# MULTI-PARTICLE DYNAMICAL LOCALIZATION IN A CONTINUOUS ANDERSON MODEL WITH AN ALLOY-TYPE POTENTIAL 

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#### Abstract

This paper is a complement to our earlier work 4. With the help of the multi-scale analysis, we derive, from estimates obtained in [4], dynamical localization for a multi-particle Anderson model in a Euclidean space $\mathbb{R}^{d}, d \geq 1$, with a short-range interaction, subject to a random alloy-type potential.


## 1. Introduction

1.1. The model. In this paper we continue our study of a multi-particle Anderson model in $\mathbb{R}^{d}$ with interaction and in an external random potential of alloy type. The Hamiltonian $\mathbf{H}\left(=\mathbf{H}^{(N)}(\omega)\right)$ is a random Schrödinger operator of the form

$$
\begin{equation*}
\mathbf{H}=-\frac{1}{2} \boldsymbol{\Delta}+\mathbf{U}(\mathbf{x})+\mathbf{V}(\omega ; \mathbf{x}) \tag{1.1}
\end{equation*}
$$

acting in $L^{2}\left(\mathbb{R}^{N d}\right)$. This means that we consider a system of $N$ interacting quantum particles in $\mathbb{R}^{d}$. Here $\mathbf{x}=\left(x_{1}, \ldots, x_{N}\right) \in \mathbb{R}^{N d}$ is for the joint position vector, where each component $x_{j} \in \mathbb{R}^{d}$ represents the position of the $j$ th particle, $1 \leq j \leq N$. Next, $\boldsymbol{\Delta}$ stands for the Laplacian in $\mathbb{R}^{N d}$. The interaction energy operator $\mathbf{U}(\mathbf{x})$ acts as multiplication by a function $U(\mathbf{x})$. Finally, the term $\mathbf{V}(\omega ; \mathbf{x})$ represents the operator of multiplication by a function

$$
\begin{equation*}
\mathbf{x} \mapsto V\left(x_{1} ; \omega\right)+\cdots+V\left(x_{N} ; \omega\right), \tag{1.2}
\end{equation*}
$$

where $x \in \mathbb{R}^{d} \mapsto V(x ; \omega)$ is a random external field potential assumed to be of the form

$$
\begin{equation*}
V(x ; \omega)=\sum_{s \in \mathbb{Z}^{d}} \mathrm{~V}_{s}(\omega) \varphi(x-s) . \tag{1.3}
\end{equation*}
$$

Here and below $\mathrm{V}_{s}, s \in \mathbb{Z}^{d}$, are i.i.d. (independent and identically distributed) real random variables on some probability space $(\Omega, \mathfrak{B}, \mathbb{P})$ and $\varphi: \mathbb{R}^{d} \rightarrow \mathbb{R}$ is usually referred to as a "bump" function.
1.2. Basic geometric notations. Throughout this paper, we will fix an integer $N \geq 2$ and work in Euclidean spaces of the form $\mathbb{R}^{l d} \cong \mathbb{R}^{d} \times \ldots \times \mathbb{R}^{d}(l$ times $)$ associated with $l$-particle sub-systems where $1 \leq l \leq N$. Correspondingly, the notations $\mathbf{x}$, $\mathbf{y}, \ldots$ will be used for vectors from $\mathbb{R}^{l d}$, depending on the context. Given a vector $\mathbf{x} \in \mathbb{R}^{l d}$, we will consider "sub-configurations" $\mathrm{x}^{\prime}$ and $\mathrm{x}^{\prime \prime}$ generated by x for a given partition of an $l$-particle system into disjoint sub-systems with $l^{\prime}$ and $l^{\prime \prime}$ particles, where $l^{\prime}+l^{\prime \prime}=l, l^{\prime}, l^{\prime \prime} \geq 1$; the vectors $\mathrm{x}^{\prime}$ and $\mathrm{x}^{\prime \prime}$ are identified with points from $\mathbb{R}^{l^{\prime} d}$ and $\mathbb{R}^{l^{\prime \prime} d}$, respectively, by re-labelling the particles accordingly.

All Euclidean spaces will be endowed with the max-norm denoted by $|\cdot|$. We will consider $l d$-dimensional cubes of integer size in $\mathbb{R}^{l d}$ centered at lattice points $\mathbf{u} \in$ $\mathbb{Z}^{l d} \subset \mathbb{R}^{l d}$ and with edges parallel to the co-ordinate axes. The cube of edge length $2 L$ centered at $\mathbf{u}$ is denoted by $\boldsymbol{\Lambda}_{L}(\mathbf{u})$; in the max-norm it represents the ball of radius $L$ centered at $\mathbf{u}$ :

$$
\begin{equation*}
\boldsymbol{\Lambda}_{L}(\mathbf{u})=\left\{\mathbf{x} \in \mathbb{R}^{l d}:|\mathbf{x}-\mathbf{u}|<L\right\} . \tag{1.4}
\end{equation*}
$$

The lattice counterpart for $\boldsymbol{\Lambda}_{L}(\mathbf{u})$ is denoted by $\mathbf{B}_{L}(\mathbf{u})$ :

$$
\mathbf{B}_{L}(\mathbf{u})=\overline{\boldsymbol{\Lambda}}_{L}(\mathbf{u}) \cap \mathbb{Z}^{l d} ; \quad \mathbf{u} \in \mathbb{Z}^{l d}
$$

Finally, we consider "cells" (cubes of radius 1) centered at lattice points $\mathbf{u} \in \mathbb{Z}^{l d}$ :

$$
\mathbf{C}(\mathbf{u})=\boldsymbol{\Lambda}_{1}(\mathbf{u}) \subset \mathbb{R}^{l d}
$$

The union of all cells $\mathbf{C}(\mathbf{u}), \mathbf{u} \in \mathbb{Z}^{l d}$, covers the entire Euclidean space $\mathbb{R}^{l d}$. For each $i \in\{1, \ldots, l\}$ we introduce the projection $\Pi_{i}: \mathbb{R}^{l d} \rightarrow \mathbb{R}^{d}$ defined by

$$
\Pi_{i}:\left(x_{1}, \ldots, x_{l}\right) \longmapsto x_{i}, \quad 1 \leq i \leq l .
$$

1.3. Interaction potential. The interaction within the system of particles is represented by the term $\mathbf{U}(\mathbf{x})$ in the expression (1.1) of the Hamiltonian H. As was said, it is the operator of multiplication by a function $\mathbf{x} \in \mathbb{R}^{l d} \mapsto U(\mathbf{x}) \in \mathbb{R}, 1 \leq l \leq N$. A usual assumption is that $U(\mathbf{x})$ (considered for $\mathbf{x} \in \mathbb{R}^{l d}$ with $1 \leq l \leq N$ ) is a sum of $k$-body potentials

$$
U(\mathbf{x})=\sum_{k=1}^{l} \sum_{1 \leq i_{1}<\ldots<i_{k} \leq l} U^{(k)}\left(x_{i_{1}}, \ldots, x_{i_{k}}\right), \quad \mathbf{x}=\left(x_{1}, \ldots, x_{l}\right) \in \mathbb{R}^{l d}
$$

In this paper we do not assume isotropy, symmetry or translation invariance of this interaction. However, we use the conditions of finite range, nonnegativity and boundedness, as stated below.

Assume a partition of a configuration $\mathbf{x} \in \mathbb{Z}^{l d}$ is given, into complementary subconfigurations $\mathbf{x}_{\mathcal{J}}=\left(x_{j}\right)_{j \in \mathcal{J}}$ and $\mathbf{x}_{\mathcal{J} c}=\left(x_{j}\right)_{j \in\{1, \ldots, l\} \backslash \mathcal{J}}$, where $\varnothing \neq \mathcal{J} \subsetneq\{1,2, \ldots, l\}$. The energy of interaction between $\mathbf{x}_{\mathcal{J}}$ and $\mathbf{x}_{\mathcal{J}}$ is defined by

$$
\begin{equation*}
U\left(\mathbf{x}_{\mathcal{J}} \mid \mathbf{x}_{\mathcal{J}^{c}}\right):=U(\mathbf{x})-U\left(\mathbf{x}_{\mathcal{J}}\right)-U\left(\mathbf{x}_{\mathcal{J}^{c}}\right) \tag{1.5}
\end{equation*}
$$

Next, define

$$
\begin{equation*}
\rho\left(\mathbf{x}_{\mathcal{J}}, \mathbf{x}_{\mathcal{J}^{c}}\right):=\min \left[\left|x_{i}-x_{j}\right|: i \in \mathcal{J}, j \in \mathcal{J}^{\mathrm{c}}\right] . \tag{1.6}
\end{equation*}
$$

We say that this interaction has range $\mathrm{r}_{0} \in(0, \infty)$ if, for all $l=1, \ldots, N$ and $\mathbf{x} \in \mathbb{R}^{l d}$,

$$
\begin{equation*}
\rho\left(\mathbf{x}_{\mathcal{J}}, \mathbf{x}_{\mathcal{J}^{c}}\right)>\mathrm{r}_{0} \Longrightarrow U\left(\mathbf{x}_{\mathcal{J}} \mid \mathbf{x}_{\mathcal{J}^{c}}\right)=0 \tag{1.7}
\end{equation*}
$$

Finally, we say that the interaction is non-negative and bounded if

$$
\begin{equation*}
\inf _{\mathbf{x} \in \mathbb{R}^{l d}} U(\mathbf{x}) \geq 0 \quad \text { and } \sup _{\mathbf{x} \in \mathbb{R}^{l d}} U(\mathbf{x})<+\infty, \quad 1 \leq l \leq N \tag{1.8}
\end{equation*}
$$

The boundedness condition can be relaxed to include hard-core interactions where $U(\mathbf{x})=+\infty$ if $\left|x_{i}-x_{j}\right| \leq a$, for some given $a \in\left(0, \mathrm{r}_{0}\right)$.
1.4. Assumptions. Our assumptions on the interaction potential $U$ are borrowed from (4):
(E1) $U$ is non-negative, bounded and has a finite range $\mathrm{r}_{0} \geq 0$.
Similarly, we use assumptions on the i.i.d. random variables $\mathrm{V}_{s}, s \in \mathbb{Z}^{d}$, and the bump function $\varphi$ introduced in [4]:
(E2) There exists a constant $\mathrm{v} \in(0, \infty)$ such that

$$
\begin{equation*}
\mathbb{P}\left\{0 \leq \mathrm{V}_{0} \leq \mathrm{v}\right\}=1 \tag{1.9}
\end{equation*}
$$

and

$$
\begin{equation*}
\forall \epsilon>0 \quad \mathbb{P}\left\{\mathrm{~V}_{0} \leq \epsilon\right\}>0 . \tag{1.10}
\end{equation*}
$$

(E3) Uniform Hölder continuity There exist constants a, b $>0$ such that for all $\epsilon \in[0,1]$, the common distribution function $F$ of the random variables $\mathrm{V}_{s}$ satisfies

$$
\begin{equation*}
\sup _{\mathrm{y} \in \mathbb{R}}[F(\mathrm{y}+\epsilon)-F(\mathrm{y})] \leq \mathrm{a} \epsilon^{\mathrm{b}} . \tag{1.11}
\end{equation*}
$$

(E4) The function $\varphi: \mathbb{R}^{d} \rightarrow \mathbb{R}$ is bounded, nonnegative and compactly supported:

$$
\begin{equation*}
\operatorname{diam}(\operatorname{supp} \varphi) \leq \mathrm{r}_{1}<\infty \tag{1.12}
\end{equation*}
$$

(E5) For all $L \geq 1$ and $u \in \mathbb{Z}^{d}$,

$$
\begin{equation*}
\sum_{s \in \Lambda_{L}(u) \cap \mathbb{Z}^{d}} \varphi(x-s) \geq \mathbf{1}_{\Lambda_{L}(u)}(x) . \tag{1.13}
\end{equation*}
$$

Here and below, $\mathbf{1}_{A}$ stands for the indicator function of a set $A$.
Henceforth, we suppose that $d$ and $N$ are fixed, as well as the interaction $\mathbf{U}$ and the structure of the external potential (i.e., the distribution function $F$ and the bump function $\varphi$ ). All constants emerging in various bounds below are introduced under this assumption.
1.5. Dynamical localization. The main result of this paper, Theorem 1.1, establishes the so-called "strong dynamical localization" for the operator $\mathbf{H}(\omega)$ defined in (1.1) near the lower edge $E^{0}$ of its spectrum. More precisely, let $E^{0}$ be the lower edge of the spectrum $\operatorname{spec}\left(\mathbf{H}^{0}\right)$ of the $N$-particle operator without interaction,

$$
\begin{equation*}
\mathbf{H}^{0}=-\frac{1}{2} \boldsymbol{\Delta}+\sum_{j=1}^{N} V\left(x_{j} ; \omega\right) . \tag{1.14}
\end{equation*}
$$

Actually, it follows from our conditions (1.9) and (1.10) that $E^{0}=0$. Owing to the non-negativity of the interaction potential $U$, the lower edge of the spectrum of $\mathbf{H}$ is bounded from below by $E^{0}$. Moreover, $\mathbf{H}$ has a non-empty spectrum in the interval [ $\left.E^{0}, E^{0}+\epsilon\right]$, for any $\epsilon>0$. This follows, e.g., from a result by Klopp and Zenk $[8$ which says that the integrated density of states for a multi-particle system with a decaying interaction is the same as for the system without interaction.

Denote by $\mathbf{X}$ the operator of multiplication by the norm of $\mathbf{x}$, i.e.,

$$
\begin{equation*}
\mathbf{X} f(\mathbf{x})=|\mathbf{x}| f(\mathbf{x}), \quad \mathbf{x} \in \mathbb{R}^{N d} . \tag{1.15}
\end{equation*}
$$

The main result of this paper is the following

[^0]Theorem 1.1. Consider the operator $\mathbf{H}$ from (1.1) and assume that conditions (E1)(E5) are fulfilled. Then for any $Q>0$ there exists a nonrandom number $\eta=\eta(Q)>0$ such that for any compact subset $\mathbf{K} \subset \mathbb{R}^{N d}$ the following bound holds:

$$
\begin{equation*}
\mathbb{E}\left[\sup _{t \in \mathbb{R}}\left\|\mathbf{X}^{Q} \mathrm{e}^{-\mathrm{i} t \mathbf{H}(\omega)} P_{I(\eta)}(\mathbf{H}(\omega)) \mathbf{1}_{\mathbf{K}}\right\|_{L^{2}\left(\mathbb{R}^{N d}\right)}\right]<\infty \tag{1.16}
\end{equation*}
$$

where $P_{I(\eta)}(\mathbf{H})$ is the spectral projection of the Hamiltonian $\mathbf{H}$ on the interval $I(\eta)=$ $\left[E^{0}, E^{0}+\eta\right]$.

Remark 1.2. The interval $I(\eta)$ is a sub-interval of the interval of energies $\left[E^{0}, E^{0}+\eta^{*}\right]$ for which the spectrum of $\mathbf{H}$ was proven to be pure point (and the eigenfunctions to be decaying exponentially); see (4].

## 2. Results of the multi-Particle MSA

The MSA works with the finite-volume approximations $\mathbf{H}_{\boldsymbol{\Lambda}_{L}(\mathbf{u})}$ of $\mathbf{H}$, relative to the cubes $\boldsymbol{\Lambda}_{L}(\mathbf{u})$. More precisely, $\mathbf{H}_{\boldsymbol{\Lambda}_{L}(\mathbf{u})}$ is an operator in $L^{2}\left(\boldsymbol{\Lambda}_{L}(\mathbf{u})\right)$, given by the same expression as in (1.1) (for $\mathbf{x} \in \boldsymbol{\Lambda}_{L}(\mathbf{u})$ ), with Dirichlet's boundary conditions on $\partial \boldsymbol{\Lambda}_{L}(\mathbf{u})$; see [4]. Specifically, the Green operator $\mathbf{G}_{\boldsymbol{\Lambda}_{L}(\mathbf{u})}(E)$ is of particular interest:

$$
\begin{equation*}
\mathbf{G}_{\boldsymbol{\Lambda}_{L}(\mathbf{u})}(E)=\left(\mathbf{H}_{\boldsymbol{\Lambda}_{L}(\mathbf{u})}-E\right)^{-1} \tag{2.1}
\end{equation*}
$$

defined for $E \in \mathbb{R} \backslash \operatorname{spec}\left(\mathbf{H}_{\boldsymbol{\Lambda}_{L}(\mathbf{u})}\right)$.
Let [.] denote the integer part. For a cube $\boldsymbol{\Lambda}_{L}(\mathbf{u})$ we denote

$$
\begin{equation*}
\boldsymbol{\Lambda}_{L}^{\mathrm{int}}(\mathbf{u})=\boldsymbol{\Lambda}_{[L / 3]}(\mathbf{u}), \quad \boldsymbol{\Lambda}_{L}^{\text {out }}(\mathbf{u})=\boldsymbol{\Lambda}_{L}(u) \backslash \boldsymbol{\Lambda}_{L-2}(u) \tag{2.2}
\end{equation*}
$$

Next, given two points $\mathbf{v}, \mathbf{w} \in \mathbf{B}_{L}(\mathbf{u})$ such that $\mathbf{C}(\mathbf{v}), \mathbf{C}(\mathbf{w}) \subset \boldsymbol{\Lambda}_{L}(\mathbf{u})$, set

$$
\begin{equation*}
\mathbf{G}_{\mathbf{v}, \mathbf{w}}^{\boldsymbol{\Lambda}_{L}(\mathbf{u})}(E):=\mathbf{1}_{\mathbf{C}(\mathbf{v})} \mathbf{G}_{\boldsymbol{\Lambda}_{L}(\mathbf{u})}(E) \mathbf{1}_{\mathbf{C}(\mathbf{w})} \tag{2.3}
\end{equation*}
$$

Following a long-standing tradition, we use a parameter $\alpha \in(1,2)$ in the definition of a sequence of scales $L_{k}$ (cf. Eqn (2.5) ); For our purposes, it suffices to set $\alpha=3 / 2$; this will be always assumed below.
Definition 2.1. A cube $\boldsymbol{\Lambda}_{L}(\mathbf{u})$ is called ( $E, m$ )-non-singular ( $(E, m)$-NS, in short) if for any $\mathbf{v} \in \mathbf{B}_{\left[L^{1 / \alpha}\right]}(\mathbf{u})$ and $\mathbf{y} \in \boldsymbol{\Lambda}_{L}^{\text {out }}(\mathbf{u}) \cap \mathbb{Z}^{N d}$ the norm of the operator $\mathbf{G}_{\mathbf{v}, \mathbf{y}}^{\boldsymbol{\Lambda}_{L}(\mathbf{u})}(E)$ satisfies

$$
\begin{equation*}
\left\|\mathbf{G}_{\mathbf{v}, \mathbf{y}}^{\boldsymbol{\Lambda}_{L}(\mathbf{u})}(E)\right\|_{L^{2}\left(\boldsymbol{\Lambda}_{L}(\mathbf{u})\right)} \leq \mathrm{e}^{-m L} \tag{2.4}
\end{equation*}
$$

Otherwise, it is called $(E, m)$-singular $((E, m)-\mathrm{S})$.
We will work with a sequence of "scales" $L_{k}$ (positive integers) defined recursively by

$$
\begin{equation*}
L_{k}:=\left[L_{k-1}^{\alpha}\right]+1, \quad \text { where } \alpha=\frac{3}{2} \tag{2.5}
\end{equation*}
$$

The sequence $L_{k}$ is determined by an initial scale $L_{0} \geq 2$. Most of arguments in Sect. 3 require $L_{0}$ to be large enough, to fulfill some specific numerical inequalities. In addition, we also assume that $L_{0} \geq \mathrm{r}_{1}$ (defined in (1.12)) in order to simplify some cumbersome technicalities.

We will use a well-known property of generalized eigenfunctions of the operator $\mathbf{H}$ which can be found, e.g., in [9, Lemma 3.3.2]:

Lemma 2.2. For every bounded set $I_{0} \subset \mathbb{R}$ there exists a constant $C^{(0)}=C^{(0)}\left(I_{0}\right)$ such that, for any cube $\boldsymbol{\Lambda}_{L}(\mathbf{u})$ with $L>7$, any point $\mathbf{w} \in \mathbf{B}_{L}(\mathbf{u})$ with $\mathbf{C}(\mathbf{w}) \subseteq \boldsymbol{\Lambda}_{L}^{\mathrm{int}}(\mathbf{u})$ and every generalized eigenfunction $\Psi$ of $\mathbf{H}$ with eigenvalue $E \in I_{0}$, the norm of the vector $\mathbf{1}_{\mathbf{C}(\mathbf{w})} \boldsymbol{\Psi}$ satisfies

$$
\begin{equation*}
\left\|\mathbf{1}_{\mathbf{C}(\mathbf{w})} \boldsymbol{\Psi}\right\| \leq C^{(0)}\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L}^{\text {out }}(u)} \mathbf{G}_{\boldsymbol{\Lambda}_{L}(\mathbf{u})}(E) \mathbf{1}_{\mathbf{C}(\mathbf{w})}\right\| \cdot\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L}^{\text {out }}(\mathbf{u})} \boldsymbol{\Psi}\right\| . \tag{2.6}
\end{equation*}
$$

(From now on we omit the subscript indicating the $L^{2}$-space where a given norm is considered, as this will be clear in the context of the argument.)

The following geometric notion is used in the forthcoming analysis.
Definition 2.3. (see [4). Let $\mathcal{J}$ be a non-empty subset of $\{1, \ldots, N\}$.
We say that the cube $\boldsymbol{\Lambda}_{L}(\mathbf{y})$ is $\mathcal{J}$-separable from the cube $\boldsymbol{\Lambda}_{L}(\mathbf{x})$ if

$$
\begin{equation*}
\left(\bigcup_{j \in \mathcal{J}} \Pi_{j} \boldsymbol{\Lambda}_{L+\mathrm{r}_{1}}(\mathbf{y})\right) \bigcap\left(\bigcup_{i \notin \mathcal{J}} \Pi_{i} \boldsymbol{\Lambda}_{L+\mathrm{r}_{1}}(\mathbf{y}) \bigcup \Pi \boldsymbol{\Lambda}_{L+\mathrm{r}_{1}}(\mathbf{x})\right)=\varnothing \tag{2.7}
\end{equation*}
$$

where $\Pi \boldsymbol{\Lambda}_{L+\mathrm{r}_{1}}(\mathbf{x})=\cup_{j=1}^{N} \Pi_{j} \boldsymbol{\Lambda}_{L+\mathrm{r}_{1}}(\mathbf{x})$.
A pair of cubes $\boldsymbol{\Lambda}_{L}(\mathbf{x}), \boldsymbol{\Lambda}_{L}(\mathbf{y})$ is separable if, for some $\mathcal{J} \subseteq\{1, \ldots, N\}$, either $\boldsymbol{\Lambda}_{L}(\mathbf{y})$ is $\mathcal{J}$-separable from $\boldsymbol{\Lambda}_{L}(\mathbf{x})$, or $\boldsymbol{\Lambda}_{L}(\mathbf{x})$ is $\mathcal{J}$-separable from $\boldsymbol{\Lambda}_{L}(\mathbf{y})$.

We will use the following easy assertion (see [4]):
Lemma 2.4. For any $L>1$ and $\mathbf{x} \in \mathbb{R}^{N d}$, there exists a collection of $N$-particle cubes $\boldsymbol{\Lambda}_{2 N\left(L+\mathrm{r}_{1}\right)}\left(\mathbf{x}^{(l)}\right), l=1, \ldots, K(\mathbf{x}, N)$, with $K(\mathbf{x}, N) \leq N^{N}$, such that if a vector $\mathbf{y} \in \mathbb{Z}^{N d}$ satisfies ${ }^{2}$

$$
\begin{equation*}
\mathbf{y} \notin \bigcup_{\ell=1}^{K(\mathbf{x}, N)} \boldsymbol{\Lambda}_{2 N\left(L+\mathrm{r}_{1}\right)}\left(\mathbf{x}^{(l)}\right) \tag{2.8}
\end{equation*}
$$

then two cubes $\boldsymbol{\Lambda}_{L}(\mathbf{x})$ and $\boldsymbol{\Lambda}_{L}(\mathbf{y})$ with $\operatorname{dist}\left(\boldsymbol{\Lambda}_{L}(\mathbf{x}), \boldsymbol{\Lambda}_{L}(\mathbf{y})\right)>2 N\left(L+\mathrm{r}_{1}\right)$ are separable. In particular, assuming $L \geq \mathrm{r}_{1}$, a pair of cubes $\boldsymbol{\Lambda}_{L}(\mathbf{x}), \boldsymbol{\Lambda}_{L}(\mathbf{y})$ is separable if

$$
\begin{equation*}
|\mathbf{y}|>|\mathbf{x}|+(4 N+2) L \tag{2.9}
\end{equation*}
$$

Since $N \geq 2$, one can replace the condition (2.9) by

$$
\begin{equation*}
|\mathbf{y}|>|\mathbf{x}|+5 N L \tag{2.10}
\end{equation*}
$$

In particular, two cubes of the form $\boldsymbol{\Lambda}_{L}(\mathbf{0}), \boldsymbol{\Lambda}_{L}(\mathbf{y})$ with $|\mathbf{y}|>5 N L$ are always separable.
The main outcome of [4] is summarized in the following Theorem [2.5.
Theorem 2.5 (see [4]). For any large enough $p>0$ there exist $m^{*}(p)>0, \eta^{*}(p)>0$, and $L_{0}^{*}(p)>0$ such that
(i) if $L_{0} \geq L_{0}^{*}(p)$ then for all $k \geq 0$ and for any pair of separable cubes $\boldsymbol{\Lambda}_{L_{k}}(\mathbf{x})$, $\boldsymbol{\Lambda}_{L_{k}}(\mathbf{y})$ with $\mathbf{x}, \mathbf{y} \in \mathbb{Z}^{N d}$,

$$
\begin{equation*}
\mathbb{P}\left\{\exists E \in\left[E^{0}, E^{0}+\eta^{*}\right]: \boldsymbol{\Lambda}_{L_{k}}(\mathbf{x}) \text { and } \boldsymbol{\Lambda}_{L_{k}}(\mathbf{y}) \text { are }(E, m)-\mathrm{S}\right\} \leq L_{k}^{-2 p} \tag{2.11}
\end{equation*}
$$

(ii) with probability one, the spectrum of $\mathbf{H}$ in the interval $I=\left[E^{0}, E^{0}+\eta^{*}(p)\right]$ is pure point, and the eigenfunctions $\boldsymbol{\Phi}_{n}$ of $\mathbf{H}$ with eigenvalues $E_{n} \in I$ satisfy

$$
\begin{equation*}
\left\|\boldsymbol{\Phi}_{n} \mathbf{1}_{\mathbf{C}(\mathbf{w})}\right\| \leq C_{n}(\omega) \mathrm{e}^{-m^{*}(p)|\mathbf{w}|}, \quad \mathbf{w} \in \mathbb{Z}^{N d}, \quad C_{n}(\omega)<\infty \tag{2.12}
\end{equation*}
$$

[^1]
## 3. Derivation of dynamical localization from MSA estimates

In this section we prove a statement that is slightly more general than Theorem 1.1 , Namely, given $Q>0$, the interval $I=I(\eta)=\left[E^{0}, E^{0}+\eta\right]$ with $\eta=\eta(Q)$, and a compact subset $\mathbf{K} \subset \mathbb{R}^{d}$, there exists a constant $C(Q, \mathbf{K}) \in(0, \infty)$ such that for any bounded measurable function $\xi: \mathbb{R} \rightarrow \mathbb{C}$ with $\operatorname{supp} \xi \subset I(\eta)$,

$$
\begin{equation*}
\mathbb{E}\left[\left\|\mathbf{X}^{Q} \xi(\mathbf{H}(\omega)) \mathbf{1}_{\mathbf{K}}\right\|\right]<C(Q, \mathbf{K})\|\xi\|_{\infty}<\infty \tag{3.1}
\end{equation*}
$$

Moreover, $Q>0$ can be made arbitrarily large, by choosing $\eta=\eta(Q)>0$ sufficiently small. Theorem 1.1 follows from (3.1) applied to the functions $\xi(s)=\mathrm{e}^{-\mathrm{i} t s} \mathbf{1}_{I(\eta)}(s)$, parametrised by $t \in \mathbb{R}$.

Throughout the section, we assume that the parameter $p$ from (2.11) satisfies

$$
\begin{equation*}
2 p>3 N d \alpha+\alpha Q \tag{3.2}
\end{equation*}
$$

More precisely, given $Q>0$ and $p$ satisfying (3.2), we work with

$$
\begin{equation*}
\eta=\eta(Q) \in\left(0, \eta^{*}(p)\right) \text { and } m=m^{*}(p)>0 \tag{3.3}
\end{equation*}
$$

where $\eta^{*}(p)$ and $m^{*}(p)$ are specified in Theorem 2.5. Further, for $p$ satisfying (3.2) we introduce the event $\Omega_{1}=\Omega_{1}(p) \subseteq \Omega$ of probability $\mathbb{P}\left(\Omega_{1}\right)=1$, defined by

$$
\begin{equation*}
\Omega_{1}=\left\{\omega \in \Omega: \text { the spectrum of } \mathbf{H}(\omega) \text { in }\left[E^{0}, E^{0}+\eta^{*}(p)\right] \text { is pure point }\right\} . \tag{3.4}
\end{equation*}
$$

3.1. Probability of "bad samples". Given $j \geq 1$, consider the event

$$
\begin{array}{r}
\mathcal{S}_{j}=\left\{\omega: \text { there exists } E \in I \text { and } \mathbf{y}, \mathbf{z} \in \mathbf{B}_{5 N L_{j+1}}(\mathbf{0})\right. \text { such that } \\
\left.\qquad \boldsymbol{\Lambda}_{L_{j}}(\mathbf{y}), \boldsymbol{\Lambda}_{L_{j}}(\mathbf{z}) \text { are separable and }(m, E) \text {-S }\right\} .
\end{array}
$$

Further, for $k \geq 1$ we denote

$$
\begin{equation*}
\Omega_{k}^{\mathrm{bad}}=\bigcup_{j \geq k} \mathcal{S}_{j} \tag{3.5}
\end{equation*}
$$

Lemma 3.1. There exists a constant $c_{1} \in(0, \infty)$ such that for all $k \geq 1$,

$$
\mathbb{P}\left\{\Omega_{k}^{\mathrm{bad}}\right\} \leq c_{1} L_{k}^{-(2 p-2 N d \alpha)}
$$

Proof. The number of separable pairs $\boldsymbol{\Lambda}_{L_{j}}(\mathbf{x}), \boldsymbol{\Lambda}_{L_{j}}(\mathbf{y})$ with $\mathbf{x}, \mathbf{y} \in \mathbf{B}_{5 N L_{j+1}}(\mathbf{0})$ is bounded by $\left(10 N L_{j+1}+1\right)^{2}<C(N) L_{j+1}^{2}$. We can apply the bound (2.11) and write

$$
\mathbb{P}\left\{\mathcal{S}_{j}\right\} \leq C(N) L_{j+1}^{2 N d} L_{j}^{-2 p} \leq L_{j}^{-2 p+2 N d \alpha}
$$

Therefore,

$$
\Omega_{k}^{\mathrm{bad}} \leq L_{k}^{-2 p+2 N d \alpha} \sum_{i \geq 0}\left(\frac{L_{k+i}}{L_{k}}\right)^{-2 p+2 N d \alpha}
$$

For $2 p>2 N d \alpha$ and $L_{0} \geq 2$ the claim follows from the inequality

$$
\frac{L_{k+i}}{L_{k}} \geq\left[L_{k}^{\alpha^{i}-1}\right]
$$

3.2. Centers of localization. Denote by $\boldsymbol{\Phi}_{n}=\boldsymbol{\Phi}_{n}(\omega)$ the normalized eigenfunctions of $\mathbf{H}(\omega), \omega \in \Omega_{1}$, with corresponding eigenvalues $E_{n}=E_{n}(\omega) \in I$. For each $n$ we define a center of localization for $\boldsymbol{\Phi}_{n}$ as a point $\hat{\mathbf{x}} \in \mathbb{Z}^{d}$ such that

$$
\begin{equation*}
\left\|\mathbf{1}_{\mathbf{C}(\hat{\mathbf{x}})} \boldsymbol{\Phi}_{n}\right\|=\max _{\mathbf{y} \in \mathbb{Z}^{N d}}\left\|\mathbf{1}_{\mathbf{C}(\mathbf{y})} \boldsymbol{\Phi}_{n}\right\| \tag{3.6}
\end{equation*}
$$

Since $\left\|\boldsymbol{\Phi}_{n}\right\|=1$, for any given $n$ such centers exist and their number is finite. We will assume that, for any eigenfunction $\boldsymbol{\Phi}_{n}$, the centers of localization $\hat{\mathbf{x}}_{n, a}, a=1, \ldots, \hat{C}(n)$, are enumerated in such a way that $\left|\hat{\mathbf{x}}_{n, 1}\right|=\min _{a}\left|\hat{\mathbf{x}}_{n, a}\right|$.

Lemma 3.2. There exists $k_{0}$ large enough such that, for all $\mathbf{u} \in \mathbb{Z}^{N d}, \omega \in \Omega_{1}$ and $k \geq k_{0}$, if $\hat{\mathbf{x}}_{n, a} \in \mathbf{B}_{L_{k}}(\mathbf{u})$ then the box $\boldsymbol{\Lambda}_{L_{k}}(\mathbf{u})$ is $\left(m, E_{n}\right)$-S.

Proof. Assume otherwise. Then from (2.6) it follows that

$$
\left\|\mathbf{1}_{\mathbf{C}\left(\hat{\mathbf{x}}_{n, a}\right)} \boldsymbol{\Phi}_{n}\right\| \leq C^{\prime} \mathrm{e}^{-m L_{k}}\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}^{\text {out }(\mathbf{u})}} \boldsymbol{\Phi}_{n}\right\|
$$

Since the number of cells in $\boldsymbol{\Lambda}_{L_{k}}^{\text {out }}(\mathbf{u})$ is bounded by $L_{k}^{N d}$, we conclude that

$$
\left\|\mathbf{1}_{\mathbf{C}\left(\hat{\mathbf{x}}_{n, a}\right)} \boldsymbol{\Phi}_{n}\right\| \leq C^{\prime} \mathrm{e}^{-m L_{k}} L_{k}^{N d} \cdot \max _{\mathbf{y} \in \mathbf{B}_{L_{k}}^{\text {out }}(\mathbf{u})}\left\|\mathbf{1}_{\mathbf{C}(\mathbf{y})} \boldsymbol{\Phi}_{n}\right\|
$$

If $k_{0}$ is large enough so that $C^{\prime} \mathrm{e}^{-m L_{k}} L_{k}^{N d}<1$ for $k \geq k_{0}$, the above inequality contradicts the definition of $\hat{\mathbf{x}}_{n, a}$ as center of localization.
3.3. Annular regions. From now on we work with the integer $k_{0}$ from Lemma 3.2. Given $k>k_{0}$, set:

$$
\begin{equation*}
\Omega_{k}^{\mathrm{good}}=\Omega_{1} \backslash \Omega_{k}^{\mathrm{bad}} \tag{3.7}
\end{equation*}
$$

Assume that $\omega \in \Omega_{k}^{\text {good }}$. Let $\hat{\mathbf{x}}_{n, a}, \hat{\mathbf{x}}_{n, b}$ be two centers of localization for the same eigenfunction $\boldsymbol{\Phi}_{n}$. It follows from the definition of the event $\Omega_{k}^{\text {good }}$ that the cubes $\boldsymbol{\Lambda}_{L_{i}}\left(\hat{\mathbf{x}}_{n, a}\right)$ and $\boldsymbol{\Lambda}_{L_{i}}\left(\hat{\mathbf{x}}_{n, b}\right)$ with $i \geq k-1$ cannot be separable, since they must be ( $m, E$ )-S. Further, by Lemma [2.4 if $L_{0} \geq \mathrm{r}_{1}$ then any cube of the form $\boldsymbol{\Lambda}_{L_{k}}(\mathbf{y})$ with $|\mathbf{y}|>\left|\hat{\mathbf{x}}_{n, 1}\right|+5 N L_{k}$ is separable from $\boldsymbol{\Lambda}_{L_{k}}\left(\hat{\mathbf{x}}_{n, 1}\right)$; this also applies, of course, to any localization center $\mathbf{y}=\hat{\mathbf{x}}_{n, a}$ with $a>1$, provided that such centers exist for a given $n$. Since $\omega \in \Omega_{k}^{\text {good }}$, for any eigenfunction $\boldsymbol{\Phi}_{n}$ there is no center of localization $\hat{\mathbf{x}}_{n, a}$ either outside the cube $\boldsymbol{\Lambda}_{\left|\hat{\mathbf{x}}_{n, 1}\right|+5 N L_{k}}(\mathbf{0})$ or inside $\boldsymbol{\Lambda}_{\left|\hat{\mathbf{x}}_{n, 1}\right|}(\mathbf{0})$ (since $\left.\left|\hat{\mathbf{x}}_{n, 1}\right|=\min _{a}\left|\hat{\mathbf{x}}_{n, a}\right|\right)$. In other words, within the event $\Omega_{k}^{\text {good }}$, all centers of localization $\hat{\mathbf{x}}_{n, a}$ with a fixed value of $n$ are located in the annulus

$$
\boldsymbol{\Lambda}_{\left|\hat{\mathbf{x}}_{n, 1}\right|+5 N L_{k}}(\mathbf{0}) \backslash \boldsymbol{\Lambda}_{\left|\hat{\mathbf{x}}_{n, 1}\right|}(\mathbf{0})
$$

of width $5 N L_{k}$ and of inner radius $\left|\hat{\mathbf{x}}_{n, 1}\right|$. This explains why, for our purposes, an eigenfunction $\boldsymbol{\Phi}_{n}$ can be effectively "labeled" by a single localization center.

In other words, although in this paper we cannot rule out the possibility of existence of multiple centers of localization at arbitrarily large distances (depending on $\boldsymbol{\Phi}_{n}$ through $\left|\hat{\mathbf{x}}_{n, 1}\right|$ ), such centers do not contribute to a "radial" quantum transport - away from the origin $\mathbf{0}$ - which might have lead to dynamical delocalization.

Lemma 3.3. Given $k>k_{0}$, there exists $j_{0} \geq k$ large enough such that if $j \geq j_{0}$, $\omega \in \Omega_{k}^{\text {good }}$ and $\hat{\mathbf{x}}_{n, 1} \in \mathbf{B}_{L_{j}}(\mathbf{0})$ then

$$
\left\|\left(1-\mathbf{1}_{\boldsymbol{\Lambda}_{L_{j+2}}(\mathbf{0})}\right) \boldsymbol{\Phi}_{n}\right\| \leq \frac{1}{4}
$$

Proof. By Lemma 2.4 (see also (2.10)),

$$
\forall i \geq j+1, \forall \mathbf{w} \in \mathbb{Z}^{N d} \backslash \boldsymbol{B}_{5 N L_{i}}(\mathbf{0}) \text {, the cubes } \boldsymbol{\Lambda}_{L_{i}}(\mathbf{w}) \text { and } \boldsymbol{\Lambda}_{L_{i}}(\mathbf{0}) \text { are separable. }
$$

In addition, we take $j \geq k$, as suggested in the lemma.
Next, we divide the complement $\mathbb{R}^{N d} \backslash \boldsymbol{\Lambda}_{5 N L_{j+2}}(\mathbf{0})$ into annular regions

$$
\begin{equation*}
\mathbf{M}_{i}(\mathbf{0}):=\boldsymbol{\Lambda}_{5 N L_{i+1}}(\mathbf{0}) \backslash \boldsymbol{\Lambda}_{5 N L_{i}}(\mathbf{0}), \quad i \geq j+2 \tag{3.8}
\end{equation*}
$$

and write

$$
\left\|\left(1-\mathbf{1}_{\boldsymbol{\Lambda}_{L_{j+2}}(\mathbf{0})}\right) \boldsymbol{\Phi}_{n}\right\|^{2}=\sum_{i \geq j+2}\left\|\mathbf{1}_{\mathbf{M}_{i}(\mathbf{0})} \boldsymbol{\Phi}_{n}\right\|^{2} \leq \sum_{i \geq j+2} \sum_{\mathbf{w} \in \mathbf{M}_{i}(\mathbf{0})}\left\|\mathbf{1}_{\mathbf{C}(\mathbf{w})} \boldsymbol{\Phi}_{n}\right\|^{2}
$$

Furthermore, $\hat{\mathbf{x}}_{n, 1} \in \mathbf{B}_{L_{j}}(\mathbf{0}) \subset \mathbf{B}_{L_{i-1}}(\mathbf{0})$, so that by Lemma 3.2, the cube $\boldsymbol{\Lambda}_{L_{i}}(\mathbf{0})$ must be $\left(m, E_{n}\right)$-S. Therefore, by the definition of the event $\Omega_{k}^{\text {good }}$, the cube $\boldsymbol{\Lambda}_{L_{i}}(\mathbf{w})$ is $\left(m, E_{n}\right)$-NS. Applying Lemma 2.2 to the cube $\boldsymbol{\Lambda}_{L_{i}}(\mathbf{w})$ and to the cell $\mathbf{C}(\mathbf{w})$, we obtain

$$
\left\|\mathbf{1}_{\mathbf{C}(\mathbf{w})} \boldsymbol{\Phi}_{n}\right\|^{2} \leq \mathrm{e}^{-2 m L_{i}}
$$

Since the volume $\left|\mathbf{M}_{i}(\mathbf{0})\right|$ of the annular region $\mathbf{M}_{i}(\mathbf{0})$ grows polynomially in $L_{i}$, the assertion of Lemma 3.3 follows.

### 3.4. Bounds on concentration of localization centers.

Lemma 3.4. There exists a constant $c_{2} \in(0, \infty)$ such that for $\omega \in \Omega_{k}^{\text {bad }}, j \geq k$,

$$
\begin{equation*}
\operatorname{card}\left\{n: \hat{\mathbf{x}}_{n, 1} \in \mathbf{B}_{L_{j+1}}(\mathbf{0})\right\} \leq c_{2} L_{j+1}^{\alpha N d} \tag{3.9}
\end{equation*}
$$

Proof. The left-hand-side of (3.9) is nondecreasing in $j$, so we can restrict ourselves to the case $j \geq j_{0}$. First, observe that, with $\boldsymbol{\Lambda}_{L_{j+2}}=\boldsymbol{\Lambda}_{L_{j+2}}(\mathbf{0})$

$$
\begin{equation*}
\sum_{n: \hat{\mathbf{x}}_{n, 1} \in \mathbf{B}_{L_{j+1}}(\mathbf{0})}\left(\mathbf{1}_{\boldsymbol{\Lambda}_{L_{j+2}}} P_{I} \mathbf{1}_{\boldsymbol{\Lambda}_{L_{j+2}}} \boldsymbol{\Phi}_{n}, \boldsymbol{\Phi}_{n}\right) \leq \operatorname{tr}\left(\mathbf{1}_{\boldsymbol{\Lambda}_{L_{j+2}}} P_{I}\right) \tag{3.10}
\end{equation*}
$$

Each term in the above sum is not less than $1 / 2$. Indeed,

$$
\begin{align*}
& \left(\mathbf{1}_{\boldsymbol{\Lambda}_{L_{j+2}}} P_{I} \mathbf{1}_{\boldsymbol{\Lambda}_{L_{j+2}}} \boldsymbol{\Phi}_{n}, \boldsymbol{\Phi}_{n}\right) \\
& \quad=\left(\mathbf{1}_{\boldsymbol{\Lambda}_{L_{j+2}}} P_{I} \boldsymbol{\Phi}_{n}, \boldsymbol{\Phi}_{n}\right)-\left(\mathbf{1}_{\boldsymbol{\Lambda}_{L_{j+2}}} P_{I}\left(1-\mathbf{1}_{\boldsymbol{\Lambda}_{L_{j+2}}}\right) \boldsymbol{\Phi}_{n}, \boldsymbol{\Phi}_{n}\right) \\
& \quad \geq\left(\mathbf{1}_{\boldsymbol{\Lambda}_{L_{j+2}}} \boldsymbol{\Phi}_{n}, \boldsymbol{\Phi}_{n}\right)-\frac{1}{4} \quad \text { (using Lemma 3.3) }  \tag{3.11}\\
& \quad=\left(\boldsymbol{\Phi}_{n}, \boldsymbol{\Phi}_{n}\right)-\left(\left(1-\mathbf{1}_{\boldsymbol{\Lambda}_{L_{j+2}}}\right) \boldsymbol{\Phi}_{n}, \boldsymbol{\Phi}_{n}\right)-\frac{1}{4} \\
& \quad \geq \frac{1}{2} \tag{3.12}
\end{align*}
$$

Substituting the lower bounds from (3.11) - (3.12) under the trace in Eqn (3.10), we get the desired upper bound on the LHS of Eqn (3.9).

### 3.5. Bounds for "good" samples of potential.

Lemma 3.5. There exists an integer $k_{1}=k_{1}\left(L_{0}\right)$ such that $\forall k \geq k_{1}, \omega \in \Omega_{k+1}^{\text {good }}$ and $\mathbf{x}$ from the annular region $\mathbf{M}_{k+1}$ defined in (3.8),

$$
\begin{equation*}
\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{x})} P_{I} \xi(\mathbf{H}) \mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{0})}\right\| \leq \mathrm{e}^{-m L_{k} / 2}\|\xi\|_{\infty} \tag{3.13}
\end{equation*}
$$

Proof. It suffices to prove (3.13) in the particular case where $\|\xi\|_{\infty} \leq 1$, which we assume below. First, we bound the LHS of (3.13) as follows:

$$
\begin{align*}
\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{x})} P_{I} \xi(\mathbf{H}) \mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{0})}\right\| & \leq \sum_{E_{n} \in I}\left|\xi\left(E_{n}\right)\right|\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{x})} \boldsymbol{\Phi}_{n}\right\|\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{0})} \boldsymbol{\Phi}_{n}\right\| \\
& \leq \sum_{E_{n} \in I}\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{x})} \boldsymbol{\Phi}_{n}\right\|\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{0})} \boldsymbol{\Phi}_{n}\right\| \tag{3.14}
\end{align*}
$$

since $\|\eta\|_{\infty} \leq 1$. Now divide the sum according to where $\hat{\mathbf{x}}_{n, 1}$ are located and write

$$
\begin{aligned}
\sum_{E_{n} \in I}\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{x})} \boldsymbol{\Phi}_{n}\right\|\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{0})} \boldsymbol{\Phi}_{n}\right\|= & \sum_{\substack{E_{n} \in I \\
\hat{\mathbf{x}}_{n, 1} \in \boldsymbol{\Lambda}_{5 N L_{k+1}}(\mathbf{0})}}\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{x})} \boldsymbol{\Phi}_{n}\right\|\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{0})} \boldsymbol{\Phi}_{n}\right\| \\
& +\sum_{j=k+1}^{\infty} \sum_{\substack{E_{n} \in I \\
\hat{\mathbf{x}}_{n, 1} \in \mathrm{M}_{j}(\mathbf{0})}}\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{x})} \boldsymbol{\Phi}_{n}\right\|\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{0})} \boldsymbol{\Phi}_{n}\right\|,
\end{aligned}
$$

with $\mathbf{M}_{i}(\mathbf{0})$ defined in (3.8). Since $\mathbf{x} \in \mathbf{M}_{k+1}(\mathbf{0})$, we have $\mathbf{B}_{L_{k}}(\mathbf{x}) \cap \mathbf{B}_{L_{k}}(\mathbf{0})=\varnothing$. Then, by Lemma 2.4, the two cubes $\mathbf{B}_{L_{k}}(\mathbf{x})$ and $\mathbf{B}_{L_{k}}(\mathbf{0})$ are separable. In turn, this implies that one of these cubes is $\left(m, E_{n}\right)-$ NS. Therefore, by Lemma 3.4

$$
\sum_{\substack{E_{n} \in I \\ \hat{\mathbf{x}}_{n, 1} \in \boldsymbol{\Lambda}_{L_{k+1}}(\mathbf{0})}}\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{x})} \boldsymbol{\Phi}_{n}\right\|\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{0})} \boldsymbol{\Phi}_{n}\right\| \leq c_{2} C^{\prime} L_{k+1}^{\alpha N d} \mathrm{e}^{-m L_{k}}
$$

Furthermore, for $k>k_{0}$ large enough,

$$
\begin{equation*}
\sum_{\substack{E_{n} \in I \\ \hat{\mathbf{x}}_{n, 1} \in \boldsymbol{\Lambda}_{L_{k+1}}(\mathbf{0})}}\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{x})} \boldsymbol{\Phi}_{n}\right\|\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{0})} \boldsymbol{\Phi}_{n}\right\| \leq \frac{1}{2} \mathrm{e}^{-m L_{k} / 2} . \tag{3.15}
\end{equation*}
$$

For any $j \geq k+2$ and $\hat{\mathbf{x}}_{n, 1} \in \mathbf{M}_{j}(\mathbf{0})$, by Lemma 3.2, the cube $\mathbf{B}_{L_{j}}\left(\hat{\mathbf{x}}_{n, 1}\right)$ must be $\left(m, E_{n}\right)$-S, so that $\mathbf{B}_{L_{j}}(\mathbf{0})$ has to be $\left(m, E_{n}\right)$-NS:

$$
\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{0})} \boldsymbol{\Phi}_{n}\right\| \leq\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{j}}(\mathbf{0})} \boldsymbol{\Phi}_{n}\right\| \leq C^{\prime} \mathrm{e}^{-m L_{j}}
$$

Applying again Lemma 3.4 we see that, if $k \geq k_{1}$, then

$$
\begin{aligned}
\sum_{j=k+1}^{\infty} \sum_{\substack{E_{n} \in I \\
\hat{\mathbf{x}}_{n, 1} \in \mathrm{M}_{j}(\mathbf{0})}}\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{x})} \boldsymbol{\Phi}_{n}\right\|\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{0})} \boldsymbol{\Phi}_{n}\right\| & \leq C \sum_{j=k+2}^{\infty} \mathrm{e}^{-m L_{j}} L_{j}^{\alpha N d} \\
& \leq \frac{1}{2} \mathrm{e}^{-m L_{k} / 2}
\end{aligned}
$$

Combining this estimate with (3.14) and (3.15), the assertion of Lemma 3.5 follows.

### 3.6. Bounds for "bad" samples of potential.

Lemma 3.6. Let $k_{1}$ be as in Lemma 3.5 and assume that $k \geq k_{1}$ and $\mathbf{x} \in \mathbf{M}_{k+1}(\mathbf{0})$. We have:

$$
\mathbb{E}\left[\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{x})} P_{I} \xi(\mathbf{H}) \mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{0})}\right\|\right] \leq\|\xi\|_{\infty}\left(C L_{k}^{-2 p+2 N d \alpha}+\mathrm{e}^{-m L_{k} / 2}\right)
$$

Proof. We again assume $\|\xi\|_{\infty} \leq 1$. For $\omega \in \Omega_{k}^{\text {bad }}$ we can use Sect. 3.1 while for $\omega \in \Omega_{k}^{\text {good }}$ we can use Sect. 3.5. More precisely, the above expectation is bounded by

$$
\mathbb{P}\left\{\Omega_{k}^{\mathrm{bad}}\right\}+\mathrm{e}^{-m L_{k} / 2} \mathbb{P}\left\{\Omega_{k}^{\mathrm{good}}\right\} \leq C L_{k}^{-2 p+2 N d \alpha}+\mathrm{e}^{-m L_{k} / 2}
$$

3.7. Conclusion. For a compact subset $\mathbf{K} \subset \mathbb{R}^{N d}$ we find an integer $k \geq k_{1}$ such that $\mathbf{K} \subset \boldsymbol{\Lambda}_{L_{k}}(\mathbf{0})$. Then, with the annular regions $\mathbf{M}_{j}(\mathbf{0})$,

$$
\begin{aligned}
\mathbb{E}\left[\left\|\mathbf{X}^{Q} P_{I} \xi(\mathbf{H}(\omega)) \mathbf{1}_{\mathbf{K}}\right\|\right] & \leq L_{k}^{Q}+\sum_{j \geq k+1} \mathbb{E}\left[\left\|\mathbf{X}^{Q} \mathbf{1}_{\mathbf{M}_{j}(\mathbf{0})} P_{I} \xi(\mathbf{H}) \mathbf{1}_{\mathbf{K}}\right\|\right] \\
& \leq L_{k}^{Q}+\sum_{j \geq k+1} L_{j+1}^{Q}\left(\sum_{\mathbf{w} \in \mathbf{M}_{j}(\mathbf{0})} \mathbb{E}\left[\left\|\mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{w})} P_{I} \xi(\mathbf{H}) \mathbf{1}_{\boldsymbol{\Lambda}_{L_{k}}(\mathbf{0})}\right\|\right]\right) \\
& \leq L_{k}^{Q}+\sum_{j \geq k+1} L_{j}^{\alpha Q} L_{j}^{N d \alpha}\left(L_{j}^{-2 p+2 N d \alpha}+\mathrm{e}^{-m L_{j} / 2}\right)<\infty
\end{aligned}
$$

since $2 p>3 N d \alpha+\alpha Q$ by assumption (3.2), and $L_{j} \sim\left[L_{0}^{\alpha^{j}}\right]$ grow fast enough.
This completes the proof of dynamical localization.

## References

[1] M. Aizenman and S. Warzel, Localization bounds for multiparticle systems, Comm. Math. Phys. 290 (2009), no. 3, 903-934.
[2] _, Complete dynamical localization in disordered quantum multi-particle systems, 2009, arXiv:math-ph/0909:5432 (2009).
[3] A. Boutet de Monvel, V. Chulaevsky, P. Stollmann, and Y. Suhov, Wegner-type bounds for a multi-particle continuous Anderson model with an alloy-type external potential, J. Stat. Phys. 138 (2010), no. 4-5, 553-566.
[4] , Anderson localization for a multi-particle alloy-type model, 2010, arXiv:mathph/1004.1300 (2010).
[5] V. Chulaevsky and Y. Suhov, Wegner bounds for a two-particle tight binding model, Comm. Math. Phys. 283 (2008), no. 2, 479-489.
[6] _, Eigenfunctions in a two-particle Anderson tight binding model, Comm. Math. Phys. 289 (2009), no. 2, 701-723.
[7] , Multi-particle Anderson localisation: induction on the number of particles, Math. Phys. Anal. Geom. 12 (2009), no. 2, 117-139.
[8] F. Klopp and H. Zenk, The integrated density of states for an interacting multielectron homogeneous model, 2003, arXiv:math-ph/0310031.
[9] P. Stollmann, Caught by disorder, Progress in Mathematical Physics, vol. 20, Birkhäuser Boston Inc., Boston, MA, 2001. Bound states in random media.
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[^0]:    ${ }^{1}$ The Hölder continuity can be relaxed to the log-Hölder continuity.

[^1]:    ${ }^{2}$ The constant $r_{1}$ is defined in (1.12).

