# Generalized eigenvalue-counting estimates for some random acoustic operators 

Yoshihiko Kitagaki


#### Abstract

For some discrete random acoustic operators, we prove Wegner estimates. These estimates are applied to show some regularity of the integrated density of states. Moreover, we prove the generalized eigenvalue-counting estimates by using Combes, Germinet, and Klein's method. As an application, the multiplicity of the eigenvalues in some interval where the Anderson localization occurs is proven to be finite. For certain models, Poisson statistics for eigenvalues and Lifshitz tails are also studied.


## 1. Introduction

In this paper, we consider the following random acoustic operator in the discrete setting:

$$
\begin{equation*}
A_{\omega}:=\nabla^{*}\left(\rho_{0}+\rho_{\omega}\right) \nabla=A_{0}+\sum_{\gamma \in \mathbf{Z}^{d}} \omega_{\gamma} \nabla^{*} \Pi_{\gamma} \nabla \tag{1.1}
\end{equation*}
$$

where $\nabla$ is defined by the discrete analogue of the partial derivative, $\nabla^{*}$ is its adjoint, $\rho_{0}$ is a periodic term, and $\rho_{\omega}:=\sum_{\gamma \in \mathbf{Z}^{d}} \omega_{\gamma} \Pi_{\gamma}$ is an Anderson-type random perturbation. The operator $A_{\omega}$ is defined on $\ell^{2}\left(\mathbf{Z}^{d}\right)$ (for details, see Section 2). In that section, we prove the Wegner-type estimate and the generalized eigenvalue-counting estimate for this operator. We apply the Wegner estimate to obtain some regularity for the integrated density of states (IDS). Moreover, we use the generalized eigenvalue-counting estimate to show that the multiplicity of the eigenvalues of (1.1) in some interval where the Anderson localization occurs is finite.

Recently, Combes, Germinet, and Klein [9] have given a nice proof of the generalized eigenvalue-counting estimate for discrete random Schrödinger operators, including the Anderson [4] model. This operator is defined by $H_{\omega}=-\Delta+V_{\omega}$, where $-\Delta$ is the discrete Laplacian and $V_{\omega}:=\sum_{\gamma \in \mathbf{Z}^{d}} \omega_{\gamma} \Pi_{\gamma}$ is the random potential. Their eigenvalue-counting estimates give an upper bound of the probability that the restricted random operator $H_{\omega}^{\Lambda}$ has $n$ eigenvalues in an interval $I$. This is an extension of the Wegner and Minami estimates. In this article, we use their ideas to study random acoustic operators under the discrete setting.

The disorder associated with $H_{\omega}$ is a type of diagonal disorder, while our acoustic operator $A_{\omega}$ has a type of off-diagonal disorder. The random potential $V_{\omega}$ in $H_{\omega}$ is the sum of independent random rank-one operators; that is, the rank of $\Pi_{\gamma}$ is just one. However, one difficulty in $A_{\omega}$ is that the rank of $\nabla^{*} \Pi_{\gamma} \nabla$ in (1.1) is more than two if $d \geq 2$. Therefore, we cannot use the method in [9] directly. To treat our random operator $A_{\omega}$, by applying the general spectral averaging results in [11], [12], and [31], we generalize the method for the rank-one case to the case when the rank is more than two.

It is well known that there are many works on the Anderson model $H_{\omega}$. Wegner [43] estimated the expectation of the number of eigenvalues of $H_{\omega}^{\Lambda}$ in $I$. This gives an upper bound of the probability that $H_{\omega}^{\Lambda}$ has at least one eigenvalue in $I$. First, this was used to show Lipschitz continuity of the IDS for $H_{\omega}$ or to give a bound of the density of states, the derivative of the IDS. Next, this was also used as a key estimate to prove the Anderson localization rigorously. This method is called the multiscale analysis (see, e.g., [8], [24], [28], [40]).

Figotin and Klein [21] proved the Wegner estimate for $A_{\omega}$ and showed the localization via the multiscale analysis when the probability distribution $\mu_{\gamma}$ of $\omega_{\gamma}$ has a bounded density. Our Wegner estimate is an extension to the case when $\mu_{\gamma}$ is general.

Minami [34] estimated the probability that $H_{\omega}^{\Lambda}$ has at least two eigenvalues in $I$ and proved that the properly rescaled eigenvalues of $H_{\omega}^{\Lambda}$ behave according to a Poisson point process. This was first shown by Molchanov [33] for 1-dimensional continuous random Schrödinger operators. It is natural that we consider the application of our generalized eigenvalue-counting estimates. However, it is uncertain whether our estimate for $A_{\omega}$ is useful to study some behavior for eigenvalues of $A_{\omega}$ since our result is a weaker estimate if $d \geq 2$ (see Section 2 for details).

Klein and Molchanov [29] used the Minami estimate as an important probabilistic estimate to study the multiplicity of eigenvalues of $H_{\omega}$. Moreover, the Minami estimate was generalized by [5], [9], and [25] to the probability that $H_{\omega}^{\Lambda}$ has at least $n$ eigenvalues in $I$ for all $n \in \mathbf{N}$.

As relating works, Faris [18]-[20] studied a simpler random acoustic model:

$$
\begin{equation*}
\Gamma_{\omega}:=\frac{1}{\sqrt{\rho_{\omega}}}(-\Delta) \frac{1}{\sqrt{\rho_{\omega}}} \tag{1.2}
\end{equation*}
$$

To show the Anderson localization, Faris gave the Wegner estimate for this operator. In this article, this random operator is also discussed (see Section 3 for details).

Our random operator $A_{\omega}$ in this article has the continuous version, which is formally defined by $-\nabla\left(\rho_{0}+\rho_{\omega}\right) \nabla=A_{0}+\sum_{\gamma \in \mathbf{Z}^{d}} \omega_{\gamma}\left(-\nabla u_{\gamma} \nabla\right)$ on $L^{2}\left(\mathbf{R}^{d}\right)$, where $\rho_{0}(x)$ is a $\mathbf{Z}^{d}$-periodic function and $\rho_{\omega}(x):=\sum_{\gamma \in \mathbf{Z}^{d}} \omega_{\gamma} u_{\gamma}(x)$ is an Anderson-type random perturbation. In [22] Figotin and Klein gave the Wegner estimate and applied this to prove the localization of acoustic waves. Unfortunately, the techniques in this article do not work for this continuous model because of the following. Although the operator $-\nabla u_{\gamma} \nabla$ is nonnegative as in the discrete model $A_{\omega}$,
this is not a bounded operator. Thus, we cannot use the general spectral averaging results in [11], [12], and [31] (see Proposition 2.1 in Section 2).

In [21], the method for the acoustic model was also applied to random Maxwell operators $M_{\omega}$, formally defined by $\nabla^{*} \times\left(1 / \varepsilon_{\omega}\right) \nabla \times$ on $\ell^{2}\left(\mathbf{Z}^{3} ; \mathbf{C}^{3}\right)$, where $\varepsilon_{\omega}$ is the position-dependent random dielectric constant. The Wegner estimate for $M_{\omega}$ was given to prove the localization of electromagnetic waves (see [21] for the discrete case and [13], [23] for the continuous case). In this article, we only treat random acoustic models; however, the method in this article works also for $M_{\omega}$.

This article is organized as follows. In Section 2, we first prove the Wegner estimate for $A_{\omega}$ by using the general spectral averaging result. Next, we extend the method of the article by Combes, Germinet, and Klein [9] for the operator $A_{\omega}$ and prove the generalized eigenvalue-counting estimate. In Section 3, we study the random operator $\Gamma_{\omega}$, including (1.2) as another acoustic model. To consider the same problem for $\Gamma_{\omega}$, we introduce an auxiliary random Schrödinger operator. In Section 3, we also state Poisson statistics for $\Gamma_{\omega}$. In Appendix A, we present a simple proof of Lifshitz tails for $\Gamma_{\omega}$. In Appendix B, we discuss some relations between $A_{\omega}$ and $\Gamma_{\omega}$, which are especially studied in the one-dimensional case. In Appendix C, we verify a condition used in Section 3.

## 2. Main results

### 2.1. Generalized eigenvalue-counting estimate for $A_{\omega}$

In general, acoustic operators are defined by

$$
\begin{equation*}
A_{a b}:=\frac{1}{\sqrt{a}} \nabla^{*} \frac{1}{b} \nabla \frac{1}{\sqrt{a}}=\sum_{j=1}^{d} \frac{1}{\sqrt{a}} \partial_{j}^{*} \frac{1}{b} \partial_{j} \frac{1}{\sqrt{a}} \tag{2.1}
\end{equation*}
$$

acting on the Hilbert space $\ell^{2}\left(\mathbf{Z}^{d}\right)$. The multiplicative operators $1 / \sqrt{a}$ and $1 / b$ are defined by

$$
\begin{equation*}
\frac{1}{\sqrt{a}}:=\sum_{\gamma \in \mathbf{Z}^{d}} \frac{1}{\sqrt{a_{\gamma}}} \Pi_{\gamma} \quad \text { and } \quad \frac{1}{b}:=\sum_{\gamma \in \mathbf{Z}^{d}} \frac{1}{b_{\gamma}} \Pi_{\gamma} \tag{2.2}
\end{equation*}
$$

for some positive real numbers $\left\{a_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}}$ and $\left\{b_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}}$, where $\Pi_{\gamma}$ is an orthogonal projector to the Kronecker delta function $\delta_{\gamma}(x) \in \ell^{2}\left(\mathbf{Z}^{d}\right)$, that is, $\Pi_{\gamma}:=$ $\left|\delta_{\gamma}\right\rangle\left\langle\delta_{\gamma}\right|=\left\langle\delta_{\gamma}, \cdot\right\rangle \delta_{\gamma}$. The discrete analogue of the partial derivative $\partial_{j}$ is defined by $\left(\partial_{j} f\right)(x):=f(x)-f\left(x-e_{j}\right)$ for $f(x) \in \ell^{2}\left(\mathbf{Z}^{d}\right)$ and $j=1, \ldots, d$, where $\left\{e_{j}\right\}_{j=1}^{d}$ are the standard basis vectors in the lattice $\mathbf{Z}^{d}$. The adjoint of $\partial_{j}$ is given by $\left(\partial_{j}^{*} f\right)(x)=f(x)-f\left(x+e_{j}\right)$. The acoustic operator $A_{a b}$ appears in the discrete wave equation

$$
\frac{\partial^{2} f(x, t)}{\partial t^{2}}+A_{a b} f(x, t)=0 .
$$

We note that the coefficient $a=\left\{a_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}}$ corresponds to the local propagation speed and $b=\left\{b_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}}$ corresponds to the mass density. Refer to [13], [18]-[22] for the physical interpretation.

In this section, we consider the operator defined by

$$
\begin{equation*}
A_{\omega}:=\nabla^{*}\left(\rho_{0}+\rho_{\omega}\right) \nabla=\sum_{j=1}^{d} \partial_{j}^{*}\left(\rho_{0}+\rho_{\omega}\right) \partial_{j} \tag{2.3}
\end{equation*}
$$

on $\ell^{2}\left(\mathbf{Z}^{d}\right)$. We set $\rho_{0}$ is a periodic term defined by $\rho_{0}:=\sum_{\gamma \in \mathbf{Z}^{d}} q_{\gamma} \Pi_{\gamma}$, where $\gamma \mapsto q_{\gamma}$ is a nonnegative bounded periodic function on $\mathbf{Z}^{d}$. More precisely, assume that there exists $p=\left(p_{1}, \ldots, p_{d}\right) \in \mathbf{Z}^{d}$ such that $q_{\gamma}=q_{\gamma+n p}$ for any $\gamma, n=\left(n_{1}, \ldots, n_{d}\right) \in \mathbf{Z}^{d}$, where $n p:=\left(n_{1} p_{1}, \ldots, n_{d} p_{d}\right)$. Let $M_{0}$ be the upper bound of $\left\{q_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}}$, and set $A_{0}:=\nabla^{*} \rho_{0} \nabla$. We next set $\rho_{\omega}$ as a positive Andersontype random perturbation defined by $\rho_{\omega}:=\sum_{\gamma \in \mathbf{Z}^{d}} \omega_{\gamma} \Pi_{\gamma}$, where $\omega=\left\{\omega_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}}$ are nontrivial random variables taking values in the finite interval $[m, M]$ for some $0<m<M<\infty$. For $\gamma \in \mathbf{Z}^{d}$, let $\mu_{\gamma}$ be the probability distribution of the random variable $\omega_{\gamma}$. We always assume that $\left\{\omega_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}}$ are independent. $A_{\omega}$ is a nonnegative, bounded self-adjoint operator. If we assume that

$$
\begin{equation*}
\left\{\omega_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}} \text { are identically distributed } \tag{2.4}
\end{equation*}
$$

with the common probability distribution $\mu$, then $A_{\omega}$ is a $\mathbf{Z}^{d}$-ergodic operator. It is known that there exist closed and nonrandom sets $\Sigma, \Sigma_{p p}, \Sigma_{a c}$, and $\Sigma_{s c}$ of $\mathbf{R}$ such that $\Sigma=\sigma\left(A_{\omega}\right), \Sigma_{p p}=\sigma_{p p}\left(A_{\omega}\right), \Sigma_{a c}=\sigma_{a c}\left(A_{\omega}\right)$, and $\Sigma_{s c}=\sigma_{s c}\left(A_{\omega}\right)$ almost surely (see [21, Theorem 1], [8], [28], [37]).

To state the main theorems for $A_{\omega}$, we define $S_{\nu}(t)$ by

$$
S_{\nu}(t):=\sup _{c \in \mathbf{R}} \nu([c, c+t])
$$

for the probability measure $\nu$ on $\mathbf{R}$ and $t>0$. For a Borel set $J \subset \mathbf{R}$ and a self-adjoint operator $T$, we denote the associated spectral projection by $\chi_{J}(T)$. For $L \in \mathbf{N}$, set a cube $\Lambda=\Lambda_{L}(0)$ centered at zero with side length $L$, that is, $\Lambda:=\left\{k \in \mathbf{Z}^{d}\left|\max _{i=1,2, \ldots, d}\right| k_{i} \mid \leq L / 2\right\} \subset \mathbf{Z}^{d}$. We denote by $A_{\omega}^{\Lambda}$ some self-adjoint restriction of $A_{\omega}$ to the cube $\Lambda$. We realize $A_{\omega}^{\Lambda}$ simply by taking the matrix elements $A_{\omega}^{\Lambda}(n, m)$ as $\left\langle\delta_{n}, A_{\omega} \delta_{m}\right\rangle$ whenever both $n$ and $m$ belong to $\Lambda$. If we let $\chi_{\Lambda}$ be the characteristic function on $\Lambda$, that is, $\chi_{\Lambda}=\sum_{\gamma \in \Lambda} \Pi_{\gamma}$, then we can also write $A_{\omega}^{\Lambda}=\chi_{\Lambda} A_{\omega} \chi_{\Lambda}$. In Appendix B , we also use the periodic boundary condition. We write $\mathbf{E}=\mathbf{E}_{\omega}$ for the expectation with respect to $\omega=\left\{\omega_{j}\right\}_{j \in \mathbf{Z}^{d}}$. We first state the following estimate.

THEOREM 2.1
Fix $E_{0} \in(0, \infty)$. Let $I$ be any bounded interval in $\left[E_{0}, \infty\right)$. Then we have the bound

$$
\begin{equation*}
\mathbf{E} \operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right) \leq C_{W} Q_{A}(|I|)|\Lambda| \tag{2.5}
\end{equation*}
$$

where $Q_{A}(|I|)=\sup _{\gamma} Q_{A, \gamma}(|I|)$,

$$
Q_{A, \gamma}(|I|)= \begin{cases}2 d\left|g_{\gamma}\right|_{\infty}|I| & \text { if } \mu_{\gamma} \text { has a bounded density } g_{\gamma},  \tag{2.6}\\ 8 d(1+2 d) S_{\mu_{\gamma}}(|I|) & \text { otherwise }\end{cases}
$$

and $C_{W}$ is a constant which depends only on $d, M_{0}, M$, and $E_{0}$.

## REMARK 1

The above constant $C_{W}$ is given by $(d+1)^{2}\left(M_{0}+M\right)^{2} / E_{0}^{2}$.
We next state the generalized eigenvalue-counting estimate for $A_{\omega}$.

## THEOREM 2.2

Fix $E_{0} \in(0, \infty)$. Let $I$ be any bounded interval in $\left[E_{0}, \infty\right)$. For $n=0,1,2, \ldots$, we have the bound

$$
\begin{equation*}
\mathbf{P}\left(\operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right) \geq n d+1\right) \leq C_{G}\left(Q_{A}(|I|)|\Lambda|\right)^{n+1} \tag{2.7}
\end{equation*}
$$

where $C_{G}$ is a constant which depends only on $n, d, M_{0}, M, m$, and $E_{0}$.

## REMARK 2

The above constant $C_{G}$ is given by

$$
\frac{(d+1)^{2(n+1)} \prod_{j=0}^{n}\left(M_{0}+2^{j} M-2^{j} m+m\right)^{2}}{E_{0}^{2(n+1)} \prod_{j=0}^{n}(j d+1)} .
$$

Let us explain that (2.7) is weaker than the corresponding estimate given in [9]. The generalized eigenvalue-counting estimate for $H_{\omega}^{\Lambda}$ in [9] is the following: for any $n \in \mathbf{N}$,

$$
\begin{equation*}
\mathbf{P}\left(\operatorname{Tr} \chi_{I}\left(H_{\omega}^{\Lambda}\right) \geq n\right) \leq \frac{1}{n!}(Q(|I|)|\Lambda|)^{n} \tag{2.8}
\end{equation*}
$$

where $Q(|I|)=\sup _{\gamma} Q_{\gamma}(|I|)$ and

$$
Q_{\gamma}(|I|)= \begin{cases}\left|g_{\gamma}\right|_{\infty}|I| & \text { if } \mu_{\gamma} \text { has a bounded density } g_{\gamma} \\ 8 S_{\mu_{\gamma}}(|I|) & \text { otherwise }\end{cases}
$$

Since we have $\mathbf{P}\left(\operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right) \geq n+1\right) \geq \mathbf{P}\left(\operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right) \geq n d+1\right)$, (2.7) is a weaker version of (2.8). However, (2.7) is the same type of estimate as (2.8) if $d=1$. Moreover, (2.7) always gives the Wegner estimate by taking $n=0$.

To state a corollary of Theorem 2.1, we introduce the IDS for $A_{\omega}$. This function is defined as follows. Since the spectrum of $A_{\omega}^{\Lambda}$ is discrete, we can define the eigenvalue-counting function by $N_{\omega}^{\Lambda}(E):=\operatorname{Tr} \chi_{(-\infty, E]}\left(A_{\omega}^{\Lambda}\right)$. If we assume (2.4), that is, that $\left\{\omega_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}}$ are independent and identically distributed (i.i.d.), we can show the existence of a nonrandom, right continuous, and increasing function $N_{A}$ such that

$$
N_{A}(E)=\lim _{|\Lambda| \rightarrow \infty} \frac{N_{\omega}^{\Lambda}(E)}{|\Lambda|}
$$

for any point of continuity of $N_{A}(E)$ and for almost all $\omega$. We call $N_{A}(E)$ the IDS for $A_{\omega}$ (see, e.g., [8], [21], [28], [37]). We say that the common probability measure $\mu$ of $\left\{\omega_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}}$ is $\alpha$-Hölder continuous if there are $c>0$ and $\alpha \in(0,1]$ such that

$$
\begin{equation*}
S_{\mu}(t) \leq c t^{\alpha} \tag{2.9}
\end{equation*}
$$

for small $t>0$. From the definition of the IDS and Fatou's lemma, it follows that

$$
\begin{aligned}
N_{A}(E)-N_{A}\left(E^{\prime}\right) & =\mathbf{E}\left[\lim _{|\Lambda| \rightarrow \infty} \frac{N_{\omega}^{\Lambda}(E)-N_{\omega}^{\Lambda}\left(E^{\prime}\right)}{|\Lambda|}\right] \\
& \leq \liminf _{|\Lambda| \rightarrow \infty} \frac{1}{|\Lambda|} \mathbf{E}\left[N_{\omega}^{\Lambda}(E)-N_{\omega}^{\Lambda}\left(E^{\prime}\right)\right] \\
& =\liminf _{|\Lambda| \rightarrow \infty} \frac{1}{|\Lambda|} \mathbf{E}\left[\operatorname{Tr} \chi_{\left(E^{\prime}, E\right]}\left(A_{\omega}^{\Lambda}\right)\right]
\end{aligned}
$$

for any point of continuity $E^{\prime}<E \in \mathbf{R}$. By Theorem 2.1 and $\alpha$-Hölder continuity of $\mu$, we have the following.

COROLLARY 2.1
Assume (2.4) and (2.9). Then the IDS for $A_{\omega}$ is $\alpha$-Hölder continuous.
We next state an application of Theorem 2.2. We say that
$A_{\omega}$ exhibits the Anderson localization in some interval $I$
if, with probability one, $A_{\omega}$ has only pure point spectrum in $I$ and the corresponding eigenfunctions decay exponentially. For $x \geq 0$, we let $[x]$ be the integral part of $x$. We prove the following result.

COROLLARY 2.2
Assume (2.4), (2.9), and (2.10). Then, with probability one, every eigenvalue of $A_{\omega}$ in I has multiplicity less than or equal to $\left[\alpha^{-1}\right] d$. In particular, if $\alpha>1 / 2$ and $d=1$, with probability one, every eigenvalue of $A_{\omega}$ in $I$ is simple.

REMARK 3
Figotin and Klein [21] showed the localization for $A_{\omega}$ near the band edge of the spectrum under some assumptions on the periodic operator $A_{0}$ and on the distribution of $\left\{\omega_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}}$.

### 2.2. Proofs of Theorems 2.1 and 2.2

We quote an essential proposition to show the main theorems.

## PROPOSITION 2.1

Let $\nu$ denote the probability distribution of the random variable $\omega$. Let I be a finite interval in $\mathbf{R}$. Let $H$ and $B$ be two self-adjoint operators on a separable Hilbert space $\mathcal{H}$, and suppose that $B$ is bounded and nonnegative. Then, for any $\varphi \in \mathcal{H}$ and any $k>0$, we have the bound

$$
\begin{align*}
& \int\left\langle\varphi, B \chi_{I}(H+\omega B) B \varphi\right\rangle \nu(d \omega)  \tag{2.11}\\
& \quad \leq \begin{cases}\|B\|_{o p}|g|_{\infty}|I|\|\varphi\|^{2} & \text { if } \nu \text { has a bounded density } g \\
4\|B\|_{o p}\left(k+\|B\|_{o p}\right) S_{\nu}(|I| / k)\|\varphi\|^{2} & \text { otherwise }\end{cases}
\end{align*}
$$

where $\|\cdot\|_{\text {op }}$ is the operator norm and $\|\cdot\|$ is the norm in $\mathcal{H}$.
This is often called the spectral averaging (see, e.g., [11, 12, 31, 42]). Combes and Hislop [11] and Kotani and Simon [31] studied the case when the probability measure $\nu$ has a bounded density. Combes and colleagues [12] proved the general case. Since [12] considered the case $k=1,(2.11)$ is slightly generalized (see [9, Appendix] of for this extension). Now, we prepare the following lemma.

LEMMA 2.1
Let $I=[a, b]$ be any finite interval in $(0, \infty)$; then we have

$$
\begin{equation*}
\operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right) \leq C \sum_{x \in \Lambda} \sum_{\gamma \in K_{x}}\left\langle\delta_{x}, B_{\gamma} \chi_{I}\left(A_{\omega}^{\Lambda}\right) B_{\gamma} \delta_{x}\right\rangle, \tag{2.12}
\end{equation*}
$$

where $C=C\left(a, d, M_{0}, M\right)=(d+1)\left(M_{0}+M\right)^{2} / a^{2}, K_{x}=\left\{x, x+e_{1}, \ldots, x+e_{d}\right\} \subset$ $\mathbf{Z}^{d}$, and $B_{\gamma}$ is the restriction of $\nabla^{*} \Pi_{\gamma} \nabla$ to the cube $\Lambda$.

## Proof

We first note that $A_{\omega}^{\Lambda}=\sum_{\gamma \in \tilde{\Lambda}}\left(q_{\gamma}+\omega_{\gamma}\right) B_{\gamma}$, where $\tilde{\Lambda} \supset \Lambda$ is defined by $\left\{k \in \mathbf{Z}^{d}\left|\max _{i=1, \ldots, d}\right| k_{i} \mid \leq L / 2+1\right\}$. We next expand the trace as $\operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right)=$ $\sum_{x \in \Lambda}\left\langle\delta_{x}, \chi_{I}\left(A_{\omega}^{\Lambda}\right) \delta_{x}\right\rangle$. By using the spectral measure $\mu_{A}^{\delta_{x}}(d \lambda)$ for $A_{\omega}^{\Lambda}$ with respect to $\delta_{x}$, we have

$$
\left\langle\delta_{x}, \chi_{I}\left(A_{\omega}^{\Lambda}\right) \delta_{x}\right\rangle=\int_{I} \mu_{A}^{\delta_{x}}(d \lambda) \leq \frac{1}{a^{2}} \int_{a}^{b} \lambda^{2} \mu_{A}^{\delta_{x}}(d \lambda) .
$$

Thus, since $\rho_{0}$ and $\rho_{\omega}$ are bounded, we get

$$
\operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right) \leq \frac{\left(M_{0}+M\right)^{2}}{a^{2}} \sum_{x \in \Lambda} \sum_{\gamma_{1}, \gamma_{2} \in \tilde{\Lambda}}\left\langle\delta_{x}, B_{\gamma_{1}} \chi_{I}\left(A_{\omega}^{\Lambda}\right) B_{\gamma_{2}} \delta_{x}\right\rangle .
$$

Now, we note that

$$
\nabla^{*} \Pi_{\gamma} \nabla=\sum_{j=1}^{d}\left|\partial_{j}^{*} \delta_{\gamma}\right\rangle\left\langle\partial_{j}^{*} \delta_{\gamma}\right|=\sum_{j=1}^{d}\left|\delta_{\gamma}-\delta_{\gamma-e_{j}}\right\rangle\left\langle\delta_{\gamma}-\delta_{\gamma-e_{j}}\right| .
$$

Then, it is easy to see that for fixed $x \in \Lambda$,

$$
\nabla^{*} \Pi_{\gamma} \nabla \delta_{x}= \begin{cases}d \delta_{x}-\sum_{j=1}^{d} \delta_{x-e_{j}} & \text { if } \gamma=x \\ \delta_{x}-\delta_{x+e_{j}} & \text { if } \gamma=x+e_{j}, j=1,2, \ldots, d \\ 0 & \text { otherwise }\end{cases}
$$

Thus, we have $B_{\gamma} \delta_{x}=0$ if $\gamma \notin K_{x}$. By using

$$
\begin{aligned}
& \left\langle\delta_{x}, B_{\gamma_{1}} \chi_{I}\left(A_{\omega}^{\Lambda}\right) B_{\gamma_{2}} \delta_{x}\right\rangle \\
& \quad \leq \frac{1}{2}\left(\left\langle\delta_{x}, B_{\gamma_{1}} \chi_{I}\left(A_{\omega}^{\Lambda}\right) B_{\gamma_{1}} \delta_{x}\right\rangle+\left\langle\delta_{x}, B_{\gamma_{2}} \chi_{I}\left(A_{\omega}^{\Lambda}\right) B_{\gamma_{2}} \delta_{x}\right\rangle\right),
\end{aligned}
$$

we complete the proof.

For $j \in \mathbf{Z}^{d}$, we write $\mathbf{E}_{\omega_{j}}$ for the expectation with respect to the random variable $\omega_{j}$. We also write $\omega_{k}^{\perp}:=\left\{\omega_{j}\right\}_{j \in \mathbf{Z}^{d} \backslash\{k\}}$ and write $\mathbf{E}_{\omega_{k}^{\perp}}$ for the corresponding expectation.

Proof of Theorem 2.1
Let $I=[a, b] \subset\left[E_{0}, \infty\right)$. Taking the expectation for (2.12), we get

$$
\begin{equation*}
\mathbf{E} \operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right) \leq C \sum_{x \in \Lambda} \sum_{\gamma \in K_{x}} \mathbf{E}_{\omega_{\gamma}^{\perp}} \mathbf{E}_{\omega_{\gamma}}\left\langle\delta_{x}, B_{\gamma} \chi_{I}\left(A_{\omega}^{\Lambda}\right) B_{\gamma} \delta_{x}\right\rangle . \tag{2.13}
\end{equation*}
$$

Now, note that $\left\|B_{\gamma}\right\|_{o p} \leq 2 d$ and $\left\|\delta_{x}\right\|=1$. By using Proposition 2.1 with $k=1$ in (2.11), we can show that the right-hand side of (2.13) is less than or equal to

$$
C \sum_{x \in \Lambda} \sum_{\gamma \in K_{x}} \mathbf{E}_{\omega_{\gamma}^{\perp}} Q_{A, \gamma}(|I|) \leq C(d+1) Q_{A}(|I|)|\Lambda|,
$$

where $Q_{A}(|I|)=\sup _{\gamma} Q_{A, \gamma}(|I|)$ and $Q_{A, \gamma}(|I|)$ is defined by (2.6). This completes the proof.

To prove Theorem 2.2, we prepare the following lemma, which is a generalization of [9, Lemma 4.1].

LEMMA 2.2
Consider the self-adjoint operator $H_{s}=H_{0}+s B$ on a separable Hilbert space $\mathcal{H}$, where $H_{0}$ and $B$ are two self-adjoint operators on $\mathcal{H}, B$ is a nonnegative operator satisfying $\operatorname{rank} B \leq d$, and $s \in \mathbf{R}$. Let $P_{s}(J)=\chi_{J}\left(H_{s}\right)$ for an interval $J$, and suppose that $\operatorname{Tr} P_{0}((-\infty, c])<\infty$ for all $c \in \mathbf{R}$. Then, given $a, b \in \mathbf{R}$ with $a<b$, we have

$$
\operatorname{Tr} P_{s}((a, b]) \leq d+\operatorname{Tr} P_{t}((a, b])
$$

for all $0 \leq s \leq t$.

Proof
Let $0 \leq s \leq t$. Note that for any $c \in \mathbf{R}$ we always have

$$
\begin{equation*}
0 \leq \operatorname{Tr} P_{s}((-\infty, c])-\operatorname{Tr} P_{t}((-\infty, c]) \leq d \tag{2.14}
\end{equation*}
$$

The last inequality is a consequence of the min-max principle, proved similarly in [28, Lemma 5.22]. We can also consider (2.14) as a bound of the spectral shift function (SSF) $\xi\left(\lambda ; H_{s}, H_{t}\right)$ for the pair $\left(H_{s}, H_{t}\right)$. Since $H_{s}-H_{t}$ is an operator of finite rank, we have the bound $0 \leq \xi\left(\lambda ; H_{s}, H_{t}\right) \leq \operatorname{rank}\left(H_{s}-H_{t}\right) \leq d$. For the general theory of the SSF, refer to [6]. Thus, we obtain

$$
\begin{aligned}
\operatorname{Tr} P_{s}((a, b]) & =\operatorname{Tr} P_{s}((-\infty, b])-\operatorname{Tr} P_{s}((-\infty, a]) \\
& \leq \operatorname{Tr} P_{s}((-\infty, b])-\operatorname{Tr} P_{t}((-\infty, a]) \\
& =\operatorname{Tr} P_{s}((-\infty, b])-\operatorname{Tr} P_{t}((-\infty, b])+\operatorname{Tr} P_{t}((a, b]) \\
& \leq d+\operatorname{Tr} P_{t}((a, b]) .
\end{aligned}
$$

## Proof of Theorem 2.2

Let $\mathbf{1}_{A}(\omega)$ be the indicator function of an event $A$. For $n \in \mathbf{N}$, we set the events $E_{n}$ and $A_{n}$ by $E_{n}:=\left\{\omega \mid \operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right) \geq n d+1\right\}$ and $A_{n}:=\left\{\omega \mid \prod_{k=0}^{n}\left(\operatorname{Tr} \chi_{I}\left(A_{\Lambda}^{\omega}\right)-\right.\right.$ $\left.k d) \geq C_{n, d}\right\}$, where $C_{n, d}=\prod_{k=0}^{n}(k d+1)$. By $E_{n} \subset A_{n}$ and Chebyshev's inequality, we have

$$
\begin{aligned}
\mathbf{P}\left(\operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right) \geq n d+1\right) & =\int_{E_{n}} \mathbf{1}_{A_{n}}(\omega) \mathbf{P}(d \omega) \\
& \leq \mathbf{E}\left[\prod_{k=0}^{n}\left(\operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right)-k d\right) \mathbf{1}_{E_{n}}(\omega)\right] / C_{n, d}
\end{aligned}
$$

When $n=0$, this is just the Wegner estimate (2.5). For the reader's convenience, we next consider the case $n=1$. By using Lemma 2.1, we first get

$$
\begin{align*}
& \operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right)\left(\operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right)-d\right) \mathbf{1}_{E_{1}}(\omega) \\
& \quad \leq C \sum_{x \in \Lambda} \sum_{\gamma \in K_{x}}\left\langle\delta_{x}, B_{\gamma} \chi_{I}\left(A_{\omega}^{\Lambda}\right) B_{\gamma} \delta_{x}\right\rangle\left(\operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right)-d\right) \mathbf{1}_{E_{1}}(\omega) \tag{2.15}
\end{align*}
$$

where $C=C\left(a, d, M_{0}, M\right)$ in Lemma 2.1. For $\omega=\left(\omega_{\gamma}^{\perp}, \omega_{\gamma}\right)$ and $\gamma \in \tilde{\Lambda}$, we define $A_{\left(\omega_{\gamma}^{\perp}, s\right)}^{\Lambda}$ by $A_{\left(\omega_{\gamma}^{\perp}, s\right)}^{\Lambda}=\sum_{k} q_{k} B_{k}+\sum_{k \neq \gamma} \omega_{k} B_{k}+s B_{\gamma}$. For each $\gamma \in \mathbf{Z}^{d}$, we take a random variable $\tilde{\omega}_{\gamma}$ with the probability distribution $\mu_{\gamma}$, so that $\tilde{\omega}_{\gamma}$ is independent from all of the other random variables, and set $\tau_{\gamma}=\tilde{\omega}_{\gamma}+M-m$. Write $\tilde{\omega}=\left\{\tilde{\omega}_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}}$. By Lemma 2.2 with $\tau_{\gamma} \geq M \geq \omega_{\gamma}$, we have that the right-hand side of (2.15) is less than or equal to

$$
C \sum_{x \in \Lambda} \sum_{\gamma \in K_{x}}\left\langle\delta_{x}, B_{\gamma} \chi_{I}\left(A_{\omega}^{\Lambda}\right) B_{\gamma} \delta_{x}\right\rangle \operatorname{Tr} \chi_{I}\left(A_{\left(\omega_{\gamma}^{\perp}, \tau_{\gamma}\right)}^{\Lambda}\right) \mathbf{1}_{E_{1}^{\prime}}\left(\omega_{\gamma}^{\perp}, \tau_{\gamma}\right),
$$

where we set $E_{1}^{\prime}=\left\{\left(\omega_{\gamma}^{\perp}, \tau_{\gamma}\right) \mid \operatorname{Tr} \chi_{I}\left(A_{\left(\omega_{\gamma}^{\perp}, \tau_{\gamma}\right)}^{\Lambda}\right) \geq 1\right\}$ and used

$$
\begin{equation*}
\mathbf{1}_{E_{1}}(\omega) \leq \mathbf{1}_{E_{1}^{\prime}}\left(\omega_{\gamma}^{\perp}, \tau_{\gamma}\right) \tag{2.16}
\end{equation*}
$$

for any $\tilde{\omega}_{\gamma}$ and $\omega$. Since $\operatorname{Tr} \chi_{I}\left(A_{\left(\omega_{\gamma}, \tau_{\gamma}\right)}^{\Lambda}\right) \mathbf{1}_{E_{1}^{\prime}}\left(\omega_{\gamma}^{\perp}, \tau_{\gamma}\right)$ is independent of $\omega_{\gamma}$, by taking the expectation with respect to $\omega$ and $\tilde{\omega}$, and using Proposition 2.1, we get

$$
\begin{aligned}
& \mathbf{E}_{\tilde{\omega}} \mathbf{E}_{\omega}\left[\operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right)\left(\operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right)-d\right) \mathbf{1}_{E_{1}}(\omega)\right] \\
& \quad \leq C \sum_{x \in \Lambda} \sum_{\gamma \in K_{x}} \mathbf{E}_{\tilde{\omega}_{\gamma}} \mathbf{E}_{\omega_{\gamma}^{\perp}}\left[\mathbf{E}_{\omega_{\gamma}}\left[\left\langle\delta_{x}, B_{\gamma} \chi_{I}\left(A_{\omega}^{\Lambda}\right) B_{\gamma} \delta_{x}\right\rangle\right] \operatorname{Tr} \chi_{I}\left(A_{\left(\omega_{\gamma}^{\perp}, \tau_{\gamma}\right)}^{\Lambda}\right) \mathbf{1}_{E_{1}^{\prime}}\left(\omega_{\gamma}^{\perp}, \tau_{\gamma}\right)\right] \\
& \quad \leq C Q_{A}(|I|) \sum_{x \in \Lambda} \sum_{\gamma \in K_{x}} \mathbf{E}_{\tilde{\omega}_{\gamma}} \mathbf{E}_{\omega_{\gamma}^{\perp}}\left[\operatorname{Tr} \chi_{I}\left(A_{\left(\omega_{\gamma}^{\perp}, \tau_{\gamma}\right)}^{\Lambda}\right) \mathbf{1}_{E_{1}^{\prime}}\left(\omega_{\gamma}^{\perp}, \tau_{\gamma}\right)\right],
\end{aligned}
$$

where $Q_{A}(|I|)=\sup _{\gamma} Q_{A, \gamma}(|I|)$ and $Q_{A, \gamma}(|I|)$ is defined in (2.6). By using the Wegner estimate (2.5), we get

$$
\mathbf{E}_{\tilde{\omega}_{\gamma}} \mathbf{E}_{\omega_{\gamma}^{\perp}}\left[\operatorname{Tr} \chi_{I}\left(A_{\left(\omega_{\gamma}^{\perp}, \tau_{\gamma}\right)}^{\Lambda}\right) \mathbf{1}_{E_{1}^{\prime}}\left(\omega_{\gamma}^{\perp}, \tau_{\gamma}\right)\right] \leq C^{\prime}(d+1) Q_{A}(|I|)|\Lambda|,
$$

where $C^{\prime}=C\left(a, d, M_{0}, 2 M-m\right)$. We remark that $\sup _{\gamma}\left\{\omega_{\gamma}^{\perp}, \tau_{\gamma}\right\}=2 M-m$. This completes the proof for $n=1$.

For the general $n \in \mathbf{N}$, we can proceed inductively as in [9]. Let us assume that for all $m, M$ with $0<m<M$ and for all possible probability distributions $\mu_{\gamma}$ of $\omega_{\gamma}$ satisfying $\operatorname{supp} \mu_{\gamma} \subset[m, M]$ for all $\gamma \in \mathbf{Z}^{d}$,

$$
\begin{equation*}
\mathbf{E}\left[\prod_{k=0}^{n-1}\left(\operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right)-k d\right) \mathbf{1}_{E_{n-1}}(\omega)\right] \leq C_{n}\left(Q_{A}(|I|)|\Lambda|\right)^{n} \tag{2.17}
\end{equation*}
$$

holds for $n \in \mathbf{N}$, where $C_{n}=C_{n}\left(a, d, M_{0}, M, m\right)$ is defined by

$$
\prod_{j=0}^{n-1}\left(M_{0}+2^{j} M-2^{j} m+m\right)^{2} \frac{(d+1)^{2 n}}{a^{2 n}} .
$$

We remark that $M \geq \sup _{\gamma} \omega_{\gamma}$ and $m \leq \inf _{\gamma} \omega_{\gamma}$. For each $\gamma \in \mathbf{Z}^{d}$, we take $\tilde{\omega}_{\gamma}$ as in the case for $n=1$. Then, from Lemmas 2.1 and 2.2, it follows that

$$
\begin{equation*}
\prod_{k=0}^{n}\left(\operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right)-k d\right) \mathbf{1}_{E_{n}}(\omega) \tag{2.18}
\end{equation*}
$$

is dominated by

$$
\begin{aligned}
C & \sum_{x \in \Lambda} \sum_{\gamma \in K_{x}}\left\langle\delta_{x}, B_{\gamma} \chi_{I}\left(A_{\omega}^{\Lambda}\right) B_{\gamma} \delta_{x}\right\rangle \prod_{k=1}^{n}\left(\operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right)-k d\right) \mathbf{1}_{E_{n}}(\omega) \\
& \leq C \sum_{x \in \Lambda} \sum_{\gamma \in K_{x}}\left\langle\delta_{x}, B_{\gamma} \chi_{I}\left(A_{\omega}^{\Lambda}\right) B_{\gamma} \delta_{x}\right\rangle \prod_{k=0}^{n-1}\left(\operatorname{Tr} \chi_{I}\left(A_{\left(\omega_{\gamma}, \tau_{\gamma}\right)}^{\Lambda}\right)-k d\right) \mathbf{1}_{E_{n}^{\prime}}\left(\omega_{\gamma}^{\perp}, \tau_{\gamma}\right)
\end{aligned}
$$

where we set $E_{n}^{\prime}=\left\{\left(\omega_{\gamma}^{\perp}, \tau_{\gamma}\right) \mid \operatorname{Tr} \chi_{I}\left(A_{\left(\omega_{\gamma}, \tau_{\gamma}\right)}^{\Lambda}\right) \geq(n-1) d+1\right\}$ and used

$$
\begin{equation*}
\mathbf{1}_{E_{n}}(\omega) \leq \mathbf{1}_{E_{n}^{\prime}}\left(\omega_{\gamma}^{\perp}, \tau_{\gamma}\right) \tag{2.19}
\end{equation*}
$$

for any $\tilde{\omega}_{\gamma}$ and $\omega$. Take the expectation with respect to $\omega=\left\{\omega_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}}$ and $\tilde{\omega}=\left\{\tilde{\omega}_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}}$. From Proposition 2.1, it follows that the expectation of (2.18) is dominated by

$$
\begin{aligned}
& C Q_{A}(|I|) \sum_{x \in \Lambda} \sum_{\gamma \in K_{x}} \mathbf{E}_{\tilde{\omega}_{\gamma}} \mathbf{E}_{\omega_{\gamma}}\left[\prod_{k=0}^{n-1}\left(\operatorname{Tr} \chi_{I}\left(A_{\left(\omega_{\gamma}, \tau_{\gamma}\right)}^{\Lambda}\right)-k d\right) \mathbf{1}_{E_{n-1}}\left(\omega_{\gamma}^{\perp}, \tau_{\gamma}\right)\right] \\
& \quad \leq C Q_{A}(|I|) \sum_{x \in \Lambda} \sum_{\gamma \in K_{x}} C_{n}^{\prime}\left(Q_{A}(|I|)|\Lambda|\right)^{n},
\end{aligned}
$$

where $C_{n}^{\prime}=C_{n}\left(a, d, M_{0}, 2 M-m, m\right)$. Note that $2 M-m \geq \sup _{\gamma}\left\{\omega_{\gamma}^{\perp}, \tau_{\gamma}\right\}$ and $m \leq \inf _{\gamma}\left\{\omega_{\gamma}^{\perp}, \tau_{\gamma}\right\}$. In the last line, we used the induction hypothesis (2.18) and the fact that $E_{n}^{\prime}$ corresponds to $E_{n-1}$. By $C(d+1) C_{n}^{\prime}=C_{n+1}$, we complete the proof.

### 2.3. Proof of Corollary 2.2

In this subsection, we prove that the multiplicity of eigenvalues for $A_{\omega}$ in some localization region is finite by using Klein and Molchanov's method [29].

## Proof of Corollary 2.2

We prove only the case for $\alpha=1$ (see [9] for the general case). The proof is based on two results (see [29, Lemmas 1, 2]). The first one is a probabilistic estimate given by Theorem 2.2. The other is a deterministic result.

We say that $\psi \in \ell^{2}\left(\mathbf{Z}^{d}\right)$ is $\beta$-decaying if $|\psi(x)| \leq C_{\psi}(1+|x|)^{-\beta}$ for some $C_{\psi}>$ 0 . If $\alpha=1$ and $\beta>5 d / 2$, we show that, with probability one, $A_{\omega}$ cannot have an eigenvalue with $d+1$ linearly independent $\beta$-decaying eigenfunctions. Pick $q>$ $2 d$. For a scale $L>0$, we cover an open interval $I$ by $2\left(\left[L^{q}|I| / 2\right]+1\right) \leq L^{q}|I|+2$ intervals of length $2 L^{-q}$ so that any subinterval $J \subset I$ with length $|J| \leq L^{-q}$ is contained in one of these intervals. We consider the event $\mathcal{B}_{L, I, q}$, which occurs if there exists an interval $J \subset I$ with $|J| \leq L^{-q}$ such that $\operatorname{Tr} \chi_{J}\left(A_{\omega}^{\Lambda}\right) \geq d+1$. By Theorem 2.2 with $n=1$, its probability can be estimated by

$$
\begin{align*}
\mathbf{P}\left(\mathcal{B}_{L, I, q}\right) & \leq\left(L^{q}|I|+2\right) C_{G}\left(C\left(2 L^{-q}\right) L^{d}\right)^{2} \\
& \leq C^{\prime}(|I|+1) L^{-q+2 d}, \tag{2.20}
\end{align*}
$$

where $C_{G}$ is the constant given in (2.7) and $C, C^{\prime}$ are some constants independent of $L$. Take $L_{k}=2^{k}$, and use (2.20) and the Borel-Cantelli lemma. It follows that if $q>2 d$, then for $\mathbf{P}$ almost every $\omega$ there exists $k(q, \omega)<\infty$ such that the event $\mathcal{B}_{L_{k}, I, q}$ does not occur if $k \geq k(q, \omega)$.

Now, suppose that for some $\omega$ there exists $E \in I$ which is an eigenvalue of $A_{\omega}$ with $d+1$ linearly independent $\beta$-decaying eigenfunctions. As in the proof of [29, Lemma 1], it follows that for $L$ large enough $A_{\omega}^{\Lambda}$ has at least $d+1$ eigenvalues in the interval $J_{E, L}=\left[E-\epsilon_{L}, E+\epsilon_{L}\right]$, where $\epsilon_{L}=C_{0} L^{-\beta+d / 2}$ for an appropriate constant $C_{0}$ independent of $L$. By $\beta>5 d / 2$, we can pick $q$ satisfying $2 d<q<\beta-d / 2$. Hence, we have $\epsilon_{L} \leq L^{-q}$ for all large $L$. However, with probability one, this is not possible since the event $\mathcal{B}_{L_{k}, I, q}$ does not occur for large $L_{k}$.

## 3. Random acoustic operator $\Gamma_{\omega}$

### 3.1. Generalized eigenvalue-counting estimate for $\Gamma_{\omega}$

In this section, we consider a simple acoustic model including (1.2). Our second model $\Gamma_{\omega}$ is the random self-adjoint operator defined by

$$
\Gamma_{\omega}:=\frac{1}{\sqrt{\rho_{\omega}}} H_{0} \frac{1}{\sqrt{\rho_{\omega}}}
$$

on $\ell^{2}\left(\mathbf{Z}^{d}\right)$, where $H_{0}$ is a nonnegative bounded self-adjoint operator. The multiplicative operator $1 / \sqrt{\rho_{\omega}}$ is defined as in (2.2). We also take the random variables $\omega=\left\{\omega_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}}$ as before. Since $H_{0}$ is any bounded operator, we may choose $H_{0}=\nabla^{*}(1 / b) \nabla$ as in (2.1). Of course, this model includes the case

$$
\begin{equation*}
H_{0}=-\Delta \tag{3.1}
\end{equation*}
$$

where $-\Delta$ is the discrete Laplacian defined by

$$
\begin{equation*}
(-\Delta f)(x):=\sum_{y:|x-y|=1}(f(x)-f(y)) . \tag{3.2}
\end{equation*}
$$

Faris [18]-[20] studied this random operator and Combes, Hislop, and Tip [13] studied the continuous version of $\Gamma_{\omega}$.

We define the restriction of $\Gamma_{\omega}$ to the cube $\Lambda$ by $\Gamma_{\omega}^{\Lambda}=\chi_{\Lambda} \Gamma_{\omega} \chi_{\Lambda}$. We also denote the IDS for $\Gamma_{\omega}$ by $N_{\Gamma}$ as in the definition of $N_{A}$ for $A_{\omega}$. We first have the following Wegner estimate and its corollary. This is an extension of the result by Faris [18]-[20] to the case for general probability measures.

THEOREM 3.1
Fix $E_{0} \in(0, \infty)$. Let I be any bounded interval in $\left[E_{0}, \infty\right)$. Then we have the bound

$$
\begin{equation*}
\mathbf{E} \operatorname{Tr} \chi_{I}\left(\Gamma_{\omega}^{\Lambda}\right) \leq Q_{\Gamma}(|I|)|\Lambda|, \tag{3.3}
\end{equation*}
$$

where $Q_{\Gamma}(|I|)=\sup _{\gamma} Q_{\Gamma, \gamma}(|I|)$ and

$$
Q_{\Gamma, \gamma}(|I|)= \begin{cases}\left|g_{\gamma}\right|_{\infty}|I| M / E_{0} & \text { if } \mu_{\gamma} \text { has a bounded density } g_{\gamma}  \tag{3.4}\\ 4\left(1+M / E_{0}\right) S_{\mu_{\gamma}}(|I|) & \text { otherwise. }\end{cases}
$$

COROLLARY 3.1
Assume (2.4), (2.9), and (3.1) for $\Gamma_{\omega}$. Then the IDS for $\Gamma_{\omega}$ is $\alpha$-Hölder continuous.

Next, we obtain the following generalized eigenvalue-counting estimate for $\Gamma_{\omega}$.

THEOREM 3.2
Fix $E_{0} \in(0, \infty)$. Let $I$ be any bounded interval in $\left[E_{0}, \infty\right)$. For $n=1,2,3, \ldots$, we have the bound

$$
\begin{equation*}
\mathbf{P}\left(\operatorname{Tr} \chi_{I}\left(\Gamma_{\omega}^{\Lambda}\right) \geq n\right) \leq \frac{1}{n!}\left(Q_{\Gamma}(|I|)|\Lambda|\right)^{n} \tag{3.5}
\end{equation*}
$$

From Theorem 3.2, we have the following, which is proven as in Corollary 2.2.

COROLLARY 3.2
Assume (2.4), (2.9), (2.10), and (3.1) for $\Gamma_{\omega}$. Then, with probability one, every eigenvalue of $\Gamma_{\omega}$ in I has multiplicity less than or equal to $\left[\alpha^{-1}\right]$. In particular, if $\alpha>1 / 2$, with probability one, every eigenvalue of $\Gamma_{\omega}$ in $I$ is simple.

REMARK 4
By using the multiscale analysis, Faris [18]-[20] showed the localization for $\Gamma_{\omega}$ near the top of the spectrum under some assumptions on the distribution of $\left\{\omega_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}}$ (see also our example in Appendix C).

### 3.2. Auxiliary random Schrödinger operator

In this subsection, to prove Theorems 3.1 and 3.2 , we introduce an auxiliary Schrödinger-type operator defined by $H_{\omega}(E)=H_{0}-E \rho_{\omega}$. In fact, $H_{\omega}(E)$ is
an Anderson model if (3.1) holds. The main statement is the following proposition.

## PROPOSITION 3.1

Suppose $0 \notin I_{\epsilon}:=[E-\epsilon, E+\epsilon]$. Then we have

$$
\begin{equation*}
\operatorname{Tr} \chi_{[-m \epsilon, m \epsilon]}\left(H_{\omega}^{\Lambda}(E)\right) \leq \operatorname{Tr} \chi_{I_{\epsilon}}\left(\Gamma_{\omega}^{\Lambda}\right) \leq \operatorname{Tr} \chi_{[-M \epsilon, M \epsilon]}\left(H_{\omega}^{\Lambda}(E)\right) . \tag{3.6}
\end{equation*}
$$

## REMARK 5

Let $N_{H}$ be the IDS for $H_{\omega}(E)$. Proposition 3.1 implies

$$
\begin{equation*}
N_{H}(m \epsilon)-N_{H}(-m \epsilon) \leq N_{\Gamma}(E+\epsilon)-N_{\Gamma}(E-\epsilon) \leq N_{H}(M \epsilon)-N_{H}(-M \epsilon), \tag{3.7}
\end{equation*}
$$

which shows upper and lower bounds of the density of states for $\Gamma_{\omega}$ in terms of that for $H_{\omega}(E)$.

This idea was already used to study the IDS for continuous random Schrödinger operators. For example, Hislop and Klopp [26] used the Birman-Schwinger-type operator $\Xi_{\omega}(E)$ as an auxiliary operator. Let us consider $H_{\omega}:=-\triangle+V_{\omega}$ on $L^{2}\left(\mathbf{R}^{d}\right)$, where $V_{\omega}(x):=\sum_{\gamma \in \mathbf{Z}^{d}} \omega_{\gamma} u_{\gamma}(x)$ is a standard Anderson-type random potential and the sign of single site potential $u_{\gamma}(x)$ is indefinite. For $E<0$, they used the following inequality:

$$
\operatorname{Tr} \chi_{[E-\epsilon, E+\epsilon]}\left(H_{\omega}^{\Lambda}\right) \leq \operatorname{Tr} \chi_{[-1-c \epsilon,-1+c \epsilon]}\left(\Xi_{\omega}^{\Lambda}(E)\right),
$$

where $c$ is some constant independent of $\Lambda, \epsilon$, and $\omega$. The Birman-Schwingertype operator $\Xi_{\omega}(E)$ is defined by $\Xi_{\omega}(E):=(-\triangle-E)^{-1 / 2} V_{\omega}(-\triangle-E)^{-1 / 2}$. Thus, they estimated the expectation of the number of eigenvalues near -1 (for some application to random magnetic fields, see, e.g., [41, Lemma 2.1]).

## Proof of Proposition 3.1

We prove only the second inequality. The same argument also applies to the first inequality.

To estimate the number of eigenvalues of $\Gamma_{\omega}^{\Lambda}$ in the interval $I_{\epsilon}$, we use the Birman-Schwinger principle (see, e.g., [38]). Suppose that $\phi_{\lambda}$ is an eigenfunction of $\Gamma_{\omega}^{\Lambda}$ with the eigenvalue $\lambda \neq 0$ and $\lambda \in I_{\epsilon}$. Then the function $\psi$ defined by $\psi=$ $\left(1 / \sqrt{\rho_{\omega}}\right) \phi_{\lambda}$ satisfies $H_{\omega}^{\Lambda}(\lambda) \psi=0$, where we set $H_{\omega}(\lambda)=H_{0}-\lambda \rho_{\omega}$. Therefore, we know that $0 \in \sigma_{p}\left(H_{\omega}^{\Lambda}(\lambda)\right)$ if and only if $\lambda \in \sigma_{p}\left(\Gamma_{\omega}^{\Lambda}\right)$. By the assumption of $\rho_{\omega}$, we have $H_{\omega}^{\Lambda}(\lambda)-M \epsilon \leq H_{\omega}^{\Lambda}(E) \leq H_{\omega}^{\Lambda}(\lambda)+M \epsilon$. Let $\mu_{n}(\lambda)$ be the $n$th eigenvalue of $H_{\omega}^{\Lambda}(\lambda)$. Similarly, let $k_{n}(E)$ be the $n$th eigenvalue of $H_{\omega}^{\Lambda}(E)$. By using the min-max principle, we have

$$
\begin{equation*}
\mu_{n}(\lambda)-M \epsilon \leq k_{n}(E) \leq \mu_{n}(\lambda)+M \epsilon . \tag{3.8}
\end{equation*}
$$

Now, take $\lambda_{n} \in I_{\epsilon}$ satisfying $\mu_{n}\left(\lambda_{n}\right)=0$. By (3.8), we have $-M \epsilon \leq k_{n}(E) \leq M \epsilon$. Hence, we obtain

$$
\#\left\{n \mid \mu_{n}(\lambda)=0 \quad \text { for some } \lambda \in I_{\epsilon}\right\} \leq \#\left\{n \mid k_{n}(E) \in[-M \epsilon, M \epsilon]\right\}
$$

As in Section 2, we first prove the Wegner estimate by using Proposition 3.1.
Proof of Theorem 3.1
We set $I=[E-\epsilon, E+\epsilon]$ and $\tilde{I}=[-M \epsilon, M \epsilon]$. From (3.6), it follows that $\mathbf{E} \operatorname{Tr} \chi_{I}\left(\Gamma_{\omega}^{\Lambda}\right)$ is dominated by

$$
\mathbf{E} \operatorname{Tr} \chi_{\tilde{I}}\left(H_{\omega}^{\Lambda}(E)\right)=\sum_{\gamma \in \Lambda} \mathbf{E}\left\langle\delta_{\gamma}, \chi_{\tilde{I}}\left(H_{\omega}^{\Lambda}(E)\right) \delta_{\gamma}\right\rangle
$$

Now, use Proposition 2.1 with $B=E \Pi_{\varphi}, k=M$, and $\varphi=\delta_{\gamma}$; then we have the bound

$$
\begin{align*}
& \int\left\langle\delta_{\gamma}, \chi_{\tilde{I}}\left(H_{\omega}^{\Lambda}(E)\right) \delta_{\gamma}\right\rangle \mu_{\gamma}\left(d \omega_{\gamma}\right)  \tag{3.9}\\
& \quad \leq \begin{cases}\left|g_{\gamma}\right|_{\infty}|\tilde{I}| / E & \text { if } \mu_{\gamma} \text { has a bounded density } g_{\gamma}, \\
4(1+M / E) S_{\mu_{\gamma}}(\tilde{I} \mid / M) & \text { otherwise } .\end{cases}
\end{align*}
$$

By (3.9) and $|\tilde{I}|=M|I|$, we complete the proof.
We next prove the generalized eigenvalue-counting estimate for $\Gamma_{\omega}$.
Proof of Theorem 3.2
We can prove Theorem 3.2 as we proved Theorem 2.2. Set $I=[E-\epsilon, E+\epsilon]$ and $\tilde{I}=[-M \epsilon, M \epsilon]$. We also set

$$
E_{n}=\left\{\omega \mid \operatorname{Tr} \chi_{I}\left(\Gamma_{\omega}^{\Lambda}\right) \geq n\right\} \quad \text { and } \quad \tilde{E}_{n}=\left\{\omega \mid \operatorname{Tr} \chi_{\tilde{I}}\left(H_{\omega}^{\Lambda}(E)\right) \geq n\right\}
$$

where we write $H_{\omega}(E)=H_{0}-E \rho_{\omega}$ as in Proposition 3.1 and $H_{\omega}^{\Lambda}(E)$ is its restriction to $\Lambda$. Then, by Proposition 3.1, we have $E_{n} \subset \tilde{E}_{n}$. From Chebyshev's inequality and Proposition 3.1, it follows that $\mathbf{P}\left(\operatorname{Tr} \chi_{I}\left(\Gamma_{\omega}^{\Lambda}\right) \geq n\right)$ is dominated by

$$
\begin{equation*}
\frac{1}{n!} \mathbf{E}\left[\prod_{k=0}^{n-1}\left(\operatorname{Tr} \chi_{I}\left(\Gamma_{\omega}^{\Lambda}\right)-k\right) \mathbf{1}_{E_{n}}\right] \leq \frac{1}{n!} \mathbf{E}\left[\prod_{k=0}^{n-1}\left(\operatorname{Tr} \chi_{\tilde{I}}\left(H_{\omega}^{\Lambda}(E)\right)-k\right) \mathbf{1}_{\tilde{E}_{n}}\right] \tag{3.10}
\end{equation*}
$$

Then, since $H_{\omega}(E)$ is an Anderson model, the right-hand side of (3.10) can be estimated similarly as in [9] from Lemma 2.2 and the spectral averaging (3.9). We proceed by the induction on $n \in \mathbf{N}$ and complete the proof.

### 3.3. Poisson statistics for eigenvalues of $\Gamma_{\omega}$

As an application of Theorem 3.2, we state Poisson statistics for eigenvalues of $\Gamma_{\omega}^{\Lambda}$. Following Molchanov [33] and Minami [34], we study a point process defined by the rescaled eigenvalues of $\Gamma_{\omega}^{\Lambda}$.

We first introduce an important point process. A point process $\xi_{\omega}$ on $\mathbf{R}$ is said to be the Poisson point process with intensity measure $\nu$ if it satisfies the following two conditions.
(a) For each bounded Borel set $J \subset \mathbf{R}, \xi_{\omega}(J)$ obeys the Poisson distribution with parameter $\nu(J)$, namely,

$$
\mathbf{P}\left(\xi_{\omega}(J)=r\right)=e^{-\nu(J)} \frac{\nu(J)^{r}}{r!}, \quad r \geq 0 .
$$

(b) If $J_{1}, \ldots, J_{n}$ are disjoint, then $\xi_{\omega}\left(J_{1}\right), \ldots, \xi_{\omega}\left(J_{n}\right)$ are independent random variables.

Now, let $E_{\omega, k}^{\Lambda}$ be $k$ th eigenvalue of $\Gamma_{\omega}^{\Lambda}(k=1,2, \ldots,|\Lambda|)$. For $E \in \sigma\left(\Gamma_{\omega}\right)$ and $E>0$, we define a point process $\xi_{\omega, E}^{\Lambda}$ on $\mathbf{R}$ by the rescaled eigenvalue of $\Gamma_{\omega}^{\Lambda}$ near $E$ :

$$
\begin{equation*}
\xi_{\omega, E}^{\Lambda}(d t)=\sum_{k=1}^{|\Lambda|} \delta_{|\Lambda|\left(E_{\omega, k}^{\Lambda}-E\right)}(d t) \tag{3.11}
\end{equation*}
$$

We always assume (2.4) and (3.1); that is, the random variables $\left\{\omega_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}}$ are independent and identically distributed (i.i.d.) by $\mu$ and $H_{0}=-\Delta$. Moreover, we introduce the following assumptions
(H1) The common distribution $\mu$ has a bounded density $g$.
(H2) There are finite constants $C, D>0, s \in(0,1)$, and $I \subset \mathbf{R}$ such that

$$
\begin{equation*}
\mathbf{E}\left|\left\langle\delta_{x},\left(\Gamma_{\omega}^{\Lambda}-z\right)^{-1} \delta_{y}\right\rangle\right|^{s} \leq C e^{-D|x-y|} \tag{3.12}
\end{equation*}
$$

for all $z \in \mathbf{C}$ satisfying $\Re z \in I, \Im z \neq 0$ and for all cubes $\Lambda \subset \mathbf{Z}^{d} ;(3.12)$ holds also for $\Gamma_{\omega}$.
(H3) The $\operatorname{IDS} N_{\Gamma}$ is differentiable at $E \in I$ with $n(E):=N_{\Gamma}^{\prime}(E)>0$.
By Minami's method in [34], we can show the following statement.

## THEOREM 3.3

Assume (H1), (H2), and (H3). Then, the point process $\xi_{\omega, E}^{\Lambda}(d t)$ converges weakly, as $|\Lambda| \rightarrow \infty$, to a Poisson point process with intensity measure $n(E) d t$.

## REMARK 6

The above point process $\xi_{\omega, E}^{\Lambda}$ converges weakly to a Poisson point process in the following sense. For any given disjoint bounded intervals $J_{1}, \ldots, J_{m}$ in $\mathbf{R}$, we have

$$
\lim _{|\Lambda| \rightarrow \infty} \mathbf{P}\left(\xi_{\omega, E}^{\Lambda}\left(J_{1}\right)=r_{1}, \ldots, \xi_{\omega, E}^{\Lambda}\left(J_{m}\right)=r_{m}\right)=\prod_{s=1}^{m} e^{-n(E)\left|J_{s}\right|} \frac{\left(n(E)\left|J_{s}\right|\right)^{r_{s}}}{r_{s}!}
$$

for all $r_{1}, \ldots, r_{m} \in \mathbf{N}$. For the reader's convenience, in Appendix C we present an example, and we prove that there exists an interval $I \subset \sigma\left(\Gamma_{\omega}\right)$ such that the condition (H2) is satisfied.

For the proof of Theorem 3.3, we use Minami's methods in [34] for our acoustic model $\Gamma_{\omega}$. We proceed as in [32, Appendix A]; hence, we need to prove two propositions for $\Gamma_{\omega}$.

We fix $\alpha \in(0,1)$ and divide $\Lambda=\Lambda_{L}(0)$ into $n_{L}$ boxes $\Lambda_{j}=\Lambda_{\ell}\left(k_{j}\right)$ of side $\ell \sim$ $L^{\alpha}$ centered at $k_{j} \in \mathbf{Z}^{d}$, that is, $\Lambda=\bigcup_{j=1}^{n_{L}} \Lambda_{j}$. Note that $n_{L}=|\Lambda| /\left|\Lambda_{j}\right| \sim L^{(1-\alpha) d}$.

For each $j$, we define point processes

$$
\xi_{\omega, E}^{\Lambda_{j}}(d t)=\sum_{k=1}^{\left|\Lambda_{j}\right|} \delta_{\left|\Lambda_{j}\right|\left(E_{\omega, k}^{\Lambda_{j}}-E\right)}(d t),
$$

where $E_{\omega, k}^{\Lambda_{j}}$ is the $k$ th eigenvalue of $\Gamma_{\omega}^{\Lambda_{j}}$. Note that from our assumptions it follows that $\left\{\xi_{\omega, E}^{\Lambda_{j}}\right\}_{j=1}^{n_{L}}$ are i.i.d. point processes. We consider the point process defined by the superposition

$$
\tilde{\xi}_{\omega, E}^{\Lambda}(d t)=\sum_{j=1}^{n_{L}} \xi_{\omega, E}^{\Lambda_{j}}(d t) .
$$

Moreover, we consider the random measure $\theta_{\omega, E}^{\Lambda}$ defined by

$$
\theta_{\omega, E}^{\Lambda}(J):=\operatorname{Tr} \chi_{\Lambda} \chi_{J}\left(|\Lambda|\left(\Gamma_{\omega}-E\right)\right) \chi_{\Lambda}
$$

for a Borel set $J \subset \mathbf{R}$. We first prepare the following proposition, which corresponds to [32, Proposition A.2] or [10, Lemma 6.1].

PROPOSITION 3.2
Assume (H1) and (H2). Then, for every $f \in \mathcal{C}_{0}^{\infty}(\mathbf{R}, d t)$, we have

$$
\begin{equation*}
\lim _{|\Lambda| \rightarrow \infty} \mathbf{E}\left[\left|\int f(t) \xi_{\omega, E}^{\Lambda}(d t)-\int f(t) \tilde{\xi}_{\omega, E}^{\Lambda}(d t)\right|\right]=0 \tag{3.13}
\end{equation*}
$$

and

$$
\begin{equation*}
\lim _{|\Lambda| \rightarrow \infty} \mathbf{E}\left[\left|\int f(t) \xi_{\omega, E}^{\Lambda}(d t)-\int f(t) \theta_{\omega, E}^{\Lambda}(d t)\right|\right]=0 \tag{3.14}
\end{equation*}
$$

To prove Proposition 3.2, we first recall a spectral averaging result for $\Gamma_{\omega}^{\Lambda}$. It is known that

$$
\begin{equation*}
\int\left\langle\delta_{\gamma}, \chi_{J}\left(\Gamma_{\omega}^{\Lambda}\right) \delta_{\gamma}\right\rangle \frac{d \omega_{\gamma}}{\omega_{\gamma}} \leq \int \chi_{J}(\lambda) \frac{d \lambda}{\lambda} \tag{3.15}
\end{equation*}
$$

for a Borel set $J \subset(0, \infty), \Lambda \subset \mathbf{Z}^{d}$, and $\gamma \in \Lambda$. We remark that (3.15) also gives another proof of Theorem 3.1 for the case when $\mu$ has bounded density. We also note that (3.15) holds for $\Gamma_{\omega}$ and $\gamma \in \mathbf{Z}^{d}$ (see [18], [20] for details). For every Borel set $J \subset \mathbf{R}$, it follows from (3.15) that

$$
\mathbf{E} \xi_{\omega, E}^{\Lambda}(J) \leq C|g|_{\infty}|J| .
$$

Similarly, we have

$$
\mathbf{E} \tilde{\xi}_{\omega, E}^{\Lambda}(J) \leq C|g|_{\infty}|J| \quad \text { and } \quad \mathbf{E} \theta_{\omega, E}^{\Lambda}(J) \leq C|g|_{\infty}|J| .
$$

We next recall the Helffer-Sjöstrand formula (see, e.g., [27, Appendix B] for details). For $f \in \mathcal{C}^{\infty}(\mathbf{R})$ and $n \in \mathbf{N}$, we set

$$
\{\{f\}\}_{n}:=\sum_{r=0}^{n} \int_{\mathbf{R}} d u\left|f^{(r)}(u)\right|\left(1+|u|^{2}\right)^{(r-1) / 2} .
$$

If $\{\{f\}\}_{n}<\infty$ for some $n \geq 2$, then for any self-adjoint operator $T$ we have

$$
\begin{equation*}
f(T)=\int_{\mathbf{R}^{2}} d \tilde{f}(z)(T-z)^{-1} \tag{3.16}
\end{equation*}
$$

where $z=x+i y, \tilde{f}(z)$ is an almost analytic extension of $f$ to the complex plane in the sense that it satisfies $\partial_{\bar{z}} \tilde{f}(z)=0$ for $z \in \mathbf{R}$, and $d \tilde{f}(z):=(1 / 2 \pi) \partial_{\bar{z}} \tilde{f}(z) d x d y$ with $\partial_{\bar{z}}:=\partial_{x}+i \partial_{y}$. Moreover, for all $p \geq 0$, we have

$$
\begin{equation*}
\int_{\mathbf{R}^{2}}|d \tilde{f}(z) \| \Im z|^{-p} \leq c_{p}\{\{f\}\}_{n}<\infty \tag{3.17}
\end{equation*}
$$

for $n \geq p+1$ with a constant $c_{p}$, where $|d \tilde{f}(z)|:=(1 / 2 \pi)\left|\partial_{\bar{z}} \tilde{f}(z)\right| d x d y$.
Proof of Proposition 3.2
We prove only (3.13) since (3.14) will be proven similarly. As in [32], let $v_{L}=$ $\beta \log L$, where $\beta>0$ is a fixed, big enough constant to be taken later. We set $\operatorname{int}\left(\Lambda_{j}\right):=\left\{x \in \Lambda_{j} \mid \operatorname{dist}\left(x, \partial \Lambda_{j}\right) \geq v_{L}\right\}$ and $\operatorname{wall}\left(\Lambda_{j}\right):=\Lambda_{j} \backslash \operatorname{int}\left(\Lambda_{j}\right)$. Then, we have

$$
\mathbf{E}\left[\left|\int f(t) d \xi_{\omega}^{\Lambda, E}(t)-\int f(t) d \tilde{\xi}_{\omega}^{\Lambda, E}(t)\right|\right] \leq \mathrm{I}_{1}+\mathrm{I}_{2}
$$

where

$$
\mathrm{I}_{1}=\sum_{j=1}^{n_{L}} \sum_{x \in \operatorname{wall}\left(\Lambda_{j}\right)} \mathbf{E}\left[\left|\left\langle\delta_{x}, f\left(|\Lambda|\left(\Gamma_{\omega}^{\Lambda}-E\right)\right) \delta_{x}\right\rangle-\left\langle\delta_{x}, f\left(|\Lambda|\left(\Gamma_{\omega}^{\Lambda_{j}}-E\right)\right) \delta_{x}\right\rangle\right|\right]
$$

and

$$
\mathrm{I}_{2}=\sum_{j=1}^{n_{L}} \sum_{x \in \operatorname{int}\left(\Lambda_{j}\right)} \mathbf{E}\left[\left|\left\langle\delta_{x}, f\left(|\Lambda|\left(\Gamma_{\omega}^{\Lambda}-E\right)\right) \delta_{x}\right\rangle-\left\langle\delta_{x}, f\left(|\Lambda|\left(\Gamma_{\omega}^{\Lambda_{j}}-E\right)\right) \delta_{x}\right\rangle\right|\right]
$$

To estimate $\mathrm{I}_{1}$, we use (3.15). Then, this yields that there is a constant $C>0$ independent of $L$ such that $\mathrm{I}_{1} \leq C|\Lambda|^{-1} \sum_{j=1}^{n_{L}} \mid$ wall $\left(\Lambda_{j}\right) \mid$, and this converges to zero as $L \rightarrow \infty$.

To estimate $\mathrm{I}_{2}$, we proceed as in [10]. By (3.16), $\mathrm{I}_{2}$ is dominated by

$$
\sum_{j=1}^{n_{L}} \sum_{x \in \operatorname{int}\left(\Lambda_{j}\right)}|\Lambda|^{-1} \int_{\mathbf{R}^{2}}|d \tilde{f}(z)| \mathbf{E}\left[\left|\left\langle\delta_{x},\left(\Gamma_{\omega}^{\Lambda}-z_{L}\right)^{-1} \delta_{x}\right\rangle-\left\langle\delta_{x},\left(\Gamma_{\omega}^{\Lambda_{j}}-z_{L}\right)^{-1} \delta_{x}\right\rangle\right|\right]
$$

where we set $z_{L}=E+z /|\Lambda|$. We recall the geometric resolvent identity for $\Gamma_{\omega}$ :

$$
\begin{align*}
& \left\langle\delta_{x},\left(\Gamma_{\omega}^{\Lambda}-z_{L}\right)^{-1} \delta_{x}\right\rangle-\left\langle\delta_{x},\left(\Gamma_{\omega}^{\Lambda_{j}}-z_{L}\right)^{-1} \delta_{x}\right\rangle \\
& \quad=\sum_{\left(y, y^{\prime}\right)} \frac{1}{\sqrt{\omega_{y} \omega_{y^{\prime}}}}\left\langle\delta_{x},\left(\Gamma_{\omega}^{\Lambda}-z_{L}\right)^{-1} \delta_{y}\right\rangle\left\langle\delta_{y^{\prime}},\left(\Gamma_{\omega}^{\Lambda_{j}}-z_{L}\right)^{-1} \delta_{x}\right\rangle \tag{3.18}
\end{align*}
$$

where the sum is over all pairs $\left(y, y^{\prime}\right)$ with $y \in \Lambda \backslash \Lambda_{j}, y^{\prime} \in \Lambda_{j}$, and $\left|y-y^{\prime}\right|=1$. From (3.18), $\omega_{\gamma} \geq m$, and the Hölder inequality, it follows that for $s \in(0,1), \mathrm{I}_{2}$
is dominated by

$$
\begin{align*}
& m^{-1}|\Lambda|^{-1} \sum_{j=1}^{n_{L}} \sum_{x \in \operatorname{int}\left(\Lambda_{j}\right)} \sum_{\left(y, y^{\prime}\right)} \int_{\mathbf{R}^{2}}|d \tilde{f}(z)|\left|\Im z_{L}\right|^{-2+s} \\
& \quad \times\left(\mathbf{E}\left|\left\langle\delta_{x},\left(\Gamma_{\omega}^{\Lambda}-z_{L}\right)^{-1} \delta_{y}\right\rangle\right|^{s}\right)^{1 / 2}\left(\mathbf{E}\left|\left\langle\delta_{y^{\prime}},\left(\Gamma_{\omega}^{\Lambda_{j}}-z_{L}\right)^{-1} \delta_{x}\right\rangle\right|^{s}\right)^{1 / 2}, \tag{3.19}
\end{align*}
$$

where we use the bound $\left|\left\langle\delta_{x},\left(\Gamma_{\omega}^{\Lambda}-z_{L}\right)^{-1} \delta_{y}\right\rangle\left\langle\delta_{y^{\prime}},\left(\Gamma_{\omega}^{\Lambda_{j}}-z_{L}\right)^{-1} \delta_{x}\right\rangle\right| \leq\left|\Im z_{L}\right|^{-2}$. Note that $\left|\Im z_{L}\right|^{-1}=|\Lambda||\Im z|^{-1}$ and that there are $O\left(L^{\alpha(d-1)}\right)$ pairs $\left(y, y^{\prime}\right)$ for each $\Lambda_{j}$. Then, by (H2) and (3.17), we have $\mathrm{I}_{2} \leq O\left(L^{d(2-s)+\alpha(d-1)} e^{-D v_{L}}\right)$. Hence, if we choose $v_{L}=\beta \log L$ and $\beta>D^{-1}(d(2-s)+\alpha(d-1)), \mathrm{I}_{2}$ converges to zero as $L \rightarrow \infty$.

Because of Proposition 3.2, it suffices to prove the following proposition to complete the proof of Theorem 3.3. This corresponds to [32, Proposition A.3].

PROPOSITION 3.3
Assume (H1), (H2), and (H3). Then, the point process $\tilde{\xi}_{\omega, E}^{\Lambda}(d t)$ converges weakly, as $|\Lambda| \rightarrow \infty$, to a Poisson point process with intensity measure $n(E) d t$.

## Proof

By standard results from the theory of point processes (see, e.g., [14, Theorem 9.2.V], [15, Theorem 11.2.V], [32, Theorem 2.3]), the weak convergence of $\tilde{\xi}_{\omega, E}^{\Lambda}$ to the Poisson point process is equivalent to verifying the following three conditions for all bounded intervals $I$ :

$$
\begin{gather*}
\lim _{L \rightarrow \infty} \max _{j=1, \ldots, n_{L}} \mathbf{P}\left(\xi_{\omega, E}^{\Lambda_{j}}(I) \geq 1\right)=0  \tag{3.20}\\
\lim _{L \rightarrow \infty} \sum_{j=1}^{n_{L}} \mathbf{P}\left(\xi_{\omega, E}^{\Lambda_{j}}(I) \geq 1\right)=n(E)|I|  \tag{3.21}\\
\lim _{L \rightarrow \infty} \sum_{j=1}^{n_{L}} \mathbf{P}\left(\xi_{\omega, E}^{\Lambda_{j}}(I) \geq 2\right)=0 \tag{3.22}
\end{gather*}
$$

Since $\xi_{\omega, E}^{\Lambda_{j}}(I)=\operatorname{Tr} \chi_{E+I /|\Lambda|}\left(\Gamma_{\omega}^{\Lambda_{j}}\right)$, (3.20) follows immediately from Theorem 3.1. In addition, from Theorem 3.2 we have $\sum_{j=1}^{n_{L}} \mathbf{P}\left(\xi_{\omega, E}^{\Lambda_{j}}(I) \geq 2\right) \leq C|I|^{2} n_{L}^{-1}$; hence, (3.22) follows. Thus, the proof is finished if we verify the condition (3.21). By using Theorem 3.2 with $n=2$, we first have

$$
\lim _{L \rightarrow \infty} \sum_{j=1}^{n_{L}} \mathbf{P}\left(\xi_{\omega, E}^{\Lambda_{j}}(I) \geq 1\right)=\lim _{L \rightarrow \infty} \mathbf{E}\left[\tilde{\xi}_{\omega, E}^{\Lambda}(I)\right]
$$

By (3.14) in Proposition 3.2, we next have

$$
\lim _{L \rightarrow \infty} \mathbf{E}\left[\tilde{\xi}_{\omega, E}^{\Lambda}(I)\right]=\lim _{L \rightarrow \infty} \mathbf{E}\left[\theta_{\omega, E}^{\Lambda}(I)\right]=\lim _{L \rightarrow \infty}|\Lambda| \int_{E+I /|\Lambda|} n(\lambda) d \lambda
$$

Finally, by (H3) and Lebesgue differentiation theorem, we complete the proof (see [10], [32], [34] for details).

## Appendix A: Simple proof of Lifshitz tails for $\Gamma_{\omega}$

In this appendix, as an application of Proposition 3.1, we discuss an asymptotic behavior of the IDS for $\Gamma_{\omega}$. We always assume (3.1), that is, $H_{0}=-\Delta$. To guarantee the existence of the IDS, we also assume (2.4), that is, that $\left\{\omega_{\gamma}\right\}_{\gamma \in \mathbf{Z}^{d}}$ are i.i.d. by $\mu$. Recall $\operatorname{supp} \mu \subset[m, M]$ and $m, M \in \operatorname{supp} \mu$. Then we have $\sigma\left(\Gamma_{\omega}\right)=\left[0, E_{\infty}\right]$ almost surely, where $E_{\infty}=4 d / m$. Our aim is to show

$$
1-N_{\Gamma}(E)=c_{1} \exp \left(-\left(c_{2}+o(1)\right)\left(E_{\infty}-E\right)^{-d / 2}\right)
$$

as $E \uparrow E_{\infty}$, which is known as Lifshitz tails. See [8], [28], [37] for the Anderson model. Note that the internal Lifshitz tails (see, e.g., [30]) for the continuous version of $A_{\omega}$ were studied by Najar [35], [36].

THEOREM A. 1
The IDS $N_{\Gamma}$ satisfies

$$
\begin{equation*}
\limsup _{\epsilon \downarrow 0} \frac{\log \left(-\log \left(1-N_{\Gamma}\left(E_{\infty}-\epsilon\right)\right)\right)}{\log \epsilon} \leq-\frac{d}{2} . \tag{A.1}
\end{equation*}
$$

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Moreover, if we assume that $\mu([m, m+\kappa]) \geq c \kappa^{\delta}$ for some $c, \delta>0$ and sufficiently small $\kappa>0$, then we can use the lower bound in Proposition 3.1 to get

$$
\begin{equation*}
\liminf _{\epsilon \downarrow 0} \frac{\log \left(-\log \left(1-N_{\Gamma}\left(E_{\infty}-\epsilon\right)\right)\right)}{\log \epsilon} \geq-\frac{d}{2} . \tag{A.2}
\end{equation*}
$$

Hence, by (A.1) and (A.2), we can determine Lifshitz exponent

$$
\lim _{\epsilon \downarrow 0} \frac{\log \left(-\log \left(1-N_{\Gamma}\left(E_{\infty}-\epsilon\right)\right)\right)}{\log \epsilon}=-\frac{d}{2} .
$$

Proof of Theorem A. 1
Choose $E=E_{\infty}$ in (3.7). Then we have

$$
N_{H}(m \epsilon)-N_{H}(-m \epsilon) \leq N_{\Gamma}\left(E_{\infty}+\epsilon\right)-N_{\Gamma}\left(E_{\infty}-\epsilon\right) \leq N_{H}(M \epsilon)-N_{H}(-M \epsilon),
$$

where $H_{\omega}\left(E_{\infty}\right)=-\Delta-E_{\infty} \rho_{\omega}$ and $N_{H}$ is the corresponding IDS. This auxiliary Schrödinger operator $H_{\omega}\left(E_{\infty}\right)$ is a nonpositive, bounded, and $\mathbf{Z}^{d}$-ergodic selfadjoint operator, and we know that $\sigma\left(H_{\omega}\left(E_{\infty}\right)\right)=\left[-E_{\infty} M, 0\right]$ almost surely. By $N_{\Gamma}\left(E_{\infty}+\epsilon\right)=N_{H}(m \epsilon)=N_{H}(M \epsilon)=1$, we get

$$
1-N_{H}(-m \epsilon) \leq 1-N_{\Gamma}\left(E_{\infty}-\epsilon\right) \leq 1-N_{H}(-M \epsilon)
$$

This implies

$$
\begin{equation*}
\limsup _{\epsilon \downarrow 0} \frac{\log \left(-\log \left(1-N_{\Gamma}\left(E_{\infty}-\epsilon\right)\right)\right)}{\log \epsilon} \tag{A.3}
\end{equation*}
$$

$$
=\underset{\epsilon \downarrow 0}{\limsup } \frac{\log \left(-\log \left(1-N_{H}(-\epsilon)\right)\right)}{\log \epsilon} .
$$

Since $H_{\omega}\left(E_{\infty}\right)$ is an Anderson model, the right-hand side of (A.3) can be estimated by standard arguments (see, e.g., [8, Theorem VI.2.7], [28]).

## Appendix B: Relation between $A_{\omega}$ and $\Gamma_{\omega}$

In this appendix, let us consider the simplest 1-dimensional model. Our aim is to compare $A_{\omega}:=\partial^{*}\left(1 / \rho_{\omega}\right) \partial$ with $\Gamma_{\omega}:=\left(1 / \sqrt{\rho_{\omega}}\right)\left(\partial^{*} \partial\right)\left(1 / \sqrt{\rho_{\omega}}\right) ; 1 / \rho_{\omega}$ and $1 / \sqrt{\rho_{\omega}}$ are defined as before. For the sake of simplicity, we assume (2.4), that is, that $\left\{\omega_{\gamma}\right\}_{\gamma \in \mathbf{Z}}$ are i.i.d. by $\mu$. Hence, we know that $\sigma\left(A_{\omega}\right)=\sigma\left(\Gamma_{\omega}\right)=\left[0, E_{\infty}\right]$ almost surely, where $E_{\infty}=4 / \mathrm{m}$.

In the 1-dimensional case, each operator is represented by a matrix: $A_{\omega}=$ $T^{*} T$ and $\Gamma_{\omega}=T T^{*}$. The matrix $T$ is defined by $T=\left(1 / \sqrt{\rho_{\omega}}\right) \partial$, where

$$
\frac{1}{\sqrt{\rho_{\omega}}}=\left(\begin{array}{llll}
\ddots & & & 0 \\
& 1 / \sqrt{\omega_{0}} & & \\
& & 1 / \sqrt{\omega_{1}} & \\
0 & & & \ddots
\end{array}\right) \quad \text { and } \quad \partial=\left(\begin{array}{llll}
\ddots & & & 0 \\
\ddots & 1 & & \\
& -1 & 1 & \\
& & -1 & \ddots \\
0 & & & \ddots
\end{array}\right) .
$$

We note that the 1-dimensional discrete Laplacian (3.2) satisfies $-\Delta=\partial^{*} \partial=$ $\partial \partial^{*}$. For $L \in \mathbf{N}$, set $\Lambda:=\{k \in \mathbf{Z}| | k \mid \leq L\} \subset \mathbf{Z}$. We first define the $((2 L+1) \times$ $(2 L+1))$-matrix by $T^{\Lambda}=\left(1 / \sqrt{\rho_{\omega}^{\Lambda}}\right) \partial^{\Lambda}$, where
$\frac{1}{\sqrt{\rho_{\omega}^{\Lambda}}}=\left(\begin{array}{ccc}1 / \sqrt{\omega_{-L}} & & 0 \\ & \ddots & \\ 0 & & 1 / \sqrt{\omega_{L}}\end{array}\right) \quad$ and $\quad \partial^{\Lambda}=\left(\begin{array}{cccc}1 & & & -1 \\ -1 & \ddots & 0 & \\ & \ddots & \ddots & \\ 0 & & -1 & 1\end{array}\right)$.
We imposed the periodic boundary condition since we need to guarantee $\partial^{\Lambda}\left(\partial^{\Lambda}\right)^{*}=\left(\partial^{\Lambda}\right)^{*} \partial^{\Lambda}=-\Delta^{\Lambda}$, where $-\Delta^{\Lambda}$ is the discrete Laplacian (3.2) with the periodic boundary condition, that is,

$$
-\Delta^{\Lambda}=\left(\begin{array}{ccccc}
2 & -1 & & & -1 \\
-1 & 2 & -1 & 0 & \\
& \ddots & \ddots & \ddots & \\
& 0 & -1 & 2 & -1 \\
-1 & & & -1 & 2
\end{array}\right)
$$

Now, we quote a lemma in [16, Theorem 2].

## LEMMA B. 1

Let $\mathcal{H}_{1}$ and $\mathcal{H}_{2}$ be two separable Hilbert spaces. Let $T$ be a densely defined closed operator from $\mathcal{H}_{1}$ to $\mathcal{H}_{2}$. Then, $T^{*} T$ on $\mathcal{H}_{1}$ and $T T^{*}$ on $\mathcal{H}_{2}$ are nonnegative
self-adjoint operators. These satisfy

$$
\sigma_{p}\left(T^{*} T\right) \backslash\{0\}=\sigma_{p}\left(T T^{*}\right) \backslash\{0\} .
$$

Moreover, for each $\lambda \in \sigma_{p}\left(T^{*} T\right) \backslash\{0\}$, it holds that

$$
\operatorname{dim} \operatorname{ker}\left(T^{*} T-\lambda\right)=\operatorname{dim} \operatorname{ker}\left(T T^{*}-\lambda\right)
$$

We first have the following statements.

## THEOREM B. 1

Fix $E_{0} \in(0, \infty)$. Let I be any bounded interval in $\left[E_{0}, \infty\right)$. Then we have

$$
\mathbf{E} \operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right)=\mathbf{E} \operatorname{Tr} \chi_{I}\left(\Gamma_{\omega}^{\Lambda}\right) \leq Q(|I|)|\Lambda|
$$

and, for all $n \in \mathbf{N}$,

$$
\mathbf{P}\left(\operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right) \geq n\right)=\mathbf{P}\left(\operatorname{Tr} \chi_{I}\left(\Gamma_{\omega}^{\Lambda}\right) \geq n\right) \leq \frac{1}{n!}(Q(|I|)|\Lambda|)^{n},
$$

where

$$
Q(|I|)= \begin{cases}|g|_{\infty}|I| M / E_{0} & \text { if } \mu \text { has a bounded density } g \\ 4\left(1+M / E_{0}\right) S_{\mu}(|I|) & \text { otherwise. }\end{cases}
$$

Proof
Suppose $0 \notin I$, and use Lemma B.1. Then we have

$$
\begin{equation*}
\operatorname{Tr} \chi_{I}\left(A_{\omega}^{\Lambda}\right)=\operatorname{Tr} \chi_{I}\left(\Gamma_{\omega}^{\Lambda}\right), \tag{B.1}
\end{equation*}
$$

where $A_{\omega}^{\Lambda}=\left(T^{\Lambda}\right)^{*} T^{\Lambda}$ and $\Gamma_{\omega}^{\Lambda}=T^{\Lambda}\left(T^{\Lambda}\right)^{*}$. Thus, because of (B.1), as in the proof of Theorem 3.1 and Theorem 3.2, we can analyze $\Gamma_{\omega}^{\Lambda}$.

We next have the following Lifshitz tail for $A_{\omega}$.

## THEOREM B. 2

The IDS $N_{A}$ in the 1-dimensional case satisfies

$$
\limsup _{\epsilon \downarrow 0} \frac{\log \left(-\log \left(1-N_{A}\left(E_{\infty}-\epsilon\right)\right)\right)}{\log \epsilon} \leq-\frac{1}{2} .
$$

Proof of Theorem B. 2
From (B.1), it follows that

$$
N_{A}(E+\epsilon)-N_{A}(E-\epsilon)=N_{\Gamma}(E+\epsilon)-N_{\Gamma}(E-\epsilon) .
$$

Therefore, we choose $E=E_{\infty}$, and we obtain

$$
\begin{equation*}
\frac{\log \left(-\log \left(1-N_{A}\left(E_{\infty}-\epsilon\right)\right)\right)}{\log \epsilon}=\frac{\log \left(-\log \left(1-N_{\Gamma}\left(E_{\infty}-\epsilon\right)\right)\right)}{\log \epsilon} \tag{B.2}
\end{equation*}
$$

Use Theorem A. 1 and (B.2) to complete the proof.

For the multidimensional case, it seems to be more difficult to apply the same approach via Lemma B. 1 since we need to study a random operator on the space of vector-valued functions. Recall that the Maxwell operator $M_{\omega}$ is a random operator defined on $\ell^{2}\left(\mathbf{Z}^{3} ; \mathbf{C}^{3}\right)$. As $T=\left(1 / \sqrt{\rho_{\omega}}\right) \nabla$ is the operator from $\ell^{2}\left(\mathbf{Z}^{d} ; \mathbf{C}\right)$ to $\ell^{2}\left(\mathbf{Z}^{d} ; \mathbf{C}^{d}\right)$, then we will have to analyze the following random operator: $T T^{*}=$ $\left(1 / \sqrt{\tilde{\rho}_{\omega}}\right) \tilde{H}_{0}\left(1 / \sqrt{\tilde{\rho}_{\omega}}\right)$, where

$$
\tilde{H}_{0}=\nabla \nabla^{*}=\left(\begin{array}{cccc}
\partial_{1} \partial_{1}^{*} & \partial_{1} \partial_{2}^{*} & \ldots & \partial_{1} \partial_{d}^{*} \\
\partial_{2} \partial_{1}^{*} & \partial_{2} \partial_{2}^{*} & \ldots & \partial_{2} \partial_{d}^{*} \\
& & \ddots & \\
\partial_{d} \partial_{1}^{*} & \partial_{d} \partial_{2}^{*} & \ldots & \partial_{d} \partial_{d}^{*}
\end{array}\right)
$$

and $1 / \sqrt{\tilde{\rho}_{\omega}}:=\sum_{\gamma}\left(1 / \sqrt{\omega_{\gamma}}\right) \tilde{\Pi}_{\gamma}$, where $\tilde{\Pi}_{\gamma}=\sum_{j=1}^{d}\left|\delta_{\gamma}^{j}\right\rangle\left\langle\delta_{\gamma}^{j}\right|$ for the standard orthonormal basis $\left\{\delta_{\gamma}^{j}\right\}_{\gamma \in \mathbf{Z}^{d}}^{j=1, \ldots, d}$ of $\ell^{2}\left(\mathbf{Z}^{d} ; \mathbf{C}^{d}\right)$. Hence, if we proceed as in Proposition 3.1, we need to study the random operator $\tilde{H}_{\omega}(E)=\tilde{H}_{0}-E \tilde{\rho}_{\omega}$ on $\ell^{2}\left(\mathbf{Z}^{d} ; \mathbf{C}^{d}\right)$. If $d \geq 2$, it is not clear how the $\operatorname{IDS}$ for $\tilde{H}_{\omega}(E)$ behaves near the edge of its spectrum. Note that $\tilde{H}_{0}$ is not elliptic if $d \geq 2$. However, our techniques in this paper work for $\tilde{H}_{\omega}^{\Lambda}(E):=\tilde{\chi}_{\Lambda} \tilde{H}_{\omega}(E) \tilde{\chi}_{\Lambda}$, where $\tilde{\chi}_{\Lambda}:=\sum_{\gamma \in \Lambda} \tilde{\Pi}_{\gamma}$, since the rank of $\tilde{\Pi}_{\gamma}$ is just $d$ and we can use Proposition 2.1. Therefore, we can obtain the eigenvalue-counting estimates for $A_{\omega}^{\Lambda}:=\nabla^{*}\left(1 / \sqrt{\rho_{\omega}}\right) \chi_{\Lambda}\left(1 / \sqrt{\rho_{\omega}}\right) \nabla$ corresponding to $\tilde{H}_{\omega}^{\Lambda}(E)$.

## Appendix C: Fractional moment bound and Localization

In this appendix, we check the condition (H2) in Theorem 3.3. Let us consider (1.2) with (3.1), that is, $\Gamma_{\omega}$ with $H_{0}=-\Delta$. We assume that (2.4) and its common probability measure $\mu$ are distributed uniformly on $[m, m+1]$ for some small parameter $m>0$. We know that $\sigma\left(\Gamma_{\omega}\right)=\left[0, E_{\infty}\right]$ almost surely, where $E_{\infty}=$ $4 d / m$. Since we proved the Wegner estimate (3.3) and the Lifshitz tail (A.1) near $E_{\infty}=4 d / m$, we can show that Anderson localization occurs near $E_{\infty}$ by using the multiscale analysis as in [18]-[20]. However, another method is well known. This is called fractional moment analysis, which was started by Aizenman and Molchanov [1], [3]. For continuous models, refer to [2] and [7]. By this approach, we will show the following exponential decay estimate for the fractional moment of the Green's function $G:=\left(\Gamma_{\omega}-z\right)^{-1}(x, y)$. The first key idea is to look at powers of $|G|^{s}$ with $0<s<1$. For the sake of simplicity, we take $s=1 / 2$. We proceed as in [38, Chapter 13] (see also [17, Section 5.1.2]).

## THEOREM C. 1

There exist finite constants $c, C, D>0$ such that

$$
\begin{equation*}
\mathbf{E}\left|\left\langle\delta_{x},\left(\Gamma_{\omega}^{\Lambda}-z\right)^{-1} \delta_{y}\right\rangle\right|^{1 / 2} \leq C e^{-D|x-y|} \tag{C.1}
\end{equation*}
$$

for all $z \in \mathbf{C}$ satisfying $|z| \geq 16 d^{2} / c^{2}$ and $\Im z \neq 0$, for all cube $\Lambda \subset \mathbf{Z}^{d}$, and for all $x, y \in \Lambda$. Note that $c$ is a universal constant, and $C, D$ depend only on $d, m$.

## REMARK 8

(C.1) holds also for $\Gamma_{\omega}$.

## Proof of Theorem C. 1

By

$$
\begin{equation*}
\left\langle\delta_{x},\left(\Gamma_{\omega}^{\Lambda}-z\right)^{-1} \delta_{y}\right\rangle=\sqrt{\omega_{x} \omega_{y}}\left\langle\delta_{x},\left(-\Delta-z \rho_{\omega}^{\Lambda}\right)^{-1} \delta_{y}\right\rangle, \tag{C.2}
\end{equation*}
$$

it is enough to study $\left\langle\delta_{x},\left(-\Delta-z \rho_{\omega}^{\Lambda}\right)^{-1} \delta_{y}\right\rangle$. Since we can write $-\Delta=2 d I-h_{0}$, where $I$ is the identity operator and $h_{0}$ is defined by $\left(h_{0} f\right)(i)=\sum_{|j|=1} f(i+j)$, then we get

$$
\begin{equation*}
\left\langle\delta_{x}, \delta_{y}\right\rangle=\left(2 d-z \omega_{x}\right)\left\langle\delta_{x},\left(-\Delta-z \rho_{\omega}^{\Lambda}\right)^{-1} \delta_{y}\right\rangle-\left\langle h_{0} \delta_{x},\left(-\Delta-z \rho_{\omega}^{\Lambda}\right)^{-1} \delta_{y}\right\rangle . \tag{C.3}
\end{equation*}
$$

By the resolvent identity, there are $\alpha, \beta \in \mathbf{C}$, which are independent of $\omega_{\gamma}$, such that

$$
\begin{equation*}
\left\langle\delta_{x},\left(-\Delta-z \rho_{\omega}^{\Lambda}\right)^{-1} \delta_{y}\right\rangle=\frac{\alpha}{\beta-z \omega_{x}} . \tag{C.4}
\end{equation*}
$$

By the decoupling lemma (see [38, Lemma 13.3], [17, Lemma 5.1.14]), there exists a universal constant $c>0$ such that

$$
\begin{align*}
& \int_{m}^{m+1}\left|\left(2 d-z \omega_{x}\right) \frac{\alpha}{\beta-z \omega_{x}}\right|^{1 / 2} d \omega_{x} \\
& \quad \geq c|z|^{1 / 2} \int_{m}^{m+1}\left|\frac{\alpha}{\beta-z \omega_{x}}\right|^{1 / 2} d \omega_{x} \tag{C.5}
\end{align*}
$$

Note that $(a+b)^{s} \leq a^{s}+b^{s}$ if $s \in(0,1)$. By (C.3), (C.4), and (C.5), we get

$$
\begin{align*}
& \mathbf{E}\left|\left\langle\delta_{x},\left(-\Delta-z \rho_{\omega}^{\Lambda}\right)^{-1} \delta_{y}\right\rangle\right|^{1 / 2} \\
& \quad \leq\left(c|z|^{1 / 2}\right)^{-1}\left(\left\langle\delta_{x}, \delta_{y}\right\rangle+\sum_{|j|=1} \mathbf{E}\left|\left\langle\delta_{x+j},\left(-\Delta-z \rho_{\omega}^{\Lambda}\right)^{-1} \delta_{y}\right\rangle\right|^{1 / 2}\right) . \tag{C.6}
\end{align*}
$$

Now, we set $\delta_{y}(x):=\left\langle\delta_{x}, \delta_{y}\right\rangle$ and $\eta(x):=\mathbf{E}\left|\left\langle\delta_{x},\left(-\Delta-z \rho_{\omega}^{\Lambda}\right)^{-1} \delta_{y}\right\rangle\right|^{1 / 2}$. By (C.6), we have

$$
\left(1-\left(c|z|^{1 / 2}\right)^{-1} h_{0}\right) \eta(x) \leq\left(c|z|^{1 / 2}\right)^{-1} \delta_{y}(x) .
$$

Note that $0<2 d\left(c|z|^{1 / 2}\right)^{-1}<1$. Since $\left(1-\left(c|z|^{1 / 2}\right)^{-1} h_{0}\right)^{-1}$ is positivity preserving (see [38, Lemma 13.4]), thus, we get

$$
\begin{aligned}
\eta(x) & \leq\left(c|z|^{1 / 2}\right)^{-1}\left(1-\left(c|z|^{1 / 2}\right)^{-1} h_{0}\right)^{-1} \delta_{y}(x) \\
& \leq\left(c|z|^{1 / 2}-2 d\right)^{-1}\left(2 d\left(c|z|^{1 / 2}\right)^{-1}\right)^{|x-y|} .
\end{aligned}
$$

By (C.2) and $\omega_{\gamma} \leq m+1$, we obtain

$$
\mathbf{E}\left|\left\langle\delta_{x},\left(\Gamma_{\omega}^{\Lambda}-z\right)^{-1} \delta_{y}\right\rangle\right|^{1 / 2} \leq C e^{-D|x-y|} .
$$

Finally, let us state a relation between the localization and Theorem C.1.

## THEOREM C. 2

Assume that the disorder parameter $m$ satisfies $0<m<c^{2} / 4 d$. Then $\Gamma_{\omega}$ has only pure point spectrum in $\left[16 d^{2} / c^{2}, E_{\infty}\right]$ almost surely.

Proof
We set $\left\langle\delta_{0},\left(\Gamma_{\omega}-E-i 0\right)^{-2} \delta_{0}\right\rangle:=\lim _{\epsilon \backslash 0}\left\langle\delta_{0},\left(\Gamma_{\omega}-E-i \epsilon\right)^{-2} \delta_{0}\right\rangle$. Theorem 3 in [18], which is called the Simon-Wolff criterion (see [39]), implies that if, with probability one,

$$
\begin{equation*}
\left\langle\delta_{0},\left(\Gamma_{\omega}-E-i 0\right)^{-2} \delta_{0}\right\rangle<\infty \tag{C.7}
\end{equation*}
$$

for almost every $E \in I$, then, with probability one, $\Gamma_{\omega}$ has only pure point spectrum in $I$. By using (C.1), we can check (C.7). Set $I=\left[16 d^{2} / c^{2}, E_{\infty}\right]$. For $z=E+i \epsilon$ with $E \in I$ and $\epsilon>0$, we get

$$
\begin{align*}
& \mathbf{E}\left|\left\langle\delta_{0},\left(\Gamma_{\omega}-E-i \epsilon\right)^{-2} \delta_{0}\right\rangle\right|^{1 / 4}  \tag{C.8}\\
& \quad \leq \sum_{x \in \mathbf{Z}^{d}} \mathbf{E}\left[\left|\left\langle\delta_{0},\left(\Gamma_{\omega}-E-i \epsilon\right)^{-1} \delta_{x}\right\rangle\right|^{1 / 4}\left|\left\langle\delta_{x},\left(\Gamma_{\omega}-E-i \epsilon\right)^{-1} \delta_{0}\right\rangle\right|^{1 / 4}\right] .
\end{align*}
$$

By using the Hölder inequality, (C.8) is dominated by

$$
\sum_{x \in \mathbf{Z}^{d}}\left(\mathbf{E}\left|\left\langle\delta_{0},\left(\Gamma_{\omega}-E-i \epsilon\right)^{-1} \delta_{x}\right\rangle\right|^{1 / 2}\right)^{1 / 2}\left(\mathbf{E}\left|\left\langle\delta_{x},\left(\Gamma_{\omega}-E-i \epsilon\right)^{-1} \delta_{0}\right\rangle\right|^{1 / 2}\right)^{1 / 2}
$$

By (C.1) for $\Gamma_{\omega}$, we know that $\mathbf{E}\left|\left\langle\delta_{0},\left(\Gamma_{\omega}-E-i \epsilon\right)^{-2} \delta_{0}\right\rangle\right|^{1 / 4}$ is bounded. This completes the proof.

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Graduate School of Human and Environmental Studies, Kyoto University, Kyoto 606-8501, Japan; ykitagak@math.h.kyoto-u.ac.jp

