq \frac{C}{1 + \|x\|^{\bar{\eta}}} \tag{3.4}$$

for all $x \in \mathbf{Z}^d$ for some $C = C(V, \bar{\eta}) < \infty$.

Theorems 3.1 and 3.2 prove Theorem 2.1. Condition (3.3) is satisfied at high disorder or low energy (see, for instance, the discussion in [9]). If in (3.4) we have $\bar{\eta} > d/2$, we can apply the results of Simon and Wolff [20] to conclude pure point spectrum with probability one, and that the corresponding eigenfunctions have at least the rate of decay of the Green's function.

Proof of Theorem 3.1. Given $L > 0$, let R_L be the statement

$$\mathbf{P}\{0 \text{ is } (\eta, E, L)\text{-regular}\} \geq 1 - \frac{1}{L^r}.$$

Theorem 3.1 follows from

Lemma 3.3. *Let $V, \alpha, K, \rho, r, \eta, A, \beta$ be as in Theorem 3.1. Suppose R_l holds, and let $L = (A + K - 1)l^\alpha$. Then R_L holds if l is large enough.*

Proof. Let $D = A + K - 1$, so $L = Dl^\alpha$. We set $\Lambda_j = \Lambda_{jl^\alpha}(0)$, $j = 1, \dots, D$. We also set $\Lambda = \Lambda_D = \Lambda_L(0)$.

Let $y \in \partial\Lambda^-$. We will estimate $G_\Lambda(E + i\varepsilon; 0, y)$ by applying (3.1) repeatedly. We have

$$\begin{aligned} |G_\Lambda(E + i\varepsilon; 0, y)| &\leq G_{\Lambda_l(0)}(E + i\varepsilon; 0, \partial) |G_\Lambda(E + i\varepsilon; v_0, y)| \\ &\leq G_{\Lambda_l(0)}(E + i\varepsilon; 0, \partial) G_{\Lambda_1}(E + i\varepsilon; v_0, \partial) |G_\Lambda(E + i\varepsilon; u_1, y)| \\ &\leq G_{\Lambda_l(0)}(E + i\varepsilon; v_0, \partial) G_{\Lambda_1}(E + i\varepsilon; v_0, \partial) G_{\Lambda_l(u_1)}(E + i\varepsilon; u_1, \partial) \\ &\quad \cdot G_{\Lambda_2}(E + i\varepsilon; v_1, \partial) G_{\Lambda_l(u_2)}(E + i\varepsilon; u_2, \partial) \\ &\quad \dots G_{\Lambda_{D-1}}(E + i\varepsilon; v_{D-1}, \partial) G_{\Lambda_l(u_{D-1})}(E + i\varepsilon; u_{D-1}, \partial) \\ &\quad \cdot |G_\Lambda(E + i\varepsilon; v_{D-1}, y)| \end{aligned}$$

for some $v_0 \in \partial\Lambda_l(0)^+$, $u_j \in \partial\Lambda_j^+$, $v_j \in \partial\Lambda_l(u_j)^+$, $j = 1, \dots, D - 1$.

We will call a cube $\Lambda_l(x)$ (E, β) -non resonant ($(E, \beta) - NR$) if

$$d(E, \sigma(H_{\Lambda_l(x)})) > \frac{1}{l^\beta}.$$

Notice also that $|\partial\Lambda_l(x)^-| \leq |\partial\Lambda_l(x)^+| \leq sl^{d-1}$ for some $s = s(d) < \infty$. Thus, if $\Lambda_l(x)$ is $(E, \beta) - NR$, we have

$$G_{\Lambda_l(x)}(E + i\varepsilon; u, \partial) \leq sl^{\beta+d-1}. \quad (3.6)$$

Let \mathcal{E} be the event defined by

- (i) All Λ_j , $j = 1, \dots, D$, and $\Lambda_l(x)$, $x \in \Lambda$, are $(E, \beta) - NR$.
- (ii) There exist $j_1 < j_2 < \dots < j_A \in \{0, 1, \dots, D - 1\}$, such that $\partial\Lambda_{j_1}^+, \dots, \partial\Lambda_{j_A}^+$ are (η, E, l) -regular, where $\partial\Lambda_0^+ = \{0\}$.

Let us now assume that the event \mathcal{E} holds. It follows from (3.5) and (3.6) that

$$\begin{aligned} G_\Lambda(E + i\varepsilon; 0, \partial) &\leq \frac{1}{l^{A\eta}} (sl^{\beta+d-1})^{K-1} (sL^{\beta+d-1})^D \\ &= s^{K-1+D} D^D l^{-[A\eta - (\beta+d-1)(K-1+\alpha D)]} \leq \frac{1}{L^\eta} = \frac{1}{D^\eta l^{\alpha\eta}} \end{aligned}$$

for l sufficiently large, since

$$\eta(A - \alpha) > (\beta + d - 1)(K - 1 + \alpha(A + K)).$$

Thus, to conclude the proof of the lemma we need only to show that

$$\mathbf{P}(\mathcal{E}) \geq 1 - \frac{1}{L^\eta}.$$

Since V is a Wegner potential, it follows from (2.1) that

$$\mathbf{P}\{\Lambda_l(x) \text{ is } (E, \beta) - NR\} \geq 1 - \frac{\sigma}{l^{\beta-d}}$$

for some $\sigma > 0$, depending only on V . Thus

$$\mathbf{P}\{\text{(i) does not hold}\} \leq \frac{D\sigma}{l^{\alpha(\beta-d)}} + \frac{(L+1)^d\sigma}{l^{\beta-d}} \leq \frac{1}{2L^r}$$

for l sufficiently large, since $\alpha r < \beta - (\alpha + 1)d$.

So it only remains to show that

$$\mathbf{P}\{\text{(ii) does not hold}\} \leq \frac{1}{2L^r}.$$

We have

$$\begin{aligned} \mathbf{P}\{\text{(ii) does not hold}\} &= \mathbf{P}\{\text{there exist } u_{j_i} \in \partial\Lambda_{j_i}^+ \text{ not } (\eta, E, l)\text{-regular} \\ &\quad i = 1, \dots, K, j_1 < j_2 < \dots < j_K \in \{0, 1, \dots, D-1\}\} \\ &\leq \sum_{j_1 < j_2 < \dots < j_K \in \{0, \dots, D-1\}} \sum_{u_i \in \partial\Lambda_{j_i}, i=1, \dots, K} \mathbf{P}\{u_1, \dots, u_K \text{ not } (\eta, E, L)\text{-regular}\}. \end{aligned}$$

Since $\|u_{i'} - u_i\| \geq \frac{l^\alpha}{2}$ for $i \neq i'$, we can use (2.2) for l sufficiently large to conclude

$$\mathbf{P}\{\text{(ii) does not hold}\} \leq \binom{D}{K} (sL^{d-1})^K \frac{C}{l^{\rho r}} \leq \frac{1}{2L^r}$$

for l sufficiently large since

$$r > \frac{\alpha(d-1)K}{\rho - \alpha}.$$

This completes the proof of Lemma 3.3 and hence of Theorem 3.1

Proof of Theorem 3.2. Let $0 < \bar{\eta} < \frac{\eta - \alpha d}{\alpha}$, we choose $\beta > d$ by $\bar{\eta} = \frac{\eta - \alpha\beta}{\alpha}$.

Let us define the events

$$A_{k+1} = \{x \text{ is } (\eta, E, L_k)\text{-regular for all } x \in \Lambda_{2L_{k+1}}(0) \text{ and } \Lambda_{L_{k+1}}(0) \text{ is } (E, \beta) - NR\}.$$

By the hypotheses of Theorem 3.2, we have

$$\mathbf{P}(A_{k+1}^c) \leq \frac{(2L_{k+1} + 1)^d}{L_k^r} + \frac{\sigma}{L_{k+1}^{\beta-d}},$$

so $\sum_{k=0}^\infty \mathbf{P}(A_{k+1}^c) < \infty$, since $r > \alpha d, \beta > d$. It follows from the Borel-Cantelli Lemma that, with probability one, we can find $\bar{k} = \bar{k}(V) < \infty$ such that A_{k+1} holds for all $k \geq \bar{k}$.

So let us fix a potential V and $\bar{k} < \infty$ such that A_{k+1} holds for $k \geq \bar{k}$. Given $x \in \mathbf{Z}^d$, we pick $\tilde{k} = \tilde{k}(x)$ in the following way:

- (i) if $\|x\| \leq \frac{1}{2}L_{\tilde{k}+1}$, pick $\tilde{k} = \bar{k}$.
- (ii) if $\|x\| > \frac{1}{2}L_{\tilde{k}+1}$, pick \tilde{k} such that $\frac{1}{2}L_{\tilde{k}} < \|x\| \leq \frac{1}{2}L_{\tilde{k}+1}$.

We let $G_k(E + i\varepsilon) = G_{\Lambda_k(0)}(E + i\varepsilon)$, and let Γ_k be the operator defined by

$$\Gamma_k(u, u') = \begin{cases} 1 & \text{if either } \langle u, u' \rangle \text{ or } \langle u', u \rangle \in \partial\Lambda_k(0) \\ 0 & \text{otherwise} \end{cases}.$$

It follows from the resolvent identity, applied repeatedly, that (we omit $E + i\varepsilon$)

$$G = G_{\tilde{k}+1} + G_{\tilde{k}+1} \Gamma_{\tilde{k}+1} G_{\tilde{k}+2} + \dots + G_{\tilde{k}+1} \Gamma_{\tilde{k}+1} G_{\tilde{k}+2} \dots \Gamma_{\tilde{k}+1+l} G_{\tilde{k}+2+l} + \dots$$

If $\tilde{k} = \bar{k}$, we just take

$$|G_{\tilde{k}+1}(E + i\varepsilon; 0, x)| \leq L_{\tilde{k}+1}^\beta.$$

If $\tilde{k} > \bar{k}$, we use (3.2) m -times where

$$m = \left\lceil \frac{\|x\|}{\left\lfloor \frac{L_{\tilde{k}}}{2} \right\rfloor + 1} \right\rceil$$

(here $\lceil t \rceil$ means the largest integer $\leq t$), to obtain

$$|G_{\tilde{k}+1}(E + i\varepsilon; 0, x)| \leq L_{\tilde{k}}^{-\eta m} L_{\tilde{k}+1}^\beta \leq D^\beta L_{\tilde{k}}^{-(\eta - \beta\alpha)} \leq D^\beta \left(\frac{D}{2}\right)^{(\eta - \beta\alpha)/\alpha} \frac{1}{\|x\|^{(\eta - \beta\alpha)/\alpha}},$$

since $m \geq 1$, $\|x\| \leq \frac{1}{2} L_{\tilde{k}+1}$.

Similarly

$$|G_{\tilde{k}+1} \Gamma_{\tilde{k}+1} \dots \Gamma_{\tilde{k}+1+l} G_{\tilde{k}+2+l}(0, x)| \leq \frac{1}{L_{\tilde{k}+1}^\eta} \frac{1}{L_{\tilde{k}+1}^{\eta \gamma_{\tilde{k}+1}}} \dots \frac{1}{L_{\tilde{k}+l}^{\eta \gamma_{\tilde{k}+l}}} L_{\tilde{k}+2+l}^\beta,$$

where

$$\gamma_k = \left\lceil \frac{\left\lfloor \frac{L_{k+1}}{2} \right\rfloor - \left\lfloor \frac{L_k}{2} \right\rfloor - 1}{\left\lfloor \frac{L_k}{2} \right\rfloor + 1} \right\rceil,$$

and thus

$$\leq \frac{C}{L_{\tilde{k}+1}^{(\eta - \alpha\beta)(l+1)}} \leq \frac{C'}{\|x\|^{((\eta - \alpha\beta)/\alpha)(l+1)}}$$

for some constants C, C'' independent of x .

Thus, we have

$$|G(E + i\varepsilon; \partial, x)| \leq L_{\tilde{k}+1}^\beta + \frac{C''}{\|x\|^{(\eta - \alpha\beta)/\alpha}} \tag{3.7}$$

if $\|x\| \leq \frac{L_{\tilde{k}+1}}{2}$, some constant $C'' < \infty$ and

$$|G(E + i\varepsilon; 0, x)| \leq \frac{C''}{\|x\|^{(\eta - \alpha\beta)/\alpha}} \tag{3.8}$$

otherwise.

Recall $\bar{\eta} = \frac{\eta - \alpha\beta}{\alpha}$, it follows from (3.7) and (3.8) that

$$|G(E + i\varepsilon; 0, x)| \leq \frac{C}{1 + \|x\|^{\bar{\eta}}}$$

for all $x \in \mathbb{Z}^d$ for some constant $C < \infty$.

This proves Theorem 3.2.

4. Gaussian Potentials

The key ingredient in the proof of Theorem 2.2 is Nelson’s best possible hypercontractive estimate [30–35]. It can be reformulated as follows [35]:

Let X, Y be jointly Gaussian random variables, both with zero mean. Then, for any measurable functions f and g of a real variable, we have

$$|\mathbf{E}(f(X)g(Y))| \leq \|f(X)\|_p \|g(Y)\|_q$$

if

$$(p-1)(q-1) \geq \frac{\mathbf{E}(XY)^2}{\mathbf{E}(X^2)\mathbf{E}(Y^2)}.$$

If particular, we can take

$$p = q = 1 + \frac{|\mathbf{E}(XY)|}{(\mathbf{E}(X^2)\mathbf{E}(Y^2))^{1/2}}.$$

We will need a multivariable version of this result.

Lemma. *Let $X_1, X_2, \dots, X_n, Y_1, Y_2, \dots, Y_n$ be jointly Gaussian random variables with mean zero. Let*

$$\mathbf{X} = (X_1, \dots, X_n), \mathbf{Y} = (Y_1, \dots, Y_n), P_1 = \mathbf{E}(\mathbf{X}'\mathbf{X}), P_2 = \mathbf{E}(\mathbf{Y}'\mathbf{Y}), Q = \mathbf{E}(\mathbf{X}'\mathbf{Y}).$$

Then, for any measurable functions f, g on \mathbf{R}^n , we have

$$|\mathbf{E}(f(\mathbf{X})g(\mathbf{Y}))| \leq \|f(\mathbf{X})\|_p \|g(\mathbf{Y})\|_q$$

if

$$(p-1)(q-1) \geq \|P_1^{-1/2} Q P_2^{-1/2}\|^2,$$

where we used the operator norm for matrices. In particular, we can take

$$p = q = 1 + \|P_1^{-1/2} Q P_2^{-1/2}\|.$$

Proof. By a limiting procedure we can take the covariance matrix

$$C = \begin{pmatrix} P_1 & Q \\ Q' & P_2 \end{pmatrix}$$

to be strictly positive definite. In this case we can also take

$$C = \begin{pmatrix} I & Q \\ Q' & I \end{pmatrix}$$

without loss of generality.

Since $C > 0$, $\|Q\| < 1$, and we have

$$C^{-1} = \begin{bmatrix} (I - QQ')^{-1} & -(I - QQ')^{-1}Q \\ -(I - Q'Q)^{-1}Q' & (I - Q'Q)^{-1} \end{bmatrix}.$$

Let us consider the operator on $\mathcal{X} = L_2(\mathbf{R}^n, (2\pi)^{-n/2} e^{-(1/2)(\mathbf{x}^2)} d\mathbf{x})$ given by the

multivariable Mehler’s formula [34; (I.34)]:

$$(\Gamma(Q)f)(\mathbf{x}) = (2\pi)^{-n/2} [\det(1 - Q^t Q)]^{1/2} \int e^{(-1/2)\langle(Q^t \mathbf{x} - \mathbf{y}), (1 - Q^t Q)^{-1}(Q^t \mathbf{x} - \mathbf{y})\rangle} f(\mathbf{y}) d\mathbf{y}$$

It is not hard to check that

$$\mathbf{E}(f(X)g(Y)) = \langle f, \Gamma(Q)g \rangle_{\mathcal{X}}.$$

The lemma now follows from Nelson’s best possible hypercontractive estimate (see [34; Theorem I.17]).

We are now ready to prove Theorem 2.2. We first show that it follows from (2.3) that V is a Wegner potential. Let $0 \in J \subset \mathbf{Z}^d$ finite, an explicit computation shows that the conditional probability distribution of $V(0)$ given $V_{J \setminus \{0\}} = \{V(x); x \in J \setminus \{0\}\}$ is absolutely continuous with density

$$f_J(V(0) | V_{J \setminus \{0\}}) = \sqrt{\frac{a_J}{2\pi}} e^{(-1/2)a_J(V(0) + (1/a_J) \sum_{x \in J \setminus \{0\}} c_J^{-1}(0, x)V(x))^2},$$

where $a_J = C_J^{-1}(0, 0)$. By (2.3) we have

$$\|f_J(V(0) | V_{J \setminus \{0\}})\|_{\infty} \leq \sqrt{\frac{a}{2\pi}} < \infty.$$

We can conclude that the conditional probability distribution of $V(0)$ given $V(0)^{\perp}$ is absolutely continuous with a density bounded by $\sqrt{\frac{a}{2\pi}}$.

Now let $0 \leq \tau \leq d$ be such that $\|C_J^{-1}\| \leq c_1 |J|^{\tau/d}$ for all $J \subset \mathbf{Z}^d$ finite, some $c_1 < \infty$. This can always be done since

$$\|C_J^{-1}\| \leq \text{Tr} C_J^{-1} \leq a |J|$$

by (2.3).

So let us pick α, ρ , and an integer k such that $\max\left\{1, \frac{d+r}{\theta}\right\} < \alpha < \rho < 2^k$. It will follow from the following lemma that V is of type $(\alpha, \rho, 2^k)$.

Lemma 4.2. *Let $L > 2l > 0$, $x_i \in \mathbf{Z}^d$, $i = 1, \dots, 2^k$, with $\|x_i - x_j\| \geq L$ for $i \neq j$, and let A_i be events on $\Lambda_l(x_i)$, $i = 1, \dots, 2^k$. Then*

$$\mathbf{P}\left(\bigcap_{i=1}^{2^k} A_i\right) \leq \prod_{i=1}^{2^k} \mathbf{P}(A_i)^{(1/p)^k}$$

with

$$p \geq 1 + bc_1 4^{k-1} \frac{(l+1)^{r+d}}{1+(L-l)^\theta}.$$

In particular, if A is an event on $\Lambda_l(0)$, we have

$$\mathbf{P}\left(\bigcap_{i=1}^{2^k} A(x_i)\right) \leq \mathbf{P}(A)^{(2/p)^k}.$$

Proof. Let $I_1 = \{1, \dots, 2^{k-1}\}$, $I_2 = \{2^{k-1} + 1, \dots, 2^k\}$. We apply Lemma 4.1 with

$$\begin{aligned} \mathbf{X} &= \{V(y); y \in \Lambda_l(x_i), i \in I_1\}, \\ \mathbf{Y} &= \{V(y); y \in \Lambda_l(x_i), i \in I_2\}. \end{aligned}$$

Then

$$\mathbf{P}\left(\bigcap_{i=1}^{2^k} A_i\right) \leq \mathbf{P}\left(\bigcap_{i \in I_1} A_i\right)^{1/p} \mathbf{P}\left(\bigcap_{i \in I_2} A_i\right)^{1/p},$$

where

$$p \geq 1 + \|P_1^{-1/2} Q P_2^{-1/2}\|,$$

P_1, P_2, Q being defined as in Lemma 4.1. But

$$\|P_j^{-1}\| \leq c_1 (2^{k-1} (l+1)^d)^{\tau/d} = c_1 2^{(k-1)(\tau/d)} (l+1)^\tau,$$

and we can estimate $\|Q\|$ by its Hilbert–Schmidt norm as follows using (2.4):

$$\|Q\| \leq 2^{(k-1)} (l+1)^d \frac{b}{1 + (L-l)^\theta}.$$

Thus

$$p \geq 1 + bc_1 4^{k-1} \frac{(l+1)^{d+\tau}}{1 + (L-l)^\theta}. \quad (4.1)$$

Repeating the above procedure for a total of k -times, which we can do using always p given by (4.1), we get the lemma.

This finishes the proof of Theorem 2.4.

5. Gibbs Potentials

In this section we present a proof of Theorem 2.3, due to S. Shlosman [28].

So let V be as in Theorem 2.3. Since the DLR equations and our assumptions immediately imply that V is a Wegner potential, we have only to prove that V is of type (α, K, K) for any $1 < \alpha < K$.

Dobrushin and Shlosman [23, 24] gave several equivalent conditions that characterize completely analytical Gibbs measures. We will show that their condition III_d implies that V is of type (α, K, K) for any $1 < \alpha < K$.

First, some notation. We will denote by r the range of the interaction (by definition completely analytical interactions have finite range, i.e., $r < \infty$). Given $\Lambda \subset \mathbf{Z}^d$, we define

$$\begin{aligned} \partial\Lambda &= \{y \in \mathbf{Z}^d \setminus \Lambda; \text{dist}(y, \Lambda) \leq r\}, \\ V_\Lambda &= \{V(x); x \in \Lambda\}, \end{aligned}$$

$\mathbf{P}_\Lambda(\cdot | V_{\Lambda^c})$ = conditional probability distribution of V_Λ given V_{Λ^c} , where $\Lambda^c = \mathbf{Z}^d \setminus \Lambda$.

We can now state Dobrushin and Shlosman's condition III_d:

There exist $C < \infty$, $\gamma > 0$, such that for all $\Lambda \subset \Omega \subset \mathbf{Z}^d$ finite, any $y \in \partial\Omega$, and

any choices V_1, V_2 of V_{Ω^c} , such that $V_1(x) = V_2(x)$ for $x \in \Omega^c \setminus \{y\}$, we have

$$\left| \frac{\mathbf{P}_{\Omega}(A|V_1)}{\mathbf{P}_{\Omega}(A|V_2)} - 1 \right| \leq C e^{-\gamma \text{dist}(y, A)} \tag{5.1}$$

for any event in A in Λ .

Lemma 5.1. *For all $\Lambda \subset \Omega \subset \mathbf{Z}^d$ finite, any choices V_1, V_2 of V_{Ω^c} we have*

$$\left| \frac{\mathbf{P}_{\Omega}(A|V_1)}{\mathbf{P}_{\Omega}(A|V_2)} - 1 \right| \leq (1 + C e^{-\gamma \text{dist}(\Lambda, \partial\Omega)})^{|\partial\Omega|} - 1. \tag{5.2}$$

Proof. Without loss of generality we can always assume $V_1(x) = V_2(x)$ for $x \in \Omega^c \setminus \partial\Omega$. For such V_1, V_2 we can always find values W_0, W_1, \dots, W_k of V_{Ω^c} , such that $W_0 = V_1, W_k = V_2, k \leq |\partial\Omega|$, and for each $i = 1, \dots, k$ we can find $y_i \in \partial\Omega$ such that $W_{i-1}(x) = W_i(x)$ for $x \neq y_i$. We have

$$\frac{\mathbf{P}_{\Omega}(A|V_1)}{\mathbf{P}_{\Omega}(A|V_2)} = \prod_{i=1}^k \frac{\mathbf{P}_{\Omega}(A|W_{i-1})}{\mathbf{P}_{\Omega}(A|W_i)}.$$

Equation (5.1) applies to each factor in the right-hand-side so we can conclude that for each $i = 1, \dots, k$, we have

$$1 - K e^{-\gamma \text{dist}(\Lambda, \partial\Omega)} \leq \frac{\mathbf{P}_{\Omega}(A|W_{i-1})}{\mathbf{P}_{\Omega}(A|W_i)} \leq 1 + K e^{-\gamma \text{dist}(\Lambda, \partial\Omega)},$$

and hence

$$-(1 - (1 - K e^{-\gamma \text{dist}(\Lambda, \partial\Omega)})^{|\partial\Omega|}) \leq \frac{\mathbf{P}_{\Omega}(A|V_1)}{\mathbf{P}_{\Omega}(A|V_2)} - 1 \leq (1 + K e^{-\gamma \text{dist}(\Lambda, \partial\Omega)})^{|\partial\Omega|} - 1,$$

so (5.2) follows.

If $\text{dist}(\Lambda, \partial\Omega)$ is sufficiently large, we have

$$|\partial\Omega| \leq C_1 [\text{dist}(\Lambda, \partial\Omega)]^{d-1}$$

for some fixed constant $C_1 < \infty$. In this case, it follows from (5.2) that we can find $C' < \infty, \gamma' > 0$, such that

$$\left| \frac{\mathbf{P}_{\Omega}(A|V_1)}{\mathbf{P}_{\Omega}(A|V_2)} - 1 \right| \leq C' e^{-\gamma' \text{dist}(\Lambda, \partial\Omega)}. \tag{5.3}$$

Now, let B be an event on Ω^c . It follows from (5.3) that

$$\left| \frac{\mathbf{P}_{\Omega}(A|B)}{\mathbf{P}(A)} - 1 \right| \leq C' e^{-\gamma' \text{dist}(\Lambda, \partial\Omega)}. \tag{5.4}$$

We are now ready to show V is of type (α, K, K) for $1 < \alpha < K$. For given events A_1, A_2, \dots, A_K , we always have

$$\mathbf{P}\left(\bigcap_{i=1}^K A_i\right) = \mathbf{P}(A_1)\mathbf{P}(A_2|A_1)\mathbf{P}(A_3|A_2 \cap A_1)\dots\mathbf{P}\left(A_K \mid \bigcap_{i=1}^{K-1} A_i\right). \tag{5.5}$$

So let $1 < l < L, x_1, \dots, x_K \in \mathbf{Z}^d$, with $\|x_i - x_j\| > L$ for all $i \neq j$. Assume $L - l$ is sufficiently large. Then, if A_i is an event on $\Lambda_l(x_i), i = 1, \dots, K$, we have from (5.4)

that, for each $i = 1, \dots, K$,

$$\frac{\mathbf{P}\left(A_i \mid \bigcap_{j=1}^{i-1} A_j\right)}{\mathbf{P}(A_i)} - 1 \leq C'e^{-\gamma(L-l)},$$

so

$$\mathbf{P}\left(A_i \mid \bigcap_{j=1}^{i-1} A_j\right) \leq \mathbf{P}(A_i)(1 + C'e^{-\gamma(L-l)}).$$

Thus, it follows from (5.5) that

$$\mathbf{P}\left(\bigcap_{i=1}^K A_i\right) \leq \prod_{i=1}^K \mathbf{P}(A_i)(1 + C'e^{-\gamma(L-l)})^K \leq 2^K \prod_{i=1}^K \mathbf{P}(A_i).$$

This proves Theorem 2.3.

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