



Localization for Anderson models on metric and discrete tree graphs

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Received: 19 February 2019 / Revised: 12 September 2019 / Published online: 25 September 2019
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Abstract

We establish spectral and dynamical localization for several Anderson models on metric and discrete radial trees. The localization results are obtained on compact intervals contained in the complement of discrete sets of exceptional energies. All results are proved under the minimal hypothesis on the type of disorder: the random variables generating the trees assume at least two distinct values. This level of generality, in particular, allows us to treat radial trees with disordered geometry as well as Schrödinger operators with Bernoulli-type singular potentials. Our methods are based on an interplay between graph-theoretical properties of radial trees and spectral analysis of the associated random differential and difference operators on the half-line.

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Communicated by Loukas Grafakos.

David Damanik was supported in part by NSF Grant DMS–1700131. Jake Fillman was supported in part by an AMS-Simons travel Grant, 2016–2018. Selim Sukhtaiev was supported in part by an AMS-Simons travel Grant, 2017–2019.

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1 Introduction

1.1 Description of main results

The central theme of this paper is Anderson localization for random models on tree graphs. In the first part of this work we establish spectral and dynamical localization for *continuum* Laplace operators subject to random Kirchhoff vertex conditions on radial trees with disordered geometry. Specifically, we consider metric trees with random branching numbers and random edge lengths. The second part of this paper addresses analogous questions for random second order difference operators on *discrete* radial trees with random branching numbers. At the outset, we emphasize that our results are all proved under the minimal possible hypotheses. Namely, we assume that the random variables used to generate the trees take at least two distinct values. We will formulate this assumption more precisely as Hypothesis 3.1. In particular, we can handle the case of Bernoulli distributions, which is generally considered to be the most challenging setting.

To begin, let us describe the models. Let Γ be a metric tree with vertices \mathcal{V} , edges \mathcal{E} , and uniformly bounded edge lengths $\{\ell_e > 0 : e \in \mathcal{E}\}$. We further assume that there is a unique vertex $o \in \mathcal{V}$ with degree 1, which we call the *root* of Γ ; see, for example, Fig. 1. For each vertex v , $\text{gen}(v)$ (the generation of v) is the combinatorial distance from v to the root. One defines $\text{gen}(e)$ for $e \in \mathcal{E}$ similarly. We consider the Laplace operator $\mathbb{H} := -\frac{d^2}{dx^2}$ acting in $L^2(\Gamma)$. In order to ensure self-adjointness of \mathbb{H} , we impose a Dirichlet condition at o , that is,

$$f(o) = 0, \quad (1.1)$$

as well as Kirchhoff vertex conditions given by

$$\begin{cases} f \text{ is continuous at } v, & v \in \mathcal{V} \\ \sum_{e \in \mathcal{E}: v \in e} \partial_v^e f(v) = q(v)f(v) & v \in \mathcal{V} \setminus \{o\}, \end{cases} \quad (1.2)$$

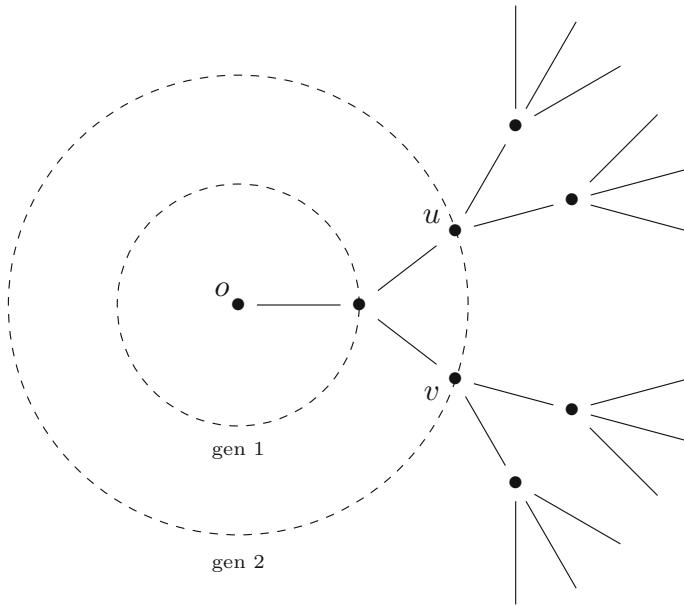


Fig. 1 $b_0 = 1, b_1 = 2, b_2 = 2, b_3 = 3$

where $q : \mathcal{V} \rightarrow \mathbb{R}$ is a real-valued function, and ∂_v^e denotes the inward-pointed derivative along the edge $e \in \mathcal{E}$. The assumption that $\deg(o) = 1$ is purely for convenience. If the root has degree 2 or higher, the Dirichlet condition (1.1) implies that the operators we study decompose into a direct sum of operators covered by the $\deg(o) = 1$ case. In the simplified case $\Gamma = \mathbb{R}_+$ the vertex conditions (1.2) provide a rigorous description of the self-adjoint realization of Schrödinger operators with zero-range potentials and coupling constants $q(v)$ (cf., e.g., [7, Section III.2.1], [9, Section 1.4.1]).

We denote the branching number of each vertex by $b(v) = \deg(v) - 1$ for $v \in \mathcal{V} \setminus \{o\}$. In this work, we assume that all quantities are *radial*. That is to say, we assume that $q(v)$ and $b(v)$ depend only on $\text{gen}(v)$ and ℓ_e depends only on $\text{gen}(e)$. The three continuum random models treated in this paper are: the random branching model (RBM), the random lengths model (RLM), and the random Kirchhoff model (RKM). In these models, the branching numbers, the Kirchhoff coupling constants, and the edge lengths are independent identically distributed Bernoulli-type random variables which depend only on the distance to the root o ; the precise description of these models is provided in Sect. 3.1. In fact, our approach can allow all three parameters to vary simultaneously; we simply single out RBM, RLM, and RKM as prominent applications of our method. Thus, these models are parameterized by a choice of a probability measure $\tilde{\mu}$ supported on a set of the form $\mathcal{A} = \{b^-, b^- + 1, \dots, b^+\} \times [\ell^-, \ell^+] \times [q^-, q^+]$, which gives the probability distribution for the branching numbers, the edge lengths, and the Kirchhoff potential at each generation. To be a little more specific, the probability space is $\Omega = \mathcal{A}^{\mathbb{N}}$ with measure $\mu = \tilde{\mu}^{\mathbb{N}}$; then, each $\omega \in \Omega$ produces a tree model with parameters dictated by

$$b(v) = \omega_1(\text{gen}(v)), \quad \ell_e = \omega_2(\text{gen}(e)), \quad q(v) = \omega_3(\text{gen}(v)), \quad v \in \mathcal{V}, \quad e \in \mathcal{E}.$$

Our approach is based on the orthogonal decomposition of $L^2(\Gamma)$ into a countable collection of reducing subspaces of the operator \mathbb{H} ; cf. [52, 53] (see also [22]). The restriction of \mathbb{H} on each subspace is unitarily equivalent to a shifted version of the model half-line operator $H := -\frac{d^2}{dx^2}$ acting in $L^2(\mathbb{R}_+)$, subject to the Dirichlet condition at 0 and self-adjoint vertex conditions of the form

$$\begin{cases} \sqrt{b_j} f(t_j^-) = f(t_j^+), & j \in \mathbb{N} \\ f'(t_j^-) + q_j f(t_j^-) = \sqrt{b_j} f'(t_j^+) & j \in \mathbb{N}, \end{cases} \quad (1.3)$$

where t_j denotes the distance from the root to vertices of generation $j \in \mathbb{N}$. Similarly b_j denotes the branching number and q_j is the Kirchhoff coupling constant at generation j .

The natural occurrence of Bernoulli models in this paper is due to random branching; in particular, the branching at each generation may only take integral values, so any randomness in the branching parameter must necessarily be discrete. Broadly speaking, the behavior of random models (at least in one spatial dimension) tends to be monotone in the randomness. In particular, increasing the randomness of the model tends to make the spectrum more singular. Thus, proving localization statements in the situation in which the single-site distribution is supported on two points (the Bernoulli case) is the most challenging task.

To prove localization for the 1D half-line operator H_ω , we adapt the approach of [19], which itself fits into the general framework of spectral analysis via transfer matrix techniques, see, e.g., [26, 55] for illuminating discussions. Recall that a *generalized eigenfunction* is an solution ψ of the eigenvalue equation $H_\omega \psi = E\psi$ that enjoys a linear upper bound; in this case, E is known as the corresponding *generalized eigenvalue*.

For the proof, we first employ Fürstenberg's Theorem to ensure positivity of the Lyapunov exponent away from a discrete set \mathfrak{D} (Theorem 3.5), and then show that almost surely all generalized eigenfunctions exhibit Lyapunov behavior in every compact interval $I \subset \mathbb{R} \setminus \mathfrak{D}$, (Theorem 3.13). This shows that the generalized eigenfunctions decay exponentially, which establishes spectral localization. At that point, the established exponential decay of generalized eigenfunctions is combined with the proof of spectral localization to bootstrap sharper bounds for the eigenfunctions in terms of their centers of localization, cf. (3.23). The latter are crucial for showing dynamical localization. We summarize this discussion by formulating the first main result of this work.

Theorem 1.1 *Suppose $\text{supp } \tilde{\mu}$ contains at least two points. Then there exists a discrete set $\mathfrak{D} \subset \mathbb{R}$ such that for every compact interval $I \subseteq \mathbb{R} \setminus \mathfrak{D}$ and every $p > 0$, there exists $\tilde{\Omega} \subset \Omega$ with $\mu(\tilde{\Omega}) = 1$ such that*

$$\sup_{t>0} \left\| |X|^p \chi_I(H_\omega) e^{-itH_\omega} \psi \right\|_{L^2(\mathbb{R}_+)} < \infty, \quad \omega \in \tilde{\Omega}, \quad (1.4)$$

whenever $\psi \in L^2(\mathbb{R}_+)$ and

$$\psi(x) \underset{x \rightarrow \infty}{=} \mathcal{O}(e^{-\log^C x}),$$

for some universal constant $C > 0$.

We prove this Theorem in Sect. 3. We deduce the second main result of the paper by combining Theorem 1.1 and the orthogonal decomposition of radial trees; see Sect. 4.

Theorem 1.2 *Suppose $\text{supp } \tilde{\mu}$ contains at least two points. Then, there exists a discrete set $\mathfrak{D} \subseteq \mathbb{R}$ such that the following two assertions hold.*

- (i) *The operator \mathbb{H}_ω exhibits Anderson localization at all energies outside of \mathfrak{D} . That is, almost surely, \mathbb{H}_ω has pure point spectrum and any eigenfunction of \mathbb{H}_ω corresponding to an energy $E \in \sigma(\mathbb{H}_\omega) \setminus \mathfrak{D}$ enjoys an exponential decay estimate of the form*

$$|f(x)| \leq \frac{C e^{-\lambda|x|}}{\sqrt{w_o(|x|)}} \quad (1.5)$$

with $C > 0$ and $\lambda > 0$, where $w_o(|x|)$ denotes the number of vertices in the generation of x , i.e., $w_o(|x|) = \#\{y \in \mathcal{V} : \text{gen}(y) = \text{gen}(x)\}$.

- (ii) *For every compact interval $I \in \mathbb{R} \setminus \mathfrak{D}$ and every $p > 0$, there exists a set $\Omega^* \subset \Omega$ with $\mu(\Omega^*) = 1$ such that for every $\omega \in \Omega^*$ and every compact set $\mathcal{K} \subset \Gamma_{b_\omega, \ell_\omega}$ one has*

$$\sup_{t>0} \left\| |X|^p \chi_I(\mathbb{H}_\omega) e^{-it\mathbb{H}_\omega} \chi_{\mathcal{K}} \right\|_{L^2(\Gamma_{b_\omega, \ell_\omega})} < \infty,$$

where $\chi_I(\mathbb{H}_\omega)$ is the spectral projection corresponding to I , and $|X|^p$ denotes the operator of multiplication by the radial function $f(x) := |x|^p$, $x \in \Gamma_{b_\omega, \ell_\omega}$, where $|x|$ denotes the distance from x to the root o .

We note that the theorem above gives localization for RBM, RLM, and RKM. We also note that the spectrum of \mathbb{H}_ω is given by a deterministic set. This is addressed in Sect. 4.1 where we also point out that the analogous question for the half-line operator H_ω presents some complications which are not typical for full-line ergodic models, see Remark 4.2.

Remark 1.3 A few remarks:

- (1) The assumption that the support of the single-generation distribution contains at least two points is clearly necessary. For, if $\text{supp } \tilde{\mu}$ consists of a single point, then there is only one operator H_ω , which is then periodic and hence does not exhibit Anderson localization.
- (2) We will refer to functions on trees obeying an estimate like (1.5) as *tree-exponentially decaying*. Since the number of vertices at the n th generation grows exponentially with n , the factor of $\sqrt{w_o(|x|)}$ in the denominator implies that the eigenfunction decay leads to square-integrability.
- (3) The transfer matrices for the half-line models can be bounded at isolated energies, and hence one cannot avoid excluding a discrete set of energies. This will be discussed in more detail in Sect. 3.

In Part 2 we address analogous questions for the discrete versions of RBM, RLM, and RKM, namely, we consider discrete Schrödinger and weighted adjacency operators on radial trees with random branching numbers, hopping parameters, and vertex potentials. Concretely, we consider rooted radial tree graphs Γ as before. Given functions $q : \mathcal{V} \rightarrow \mathbb{R}$ and $p : \mathcal{E} \rightarrow (0, \infty)$, the corresponding Schrödinger operators \mathbb{A} and \mathbb{S} are given by

$$[\mathbb{A}f](u) = - \sum_{v \sim u} p(u, v) f(v), \quad f \in \ell^2(\mathcal{V}), \quad v \in \mathcal{V}. \quad (1.6)$$

$$[\mathbb{S}f](u) = \sum_{v \sim u} (q(u) f(u) - f(v)), \quad f \in \ell^2(\mathcal{V}), \quad v \in \mathcal{V}. \quad (1.7)$$

As before, we will assume that b , p , and q are bounded radial functions, so the randomness will be encoded in a measure $\tilde{\mu}$ which gives the distribution of branching numbers, edge weights, and vertex potentials in each generation. We will define this more precisely in Part 2. Our third main result is the following theorem which is proved in Sect. 5.3. The quantity $w_y(r)$ in (1.9) below denotes the number of points in the subtree rooted at y that are at a distance r from y ; see (2.1) for the definition.

Theorem 1.4 *Assume $\text{supp } \tilde{\mu}$ contains at least two points. Let $\mathbb{J}_\omega = \mathbb{A}_\omega$ or $\mathbb{J}_\omega = \mathbb{S}_\omega$. Then there exists a set \mathcal{D} of cardinality at most one such that the following assertions hold.*

- (i) *The operator \mathbb{J}_ω exhibits Anderson localization at all energies outside of \mathcal{D} . That is, almost surely, \mathbb{J}_ω has pure point spectrum and any eigenfunction of \mathbb{J}_ω corresponding to an energy $E \in \sigma(\mathbb{J}_\omega) \setminus \mathcal{D}$ enjoys an exponential decay estimate of the form*

$$|f(x)| \leq \frac{C e^{-\lambda|x|}}{\sqrt{w_o(|x|)}}, \quad x \in \mathcal{V}, \quad (1.8)$$

where $C, \lambda > 0$ are constants.

- (ii) *For every compact interval $I \subset \mathbb{R} \setminus \mathcal{D}$ there exist $\Omega^* \subset \Omega$ with $\mu(\Omega^*) = 1$ and $\theta > 0$ such that for every $x, y \in \mathcal{V}$, $|x| \geq |y|$, $\omega \in \Omega^*$ one has*

$$\sup_{t>0} |\langle \delta_x, \chi_I(\mathbb{J}_\omega) e^{-it\mathbb{J}_\omega} \delta_y \rangle_{\ell^2(\mathcal{V})}| \leq \frac{C e^{-\theta \text{dist}(x,y)}}{\sqrt{w_y(|x| - |y|)}}, \quad (1.9)$$

for some $C = C(y, \omega, \theta) > 0$. In particular, for all $y \in \mathcal{V}$, $\omega \in \Omega^*$, $R > 0$ one has

$$\sum_{|x| \geq R} \sup_{t>0} |\langle \delta_x, \chi_I(\mathbb{J}_\omega) e^{-it\mathbb{J}_\omega} \delta_y \rangle_{\ell^2(\mathcal{V})}| \leq \gamma e^{-\kappa R}, \quad (1.10)$$

for some $\kappa = \kappa(y) > 0$ and $\gamma = \gamma(y) > 0$.

It is proved in Sect. 5.1 that the spectrum of \mathbb{A}_ω is given by a deterministic set. It is interesting to contrast this result with the work of Klein [51] (see also [2,34] for alternative proofs), which works without the radial assumption. In that model, each vertex potential is an i.i.d. random variable, and that model exhibits absolutely

continuous spectrum in suitable energy regions for small coupling; it therefore does not exhibit localization uniformly, whereas the model in this work does. In particular, the model of [51] is more random than this one, and yet the spectral type is more regular.

Our work is motivated by the paper [44], which investigated RLM and RKM, and can be viewed as a natural continuation of [27] where discrete RBM was considered. It is worth noting that the methods of [44] are not applicable in the present setting since they are based on spectral averaging and hence rely heavily on the assumption that the random variables are absolutely continuous. Of course, in the case of random branching numbers such a hypothesis cannot be made. We stress again that RBM naturally presents the most challenging case of random models, which are commonly referred to as Bernoulli–Anderson-type models. A textbook discussion of some difficulties presented by Bernoulli-type potentials is provided in the Notes sections of Chapters 4, 7, and 12 of [6].

1.2 Background

The spectral theory of Schrödinger operators on tree graphs has attracted a lot of attention cf., e.g., [1–3, 13–18, 22, 27–30, 32–35, 42–44, 49–53, 57–59]. The recurring topic in these works is the dependence of the spectrum of differential operators on the geometry of trees, in particular, on their growth rates, edge lengths, and branching numbers. For example, Ekholm, Frank, and Kovarik established Lieb–Thirring inequalities which heavily depend on the growth rate and the global dimension of underlying trees, cf. [29], and Frank and Kovarik obtained heat kernel estimates for various trees in [30]. Evans, Harris, and Pick studied Hardy inequalities on trees in the context of eigenvalue counting for the Neumann Laplacian on bounded domains with fractal boundaries cf. [31, 32]. This topic was further developed by Naimark and Solomyak [52, 53]. As far as the discrete spectrum is concerned, Solomyak also obtained Weyl’s asymptotic formula for compact metric trees with the standard power-law behavior replaced by $c(\Gamma)\sqrt{\lambda} \log \lambda$ (this hints on mixed dimensionality of the model) with $c(\Gamma)$ depending on the tree, cf. [59]. Further, the dependence of the spectral type on the geometry was investigated by Breuer et al. [13, 16]. Exponential decay of the eigenfunctions on trees (and more general graphs) was recently discussed by Harrell and Maltsev [43]. Aizenman, Sims, and Warzel studied the effects of disorder in the geometry of trees. In particular, they considered trees with edge lengths given by $\ell_e(\omega) = e^{\lambda\omega_e}$ where $\lambda \in [0, 1]$ and $\{\omega_e\}_{e \in \mathcal{E}}$ are i.i.d. random variables, and proved in [1] that the absolutely continuous spectrum of the Laplace operator is continuous (in the sense of [1, Theorem 1.1]) at $\lambda = 0$ almost surely. That such a continuity property fails in the case of radial disorder is conjectured in [1] and proved by Hislop and Post [44]. As already mentioned earlier, the existence of absolutely continuous spectrum for the Anderson Hamiltonian on the regular trees in the regime of small disorder was shown by Klein [51] (and also by Aizenman et al. [2] as well as by Froese et al. [34]). Thematically related recent results are due to Aizenman and Warzel [4, 5] showing delocalization near the spectral edges for random Schrödinger operators on discrete trees.

The structure of the paper follows. In Sect. 2, we discuss the spectral theory of deterministic continuum operators on metric tree graphs. We use this to set notation

and to give the reader relevant background on a reduction from the metric tree graphs to Schrödinger operators on a half-line with singular potentials. We prove a localization result for these half-line operators in Sect. 3, which we then use to prove our main results for metric tree graphs in Sect. 4. The case of discrete operators on random tree graphs is taken up in Part 2.

Part 1. Anderson localization for continuum radial trees

2 Spectral theory of deterministic continuum operators

In this section we introduce deterministic Laplace operators on radial tree graphs, discuss their orthogonal decomposition, and establish several auxiliary results regarding the spectral theory of the one-dimensional half-line operators arising in such a decomposition.

To set the stage, we fix a metric rooted tree $\Gamma = (\mathcal{V}, \mathcal{E})$ with vertices \mathcal{V} , edges \mathcal{E} , root $o \in \mathcal{V}$, and edge lengths $\{\ell_e\}_{e \in \mathcal{E}}$. The shortest path connecting $x \in \Gamma$ and $y \in \Gamma$ and its length are denoted by $p(x, y)$ and $d(x, y)$, respectively, and $|x| := d(o, x)$. The generation and the branching number of a vertex v are defined by

$$\text{gen}(v) := \#\{x \in \mathcal{V} \setminus \{v\} : x \in p(o, v)\}, \quad b(v) := \begin{cases} \deg(v) - 1, & v \neq o, \\ 1 & v = o. \end{cases}$$

In other words, $\text{gen}(v)$ is the combinatorial graph distance from v to the root and $b(v)$ is the number of children of v . For an edge $e = (u, v)$, we define $\text{gen}(e) = \max(\text{gen}(u), \text{gen}(v))$. Furthermore, $T_v \subset \Gamma$ denotes the “forward” subtree of Γ rooted at v , that is, $T_v := \{x \in \Gamma : v \in p(o, x), |v| \leq |x|\}$; its branching function is given by

$$w_v(t) := \#\{x \in T_v : d(v, x) = t\}, \quad t > 0. \quad (2.1)$$

For example, given a vertex v , $w_o(|v|)$ counts the number of vertices in the same generation as v .

Hypothesis 2.1 Γ is a rooted radial metric tree with bounded branching b and bounded edge lengths, ℓ , and $q : \mathcal{V} \rightarrow \mathbb{R}$ is a bounded radial potential. More precisely:

- (i) There are constants $b^\pm \in [2, \infty)$, $\ell^\pm \in (0, \infty)$ and sequences $b := \{b_n\}_{n=0}^\infty$, $\ell := \{\ell_n\}_{n=1}^\infty$ such that
 - $b(v) = b_{\text{gen}(v)} \in [b^-, b^+] \cap \mathbb{N}$ for all $v \in \mathcal{V}$ (except, $b(o) = b_0 = 1$),
 - $\ell_e = \ell_{\text{gen}(e)} \in [\ell^-, \ell^+]$ for all $e \in \mathcal{E}$.
- (ii) There are constants $q^\pm \in \mathbb{R}$ and a sequence $\{q_n\}_{n=1}^\infty$ such that $q(v) = q_{\text{gen}(v)} \in [q^-, q^+]$.

When Γ satisfies Hypothesis 2.1, we will write $\Gamma = \Gamma_{b, \ell}$ to emphasize the dependence of Γ on the branching and length sequences.

Given Γ satisfying Hypothesis 2.1, we equip \mathbb{R}_+ with a sequence of degree two vertices $\{t_j\}_{j=1}^\infty$, where t_j denotes the distance from the root to vertices at generation j , that is,

$$t_0 := 0, \quad t_j := \sum_{i=1}^j \ell_i, \quad j > 0. \quad (2.2)$$

Then, we introduce the Sobolev spaces on such a chain of intervals

$$\widehat{H}^k(\mathbb{R}_+) := \bigoplus_{j=0}^{\infty} H^k(t_j, t_{j+1}), \quad j \in \mathbb{Z}_+, \quad k = 0, 1, 2.$$

A note on notation: throughout this paper, we write \mathbb{N} for $\{1, 2, 3, \dots\}$ and \mathbb{Z}_+ for $\mathbb{N} \cup \{0\}$. Let us note that we use the notation $\widehat{H}^k(\mathbb{R}_+)$ even though the exact composition of the space depends on the vertices $\{t_j\}_{j=0}^\infty$. Similarly, on Γ , we define

$$\widehat{H}^k(\Gamma) := \bigoplus_{e \in \mathcal{E}} H^k(e), \quad \|f\|_{\widehat{H}^k(\Gamma)}^2 := \sum_{e \in \mathcal{E}} \|f|_e\|_{H^k(e)}^2, \quad k = 0, 1, 2.$$

Notice that the elements of $\widehat{H}^k(\mathbb{R}_+)$ or $\widehat{H}^k(\Gamma)$ may be discontinuous at the vertices.

2.1 Orthogonal decomposition of radial trees

Given a radial tree $\Gamma_{b,\ell}$ and a potential q satisfying Hypothesis 2.1, we consider the self-adjoint operator $\mathbb{H} = \mathbb{H}(b, \ell, q)$ defined by

$$\begin{aligned} \mathbb{H}(b, \ell, q) &:= -\frac{d^2}{dx^2}, \quad \mathbb{H}(b, \ell, q) : \text{dom}(\mathbb{H}(b, \ell, q)) \subset L^2(\Gamma_{b,\ell}) \rightarrow L^2(\Gamma_{b,\ell}), \\ \text{dom}(\mathbb{H}(b, \ell, q)) &= \left\{ f \in \widehat{H}^2(\Gamma_{b,\ell}) : f \text{ satisfies (1.1) and (1.2)} \right\}. \end{aligned} \quad (2.3)$$

Due to the radial structure of the graph, $L^2(\Gamma_{b,\ell})$ enjoys an orthogonal decomposition into \mathbb{H} -reducing subspaces; cf. [22,52,58,59]. Namely, to every vertex $v \in \mathcal{V}$ there corresponds an \mathbb{H} -reducing subspace \mathcal{S}_v such that

$$L^2(\Gamma_{b,\ell}) = \bigoplus_{v \in \mathcal{V}} \mathcal{S}_v, \quad \mathbb{H}P_{\mathcal{S}_v} = P_{\mathcal{S}_v}\mathbb{H}, \quad (2.4)$$

where $P_{\mathcal{S}_v}$ denotes the orthogonal projection onto \mathcal{S}_v in $L^2(\Gamma_{b,\ell})$. Furthermore, each subspace \mathcal{S}_v can be further decomposed into $b_{\text{gen}(v)} - 1$ subspaces, each of which is also \mathbb{H} -reducing, that is,

$$\mathcal{S}_v = \begin{cases} \bigoplus_{k=1}^{b_{\text{gen}(v)}-1} \mathcal{L}_{v,k}, & v \neq o, \\ \mathcal{L}_o, & v = o, \end{cases} \quad (2.5)$$

and $\mathbb{H}P_{\mathcal{L}_{v,k}} = P_{\mathcal{L}_{v,k}}\mathbb{H}$, $\mathbb{H}P_{\mathcal{L}_o} = P_{\mathcal{L}_o}\mathbb{H}$. Moreover, the reduced operators are unitarily equivalent to 1D Schrödinger operators acting in $L^2(\mathbb{R}_+)$. Concretely, the operators

$$\mathbb{H}(b, \ell, q)P_{\mathcal{L}_{v,k}}, \mathbb{H}(b, \ell, q)P_{\mathcal{L}_o},$$

are unitarily equivalent to the operator

$$H(T^{\text{gen}(v)}b, T^{\text{gen}(v)}\ell, T^{\text{gen}(v)}q) \text{ acting in } L^2(t_{\text{gen}(v)}, \infty), v \in \mathcal{V} \quad (2.6)$$

where T denotes the left shift $(Tx)_n := x_{n+1}$ and

$$\begin{aligned} H(T^\varkappa b, T^\varkappa \ell, T^\varkappa q) &:= -\frac{d^2}{dx^2}, \\ H(T^\varkappa b, T^\varkappa \ell, T^\varkappa q) : \text{dom}(H(T^\varkappa b, T^\varkappa \ell, T^\varkappa q)) &\subset L^2(t_\varkappa, \infty) \rightarrow L^2(t_\varkappa, \infty) \\ \text{dom}(H(T^\varkappa b, T^\varkappa \ell, T^\varkappa q)) &= \left\{ f \in \widehat{H}^2(t_\varkappa, \infty) : f(t_\varkappa^+) = 0, \begin{array}{l} f \text{ satisfies (1.3)} \\ \text{for all } j > \varkappa \end{array} \right\} \end{aligned} \quad (2.7)$$

$$(2.8)$$

for $\varkappa \in \mathbb{Z}_+$. The unitary map

$$\mathcal{U}_{v,k} : \mathcal{L}_{v,k} \rightarrow L^2(t_{\text{gen}(v)}, \infty), v \in \mathcal{V} \setminus \{o\}, 1 \leq k \leq b_{\text{gen}(v)} - 1,$$

realizing the equivalence is defined by

$$(\mathcal{U}_{v,k}^{-1}f)(x) = \begin{cases} \frac{\exp\left(\frac{2\pi i j k}{b_{\text{gen}(v)}}\right)f(|x|)}{\sqrt{w_v(|x|)}}, & x \in T_v(j), 1 \leq j \leq b_{\text{gen}(v)}, \\ 0, & \text{otherwise,} \end{cases} \quad (2.9)$$

where $T_v(j) \subset T_v$ denotes the forward subtree determined by the j th edge emanating from the vertex v . Letting $k = 0$ in (2.9), one defines \mathcal{U}_o . We point out that $(\mathcal{U}_{v,k}^{-1}f) \in \text{dom}(\mathbb{H}(b, \ell, q))$ whenever f belongs to the domain of the operator defined in (2.6). Indeed, continuity of $\mathcal{U}_{v,k}^{-1}f$ at v is ensured by the Dirichlet condition (2.8) while the Kirchhoff condition at v is satisfied due to (2.9) and the fact that the sum of roots of unity is equal to zero. At all other vertices, one has continuity and the Kirchhoff condition by (1.3).

Combining these unitary operators together, one defines

$$\Psi_{b,\ell} := \mathcal{U}_o \oplus \bigoplus_{v \in \mathcal{V} \setminus \{o\}} \bigoplus_{k=1}^{b_{\text{gen}(v)}-1} \mathcal{U}_{v,k}, \quad (2.10)$$

and has, [52, Theorem 4.1],

$$\Psi_{b,\ell} : L^2(\Gamma_{b,\ell}) \rightarrow \bigoplus_{n=0}^{\infty} \bigoplus_{k=1}^{m(n)} L^2(t_n, \infty), \quad (2.11)$$

$$\Psi_{b,\ell} \mathbb{H}(b, \ell, q) \Psi_{b,\ell}^{-1} = \bigoplus_{n=0}^{\infty} \bigoplus_{k=1}^{m(n)} H(T^n b, T^n \ell, T^n q),$$

$$m(n) := \begin{cases} b_0 \cdot b_1 \cdots b_{n-1} \cdot (b_n - 1), & n \geq 1, \\ 1, & n = 0. \end{cases} \quad (2.12)$$

Next, we turn to the spectral analysis of $H(b, \ell, q)$ for fixed admissible b, ℓ, q . First, the eigenvalue problem for this operator can be written in terms of suitable $\text{SL}(2, \mathbb{R})$ matrices. Namely, if f is a solution to the problem

$$\begin{cases} -f'' = Ef, & f(t_0) = 0, \\ f(t_j^+) = \sqrt{b_j} f(t_j^-) & j \in \mathbb{N} \\ f'(t_j^+) = \frac{f'(t_j^-) + q_j f(t_j^-)}{\sqrt{b_j}} & j \in \mathbb{N}, \\ f \in H^2(t_{j-1}, t_j) & j \in \mathbb{N}, \end{cases} \quad (2.13)$$

then one has

$$\begin{bmatrix} f(t_j^+) \\ f'(t_j^+) \end{bmatrix} = M^E(b_j, \ell_j, q_j) \begin{bmatrix} f(t_{j-1}^+) \\ f'(t_{j-1}^+) \end{bmatrix} \text{ for all } j \in \mathbb{N}, \quad (2.14)$$

where $M^E(\beta, \lambda, \varkappa) := D(\beta)S(\varkappa)R_{\sqrt{E}}(\lambda\sqrt{E})$, $\text{Im}(\sqrt{E}) \geq 0$ and

$$D(\beta) := \begin{bmatrix} \beta^{1/2} & 0 \\ 0 & \beta^{-1/2} \end{bmatrix}, \quad S(\varkappa) := \begin{bmatrix} 1 & 0 \\ \varkappa & 1 \end{bmatrix}, \quad R_{\mu}(\varphi) := \begin{bmatrix} \cos \varphi & \frac{\sin \varphi}{\mu} \\ -\mu \sin \varphi & \cos \varphi \end{bmatrix}. \quad (2.15)$$

In this case, we can interpolate between the vertices to get

$$f(x) = f(t_{j-1}^+) \cos(\sqrt{E}(x - t_{j-1})) + \frac{f'(t_{j-1}^+) \sin(\sqrt{E}(x - t_{j-1}))}{\sqrt{E}}, \quad (2.16)$$

for all $x \in (t_{j-1}, t_j)$, $j \in \mathbb{Z}_+$. Conversely, given initial data $(f(0^+), f'(0^+))^\top$, then (2.14) and (2.16) construct a solution to the problem (2.13). Furthermore, $f \in L^2(\mathbb{R}_+)$ if and only if

$$\left\{ \begin{bmatrix} f(t_j^+) \\ f'(t_j^+) \end{bmatrix} \right\}_{j=0}^{\infty} \in \ell^2(\mathbb{Z}_+, \mathbb{C}^2).$$

2.2 Quadratic form of the model half-line operator

The following proposition describes the quadratic form of $H(b, \ell, q)$ and provides prerequisites for the Weyl criteria used in the proof of later results (e.g. Theorem 4.1).

Lemma 2.2 *Assume Hypothesis 2.1 and consider the sesquilinear form $\mathfrak{h} = \mathfrak{h}(b, \ell, q)$ defined by*

$$\mathfrak{h} : \text{dom}(\mathfrak{h}) \times \text{dom}(\mathfrak{h}) \rightarrow \mathbb{C}, \quad (2.17)$$

$$\text{dom}(\mathfrak{h}) = \left\{ f \in \widehat{H}^1(t_0, \infty) : \begin{array}{l} f(0^+) = 0, \\ \sqrt{b_j} f(t_j^-) = f(t_j^+), \quad j > 0 \end{array} \right\}, \quad (2.18)$$

$$\mathfrak{h}[u, v] = \langle u', v' \rangle_{L^2(t_0, \infty)} + \sum_{j=1}^{\infty} q_j \overline{u(t_j^-)} v(t_j^-) \quad u, v \in \text{dom}(\mathfrak{h}). \quad (2.19)$$

Then \mathfrak{h} is densely defined, closed, and bounded from below (i.e. for some $\gamma \in \mathbb{R}$ one has $\mathfrak{h}[u, u] \geq \gamma \|u\|_{L^2(t_0, \infty)}^2$, $u \in \text{dom}(\mathfrak{h})$). It is uniquely associated with the operator $H = H(b, \ell, q)$, that is,

$$\mathfrak{h}[u, v] = \langle u, Hv \rangle_{L^2(t_0, \infty)}, \quad (2.20)$$

for all $u \in \text{dom}(\mathfrak{h})$ and $v \in \text{dom}(H)$. Furthermore, there exist positive constants $c, C > 0$ such that

$$c \|u\|_{\widehat{H}^1(t_0, \infty)}^2 \leq (\mathfrak{h} - \gamma + 1)[u, u] \leq C \|u\|_{\widehat{H}^1(t_0, \infty)}^2, \quad u \in \text{dom}(\mathfrak{h}), \quad (2.21)$$

where γ is a lower bound of \mathfrak{h} . In addition, the space of compactly supported functions contained in $\text{dom}(\mathfrak{h})$ is a core of the form \mathfrak{h} .

Proof Throughout this proof we will abbreviate $\mathfrak{h} := \mathfrak{h}(b, \ell, q)$ and $H := H(b, \ell, q)$ for an admissible fixed triple (b, ℓ, q) . First, we show that \mathfrak{h} is bounded from below. If $q^- \geq 0$, the form is non-negative. Suppose that $q^- < 0$. By a standard Sobolev-type inequality (cf., e.g. [21, Corollary 4.2.10], [46, IV.1.2]) one has

$$\max \left\{ |u(t_{j-1}^+)|^2, |u(t_j^-)|^2 \right\} \lesssim C \|u\|_{L^2(t_{j-1}, t_j)}^2 + \varepsilon \|u'\|_{L^2(t_{j-1}, t_j)}^2, \quad (2.22)$$

for all $\varepsilon > 0$ and $j \in \mathbb{N}$, where $C = C(\varepsilon, \ell^-, \ell^+) > 0$. Then

$$\mathfrak{h}[u, u] \gtrsim \|u'\|_{L^2(t_0, \infty)}^2 + q^- C \|u\|_{L^2(t_0, \infty)}^2 + q^- \varepsilon \|u'\|_{L^2(t_0, \infty)}^2 \quad (2.23)$$

$$\begin{aligned} &\geq (1 + q^- \varepsilon) \|u'\|_{L^2(\mathbb{R}_+)}^2 + q^- C \|u\|_{L^2(t_0, \infty)}^2 \\ &\geq \gamma \|u\|_{L^2(t_0, \infty)}^2, \end{aligned} \quad (2.24)$$

where we chose $\varepsilon > 0$ so that $1 + q^- \varepsilon > 0$ and set $\gamma := q^- C$.

Next, we prove that \mathfrak{h} is closed, i.e., that $\text{dom}(\mathfrak{h})$ is closed with respect to the topology induced by the inner product $\mathfrak{h} - \gamma + 1$. First, using (2.23), (2.24) one infers

$$(\mathfrak{h} - \gamma + 1)[u, u] \gtrsim \|u\|_{\widehat{H}^1(t_0, \infty)}^2. \quad (2.25)$$

Suppose that $\{u_k\}_{k \geq 1} \subset \text{dom}(\mathfrak{h})$ is a Cauchy sequence with respect to the inner product $\mathfrak{h} - \gamma + 1$. In that case, it is Cauchy in $\widehat{H}^1(t_0, \infty)$ and hence has a limit $u \in \widehat{H}^1(t_0, \infty)$:

$$u_k \xrightarrow{\widehat{H}^1(t_0, \infty)} u, \quad k \rightarrow \infty. \quad (2.26)$$

In order to show that \mathfrak{h} is closed, it is enough to prove that u satisfies the vertex conditions at every vertex t_j . To that end, we notice that for all $k \in \mathbb{N}$, $j > 0$ we have $\sqrt{b_j}u_k(t_j^-) = u_k(t_j^+)$. Then, by (2.22) and (2.26) we may pass to the limit as $k \rightarrow \infty$ and obtain $\sqrt{b_j}u(t_j^-) = u(t_j^+)$ for all $j > 0$. Similarly, we get $u(t_0^+) = 0$.

The first inequality in (2.21) is already proved; see (2.25). The second one follows from the Cauchy–Schwarz inequality and the Sobolev-type estimate (2.22).

Next, we prove (2.20). Notice that the subspace

$$\{v \in \text{dom}(H) : \text{supp}(v) \text{ is compact in } [t_0, \infty)\} \subset \text{dom}(H),$$

is a core of H . Hence it is sufficient to check (2.20) for arbitrary $u \in \text{dom}(\mathfrak{h})$, $v \in \text{dom}(H)$ with $\text{supp}(v) \subset [t_0, t_K)$ for some $K \in \mathbb{N}$. One has

$$\begin{aligned} \langle u, Hv \rangle_{L^2(t_0, \infty)} &= - \sum_{j=1}^K \int_{t_{j-1}}^{t_j} \overline{u(x)} v''(x) dx \\ &= \langle u', v' \rangle_{L^2(t_0, \infty)} + \overline{u(t_0^+)} v'(t_0^+) + \sum_{j=1}^K \overline{u(t_j^+)} v'(t_j^+) - \overline{u(t_j^-)} v'(t_j^-) \\ &= \langle u', v' \rangle_{L^2(t_0, \infty)} + \sum_{j=1}^K \frac{\overline{\sqrt{b_j}u(t_j^-)} v'(t_j^-) + q_j v(t_j^-)}{\sqrt{b_j}} - \overline{u(t_j^-)} v'(t_j^-) \\ &= \mathfrak{h}[u, v]. \end{aligned}$$

□

The following Weyl-type criterion holds.

Proposition 2.3 Assume Hypothesis 2.1, and denote $\mathfrak{h} = \mathfrak{h}(b, \ell, q)$ and $H = H(b, \ell, q)$ as in Lemma 2.2. Let $D \subset \text{dom}(\mathfrak{h})$ be a dense subset with respect to the $\widehat{H}^1(t_0, \infty)$ norm (or, equivalently, with respect to the norm $\|\cdot\|_{\mathfrak{h}}^2 := (\mathfrak{h} - \gamma + 1)[\cdot, \cdot]$). Then $E \in \sigma(H)$ if and only if there exist $\{\varphi_k\}_{k=1}^\infty \subset D$ and $\{m_k\}_{k=1}^\infty \subset \mathbb{N}$ such that

$$\|\varphi_k\|_{L^2(t_0, \infty)} = 1, \quad \text{supp}(\varphi_k) \subset [t_0, t_{m_k}), \quad (2.27)$$

$$\sup_{k \in \mathbb{N}} \|\varphi_k\|_{\widehat{H}^1(t_0, \infty)} < \infty, \quad (2.28)$$

$$\sup_{\substack{g \in \text{dom}(\mathfrak{h}) \\ \|g\|_{\widehat{H}^1(t_0, \infty)} \leq 1}} (\mathfrak{h} - E)[\varphi_k, g] \rightarrow 0, \quad k \rightarrow \infty. \quad (2.29)$$

Proof Since the norm $\|\cdot\|_{\widehat{H}^1(t_0, \infty)}$ is equivalent to the form domain norm $\|\cdot\|_{\mathfrak{h}}$, (2.27), (2.29), together with the standard Weyl's criterion cf., e.g. [60, Proposition 1.4.4], yield $E \in \sigma(H)$ proving the “if” part.

To prove the “only if” part we combine Weyl's criterion and the last part of Lemma 2.2 to obtain a sequence satisfying (2.27), (2.29). Without loss of generality we may assume that $\gamma \geq 0$. In that case, one has

$$\begin{aligned} \|\varphi_k\|_{\widehat{H}^1(t_0, \infty)}^2 &\lesssim |\mathfrak{h}[\varphi_k, \varphi_k]| = \sup_{\substack{g \in \text{dom}(\mathfrak{h}) \\ \|g\|_{\mathfrak{h}}=1}} |\mathfrak{h}[\varphi_k, g]| \\ &\leq \sup_{\substack{g \in \text{dom}(\mathfrak{h}) \\ \|g\|_{\mathfrak{h}}=1}} |(\mathfrak{h} - E)[\varphi_k, g]| \\ &\quad + \sup_{\substack{g \in \text{dom}(\mathfrak{h}) \\ \|g\|_{\mathfrak{h}}=1}} |E\langle \varphi_k, g \rangle_{L^2(\mathbb{R}_+)}| \underset{k \rightarrow \infty}{=} o(1) + \mathcal{O}(1). \end{aligned}$$

Thus (2.28) holds as asserted. \square

In the sequel we will refer to the Dirichlet–Neumann truncation of the half-line operator $H(b, \ell, q)$ defined as follows

$$\begin{aligned} H^k(b, \ell, q) &:= -\frac{d^2}{dx^2}, \\ H^k(b, \ell, q) : \text{dom}(H^k(b, \ell, q)) &\subset L^2(t_0, t_k) \rightarrow L^2(t_0, t_k), \\ \text{dom}(H^k(b, \ell, q)) &= \left\{ \widehat{H}^2(t_0, t_k) : \begin{array}{l} f(t_0^+) = f'(t_k^-) = 0 \\ f \text{ satisfies (1.3) for all } 0 < j < k \end{array} \right\}. \end{aligned}$$

Proposition 2.4 *Let us fix $n \geq 1$, $E \notin \sigma(H^n(b, \ell, q))$, and suppose that u_{\pm} satisfy (1.3) for all $0 < j < n$, $-u''_{\pm} = Eu_{\pm}$, $u_-(t_0^+) = u'_+(t_n^-) = 0$, and $u'_-(t_0^+) = u_+(t_n^-) = 1$. Then the Green function of the operator $H^n(b, \ell, q)$ is given by*

$$G_n^E(x, y) = G_{[t_0, t_n]}^E(x, y) := \frac{1}{W(u_+, u_-)} \begin{cases} u_+(y)u_-(x), & y \geq x, \\ u_+(x)u_-(y), & y \leq x, \end{cases}$$

where $0 \neq W(u_+, u_-) = u'_-(t_n^-) = u_+(t_0^+)$ denotes the Wronskian of linearly independent solutions u_{\pm} . That is, $(H^n(b, \ell, q) - E)^{-1}$ is an integral operator with the kernel $G_{[t_0, t_n]}^E$.

Proof For a fixed $g \in L^2(t_0, t_n)$ the unique nonzero function u satisfying all vertex conditions and solving the non-homogeneous differential equation $-u'' - Eu = g$ is given by

$$u(y) = [\mathcal{R}_E g](y) := \int_{t_0}^{t_n} G_{[t_0, t_n]}^E(x, y) g(x) dx.$$

Evidently, the operator \mathcal{R}_E is bounded and

$$(H^n(b, \ell, q) - E)\mathcal{R}_E = \mathcal{R}_E(H^n(b, \ell, q) - E) = I_{L^2(t_0, t_n)},$$

as asserted. Finally, evaluating the Wronskian at t_0 and t_{n+1} , we get

$$W(u_+, u_-) = u'_-(t_n^-) = u_+(t_0^+)$$

(see also [44, Lemma D.12]). □

3 Proof of localization for half-line random operators

The main goal of this section is to prove dynamical and spectral localization for the random half-line operators H_ω arising in the orthogonal decomposition of \mathbb{H}_ω . Theorem 3.5 ensures positivity of the Lyapunov exponent outside of a discrete set \mathfrak{D} . In Theorem 3.13 we prove spectral localization and SULE for H_ω . Finally, we conclude with the proof of Theorem 1.1, which addresses dynamical localization.

3.1 Description of random models

The *random branching model* (abbreviated RBM) is described by a family of Laplace operators subject to Neumann–Kirchhoff vertex conditions on radial metric trees with random branching numbers. In other words, we assume Hypothesis 2.1 with the following parameters

$$b = \{b_\omega(n)\}_{n \in \mathbb{N}} \subset \{2, \dots, d\}, \quad d \geq 3, \quad \ell^- = \ell^+ = 1, \quad q^- = q^+ = 0,$$

where $\{b_\omega(n)\}_{n \in \mathbb{N}}$ is a sequence of independent and identically distributed random variables whose common distribution contains at least two points in its support.

The *random lengths model* (RLM) is given by a family of the Neumann–Kirchhoff Laplace operators on radial metric trees with random edge lengths. That is, we assume Hypothesis 2.1 with

$$b^- = b^+ = d, \quad \ell = \{\ell_\omega(n)\}_{n \in \mathbb{N}} \subset [\ell^-, \ell^+], \quad q^- = q^+ = 0,$$

where $\{\ell_\omega(n)\}_{n \in \mathbb{N}}$ is a sequence of independent and identically distributed random variables whose common distribution contains at least two points in its support.

The *random Kirchhoff model* (RKM) is given by the Laplace operators subject to random δ -type vertex conditions. That is, we assume Hypothesis 2.1 with

$$b^- = b^+ = d, \quad \ell^- = \ell^+ = 1, \quad q = \{q_\omega(n)\}_{n \in \mathbb{N}} \subset [q^-, q^+],$$

where $\{q_\omega(n)\}_{n \in \mathbb{N}}$ is a sequence of independent and identically distributed random variables whose common distribution contains at least two points in its support.

In order to unify these models we consider three-dimensional random variables with common distribution $\tilde{\mu}$.

Hypothesis 3.1 *Let $\tilde{\mu}$ be a probability measure with*

$$\text{supp}(\tilde{\mu}) \subset \mathcal{A} := \{b^-, \dots, b^+\} \times [\ell^-, \ell^+] \times [q^-, q^+].$$

Suppose that $\text{supp}(\tilde{\mu})$ contains at least two distinct points, and let $(\Omega, \mu) := (\mathcal{A}^{\mathbb{N}}, \tilde{\mu}^{\mathbb{N}})$.

Remark 3.2 We notice that

- RBM arises when $\text{supp } \tilde{\mu} \subseteq \{b^-, \dots, b^+\} \times \{1\} \times \{0\}$,
- RLM arises when $\text{supp } \tilde{\mu} \subseteq \{d\} \times [\ell^-, \ell^+] \times \{0\}$,
- RKM arises when $\text{supp } \tilde{\mu} \subseteq \{d\} \times \{1\} \times [q^-, q^+]$.

For $\omega \in \Omega$ we denote the components of ω as $\omega(n) = (b_\omega(n), \ell_\omega(n), q_\omega(n))$, since we will use them to define the branching, edge lengths, and Kirchhof potential of an operator. In particular, the vertices in \mathbb{R}_+ are denoted $t_\omega(n)$. Given ω , define the operators $\mathbb{H}_\omega = \mathbb{H}(b_\omega, \ell_\omega, q_\omega)$ acting in $L^2(\Gamma_{b_\omega, \ell_\omega})$ as in (2.3). Similarly, for $j \in \mathbb{Z}_+$, define

$$H_{T^j \omega} := H(T^j b_\omega, T^j \ell_\omega, T^j q_\omega) \text{ acting in } L^2(t_\omega(j), \infty),$$

as in (2.7), (2.8) and let $\mathfrak{h}_{T^j \omega}$ denote the corresponding quadratic forms.

3.2 Positivity of Lyapunov exponents via Fürstenberg's theorem

Inspired by (2.14) and (2.15), we introduce an $\text{SL}(2, \mathbb{R})$ -cocycle over T (the left shift $\Omega \rightarrow \Omega$) as follows. First, let \mathcal{A} , b^\pm , ℓ^\pm , and q^\pm be as in Hypothesis 3.1. For each $E \in \mathbb{R}$, (2.14), (2.15) lead us to define $M^E : \mathcal{A} \rightarrow \text{SL}(2, \mathbb{R})$ by

$$\mathcal{A} \ni \alpha = (\beta, \lambda, \varkappa) \mapsto M^E(\alpha) = D(\beta)S(\varkappa)R_{\sqrt{E}}(\lambda\sqrt{E}). \quad (3.1)$$

This induces a map $M^E : \Omega \rightarrow \text{SL}(2, \mathbb{R})$ via $M^E(\omega) = M^E(\omega(1))$, and then a skew product

$$(T, M^E) : \Omega \times \mathbb{R}^2 \rightarrow \Omega \times \mathbb{R}^2, \quad (T, M^E)(\omega, v) = (T\omega, M^E(\omega)v).$$

Then denoting the n -step transfer matrix by

$$M_n^E(\omega) = \prod_{r=n-1}^0 M^E(T^r \omega) = M^E(T^{n-1} \omega) \cdots M^E(T\omega) M^E(\omega), \quad n \in \mathbb{N}, \quad (3.2)$$

we note that the iterates over the skew product are given by $(T, M^E)^n = (T^n, M_n^E)$. The Lyapunov exponent is defined by

$$L(E) := \lim_{n \rightarrow \infty} \frac{1}{n} \int_{\Omega} \log \|M_n^E(\omega)\| d\mu(\omega). \quad (3.3)$$

By Kingman's Subadditive Ergodic Theorem we have

$$L(E) = \lim_{n \rightarrow \infty} F_n(\omega, E); \quad F_n(\omega, E) := \frac{1}{n} \log \|M_n^E(\omega)\|, \quad (3.4)$$

for μ -almost every ω .

Remark 3.3 Let us note that there are two natural cocycles that one can work with here. In addition to the discrete cocycle just described, there is also the continuum cocycle \tilde{M}^E defined by

$$\tilde{M}_x^E(\omega) : \begin{bmatrix} u(0^+) \\ u'(0^+) \end{bmatrix} \mapsto \begin{bmatrix} u(x^+) \\ u'(x^+) \end{bmatrix}$$

whenever $-u'' = Eu$ and u satisfies the vertex conditions defining $\text{dom}(H_\omega)$. Evidently,

$$M_n^E(\omega) = \tilde{M}_{t_\omega(n)}^E(\omega).$$

This leads to a simple relationship between the Lyapunov exponents of M^E and \tilde{M}^E . By Birkhoff's Ergodic Theorem,

$$\lim_{n \rightarrow \infty} \frac{1}{n} t_\omega(n) = \langle \ell \rangle := \int_{\mathcal{A}} \alpha_2 d\tilde{\mu}(\alpha),$$

the average length. Then, one has

$$L(E) = \tilde{L}(E) \cdot \langle \ell \rangle. \quad (3.5)$$

Our next goal is to show that Lyapunov exponents are positive away from a discrete set of energies. To that end, we first recall Fürstenberg's Theorem and some related facts. In order to state Fürstenberg's Theorem, let us recall that a few definitions from the general theory. A group $G \subseteq \text{SL}(2, \mathbb{R})$ is called *strongly irreducible* if there does not exist a finite set $\Lambda \subseteq \mathbb{RP}^1$ such that $\{gv : v \in \Lambda\} = \Lambda$ for all $g \in G$; G is called *contracting* if there exist $g_n \in G$, $n \geq 1$ such that $\|g_n\|^{-1} g_n$ converges to a rank-one operator as $n \rightarrow \infty$. Given Borel probability measures ν_k supported in $\text{SL}(2, \mathbb{R})$, $k \geq 1$, we say $\nu_k \rightarrow \nu$ *weakly and boundedly* if

$$\int_{\|M\| \geq N} \log^+ \|M\| d\nu_k(M) + \int_{\|M\| \geq N} \log^+ \|M\| d\nu(M) \rightarrow 0$$

as $N \rightarrow \infty$, uniformly in k and

$$\int f \, dv_k \rightarrow \int f \, dv$$

for all continuous complex-valued functions $f : \mathrm{SL}(2, \mathbb{R}) \rightarrow \mathbb{C}$ having compact support.

Theorem 3.4 *Let ν be a probability measure on $\mathrm{SL}(2, \mathbb{R})$ satisfying*

$$\int \log \|M\| \, d\nu(M) < \infty.$$

Let G_ν be the smallest closed subgroup of $\mathrm{SL}(2, \mathbb{R})$ that contains $\mathrm{supp} \, \nu$.

- (i) [36, Theorem 8.6] *Assume that G_ν is not compact and that it is strongly irreducible. Then the Lyapunov exponent $L(\nu)$ associated with ν is positive.*
- (ii) [37, Theorem B] *Assume that the set*

$$\mathrm{Fix}(G_\nu) := \left\{ V \in \mathbb{RP}^1 : MV = V \text{ for every } M \in G_\nu \right\}$$

contains at most one element. If $\nu_k \rightarrow \nu$ weakly and boundedly, then $L(\nu_k) \rightarrow L(\nu)$ as $k \rightarrow \infty$.

In the present setting, we have a one-parameter family of measures induced on $\mathrm{SL}(2, \mathbb{R})$, namely, we consider ν_E , the pushforward of $\tilde{\mu}$ under the map M^E in (3.1).

Theorem 3.5 *Assume Hypothesis 3.1. Then there is a discrete set $\mathcal{D} \subseteq \mathbb{R}$ such that $G = G_{\nu(E)}$ enjoys the following properties for $E \in \mathbb{R} \setminus \mathcal{D}$.*

- (i) *G is noncompact*
- (ii) *G is strongly irreducible*
- (iii) *G is contracting*
- (iv) *$\mathrm{Fix}(G) = \emptyset$*

In particular, L is continuous and positive on $\mathbb{R} \setminus \mathcal{D}$.

Proof In view of Theorem 3.4, positivity follows from (i) and (ii), while continuity on $\mathbb{R} \setminus \mathcal{D}$ follows from (iv). Moreover, (ii) \implies (iv), so we only need to prove (i)–(iii). Write

$$\begin{aligned} M^E(\beta, \lambda, \varkappa) &= D(\beta)S(\varkappa)R_{\sqrt{E}}(\lambda\sqrt{E}) \\ &= \begin{bmatrix} \sqrt{\beta} & 0 \\ 0 & \frac{1}{\sqrt{\beta}} \end{bmatrix} \begin{bmatrix} 1 & 0 \\ \varkappa & 1 \end{bmatrix} \begin{bmatrix} \cos(\lambda\sqrt{E}) & \frac{\sin(\lambda\sqrt{E})}{\sqrt{E}} \\ -\sqrt{E} \sin(\lambda\sqrt{E}) & \cos(\lambda\sqrt{E}) \end{bmatrix} \\ &= \begin{bmatrix} \beta^{1/2} \cos(\lambda\sqrt{E}) & \beta^{1/2} \frac{\sin(\lambda\sqrt{E})}{\sqrt{E}} \\ \varkappa \beta^{-1/2} \cos(\lambda\sqrt{E}) - \beta^{-1/2} \sqrt{E} \sin(\lambda\sqrt{E}) & \frac{\varkappa \sin(\lambda\sqrt{E})}{\beta^{1/2} \sqrt{E}} + \beta^{-1/2} \cos(\lambda\sqrt{E}) \end{bmatrix} \end{aligned}$$

Now, let $(b_1, \ell_1, q_1) \neq (b_2, \ell_2, q_2)$ be distinct elements of $\mathrm{supp} \, \tilde{\mu}$, abbreviate

$$M_j = M_j(E) := M^E(b_j, \ell_j, q_j),$$

and define the commutator

$$g = g(E) = [M_1, M_2] = M_1 M_2 - M_2 M_1.$$

To conclude the proof, it suffices to show that $g(E)$ does not vanish identically. Concretely, it is easy to see that the matrices M_j are analytic functions of E with non-constant trace and that the entries of M_j are real whenever $\operatorname{tr} M_j \in [-2, 2]$. Thus, the matrices $M_j(E)$ satisfy the first three hypotheses of [20, Theorem 2.1], so, if $g(E)$ does not vanish identically, we can conclude that there is a discrete set \mathfrak{D} such that (i)–(iii) hold for $E \in \mathbb{R} \setminus \mathfrak{D}$ by [20, Theorem 2.1].

To that end, suppose for the purpose of establishing a contradiction that g vanishes identically. In particular, the upper left matrix element $g_{11}(E)$ vanishes identically. One may calculate $g_{11}(E)$ directly:

$$\begin{aligned} g_{11}(E) &= b_1^{1/2} \frac{\sin(\ell_1 \sqrt{E})}{\sqrt{E}} \left(q_2 b_2^{-1/2} \cos(\ell_2 \sqrt{E}) - b_2^{-1/2} \sqrt{E} \sin(\ell_2 \sqrt{E}) \right) \\ &\quad - b_2^{1/2} \frac{\sin(\ell_2 \sqrt{E})}{\sqrt{E}} \left(q_1 b_1^{-1/2} \cos(\ell_1 \sqrt{E}) - b_1^{-1/2} \sqrt{E} \sin(\ell_1 \sqrt{E}) \right). \end{aligned}$$

For ease of notation, write $r_1 = b_2^{1/2}/b_1^{1/2}$, $r_2 = b_1^{1/2}/b_2^{1/2}$, and $w = \sqrt{E}$. Expanding the trigonometric functions, we get

$$\begin{aligned} g_{11} &= \frac{q_2 r_2}{4i w} (e^{i \ell_2 w} + e^{-i \ell_2 w}) (e^{i \ell_1 w} - e^{-i \ell_1 w}) \\ &\quad - \frac{q_1 r_1}{4i w} (e^{i \ell_1 w} + e^{-i \ell_1 w}) (e^{i \ell_2 w} - e^{-i \ell_2 w}) \\ &\quad - \frac{r_1 - r_2}{4} (e^{i \ell_1 w} - e^{-i \ell_1 w}) (e^{i \ell_2 w} - e^{-i \ell_2 w}). \end{aligned}$$

Thus,

$$\begin{aligned} 4i w^2 g_{11} &= \left(q_2 r_2 w - q_1 r_1 w - i w^2 (r_1 - r_2) \right) e^{i(\ell_1 + \ell_2)w} \\ &\quad + \left(q_1 r_1 w - q_2 r_2 w - i w^2 (r_1 - r_2) \right) e^{-i(\ell_1 + \ell_2)w} \\ &\quad + \left(q_2 r_2 w + q_1 r_1 w + i w^2 (r_1 - r_2) \right) e^{i(\ell_1 - \ell_2)w} \\ &\quad + \left(-q_1 r_1 w - q_2 r_2 w + i w^2 (r_1 - r_2) \right) e^{-i(\ell_1 - \ell_2)w}. \end{aligned} \tag{3.6}$$

Since g_{11} vanishes identically and $\ell_1, \ell_2 > 0$, this forces

$$\begin{aligned} q_2 r_2 w - q_1 r_1 w - i w^2 (r_1 - r_2) &\equiv 0 \\ q_1 r_1 w - q_2 r_2 w - i w^2 (r_1 - r_2) &\equiv 0 \end{aligned}$$

It is easy to see that this yields $r_1 = r_2$ (hence $b_1 = b_2$) and $q_1 = q_2$. Since $(b_1, \ell_1, q_1) \neq (b_2, \ell_2, q_2)$, we must have $\ell_1 \neq \ell_2$. Going back to (3.6), this implies

$$\begin{aligned} q_2 r_2 w + q_1 r_1 w + i w^2 (r_1 - r_2) &\equiv 0 \\ -q_1 r_1 w - q_2 r_2 w + i w^2 (r_1 - r_2) &\equiv 0. \end{aligned}$$

and hence $q_1 = q_2 = 0$. Writing $b_1 = b_2 =: b$, and substituting $q_1 = q_2 = 0$, we may directly calculate g :

$$g(E) = \begin{bmatrix} 0 & \frac{b-1}{\sqrt{E}} \sin((\ell_2 - \ell_1)\sqrt{E}) \\ \frac{b-1}{b} \sqrt{E} \sin((\ell_2 - \ell_1)\sqrt{E}) & 0 \end{bmatrix} \quad (3.7)$$

which clearly only vanishes on the discrete set

$$\mathfrak{D} = \{(\ell_1 - \ell_2)^{-2} \pi^2 k^2 : k \in \mathbb{Z}_+\},$$

a contradiction. \square

The proof above implicitly uses the following statement.

Lemma 3.6 *Suppose $\{a_j : j = 0, \dots, n\}$ is a set of $n + 1$ distinct complex numbers and $\{p_j : j = 0, \dots, n\}$ are polynomials in z . Then, the function*

$$Q(z) := \sum_{j=0}^n p_j(z) e^{a_j z}$$

vanishes identically if and only if $p_j \equiv 0$ for each j .

Proof Write $D = d/dz$ and $M = \max(\deg(p_j))$. Suppose on the contrary that

$$p_0(z) e^{a_0 z} \equiv \sum_{j=1}^n p_j(z) e^{a_j z}$$

with $p_0 \not\equiv 0$. Notice that $\prod_{j=1}^n (D - a_j)^{M+1}$ annihilates the right hand side. However, if $b \neq a_0$, one readily verifies that

$$(D - b)[p_0(z) e^{a_0 z}] = \tilde{p}_0(z) e^{a_0 z},$$

where \tilde{p}_0 has the same degree as p_0 . Consequently, a straightforward induction implies that

$$\prod_{j=1}^n (D - a_j)^{M+1} [p_0(z) e^{a_0 z}]$$

does not vanish identically, a contradiction. \square

Remark 3.7 Let us make a few comments about the proof of Theorem 3.5.

- (1) Since the argument above is soft, we do not get any information about \mathfrak{D} , except that \mathfrak{D} is discrete. However, in concrete situations in which one has more information, one can say more. For example, the g from (3.7) corresponds to the RLM; we can explicitly see that $\mathfrak{D} = \{(\ell_1 - \ell_2)^{-2} \pi^2 k^2 : k \in \mathbb{Z}_+\}$. For another example, in the RBM, one has $\text{supp } \tilde{\mu} \subseteq \{b_-, \dots, b_+\} \times \{1\} \times \{0\}$, so one can choose $(b_1, 1, 0) \neq (b_2, 1, 0) \in \text{supp } \tilde{\mu}$. After some calculations, one obtains

$$\det g = -\frac{(b_1 - b_2)^2}{b_1 b_2} \sin^2(\sqrt{E}),$$

so Fürstenberg's Theorem holds away from $\mathfrak{D} = \{\pi^2 k^2 : k \in \mathbb{Z}_+\}$. In this setting there exists a finite set of invariant directions at these special energies. That said, we note that the Lyapunov exponent is still positive by direct calculation.

- (2) Let us also remark that the transfer matrices may be bounded at a discrete set of energies (compare [25]). For example, take parameters $(b_1, \ell_1, q_1) = (2, 1, 0)$ and $(b_2, \ell_2, q_2) = (2, 3, 0)$. Then, at energies $E = \frac{1}{4} \pi^2 (2k + 1)^2$ with $k \in \mathbb{Z}_+$, M_1 and M_2 are commuting and elliptic.¹ In particular, the transfer matrices at these energies are uniformly bounded, so [24, Corollaries 2.1 and 2.2] suggest that dynamical localization as formulated in Theorem 1.2.(ii) cannot hold without excluding these energies.

Remark 3.8 As far as spectral localization is concerned, it suffices to ensure that for every compact interval $I \in \mathbb{R} \setminus \mathfrak{D}$, almost surely all generalized eigenvalues exhibit Lyapunov behavior. We will construct a full measure set $\Omega^* \subset \Omega$ such that one has

$$0 < L(E) = \lim_{n \rightarrow \infty} \frac{1}{n} \log \|M_n^E(\omega)\|$$

for every generalized eigenvalue $E \in I$ of H_ω [$M_n^E(\omega)$ is defined in (3.2)]. As discussed in [19], the work of Gorodetski and Kleptsyn [40] shows that dropping the assumption that E is a generalized eigenvalue invalidates the above assertion.

3.3 Dynamical localization for half-line operators

Our approach relies on the Large Deviation Theorem (LDT) [19, Theorem 3.1]. Although this is not stated explicitly in [19], the LDT and its corollaries [19, Theorem 4.1, Corollary 5.3, (5.13)] are applicable whenever the conditions of the Fürstenberg Theorem are met, the corresponding subgroup is contracting and the transfer matrices satisfy Lipschitz estimates which are supplied by the following lemma.

Lemma 3.9 *Fix a compact interval $I \subseteq \mathbb{R}$. There are constants $C > 0$, $\rho > 0$ such that*

$$\|M_n^E(\omega) - M_n^{E'}(\omega')\| \leq C n \rho^{n-1} (|E - E'| + \|\omega - \omega'\|_\infty)$$

¹ I.e., $|\text{tr } M_j| < 2$.

for all $\omega, \omega' \in \Omega$, $E, E' \in I$, and $n \in \mathbb{Z}_+$. The constants depend only on I and $\text{supp } \tilde{\mu}$. Consequently,

$$|F_n(\omega, E) - F_n(\omega', E')| \leq C\rho^{n-1}(|E - E'| + \|\omega - \omega'\|_\infty), \quad (3.8)$$

where F_n is defined as in (3.4).

Proof Let $n, E, E', \alpha = (\beta, \varkappa, \lambda) \in \mathcal{A}$, and $\alpha' = (\beta', \varkappa', \lambda') \in \mathcal{A}$ be given. One immediately has

$$\|S(\varkappa) - S(\varkappa')\| = |\varkappa - \varkappa'| \quad (3.9)$$

and

$$\|D(\beta) - D(\beta')\| = \left| \sqrt{\beta} - \sqrt{\beta'} \right| \leq \frac{1}{2\sqrt{2}} |\beta - \beta'| \quad (3.10)$$

since $\beta, \beta' \geq 2$. Writing $\kappa = \sqrt{E}$, and $\kappa' = \sqrt{E'}$, we get

$$\begin{aligned} \|R_\kappa(\lambda\kappa) - R_{\kappa'}(\lambda'\kappa')\| &\leq \|R_\kappa(\kappa\lambda) - R_{\kappa'}(\kappa'\lambda)\| + \|R_{\kappa'}(\kappa'\lambda) - R_{\kappa'}(\kappa'\lambda')\| \\ &\leq C(\ell^\pm, I)(|E - E'| + |\lambda - \lambda'|). \end{aligned} \quad (3.11)$$

Using the triangle inequality to change a single one-step transfer matrix at a time, one has

$$\begin{aligned} \|M_n^E(\omega) - M_n^{E'}(\omega')\| &\leq \sum_{k=0}^{n-1} \left\| M_{n-k-1}^{E'}(T^{k+1}\omega')(M_1^E(T^k\omega) - M_1^{E'}(T^k\omega'))M_k^E(\omega) \right\|, \end{aligned}$$

where T is the left shift operator. Writing

$$\rho = \sup \left\{ \|M_1^E(\omega)\| : E \in I, \omega \in \Omega \right\}, \quad (3.12)$$

we can estimate the first and third factors by ρ^{n-k-1} and ρ^k respectively. On other hand, (3.9), (3.10), and (3.11) yield

$$\|M_1^E(T^k\omega) - M_1^{E'}(T^k\omega')\| \leq C(|E - E'| + \|\omega - \omega'\|_\infty),$$

so, putting everything together, we have

$$\begin{aligned} \|M_n^E(\omega) - M_n^{E'}(\omega')\| &\leq \sum_{k=0}^{n-1} C\rho^{n-1}(|E - E'| + \|\omega - \omega'\|_\infty) \\ &= Cn\rho^{n-1}(|E - E'| + \|\omega - \omega'\|_\infty), \end{aligned}$$

proving the first inequality. The second follows from this and the statement $|\log a - \log b| \leq |a - b|$ for $a, b \geq 1$. \square

Having established Theorem 3.5 and Lemma 3.9, we may utilize the LDT in our setting. In particular, we have the following:

Theorem 3.10 Assume Hypothesis 3.1 holds true.

(i) [19, Theorem 3.1] For any $\varepsilon > 0$, there exist $C, \eta > 0$ such that

$$\mu \left\{ \omega \in \Omega : \left| L(E) - \frac{1}{n} \log \|M_n^E(\omega)\| \right| \geq \varepsilon \right\} \leq C e^{-\eta n}, \quad (3.13)$$

for all $n \geq 0$ and all $E \in I$.

(ii) [19, Theorem 4.1] There exist constants $C = C(I, \tilde{\mu})$, $\beta = \beta(I, \tilde{\mu}) > 0$ such that

$$|L(E) - L(E')| \leq C |E - E'|^\beta, \quad E, E' \in I. \quad (3.14)$$

(iii) [19, Corollary 5.3] For every $\varepsilon \in (0, 1)$ there exists a full measure set $\Omega_1(\varepsilon)$ with $\mu(\Omega_1(\varepsilon)) = 1$ such that for every $\omega \in \Omega_1(\varepsilon)$ there exists $n_1 = n_1(\varepsilon, \omega)$ such that

$$\frac{1}{n} \log \|M_n^E(T^{\zeta_0} \omega)\| \leq L(E) + \varepsilon, \quad (3.15)$$

for any $\zeta_0 \in \mathbb{Z}_+$ and $n \geq \max(n_1, \log^2(\zeta_0 + 1))$.

(iv) For every $\varepsilon \in (0, 1)$ there exists $\Omega_2(\varepsilon) \subseteq \Omega$, $\mu(\Omega_2(\varepsilon)) = 1$ with the following property: For every $\omega \in \Omega_2(\varepsilon)$, there exists $n_2 = n_2(\omega, \varepsilon)$ such that

$$\left| L(E) - \frac{1}{n^2} \sum_{s=0}^{n^2-1} \frac{\log \|M_n^E(T^{\zeta+sn} \omega)\|}{n} \right| < \varepsilon, \quad (3.16)$$

for all $\zeta \in \mathbb{Z}_+$, $n \geq \max(n_2, \log^{\frac{2}{3}}(\zeta + 1))$, and $E \in I$.

Part (iii) yields

$$\mu \left\{ \omega : \text{for all } E \in I, \limsup_{n \rightarrow \infty} \frac{1}{n} \log \|M_n^E(\omega)\| \leq L(E) \right\} = 1. \quad (3.17)$$

This fact may also be derived from the Craig–Simon approach [23] (see also [45]). Our main focus is on showing

$$\mu \left\{ \omega : \begin{array}{l} \liminf_{n \rightarrow \infty} \frac{1}{n} \log \|M_n^E(\omega)\| \geq L(E) \\ \text{for all generalized eigenvalues } E \in I \end{array} \right\} = 1.$$

The next key step is an analog of the elimination of double resonances. Let us note that we do not use the typical formulation of double resonances (cf., e.g., [47, (9.21)]), since our ultimate goal is to work with transfer matrices in order to apply the Avalanche Principle. The resonances we wish to exclude are those for which there are large disjoint intervals $I_1, I_2 \subseteq \mathbb{Z}$ so that some energy E is very close to an eigenvalue of H_ω restricted to I_1 , and the norm of the transfer matrix across I_2 at energy E deviates substantially from $\exp(|I_2|L(E))$. In particular, we would like to show that

this event occurs with very small probability, see [12]. We shall make this precise and quantitative in Theorem 3.11.

By convention, we write $\|(H_\omega^n - E)^{-1}\|_{\mathcal{B}(L^2(t_0, t_n))} = +\infty$ whenever $E \in \sigma(H_\omega^n)$. Let us recall $F_n(\omega, E)$ from (3.4), and abbreviate $\bar{K} := \lfloor K^{\log K} \rfloor$.

Theorem 3.11 *Given $\varepsilon \in (0, 1)$, $N \in \mathbb{N}$, let*

$$\mathcal{D}_N(\varepsilon) := \left\{ \omega \in \Omega : \begin{array}{l} \text{for some } \zeta \in \mathbb{Z}_+, E \in I, \\ K \geq \max\{N, \log^2(\zeta + 1)\}, 0 < n \leq K^9, \text{ one has:} \\ \left\{ \begin{array}{l} \|(H_\omega^{\zeta+n} - E)^{-1}\|_{\mathcal{B}(L^2(t_0, t_{\zeta+n}))} \geq e^{K^2} \\ \text{and } |F_m(T^{r+\zeta}\omega, E)| \leq L(E) - \varepsilon \\ \text{for some } K^{10} \leq r \leq \bar{K}, m \in \{K, 2K\} \end{array} \right. \end{array} \right\}$$

Then there exist $C = C(\varepsilon) > 0$, $\eta(\varepsilon) > 0$ such that

$$\mu(\mathcal{D}_N(\varepsilon)) \leq C e^{-\eta N}. \quad (3.18)$$

In particular, one has

$$\mu(\Omega_3(\varepsilon)) = 1 \text{ where } \Omega_3(\varepsilon) := \Omega \setminus \limsup_{N \rightarrow \infty} \mathcal{D}_N(\varepsilon). \quad (3.19)$$

Proof Let us fix

$$\zeta \in \mathbb{Z}_+, K \geq \max\{N, \log^2(\zeta + 1)\}, 0 < n \leq K^9, K^{10} \leq r \leq \bar{K}, j \in \{1, 2\}, \quad (3.20)$$

and denote

$$\mathcal{D}_j(K, n, r, \zeta) := \left\{ \omega \in \Omega : \begin{array}{l} \text{for some } E \in I, \text{ one has} \\ \|(H_\omega^{\zeta+n} - E)^{-1}\|_{\mathcal{B}(L^2(t_0, t_{\zeta+n}))} \geq e^{K^2} \text{ and} \\ |F_{jK}(T^{r+\zeta}\omega, E)| \leq L(E) - \varepsilon \end{array} \right\}$$

In order to estimate $\mu(\mathcal{D}_j(K, n, r, \zeta))$, we pick $\omega \in \mathcal{D}_j(K, n, r, \zeta)$, consider the corresponding $E \in I$, and notice that (due to the resolvent bound) E is close to an eigenvalue of the Dirichlet–Neumann truncation, that is,

$$|E - E_0| \leq e^{-K^2} \text{ for some } E_0 \in \sigma(H_\omega^{\zeta+n}). \quad (3.21)$$

Combining (3.8), (3.14), (3.21), and choosing N (hence K) sufficiently large we obtain

$$F_{jK}(T^{\zeta+r}\omega, E_0) \leq L(E_0) - \frac{\varepsilon}{2},$$

whenever $\omega \in \mathcal{D}_j(K, n, r, \zeta)$ and $E_0 = E_0(\omega_1, \dots, \omega_{\zeta+n})$ is as in (3.21). In other words

$$\mathcal{D}_j(K, n, r, \zeta) \subset \widehat{\mathcal{D}}_j(K, n, r, \zeta),$$

where

$$\widehat{\mathcal{D}}_j(K, n, r, \zeta) := \bigcup_{E_0 \in \sigma(H_\omega^{\zeta+n}) \cap \widehat{I}} \left\{ \omega \in \Omega : \frac{\varepsilon}{2} \leq L(E_0) - F_{jK}(T^{\zeta+r} \omega, E_0) \right\},$$

where $\widehat{I} := [\min I - 1, \max I + 1]$. We note that $H_\omega^{\zeta+n}$ and the standard Dirichlet Laplacian $H_D^{\zeta+n}$ on $(t_0, t_{\zeta+n})$ are self-adjoint extensions of a symmetric (minimal) operator with deficiency indices $(2(\zeta + n), 2(\zeta + n))$ (Sect. 2.1, [10]). Then the spectral shift for these two operators is at most $2(\zeta + n)$, see [8, Lemma 9.3.2 p.214, Theorem 9.3.3, p. 215]. Combining this with an explicit computation of eigenvalues of $H_D^{\zeta+n}$ we get

$$\#(\sigma(H_\omega^{\zeta+n}) \cap \widehat{I}) \leq C|\widehat{I}|(n + \zeta),$$

where $C > 0$ is a universal constant [we recall from (2.2) that $\ell^-(\zeta + n) \leq |t_{\zeta+n}| \leq \ell^+(\zeta + n)$]. Then using (3.13) and $[0, \zeta + n] \cap [\zeta + r, \zeta + r + jK] = \emptyset$, we estimate

$$\mu(\widehat{\mathcal{D}}_j(K, n, r, \zeta)) \leq C(n + \zeta)e^{-\eta K} \leq C(K^9 + e^{\sqrt{K}})e^{-\eta K} \leq Ce^{-\eta_1 K},$$

for some $\eta_1 = \eta_1(\varepsilon) > 0$. Clearly, one has

$$\mu(\mathcal{D}_N(\varepsilon)) \leq \sum_{K, n, r, \zeta, j \text{ as in (3.20)}} \mu(\widehat{\mathcal{D}}_j(K, n, r, \zeta)).$$

Then for a fixed K , the summation with respect to n, r introduces a subexponential number of terms bounded by $e^{-\eta_1 K}$, and summation with respect to ζ introduces no more than $\lceil e^{\sqrt{K}} \rceil$ terms bounded by $e^{-\eta_1 K}$ (the precise calculation is carried out in the proof of [19, Proposition 6.1]). Thus (3.18) holds as asserted, which together with the Borel–Cantelli lemma yields (3.19). \square

Let us recall the Avalanche Principle employed in the proof of Theorem 3.13.

Lemma 3.12 (Avalanche Principle) *Let $A^{(1)}, \dots, A^{(n)}$ be a finite sequence in $\mathrm{SL}(2, \mathbb{R})$ satisfying the following conditions:*

$$\begin{aligned} \min_{1 \leq j \leq n} \|A^{(j)}\| &\geq \lambda > n, \\ \max_{1 \leq j < n} \left| \log \|A^{(j+1)}\| + \log \|A^{(j)}\| - \log \|A^{(j+1)}A^{(j)}\| \right| &< \frac{1}{2} \log \lambda. \end{aligned}$$

Then for some absolute constant $C > 0$ one has

$$\left| \log \|A^{(n)} \dots A^{(1)}\| + \sum_{j=2}^{n-1} \log \|A^{(j)}\| - \sum_{j=1}^{n-1} \log \|A^{(j+1)}A^{(j)}\| \right| \leq C \frac{n}{\lambda}.$$

See [39, Proposition 2.2] for a proof of Lemma 3.12.

In order to streamline notation, we use the shorthand t_n for the point $t_\omega(n)$.

Theorem 3.13 *There exist a discrete set $\mathfrak{D} \subset \mathbb{R}$ and a set $\tilde{\Omega} \subset \Omega$ with $\mu(\tilde{\Omega}) = 1$ such that for every compact interval $I \subset \mathbb{R} \setminus \mathfrak{D}$ and every $\omega \in \tilde{\Omega}$ the following assertions hold:*

(i) *For every generalized eigenvalue $E \in I$ of the operator H_ω , one has*

$$\lim_{n \rightarrow \infty} \frac{1}{n} \log \|M_n^E(\omega)\| = L(E). \quad (3.22)$$

(ii) *The spectral subspace $\text{ran}(\chi_I(H_\omega))$ admits a basis of exponentially decaying eigenfunctions.*

(iii) *Given $\delta \in (0, 1)$ and a normalized eigenfunction*

$$f \in \ker(H_\omega - E) \setminus \{0\}, E \in I, \|f\|_{L^2(\mathbb{R}_+)} = 1,$$

there exist $\zeta = \zeta(f) \in \mathbb{N}$, $C_{\omega, \delta} > 0$, $C_\delta > 0$ such that²

$$|f(x^+)| \leq C_{\omega, \delta} e^{C_\delta \log^C(\zeta+1)} e^{-(1-\delta)\tilde{L}(E)|x-\zeta|}, \quad x \geq 0, \quad (3.23)$$

for an absolute constant $C > 0$.

Proof We will show that the statement of the theorem holds with \mathfrak{D} as in Theorem 3.5 and

$$\tilde{\Omega} := \bigcap_{\varepsilon \in (0, \tau) \cap \mathbb{Q}} \Omega_1(\varepsilon) \cap \Omega_2(\varepsilon) \cap \Omega_3(\varepsilon), \quad \tau := \frac{1}{3} \min_{E \in I} L(E),$$

where $\Omega_{1,2,3}(\varepsilon)$ are defined in Theorem 3.10 (iii), (iv) and in Theorem 3.11 respectively. Note that $\tau > 0$ by Theorem 3.5.

Proof of Part (i). Due to (3.17), it is enough to prove that for a given $\omega \in \tilde{\Omega}$ and for a generalized eigenvalue $E = E_\omega \in I$ (which are henceforth fixed) one has

$$\liminf_{n \rightarrow \infty} \frac{1}{n} \log \|M_n^E(\omega)\| \geq L(E). \quad (3.24)$$

Let u be the generalized eigenfunction of H_ω corresponding to E , that is,

$$\begin{aligned} -u'' &= Eu, u(0^+) = 0, u \text{ satisfies (1.3) for all } j > 0, \\ \max \{|u'(t_n^\pm)|, |u(t_n^\pm)|\} &\leq C_u(1+n), n \in \mathbb{Z}_+, \text{ for some } C_u > 0. \end{aligned} \quad (3.25)$$

Our goal is to show that for a given $\varepsilon \in (0, \tau)$ and for all sufficiently large K one has

$$\frac{1}{n} \log \|M_n^E(\omega)\| \geq L(E) - 6\varepsilon, \text{ for all } n \in [K^{11} + K^{10}, \overline{K}]. \quad (3.26)$$

Since these intervals cover a half-line, (3.26) yields (3.24).

² Recall that L and \tilde{L} are related via (3.5).

For a given³ $\zeta \in \mathbb{Z}_+$ let

$$K(N) := \max \left\{ N, n_1, n_2, n_3, \lceil \log^2(\zeta + 1) \rceil \right\}, \quad (3.27)$$

where $N \in \mathbb{N}$ is to be determined,⁴ n_1, n_2 are as in Theorem 3.10 (iii), (iv) correspondingly, and $n_3 = n_3(\omega, \varepsilon)$ is the smallest integer for which

$$\omega \in \bigcap_{i \geq n_3} (\Omega \setminus \mathcal{D}_i(\varepsilon)). \quad (3.28)$$

□

Step 1 There exists $N = N(C_u) > 0$ such that for all $K \geq K(N)$ there exists an integer $m \in [0, \zeta + K^9]$ such that

$$|u(t_m^-)| \leq e^{-2K^2}, \quad |u'(t_m^-)| \leq e^{-2K^2}. \quad (3.29)$$

Proof First we note that (3.16) with $n = K^3$ yields

$$L(E) - \frac{\log \|M_{K^3}^E(T^{\zeta+sK^3}\omega)\|}{K^3} < \varepsilon,$$

or, equivalently,

$$\exp((L(E) - \varepsilon)K^3) < \|M_{K^3}^E(T^{\zeta+sK^3}\omega)\|, \quad (3.30)$$

for some $s \in [0, K^6 - 1] \cap \mathbb{Z}_+$. Focusing on the s -th block we introduce the following notation

$$[\alpha, \beta] := [\zeta + sK^3, \zeta + (s+1)K^3], \quad m := \lfloor \frac{\alpha + \beta}{2} \rfloor.$$

Our argument is based on a representation of u in terms of its boundary values $u(t_\alpha^+)$, $u(t_\beta^-)$ and special solutions ψ_\pm satisfying certain boundary conditions. The choice of the boundary conditions, hence the representation of u , depends on the entry of the matrix

$$S^{-1}(q_\beta)D^{-1}(b_\beta)M_{K^3}^E(T^\alpha\omega) \quad (3.31)$$

that dominates its norm. Specifically, letting m_{ij} denote the ij th entry of (3.31) and assuming that ψ_\pm satisfy $-\psi_\pm'' = E\psi_\pm$, the interior vertex conditions in the interval $[\alpha, \beta]$, and the boundary conditions indicated below, we consider the following four cases.

Case 1. If $\|S^{-1}(q_\beta)D^{-1}(b_\beta)M_{K^3}^E(T^\alpha\omega)\| \leq 4|m_{11}|$ then we let

$$\psi_-(t_\alpha^+) = 1, \quad \psi'_-(t_\alpha^+) = 0, \quad \psi_+(t_\beta^-) = 0, \quad \psi'_+(t_\beta^-) = 1,$$

³ In the sequel ζ will be determined by the center of localization.

⁴ N will depend on u through C_u . In particular, if all generalized eigenfunctions are uniformly bounded, N is u -independent.

and observe that

$$|W(\psi_+, \psi_-)| = |\psi'_+(t_\alpha^+)| = |\psi_-(t_\beta^-)| = |m_{11}| > 0. \quad (3.32)$$

In particular, (3.32) shows that ψ_- and ψ_+ are linearly independent, which shows that we may represent

$$u(t_m^-) = u'(t_\alpha^+) \frac{\psi_+(t_m^-)}{\psi'_+(t_\alpha^+)} + u(t_\beta^-) \frac{\psi_-(t_m^-)}{\psi_-(t_\beta^-)}. \quad (3.33)$$

Case 2. If $\|S^{-1}(q_\beta)D^{-1}(b_\beta)M_{K^3}^E(T^\alpha\omega)\| \leq 4|m_{12}|$ then

$$\begin{aligned} \psi_-(t_\alpha^+) &= 0, \quad \psi'_-(t_\alpha^+) = 1, \quad \psi_+(t_\beta^-) = 0, \quad \psi'_+(t_\beta^-) = 1, \\ u(t_m^-) &= u(t_\alpha^+) \frac{\psi_+(t_m^-)}{\psi_+(t_\alpha^+)} + u(t_\beta^-) \frac{\psi_-(t_m^-)}{\psi_-(t_\beta^-)}, \\ |W(\psi_+, \psi_-)| &= |\psi_+(t_\alpha^+)| = |\psi_-(t_\beta^-)| = |m_{12}| > 0. \end{aligned}$$

Case 3. If $\|S^{-1}(q_\beta)D^{-1}(b_\beta)M_{K^3}^E(T^\alpha\omega)\| \leq 4|m_{21}|$ then

$$\begin{aligned} \psi_-(t_\alpha^+) &= 1, \quad \psi'_-(t_\alpha^+) = 0, \quad \psi_+(t_\beta^-) = 1, \quad \psi'_+(t_\beta^-) = 0, \\ u(t_m^-) &= u'(t_\alpha^+) \frac{\psi_+(t_m^-)}{\psi'_+(t_\alpha^+)} + u'(t_\beta^-) \frac{\psi_-(t_m^-)}{\psi'_-(t_\beta^-)}, \\ |W(\psi_+, \psi_-)| &= |\psi'_+(t_\alpha^+)| = |\psi'_-(t_\beta^-)| = |m_{21}| > 0. \end{aligned}$$

Case 4. If $\|S^{-1}(q_\beta)D^{-1}(b_\beta)M_{K^3}^E(T^\alpha\omega)\| \leq 4|m_{22}|$ then

$$\begin{aligned} \psi_-(t_\alpha^+) &= 0, \quad \psi'_-(t_\alpha^+) = 1, \quad \psi_+(t_\beta^-) = 1, \quad \psi'_+(t_\beta^-) = 0, \\ u(t_m^-) &= u(t_\alpha^+) \frac{\psi_+(t_m^-)}{\psi_+(t_\alpha^+)} + u'(t_\beta^-) \frac{\psi_-(t_m^-)}{\psi'_-(t_\beta^-)}, \\ |W(\psi_+, \psi_-)| &= |\psi_+(t_\alpha^+)| = |\psi'_-(t_\beta^-)| = |m_{22}| > 0. \end{aligned}$$

We proceed with Case 1; the other three cases can be handled similarly. Let us estimate each term in the right-hand side of (3.33). Combining (3.3) and (3.32), we get

$$\begin{aligned} |\psi'_+(t_\alpha^+)| = |\psi_-(t_\beta^-)| = |m_{11}| &\geq \frac{\|S^{-1}(q_\beta)D^{-1}(b_\beta)M_{K^3}^E(T^\alpha\omega)\|}{4} \\ &\geq \frac{\|M_{K^3}^E(T^\alpha\omega)\|}{4\|D(b_\beta)S(q_\beta)\|} \\ &\geq c(b^\pm, q^\pm) \exp((L(E) - \varepsilon)K^3), \end{aligned} \quad (3.34)$$

for some $c(b^\pm, q^\pm) > 0$. By (3.25) we get

$$\max \left\{ |u'(t_\alpha^+)|, |u(t_\beta^-)| \right\} \leq C_u(\beta + 1) \leq C_u(K^9 + e^{\sqrt{K}}).$$

Employing (3.15) with $n = \lfloor \frac{K^3}{2} \rfloor$, $\zeta_0 = \zeta + sK^3$, and choosing N so that $\lfloor \frac{K^3}{2} \rfloor \geq \log^2(\zeta + sK^3)$ we obtain

$$\begin{aligned} |\psi_-(t_m^-)| &\leq \left| \left\langle \begin{bmatrix} 1 \\ 0 \end{bmatrix}, S^{-1}(q_m) D^{-1}(b_m) M_{\lfloor \frac{K^3}{2} \rfloor}^E(T^{\zeta+sK^3}\omega) \begin{bmatrix} 1 \\ 0 \end{bmatrix} \right\rangle \right|, \\ &\leq C(b^\pm, q^\pm) \exp\left(\frac{(L(E) + \varepsilon)K^3}{2}\right), \end{aligned} \quad (3.35)$$

for some $C(b^\pm, q^\pm) > 0$. Similarly for N so large that $\lfloor \frac{K^3}{2} \rfloor \geq \log^2(\zeta + sK^3 + \frac{K^3}{2})$ we obtain

$$|\psi_+(t_m^-)| \leq C(b^\pm, q^\pm) \exp\left(\frac{(L(E) + \varepsilon)K^3}{2}\right), \quad C(b^\pm, q^\pm) > 0. \quad (3.36)$$

Combining (3.33), (3.34)–(3.36) one obtains

$$|u(t_m^-)| \leq 2C_u C(b^\pm, q^\pm) (K^9 + e^{\sqrt{K}}) \exp\left(\frac{-L(E)K^3 + 3\varepsilon K^3}{2}\right) \leq e^{-2K^2},$$

where the last inequality holds whenever $N = N(C_u)$ is large enough and $C(b^\pm, q^\pm) > 0$. Replacing $u(t_m^-)$ by $u'(t_m^-)$, $\psi_\pm(t_m^-)$ by $\psi'_\pm(t_m^-)$ in (3.33), and $[1, 0]^\top$ by $[0, 1]^\top$ in (3.35), (3.36) we obtain

$$|u'(t_m^-)| \leq e^{-2K^2}.$$

□

Step 2 Suppose that $|u(\tau)| = 1$ for some $\tau \in \mathbb{R}_+$, let ζ be the largest integer such that $t_\zeta \leq \tau$, and recall $m \in [0, \zeta + K^9]$ from Step 1 for such ζ . Then

$$\|(H_\omega^m - E)^{-1}\|_{B(L^2(t_0, t_m))} \geq e^{K^2}. \quad (3.37)$$

Proof It suffices to show that

$$|G_{\omega, [0, t_m]}^E(x, y)| \geq Ce^{2K^2}, \quad (x, y) \in J_1 \times (t_m - \delta, t_m), \quad (3.38)$$

for some K -independent interval $J \subset (t_\zeta, t_{\zeta+1})$, K -independent $\delta > 0$, and $C = C(\ell^\pm, I)$. Indeed, denoting the characteristic functions of J , $(t_m - \delta, t_m)$ by χ_1, χ_2 respectively, we get

$$e^{K^2} \leq \frac{|\langle \chi_1, (H_\omega^m - E)^{-1} \chi_2 \rangle_{L^2(t_0, t_m)}|}{\|\chi_1\|_{L^2(t_0, t_m)} \|\chi_2\|_{L^2(t_0, t_m)}} \leq \|(H_\omega^m - E)^{-1}\|_{B(L^2(t_0, t_m))},$$

for N in (3.27) sufficiently large (depending only on $C(\ell^\pm, I)$). To prove (3.38) we notice that

$$\begin{aligned} u(x) &= u(0^+) \frac{\psi_+(x)}{W(\psi_+, \psi_-)} + u'(t_m^-) \frac{\psi_-(x)}{W(\psi_+, \psi_-)} = u'(t_m^-) \frac{\psi_-(x)}{W(\psi_+, \psi_-)} \\ &= u'(t_m^-) G_{\omega, m}^E(x, t_m), \quad x \in (t_\zeta, t_{\zeta+1}), \end{aligned}$$

(this is similar to Case 4 in Step 1 above). By right-continuity of u and $|u(\tau)| = 1$ we have

$$1/2 \leq |u(x)|, \quad x \in J \subset (t_\zeta, t_{\zeta+1}),$$

for some K -independent interval J . Employing (3.29) one infers

$$1 \lesssim |u(x)| = |u'(t_m)| \left| \frac{\psi_-(x)}{W(\psi_+, \psi_-)} \right| \leq e^{-2K^2} \left| \frac{\psi_-(x)}{W(\psi_+, \psi_-)} \right|,$$

for all $x \in J$. That is,

$$e^{2K^2} \lesssim \left| \frac{\psi_-(x)}{W(\psi_+, \psi_-)} \right|, \quad x \in J.$$

Furthermore, noticing that

$$\psi_+(y) = \cos(\sqrt{E}(y - t_m)) \geq 1/2 \text{ for all } y \in (t_m - \delta, t_m],$$

for some K -independent sufficiently small constant $\delta > 0$, and using Proposition 2.4 we arrive at

$$|G_{\omega, [0, t_m]}^E(x, y)| = \left| \frac{\psi_-(x) \psi_+(y)}{W(\psi_+, \psi_-)} \right| \geq \left| \frac{\psi_-(x)}{2W(\psi_+, \psi_-)} \right| \gtrsim e^{2K^2},$$

for all $(x, y) \in J \times (t_m - \delta, t_m]$. Thus (3.38) holds as required. \square

Step 3 Let ζ be as in Step 2. Then there exists $N = N(C_u)$ such that for all $K \geq K(N)$ and all $n \in [K^{11} + K^{10}, \bar{K}]$ one has

$$\frac{1}{n} \log \|M_n^E(T^\zeta \omega)\| \geq L(E) - 5\varepsilon. \quad (3.39)$$

Proof Combining (3.28), (3.37) and Theorem 3.11 one infers

$$\frac{1}{mK} \log \|M_{mK}^E(T^{\zeta+r} \omega)\| \geq L(E) - \varepsilon, \quad r \in [K^{10}, \bar{K}], \quad m \in \{1, 2\}. \quad (3.40)$$

We will use (3.40) to apply the Avalanche principle, see Lemma 3.12. Concretely, choose $q \in \mathbb{Z}_+$ with $K^{10} \leq q \leq K^{-1}\bar{K} - K^9$, define

$$A^{(j)} := M_K^E(T^{\zeta+K^{10}+(j-1)K}\omega), \quad 1 \leq j \leq q.$$

With $\lambda := \exp(K(L(E) - \varepsilon))$, (3.40) gives

$$\|A^{(j)}\| \geq \lambda \geq q$$

for all j , where the second inequality holds as long as N , cf. (3.27), is sufficiently large. Since $K \geq \tilde{n}_1$ and $K \geq \log^2(|\zeta| + |\bar{K}| + 1)$ (enlarge N if necessary), we may use (3.15) to obtain

$$\|A^{(j)}\| \leq \exp(K(L(E) + \varepsilon)), \quad 1 \leq j \leq q.$$

Thus, implies

$$\begin{aligned} & \left| \log \|A^{(j+1)}\| + \log \|A^{(j)}\| - \log \|A^{(j+1)}A^{(j)}\| \right| \\ & \log \|A^{(j+1)}\| + \log \|A^{(j)}\| - \log \|A^{(j+1)}A^{(j)}\| \\ & < 2K(L(E) + \varepsilon) - 2K(L(E) - \varepsilon) \\ & = 4K\varepsilon \\ & \leq \frac{1}{2} \log \lambda, \end{aligned}$$

where the final inequality needs ε to be sufficiently small; we note that this smallness condition depends only on $\tilde{\mu}$. Thus, taking $\hat{N} = qK$ and $r_0 = K^{10}$, we have $\hat{N} \in [K^{11}, \bar{K} - K^{10}]$ and the Avalanche Principle (Lemma 3.12) yields

$$\begin{aligned} \log \|M_{\hat{N}}(T^{\zeta+r_0}\omega)\| &= \log \|A^{(q)} \cdots A^{(1)}\| \\ &\geq \sum_{j=1}^{q-1} \log \|A^{(j+1)}A^{(j)}\| - \sum_{j=2}^{q-1} \log \|A^{(j)}\| - C\frac{q}{\lambda} \\ &\geq (q-1)2K(L(E) - \varepsilon) - (q-2)K(L(E) + \varepsilon) - C \\ &\geq \hat{N}(L(E) - 4\varepsilon) \end{aligned}$$

again, by choosing N large.

Putting this together, we can control $\|M_n^E(T^\zeta\omega)\|$ for general $K^{11} + K^{10} \leq n \leq \bar{K}$ by interpolation. In particular, writing $n = qK + p$ with $0 \leq p < K$ and $q \geq K^{10} + K^9$, we have

$$\begin{aligned} \|M_n^E(T^\zeta\omega)\| &\geq \frac{\|M_{n-K^{10}}(T^{\zeta+K^{10}}\omega)\|}{\|M_{K^{10}}(T^\zeta\omega)\|} \\ &\geq \rho^{-K^{10}-p} \|M_{qK-K^{10}}(T^{\zeta+K^{10}}\omega)\| \end{aligned}$$

$$\begin{aligned} &\geq \rho^{-K^{10}-p} e^{(qK-K^{10})(L(E)-4\varepsilon)} \\ &\geq e^{n(L(E)-5\varepsilon)}, \end{aligned}$$

as long as N is sufficiently large [recall ρ from (3.12)]. \square

Picking $\tau \in (t_0, t_1)$ such that $u(\tau) \neq 0$, replacing u by $\frac{u}{u(\tau)}$, and using (3.39) one infers (3.26) which in turn yields (3.24) and (3.22).

Proof of Part (ii). By Part (i) and Ruelle's deterministic version of Oseledec' Theorem [54,56], every generalized eigenvalue is, in fact, an eigenvalue corresponding to an exponentially decaying eigenfunction. Furthermore, since the spectral measure of $H_\omega \chi_I(H_\omega)$ is supported by the generalized eigenvalues belonging to I , cf. [44, Theorem C.17], one infers that $\text{ran}(\chi_I(H_\omega))$ admits a basis of exponential decaying eigenfunctions.

Proof of Part (iii). First, we notice that

$$\begin{aligned} &\max \{ \|f\|_{L^\infty(t_j, t_{j+1})}, \|f'\|_{L^\infty(t_j, t_{j+1})} \} \\ &\leq c(\ell^-, \ell^+) (\|f\|_{L^2(t_j, t_{j+1})} + \|f''\|_{L^2(t_j, t_{j+1})}) \\ &\leq c(\ell^-, \ell^+, I) \|f\|_{L^2(\mathbb{R}_+)} = c(\ell^-, \ell^+, I), \end{aligned} \quad (3.41)$$

and

$$\begin{aligned} \|f'\|_{L^\infty(t_j, t_{j+1})} &\leq C(\ell^-, \ell^+) (\|f\|_{L^2(t_j, t_{j+1})} + \|f''\|_{L^2(t_j, t_{j+1})}) \\ &\leq C(\ell^-, \ell^+, I) \|f\|_{L^2(t_j, t_{j+1})} \leq C(\ell^-, \ell^+, I) \|f\|_{L^\infty(t_j, t_{j+1})}, \end{aligned} \quad (3.42)$$

for some $C(\ell^-, \ell^+, I) > 0$, and all $j \in \mathbb{Z}_+$ cf., e.g. [21, Corollary 4.2.10], [46, IV.1.2]. In addition we remark that f attains its maximum since

$$\left\{ \begin{bmatrix} f(t_j^+) \\ f'(t_j^+) \end{bmatrix} \right\}_{j=0}^\infty \in \ell^2(\mathbb{Z}_+, \mathbb{C}^2) \text{ and thus } \lim_{t \rightarrow \infty} (|f(t)| + |f'(t)|) = 0.$$

Therefore, we may repeat the arguments of the proof of Part (i) with

$$\begin{aligned} u &= \frac{f}{\|f\|_{L^\infty(\mathbb{R}_+)}} \quad C_u = \max \{1, C(\ell^-, \ell^+, I)\} \text{ in Step 1,} \\ \tau &= \operatorname{argmax} |f| \text{ (i.e. } \tau \text{ is chosen so that } |f(\tau)| = \|f\|_\infty \text{) in Step 2,} \end{aligned}$$

where we pick any value of argmax if there is more than one extremum. Then for a given $\varepsilon \in (0, \tau)$ there exists $N = N(\varepsilon, \omega)$ (which does not depend on f) such that for all $K \geq K(N, \log^2(\zeta + 1))$ and all $n \in [K^{11} + K^{10}, \bar{K}]$ one has

$$\frac{1}{n} \log \|M_n^E(T^\zeta \omega)\| \geq L(E) - 6\varepsilon.$$

Utilizing this with sufficiently small ε (depending on δ only) and letting

$$\varkappa := c(b^\pm, \ell^\pm I) \max \{1, C(\ell^-, \ell^+, I)\},$$

see (3.41), (3.42), we will show that

$$|f(t_{\zeta+n}^+)| \leq \varkappa e^{-(1-\delta)L(E)n}, \text{ for all } n \in \left[\frac{p}{4}, \frac{p-1}{2}\right],$$

for all $p \in [K^{11} + K^{10}, \overline{K}]$, $K \geq K(N)$. As in Step 1 our subsequent argument relies on a representation of f considered on the interval $[t_\zeta, t_{\zeta+p}]$ in terms of its boundary values. Our choice of the representation, as before, depends on the entry of

$$S^{-1}(q_{\zeta+p})D^{-1}(b_{\zeta+p})M_p^E(T^\zeta \omega)$$

that dominates its norm. We will provide the argument assuming that the maximizing entry is 11 and note that the other three cases can be treated almost identically.

One has

$$\frac{f(t_{\zeta+n}^+)}{M_f} = \frac{f'(t_\zeta^+)\psi_+(t_{\zeta+n}^+)}{M_f \psi'_+(t_\zeta^+)} + \frac{f(t_{\zeta+p}^-)\psi_-(t_{\zeta+n}^+)}{M_f \psi_-(t_{\zeta+p}^-)}, \quad (3.43)$$

where $M_f := \|f\|_{L^\infty(\mathbb{R}_+)}$, $-\psi'_\pm = E\psi_\pm$, ψ_\pm satisfies the interior vertex conditions in the interval $[t_\zeta, t_{\zeta+p}]$, and

$$\psi_-(t_\zeta^+) = 1, \quad \psi'_-(t_\zeta^+) = 0, \quad \psi_+(t_{\zeta+p}^-) = 0, \quad \psi'_+(t_{\zeta+p}^-) = 1,$$

and

$$\begin{aligned} |W(\psi_+, \psi_-)| &= |\psi'_+(t_\zeta^+)| = |\psi_-(t_{\zeta+p}^-)| \\ &\geq \frac{\|S^{-1}(q_{\zeta+p})D^{-1}(b_{\zeta+p})M_p^E(T^\zeta \omega)\|}{4} \\ &\geq \frac{\|M_p^E(T^\zeta \omega)\|}{4\|D(b_{\zeta+p})S(q_{\zeta+p})\|} \\ &\geq c(b^\pm, \ell^\pm) \exp((L(E) - 6\varepsilon)p), \end{aligned} \quad (3.44)$$

for some $c(b^\pm, \ell^\pm) > 0$. In order to estimate $\psi_-(t_{\zeta+n}^+)$, we rewrite it in terms of the transfer matrices and use (3.15) as follows

$$|\psi_-(t_{\zeta+n}^+)| = \left| \left\langle \begin{bmatrix} 1 \\ 0 \end{bmatrix}, M_n^E(T^\zeta \omega) \begin{bmatrix} 1 \\ 0 \end{bmatrix} \right\rangle \right| \leq \exp((L(E) + \varepsilon)n).$$

Similarly one can estimate $\psi_+(t_{\zeta+n}^+)$. Combining this and (3.41), (3.42), (3.43), (3.44) we get

$$\begin{aligned} |f(t_{\zeta+n}^+)| &\leq \varkappa \exp((L(E) + \varepsilon)n - (L(E) - 6\varepsilon)p) \\ &\quad + \varkappa \exp((L(E) + \varepsilon)(p - n) - (L(E) - 6\varepsilon)p) \\ &\leq \varkappa \exp(-(p - n)L(E) + (n + 6p)\varepsilon) \\ &\quad + \varkappa \exp(-nL(E) + (7p - n)\varepsilon) \\ &\leq 2\varkappa \exp(-nL(E) + 8p\varepsilon) \leq 2\varkappa \exp(-nL(E) + 32n\varepsilon) \\ &\leq 2\varkappa e^{-(1-\delta)nL(E)}, \end{aligned}$$

to facilitate the last inequality we pick $\varepsilon = \varepsilon(\delta) > 0$ sufficiently small (depending only on δ). Thus

$$|f(t_{\zeta+n}^+)| \leq 2\varkappa e^{-(1-\delta)L(E)n}, \quad (3.45)$$

for all $n \in [\frac{K^{11}+K^{10}}{4}, \frac{\bar{K}-1}{2}]$ and $K \geq K(N)$. Since these intervals cover the half-line $[\frac{K^{11}}{2}, \infty)$ for sufficiently large N , the inequality in (3.45) holds for all

$$n \geq \frac{K^{11}}{2} = \frac{1}{2} \max \left\{ N(\omega, \varepsilon), \log^2(\zeta + 1) \right\}^{11}.$$

Furthermore, estimating $f(t_{\zeta+n}^+)$ for

$$n \in \left[0, 2^{-1} \max \left\{ N(\omega, \varepsilon), \log^2(\zeta + 1) \right\}^{11} \right]$$

trivially and changing variables $k = \zeta + n$, we get

$$\begin{aligned} |f(t_k^+)| &\leq 2\varkappa e^{(1-\delta)L(E) \max \left\{ N(\omega, \varepsilon), \log^2(\zeta + 1) \right\}^{11}} e^{-(1-\delta)L(E)(k-\zeta)} \\ &\leq C_{\omega, \delta} e^{C_{\delta} \log^{22}(\zeta + 1)} e^{-(1-\delta)L(E)|k-\zeta|}, \quad k \geq \zeta. \end{aligned} \quad (3.46)$$

A similar estimate can be obtained for $k \in [0, \zeta]$: In this case, the Lyapunov behavior (3.45) is observed only for sufficiently large ζ , in which case (3.45) holds for $k \in [0, \zeta - \frac{K^{11}}{2}]$ (for small ζ , use the trivial bound).

In order to show a version of (3.46) with f replaced by f' , we employ

$$\frac{f'(t_{\zeta+n}^+)}{M_f} = \frac{f'(t_{\zeta}^+)\psi'_+(t_{\zeta+n}^+)}{M_f \psi'_+(t_{\zeta}^+)} + \frac{f(t_{\zeta+p}^-)\psi'_-(t_{\zeta+n}^+)}{M_f \psi_-(t_{\zeta+p}^-)},$$

and repeat (3.44)–(3.46). Finally, keeping in mind Remark 3.3 and interpolating between the discrete vertices, we infer (3.23). \square

Having established existence of a basis of semi-uniformly localized eigenfunctions (SULE) we turn to dynamical localization. Our argument stems from the proof of [38, Theorem 2.1].

Proof of Theorem 1.1 Our first goal is to derive an upper bound for the number of centers of localization⁵ located in a large interval $[0, L]$. Let $\{\varphi_n\}_{n=1}^\infty$ be an $L^2(\mathbb{R}_+)$ -orthonormal basis of exponentially decaying eigenfunctions of the spectral subspace $\text{ran}(\chi_I(H_\omega))$; the corresponding eigenvalues are denoted by $E_n \in I$, $n \geq 1$. Then by (3.23) with

$$\delta := 1/2, \nu := \min_{E \in I} \tilde{L}(E, 1) > 0,$$

we have

$$|\varphi_n(x)| \leq C_\omega e^{C \log^{22}(\zeta_n+1)} e^{-\frac{\nu|x-\zeta_n|}{2}}, \quad x \geq 0. \quad (3.47)$$

We claim that

$$\mathcal{N}(L) := \#\{n : \zeta_n \leq L\} \leq C(\omega, I)L, \quad L \geq L_0, \quad (3.48)$$

for sufficiently large $L_0 > 0$. For $L > 0$ let $\chi_{3L} \in \mathcal{B}(L^2(\mathbb{R}_+))$ denote the operator of multiplication by the characteristic function of $[0, 3L]$, let $R(H_\omega)$ denote the resolvent of H_ω at $\lambda = \min \sigma(H_\omega) - 1$ and note that $\|R^2(H_\omega)\|_{\mathcal{B}(L^2(\mathbb{R}_+))} \leq 1$. Next we show

$$\mathcal{N}(L) \leq C(\omega, I) \text{tr}(\chi_{3L} R^2(H_\omega) \chi_{3L}), \quad (3.49)$$

for sufficiently large L and some $C(\omega, I)$. To that end, notice that

$$\begin{aligned} \frac{1}{(E_n - \lambda)^2} &= \langle \varphi_n, R^2(H_\omega) \varphi_n \rangle_{L^2(\mathbb{R}_+)} \\ &= \langle \varphi_n, \chi_{3L} R^2(H_\omega) \chi_{3L} \varphi_n \rangle_{L^2(\mathbb{R}_+)} \\ &\quad + \langle \varphi_n, \chi_{3L} R^2(H_\omega) (1 - \chi_{3L}) \varphi_n \rangle_{L^2(\mathbb{R}_+)} \\ &\quad + \langle \varphi_n, (1 - \chi_{3L}) R^2(H_\omega) \chi_{3L} \varphi_n \rangle_{L^2(\mathbb{R}_+)} \\ &\quad + \langle \varphi_n, (1 - \chi_{3L}) R^2(H_\omega) (1 - \chi_{3L}) \varphi_n \rangle_{L^2(\mathbb{R}_+)}. \end{aligned} \quad (3.50)$$

$$(3.51)$$

Assuming that $\zeta_n \leq L$, $E_n \in I$, and $C \log^{22}(L+1) < \frac{\nu L}{4}$ and using (3.47) we obtain

$$|\varphi_n(x)| \leq C_\omega e^{\frac{\nu L}{4}} e^{-\frac{\nu|x-\zeta_n|}{2}}, \quad x \geq 0,$$

and

$$\begin{aligned} \langle \varphi_n, \chi_{3L} R^2(H_\omega) (1 - \chi_{3L}) \varphi_n \rangle_{L^2(\mathbb{R}_+)} &\leq \|(1 - \chi_{3L}) \varphi_n\|_{L^2(\mathbb{R}_+)} \\ &\leq C_\omega e^{\frac{\nu L}{4}} \left(\int_{3L}^\infty e^{-\nu|x-\zeta_n|} dx \right)^{1/2} \\ &\leq C_\omega e^{\frac{\nu L}{4}} e^{\frac{\nu \zeta_n}{2}} e^{-\frac{3\nu L}{2}} \nu^{-\frac{1}{2}} \leq C_\omega e^{-\frac{3\nu L}{4}} \nu^{-\frac{1}{2}} \underset{L \rightarrow \infty}{=} o(1). \end{aligned}$$

Similar estimates hold for (3.50) and (3.51). Therefore we have

$$\text{tr}(\chi_{3L} R^2(H_\omega) \chi_{3L})$$

⁵ ζ from (3.23) is called the center of localization of f .

$$\begin{aligned}
&\geq \sum_{n: \zeta_n \leq L} \langle \varphi_n, \chi_{3L} R^2(H_\omega) \chi_{3L} \varphi_n \rangle_{L^2(\mathbb{R}_+)} \\
&\geq \sum_{n: \zeta_n \leq L} \left(\frac{1}{(E_n - \lambda)^2} - 3C_\omega e^{-\frac{3vL}{4}} v^{-\frac{1}{2}} \right) \\
&\geq C(I, \omega) \#\{n : \zeta_n \leq L\},
\end{aligned}$$

for some $C(I, \omega) > 0$.

Next we estimate the right-hand side of (3.49). Let us recall that $AB \in \mathcal{B}_2(L^2(\mathbb{R}_+))$ (the space of Hilbert–Schmidt operators on $L^2(\mathbb{R}_+)$) and

$$\|AB\|_{\mathcal{B}_2(L^2(\mathbb{R}_+))} \lesssim \|A\|_{\mathcal{B}(L^\infty(\mathbb{R}_+), L^2(\mathbb{R}_+))} \|B\|_{\mathcal{B}(L^2(\mathbb{R}_+), L^\infty(\mathbb{R}_+))},$$

whenever $A \in \mathcal{B}(L^\infty(\mathbb{R}_+), L^2(\mathbb{R}_+))$, $B \in \mathcal{B}(L^2(\mathbb{R}_+), L^\infty(\mathbb{R}_+))$. A discussion of this fact together with related references can be found, for instance, in [60, Section 4.1.11] and [61, pp. 418–419]. This result is applicable in our case due to [44, Lemma C.12] which asserts that $R(H_\omega)$ maps (boundedly) $L^2(\mathbb{R}_+)$ into $L^\infty(\mathbb{R}_+)$. Combining these facts we infer

$$\begin{aligned}
\mathrm{tr}(\chi_{3L} R^2(H_\omega) \chi_{3L}) &= \|\chi_{3L} R(H_\omega)\|_{\mathcal{B}_2(L^2(\mathbb{R}_+))}^2 \\
&\leq \left(\sqrt{3L} \|R(H_\omega)\|_{\mathcal{B}(L^2(\mathbb{R}_+), L^\infty(\mathbb{R}_+))} \right)^2 \leq C(\omega)L,
\end{aligned} \tag{3.52}$$

for some $C(\omega) > 0$. Then (3.49) and (3.52) yield (3.48).

Next, we turn to (1.4). For brevity, denote $\gamma := 22 + \varepsilon$ and let $\kappa > 0$ be such that

$$|\log^\gamma(x + \kappa) - \log^\gamma(y + \kappa)| \leq \frac{v|x - y|}{4}, \quad x, y > 0. \tag{3.53}$$

Then we have

$$\begin{aligned}
&\left\| |X|^p \chi_I(H_\omega) e^{-itH_\omega} \psi \right\|_{L^2(\mathbb{R}_+)} \\
&\leq \sum_{n=1}^{\infty} |\langle \varphi_n, \psi \rangle_{L^2(\mathbb{R}_+)}| \| |X|^p \varphi_n \|_{L^2(\mathbb{R}_+)} \\
&\leq \sum_{n=1}^{\infty} C_{\omega, I} e^{2C \log^{22}(\zeta_n + 1)} \int_{\mathbb{R}_+} |\psi(x)| e^{-\frac{v|x - \zeta_n|}{2}} dx \left(\int_{\mathbb{R}_+} x^{2p} e^{-v|x - \zeta_n|} dx \right)^{1/2} \\
&\leq \sum_{n=1}^{\infty} C_{\omega, I, p, \psi} e^{2C \log^{22}(\zeta_n + 1)} \zeta_n^p \int_{\mathbb{R}_+} e^{-\log^\gamma(x + \kappa)} e^{-\frac{v|x - \zeta_n|}{2}} dx \\
&\leq \sum_{n=1}^{\infty} C_{\omega, I, p, \psi} e^{2C \log^{22}(\zeta_n + 1) + p \log(\zeta_n + 1) - \log^\gamma(\zeta_n + \kappa)} \\
&\quad \times \int_{\mathbb{R}_+} e^{-\log^\gamma(x + \kappa) + \log^\gamma(\zeta_n + \kappa) - \frac{v|x - \zeta_n|}{2}} dx
\end{aligned} \tag{3.54}$$

$$\begin{aligned}
 & \stackrel{(3.53)}{\leq} \sum_{n=1}^{\infty} C_{\omega, I, p, \psi} e^{2C \log^{22}(\zeta_n+1) + p \log(\zeta_n+1) - \log^{\gamma}(\zeta_n+\kappa)} \\
 & \quad \times \int_{\mathbb{R}_+} e^{-\frac{\nu|x-\zeta_n|}{4}} dx \\
 & \leq \tilde{C}_{\omega, I, p, \psi} \sum_{n=1}^{\infty} e^{2C \log^{22}(\zeta_n+1) + p \log(\zeta_n+1) - \log^{\gamma}(\zeta_n+\kappa)} \\
 & \leq \tilde{C}_{\omega, I, p, \psi} \sum_{L=0}^{\infty} \sum_{n: \zeta_n=L} e^{2C \log^{22}(\zeta_n+1) + p \log(\zeta_n+1) - \log^{\gamma}(\zeta_n+\kappa)} \\
 & \leq \tilde{C}_{\omega, I, p, \psi} \sum_{L=0}^{\infty} \mathcal{N}(L) e^{2C \log^{22}(L+1) + p \log(L+1) - \log^{22+\varepsilon}(L+\kappa)} < \infty, \quad (3.55)
 \end{aligned}$$

where we used (3.48) in the last inequality.

4 Random metric trees

4.1 The almost-sure spectrum for continuum models

Our first objective is to show that almost surely the spectrum of \mathbb{H}_{ω} is given by a deterministic set Σ .

Theorem 4.1 *There exists a full μ -measure set $\hat{\Omega} \subset \Omega$ such that*

$$\sigma(\mathbb{H}_{\omega}) = \Sigma := \overline{\bigcup_{(b, \ell, q) \text{ periodic}} \sigma(\mathbb{H}(b, \ell, q))}, \quad \omega \in \hat{\Omega}.$$

Proof Since

$$\sigma(\mathbb{H}(b, \ell, q)) = \overline{\bigcup_{k \in \mathbb{Z}_+} \sigma(H(T^k b, T^k \ell, T^k q))},$$

one has

$$\sigma(\mathbb{H}_{\omega}) = \overline{\bigcup_{k \in \mathbb{Z}_+} \sigma(H_{T^k \omega})}; \quad \Sigma = \overline{\bigcup_{(b, \ell, q) \text{ periodic}} \sigma(H(b, \ell, q))}.$$

First, we will first show that

$$\sigma(H_{\omega}) \subset \Sigma, \quad \text{for all } \omega \in \Omega,$$

and therefore $\sigma(\mathbb{H}_{\omega}) \subset \Sigma$. Let us fix $\omega \in \Omega$. Seeking a contradiction, we pick $E \in \sigma(H_{\omega}) \setminus \Sigma$. Then there exist

$$\{f_k\}_{k=1}^{\infty} \subset \text{dom}(H_{\omega}) \text{ and } \{m_k\}_{k=1}^{\infty} \subset \mathbb{N},$$

such that

$$\|f_k\|_{L^2(t_0, \infty)} = 1, \quad \text{supp}(f_k) \subset [t_0, t_{m_k}],$$

$$\sup_{\substack{g \in \text{dom}(\mathfrak{h}_\omega) \\ \|g\|_{\widehat{H}^1(t_0, \infty)} \leq 1}} (\mathfrak{h}_\omega - E)[f_k, g] \rightarrow 0, \quad k \rightarrow \infty, \quad (4.1)$$

where $\mathfrak{h}_\omega = \mathfrak{h}(b_\omega, \ell_\omega, q_\omega)$, cf. (2.17)–(2.19) [we recall that \widehat{H}^1 -norm is equivalent to the form norm, see (2.21)]. Let $(b^k, \ell^k, q^k) \in \Omega$ denote the m_k -periodic sequence whose first m_k elements are given by $\omega_1, \dots, \omega_{m_k}$. Then since $E \notin \Sigma$ one has

$$C := \sup_{k \in \mathbb{N}} \|F_k\|_{\widehat{H}^1(t_0, \infty)} < \infty, \quad F_k := (H(b^k, \ell^k, q^k) - E)^{-1} f_k,$$

where the first inequality follows from the fact that $F_k'' = -E F_k - f_k$ and Sobolev inequalities. Suitable truncations of F_k belong to $\text{dom}(\mathfrak{h}_\omega)$. Indeed, for $k \in \mathbb{N}$, let $\varphi_k \in C_0^\infty[t_0, \infty)$ be such that $\text{supp } \varphi_k \subset [t_0, t_{m_k+1}]$, $0 \leq \varphi_k(x) \leq 1$, $x \geq t_0$, and

$$\varphi_k(x) = \begin{cases} 1, & x \in [t_0, t_{m_k}], \\ 0, & x \in [t_{m_k+1}, \infty). \end{cases}$$

Then for all $k \in \mathbb{N}$ one has

$$\begin{aligned} (\varphi_k F_k) &\in \text{dom}(\mathfrak{h}_\omega) \\ \|\varphi_k F_k\|_{\widehat{H}^1(\mathbb{R}_+)} &\leq \max \left\{ 1, \|\varphi_k\|_{H^1(t_{m_k}, t_{m_k+1})} \right\} \|F_k\|_{\widehat{H}^1(\mathbb{R}_+)} \lesssim 1, \end{aligned} \quad (4.2)$$

where we used $\|\varphi_k F_k\|_{H^1(t_{m_k}, t_{m_k+1})} \lesssim \|\varphi_k\|_{H^1(t_{m_k}, t_{m_k+1})} \|F_k\|_{H^1(t_{m_k}, t_{m_k+1})}$, see [41, Theorem 4.14]. Moreover, one has

$$\begin{aligned} (\mathfrak{h}_\omega - E)[f_k, \varphi_k F_k] &= \langle \varphi_k F_k, -f_k'' - E f_k \rangle_{L^2(\mathbb{R}_+)} \\ &= \langle F_k, -f_k'' - E f_k \rangle_{L^2(\mathbb{R}_+)} \\ &= \left\langle (H(b^k, \ell^k, q^k) - E)^{-1} f_k, (H(b^k, \ell^k, q^k) - E) f_k \right\rangle_{L^2(\mathbb{R}_+)} = 1. \end{aligned} \quad (4.3)$$

Combining (4.1), (4.2) and (4.3) we obtain a contradiction.

Next we show that exists a full μ -measure set $\widehat{\Omega} \subset \Omega$ such that

$$\Sigma \subset \sigma(\mathbb{H}_\omega), \quad \omega \in \widehat{\Omega}. \quad (4.4)$$

First of all, we note that $E \in \sigma(\mathbb{H}_\omega)$ whenever there exist two sequences of natural numbers

$$\{r_k\}_{k=1}^\infty \subset \mathbb{N}, \quad \{m_k\}_{k=1}^\infty \subset \mathbb{N}, \quad (4.5)$$

and a sequence of functions $\{f_k\}_{k=1}^\infty$ such that $f_k \in \text{dom}(\mathfrak{h}_{T^{r_k}\omega})$ satisfying

$$\liminf_{k \rightarrow \infty} \|f_k\|_{L^2(t_\omega(r_k), \infty)} > 0, \quad \text{supp}(f_k) \subset [t_\omega(r_k), t_\omega(r_k + m_k)], \quad k \in \mathbb{N}, \quad (4.6)$$

and

$$\sup (\mathfrak{h}_{T^{r_k}\omega} - E)[f_k, g] \rightarrow 0, \quad k \rightarrow \infty, \quad (4.7)$$

where the supremum is taken over the set

$$\{g \in \text{dom}(\mathfrak{h}_{T^{r_k}\omega}) : \|g\|_{\widehat{H}^1(t_\omega(r_k), \infty)} \leq 1\}.$$

This is due to orthogonal decomposition (2.11) and the standard Weyl criterion for \mathbb{H}_ω . Secondly, there exists $\widehat{\Omega} \subset \Omega$, $\mu(\widehat{\Omega}) = 1$ such that for arbitrary

$$\omega \in \widehat{\Omega}, (b, \ell, q) \in \text{supp}(\mu), \{m_k\}_{k=1}^\infty \subset \mathbb{N},$$

there exists a sequence $\{r_k\}_{k=1}^\infty$ such that for all $k \in \mathbb{N}$ one has

$$b_\omega(r_k + i) = b_i \text{ for all } i \in \{1, \dots, m_k\},$$

$$\max_{1 \leq i \leq m_k} |\ell_\omega(i + r_k) - \ell_i| \leq \frac{\sqrt{\ell^-}}{k}, \quad (4.8)$$

$$\max_{1 \leq i \leq m_k} |q_\omega(i + r_k) - q_i| \leq \frac{1}{k}, \quad (4.9)$$

see, for example, [47, Proposition 3.8]. We claim that (4.4) holds with this choice of $\widehat{\Omega}$. Indeed, pick any periodic sequence (b, ℓ, q) and $E \in \sigma(H(b, \ell, q))$. Then by Proposition 2.3 there exist

$$\{\varphi_k\}_{k=1}^\infty \subset \text{dom}(\mathfrak{h}(b, \ell, q)), \{m_k\}_{k=1}^\infty \subset \mathbb{N},$$

such that

$$\sup_{k \in \mathbb{N}} \|\varphi_k\|_{\widehat{H}^1(t_0, \infty)} < \infty, \|\varphi_k\|_{L^2(t_0, \infty)} = 1, \text{supp}(\varphi_k) \subset [t_0, t_{m_k}],$$

$$\sup_{\substack{g \in \text{dom}(\mathfrak{h}(b, \ell, q)) \\ \|g\|_{\widehat{H}^1(t_0, \infty)} \leq 1}} (\mathfrak{h}(b, \ell, q) - E)[\varphi_k, g] \rightarrow 0, k \rightarrow \infty. \quad (4.10)$$

In order to produce a singular sequence for \mathbb{H}_ω we will rescale φ_k from $[t_{i-1}, t_i]$ to $[t_\omega(r_k + i - 1), t_\omega(r_k + i)]$. That is, for every $i, k \in \mathbb{N}$ we let

$$f_k(y) := \varphi_k(s_{i,k}^{-1}(y)), y \in [t_\omega(r_k + i - 1), t_\omega(r_k + i)],$$

where

$$s_{i,k}(x) := \frac{t_\omega(r_k + i) - t_\omega(r_k + i - 1)}{\ell_i}(x - t_{i-1}) + t_\omega(r_k + i - 1),$$

for $x \in [t_{i-1}, t_i]$. Then changing variables one obtains

$$\langle f'_k, g' \rangle_{L^2(t_\omega(r_k+i-1), t_\omega(r_k+i))} = \frac{\ell_i}{\ell_\omega(r_k + i)} \langle \varphi'_k, (g \circ s_{i,k})' \rangle_{L^2(t_{i-1}, t_i)}, \quad (4.11)$$

$$\langle f_k, g \rangle_{L^2(t_\omega(r_k+i-1), t_\omega(r_k+i))} = \frac{\ell_\omega(r_k+i)}{\ell_i} \langle \varphi_k, g \circ s_{i,k} \rangle_{L^2(t_{i-1}, t_i)}, \quad (4.12)$$

where $g \in \widehat{H}^1(t_\omega(r_k), \infty)$. Let us denote

$$\widetilde{g}_k(x) := (g \circ s_{i,k})(x), \quad x \in [t_{i-1}, t_i], \quad i \in \mathbb{N}, \quad k \in \mathbb{N}.$$

Then using (4.11), (4.12) with f_k replaced by g we note that there exists a constant $C > 0$ which does not depend on k such that

$$\|\widetilde{g}_k\|_{\widehat{H}^1(t_0, t_{m_k})} \leq C \text{ if } \|g\|_{\widehat{H}^1(t_\omega(r_k), \infty)} \leq 1, \quad k \in \mathbb{N}. \quad (4.13)$$

We claim that $\{f_k\}_{k=1}^\infty$ is a singular sequence satisfying (4.5)–(4.7). First, we know that $f_k \in \text{dom}(\mathfrak{h}_{T^{r_k}\omega})$ holds since the vertex conditions displayed in (2.18) are scale-invariant. Next, the conditions in (4.6) hold due to (4.10) and (4.12) (with $g = f_k$). In order to check (4.7), let us fix $k \in \mathbb{N}$ and g with $\|g\|_{\widehat{H}^1(t_\omega(r_k), \infty)} \leq 1$. Then one has

$$\begin{aligned} & |(\mathfrak{h}_{T^{r_k}\omega} - E)[f_k, g] - (\mathfrak{h}(b, \ell, q) - E)[\varphi_k, \widetilde{g}_k]| \\ & \leq \left| \sum_{i=1}^{m_k} \left(\frac{\ell_i}{\ell_\omega(r_k+i)} - 1 \right) \langle \varphi'_k, (g \circ s_{i,k})' \rangle_{L^2(t_{i-1}, t_i)} \right. \\ & \quad \left. - E \left(\frac{\ell_\omega(r_k+i)}{\ell_i} - 1 \right) \langle \varphi_k, g \circ s_{i,k} \rangle_{L^2(t_{i-1}, t_i)} \right| \\ & \quad + \left| \sum_{i=1}^{m_k} (q_i - q_\omega(r_k+i)) \overline{\varphi_k(t_i^-)} (g \circ s_{i,k})(t_i^-) \right| \\ & \lesssim \frac{\|\varphi_k\|_{\widehat{H}^1(\mathbb{R}_+)} \|\widetilde{g}_k\|_{\widehat{H}^1(\mathbb{R}_+)}}{k} \\ & \lesssim \frac{1}{k} \rightarrow 0, \quad k \rightarrow \infty. \end{aligned}$$

In the first inequality we employed (4.11) and (4.12); in the second one we used the Cauchy–Schwarz inequality, the fact that $|\varphi_k(t_i^-)| \lesssim \|\varphi_k\|_{\widehat{H}^1(t_{i-1}, t_i)}$, (4.8), and (4.9); and finally in the last inequality we used (4.10) and (4.13). Hence, (4.7) holds and $E \in \sigma(\mathbb{H}_\omega)$ as asserted. \square

Remark 4.2 It is natural to conjecture that the spectrum for the half-line operator H_ω is a deterministic set given by the union of periodic spectra of $H(b, \ell, q)$. The latter, under some spectral monotonicity assumption, in turn equals the union of constant spectra, which in certain scenarios can be computed explicitly. However, neither standard ergodicity arguments (e.g., proof of Pastur’s Theorem) nor spectral theoretical arguments (cf. [60, proof Lemma 1.4.2] and [48]) seem to be applicable to the *half-line* models in question. We note that the half-line models present both probabilistic and spectral-theoretical complications which are not typical for operators on \mathbb{R} .

4.2 Proof of dynamical and exponential localization for metric trees

We say that a function $f : \Gamma_{b,\ell} \rightarrow \mathbb{R}$ is *tree-exponentially decaying* if there exist $\lambda \geq 0$ and $C = C(f, \lambda) > 0$ such that

$$|f(x)| \leq \frac{C e^{-\lambda|x|}}{\sqrt{w_o(|x|)}},$$

where $w_o(|x|)$ denotes the number of vertices in the same generation as x ; cf. (2.1).

Proof of Theorem 1.2 (i) By Theorem 4.1 and part (ii) of Theorem 3.13, there exist full measure sets $\tilde{\Omega}, \tilde{\Omega} \subset \Omega$ such that

$$\sigma(\mathbb{H}_\omega) = \Sigma, \quad \sigma_c(H_\omega) = \emptyset, \quad \omega \in \hat{\Omega} \cap \tilde{\Omega},$$

and the operator H_ω enjoys a basis of exponentially decaying eigenfunctions. Then letting

$$\Omega^* := \bigcap_{n \in \mathbb{Z}_+} T^{-n}(\hat{\Omega} \cap \tilde{\Omega}), \quad (4.14)$$

we notice that $\mu(\Omega^*) = 1$ and that

$$\sigma(\mathbb{H}_\omega) = \Sigma, \quad \sigma_c(\mathbb{H}_\omega) = \overline{\bigcup_{n \in \mathbb{Z}_+} \sigma_c(H_{T^n \omega})} = \emptyset, \quad \omega \in \Omega^*,$$

where we used the orthogonal decomposition (2.12). Next we show that \mathbb{H}_ω admits a basis of tree-exponentially decaying eigenfunctions almost surely. To that end, let us fix $\omega \in \tilde{\Omega}$, $v \in \mathcal{V} \setminus \{o\}$, $\text{gen}(v) = n \in \mathbb{N}$, and $1 \leq k \leq b_{n-1}$. Then it suffices to construct a basis of tree-exponentially decaying eigenfunctions in $\mathcal{L}_{v,k} = \mathcal{U}_{v,k}^{-1}(L^2(t_\omega(n), \infty))$, cf. (2.4), (2.5). For a basis element $f \in \ker(H_{T^n \omega} - E)$ of $L^2(t_\omega(n), \infty)$, we define the corresponding basis element of $\mathcal{L}_{v,k}$,

$$\psi_f := \mathcal{U}_{v,k}^{-1} f, \quad \psi_f \in \text{dom}(\mathbb{H}_\omega).$$

Then (2.9) yields

$$|\psi_f(x)| \leq \frac{C_f e^{-\frac{\tilde{L}(E)|x|}{2}}}{\sqrt{w_v(|x|)}}.$$

A basis of tree-exponentially decaying eigenfunctions of \mathcal{L}_o can be constructed similarly.

(ii) Let $v \in \mathcal{V}$ and $n := \text{gen}(v)$, then by Part (iii) of Theorem 3.13, the subspace $\text{ran}(\chi_I(H_{T^n \omega}))$ is spanned by semi-uniformly localized eigenfunctions

$$f_{n,j} \in \ker(H_{T^n \omega} - E_j(n)), \quad j \in \mathbb{Z}_+, \quad E_j(n) \in I, \quad n = \text{gen}(v). \quad (4.15)$$

For $1 \leq k \leq b_{n-1}$, $j \in \mathbb{Z}_+$ we introduce

$$\psi_{v,k,j} := \mathcal{U}_{v,k}^{-1} f_{n,j} \in \text{dom}(\mathbb{H}_\omega),$$

and notice that

$$\text{supp}(\psi_{v,k,j}) \subset T_v, \quad (4.16)$$

the forward subtree rooted at v . Then for $\omega \in \Omega^*$ one has (abbreviating $\Gamma = \Gamma_{b_\omega, \ell_\omega}$):

$$\begin{aligned} & \left\| |X|^p \chi_I(\mathbb{H}_\omega) e^{-it\mathbb{H}_\omega} \chi_K \right\|_{L^2(\Gamma)} \\ & \leq \sum_{v \in \mathcal{V}} \sum_{k=1}^{b_v-1} \sum_{\left\{ j: \begin{array}{l} E_j(n) \in I, \\ E_j(n) \text{ as in (4.15)} \end{array} \right\}} |\langle \psi_{v,k,j}, \chi_K \rangle_{L^2(\Gamma)}| \| |X|^p \psi_{v,k,j} \|_{L^2(\Gamma)} \\ & \stackrel{(4.16)}{\leq} \sum_{v \in \mathcal{V}, T_v \cap K \neq \emptyset} \sum_{\substack{1 \leq k \leq b_v-1 \\ j: E_j(n) \in I}} |\langle \psi_{v,k,j}, \chi_K \rangle_{L^2(\Gamma)}| \| |X|^p \psi_{v,k,j} \|_{L^2(\Gamma)} \\ & \leq \sum_{\substack{v \in \mathcal{V}, T_v \cap K \neq \emptyset \\ 1 \leq k \leq b_v-1 \\ j: E_j(n) \in I}} \int_{K \cap T_v} |\psi_{v,k,j}(x)| dx \left(\int_{\Gamma} x^{2p} |\psi_{v,k,j}(x)|^2 dx \right)^{1/2} \\ & \leq \sum_{\substack{v \in \mathcal{V}, T_v \cap K \neq \emptyset, \\ n = \text{gen}(v), \\ 1 \leq k \leq b_v-1, \\ j: E_j(n) \in I}} \int_{|K \cap T_v|} (w_v(t))^{1/2} |f_{n,j}(t + |v|)| dt \\ & \quad \times \left(\int_{\Gamma} |x|^{2p} |\psi_{v,k,j}(x)|^2 dx \right)^{1/2} \\ & \leq \sum_{\substack{v \in \mathcal{V}, T_v \cap K \neq \emptyset, \\ n = \text{gen}(v), \\ j: E_j(n) \in I}} C_{v,K} \int_{|K \cap T_v|} |f_{n,j}(t + |v|)| dt \\ & \quad \times \left(\int_{|v|}^{\infty} |\tau|^{2p} |f_{n,j}(\tau)|^2 d\tau \right)^{1/2}, \end{aligned} \quad (4.17)$$

where $|K| := [0, \text{diam}(K)]$. Proceeding as in (3.54), (3.55) with ψ replaced by the characteristic function of the interval $[0, \text{diam}(K)]$, we deduce that (4.17) converges as asserted. \square

Remark 4.3 We notice that all eigenfunctions ψ_E (including those corresponding to energies $E \in \mathfrak{D}$) satisfy

$$|\psi_E(x)| \leq \frac{C e^{-\lambda_E |x|}}{\sqrt{w_o(|x|)}}, \quad (4.18)$$

for some $\lambda_E \geq 0$ and $C > 0$, where $w_o(|x|)$ denotes the number of vertices in the same generation as x ; cf. (2.1). Moreover, one has $\lambda_E > 0$ whenever $E \notin \mathfrak{D}$, in particular, (4.18) yields $\psi_E \in L^2(\Gamma_{b,\ell})$ in this case. Furthermore, if $E \in \mathfrak{D}$ and $\lambda_E = 0$ then ψ_E still decays exponentially, $|\psi_E(v)| \leq \frac{C}{2^{|\text{gen}(v)|/2}}$ for all $v \in \mathcal{V}$. However, this inequality alone is insufficient to deduce $L^2(\Gamma_{b,\ell})$ integrability. The analogous issue does not arise in the setting of metric graphs for which the volume of the ball centered at the root with radius r grows polynomially as $r \uparrow +\infty$, e.g., as in the metric graph spanned by \mathbb{Z}^d .

Part 2. Anderson localization for discrete radial trees

5 Random discrete trees

This part of the paper concerns Anderson localization for discrete radial trees.

Hypothesis 5.1 *Let $\Gamma = (\mathcal{V}, \mathcal{E})$ be a rooted, radial discrete tree. Assume that the branching numbers $b_v \in [b^-, b^+]$, $b^- \geq 2$, and the potential $q_v \in [q^-, q^+]$ are radial. Let*

$$p : \{(u, v) \in \mathcal{V}^2 : d(u, v) = 1\} \rightarrow [p^-, p^+],$$

be radial, symmetric, and bounded, that is,

$$p(u, v) = p_{\min(\text{gen}(u), \text{gen}(v))}, \text{ for } u, v \in \mathcal{V};$$

and $p := \{p_n\}_{n=0}^\infty \subset [p^-, p^+]$, $p_{-1} = 0$, $p^\pm \in (0, \infty)$.

Assuming this hypothesis, we introduce a bounded operator $\mathbb{J}(b, p, q) \in \mathcal{B}(\ell^2(\mathcal{V}))$ as follows

$$(\mathbb{J}(b, p, q)f)(u) := \sum_{v \sim u} p(u, v)(q(u)f(u) - f(v)), \quad f \in \ell^2(\mathcal{V}). \quad (5.1)$$

In this part, we adopt the notation of the previous sections with the convention that all edges have length one. Thus, for vertices $x, y \in \mathcal{V}$, $\text{dist}(x, y)$ is the combinatorial distance between them, and, in particular $|x| = \text{gen}(x)$ for all $x \in \mathcal{V}$.

5.1 The almost-sure spectrum for discrete models

The following hypothesis is assumed throughout this section.

Hypothesis 5.2 *Let $\tilde{\mu}$ be a probability measure with $\text{supp}(\tilde{\mu}) = \mathcal{A}$, $\#\mathcal{A} \geq 2$, and either*

$$\mathcal{A} \subseteq \{b_-, \dots, b_+\} \times \{1\} \times [q_-, q_+] \quad (5.2)$$

or

$$\begin{aligned} \mathcal{A} &\subseteq \{b_-, \dots, b_+\} \times [p_-, p_+] \times \{0\} \\ &\text{and } \exists(b, p, 0), (b', q', 0) \in \text{supp } \tilde{\mu} \text{ with } p\sqrt{b} \neq p'\sqrt{b'}. \end{aligned} \quad (5.3)$$

Let us remark that the secondary hypothesis in (5.3) is essential, for, if $\text{supp } \tilde{\mu}$ is concentrated on a set for which $q = 0$ and $p\sqrt{b} = \text{const.}$, then the Jacobi matrices arising in the orthogonal decomposition of \mathbb{J}_ω will all have constant entries.

We introduce $(\Omega, \mu) := (\mathcal{A}^{\mathbb{Z}_+}, \tilde{\mu}^{\mathbb{Z}_+})$. For $\omega \in \Omega$, define the operators $\mathbb{J}_\omega := \mathbb{J}(b_\omega, p_\omega, q_\omega)$ and Jacobi matrices $J_\omega := J(b_\omega, p_\omega, q_\omega)$, where

$$\{(b_\omega(n), p_\omega(n), q_\omega(n))\}_{n=0}^\infty,$$

is a sequence of i.i.d. random vectors with common distribution $\tilde{\mu}$. Let us notice that

$$\mathbb{J}_\omega = \begin{cases} \mathbb{S}_\omega \text{ (cf. (1.6))}, & \text{if (5.2) holds,} \\ \mathbb{A}_\omega \text{ (cf. (1.7))}, & \text{if (5.3) holds.} \end{cases}$$

In particular,

- Random Branching Model (RBM) arises when $\text{supp } \tilde{\mu} \subseteq \{b_-, \dots, b_+\} \times \{1\} \times \{1\}$,
- Random Weight Model (RWM) arises when $\text{supp } \tilde{\mu} \subseteq \{d\} \times [p_-, p_+] \times \{0\}$,
- Random Schrödinger Operator (RSO) arises when $\text{supp } \tilde{\mu} \subseteq \{d\} \times \{1\} \times [q_-, q_+]$.

Remark 5.3 We point out that RBM and RSO concern random realizations of the *discrete Laplace operator*, while RWM is focused on *the adjacency matrices*, i.e. $q \equiv 0$. Typically (e.g., for \mathbb{Z}^d models) the distinction between the discrete Laplace operator and the adjacency matrix of the graph is irrelevant as the two operators differ by a scalar multiple of the identity operator. In the setting of non-constant trees, however, the distinction is more subtle since it depends on the branching numbers. What is more, the consecutive transfer matrices for RWM are correlated unless $q \equiv 0$.

Abusing notation somewhat, we will identify a scalar with a constant sequence consisting of that scalar, for example writing $\mathbb{A}(2, 1, 0)$ to mean the adjacency operator for which all branching numbers are two and all p 's are one.

Theorem 5.4 *There exists a full μ -measure set $\widehat{\Omega} \subset \Omega$ such that*

$$\sigma(\mathbb{A}_\omega) = \Sigma := \overline{\bigcup_{(b,p) \text{ periodic}} \sigma(\mathbb{A}(b, p, 0))}, \quad \omega \in \widehat{\Omega}. \quad (5.4)$$

Proof First, we show that

$$\sigma(\mathbb{A}_\omega) \subset \Sigma, \quad \text{for all } \omega \in \Omega.$$

Seeking contradiction, we assume that $E \in \sigma(\mathbb{A}_\omega) \setminus \Sigma$ for some $\omega \in \Omega$. Then there exist

$$\{f_k\}_{k=1}^\infty \subset \ell^2(\Gamma) \text{ and } \{m_k\}_{k=1}^\infty \subset \mathbb{N},$$

such that

$$\begin{aligned}\|f_k\|_{\ell^2(\Gamma)} &= 1, \quad \text{supp}(f_k) \subset B(o; m_k), \\ \|(\mathbb{A}_\omega - E)f_k\|_{\ell^2(\Gamma)} &\rightarrow 0, \quad k \rightarrow \infty.\end{aligned}\tag{5.5}$$

where $B(o; m_k)$ denotes the ball centered at o with radius m_k . The $m_k + 2$ -periodic sequence with the first $m_k + 2$ elements given by $\omega_1, \dots, \omega_{m_k+1}$ is denoted by $(b^k, p^k, 0)$. Then since $E \notin \Sigma$ one has

$$\|(\mathbb{A}(b^k, p^k, 0) - E)^{-1}\|_{\mathcal{B}(\ell^2(\Gamma))} \leq C < \infty,$$

and thus for all k we get

$$\|(\mathbb{A}_\omega - E)f_k\|_{\ell^2(\Gamma)} = \|(\mathbb{A}(b^k, p^k, 0) - E)f_k\|_{\ell^2(\Gamma)} \geq C^{-1} > 0,$$

which contradicts (5.5).

Next, we show

$$\Sigma \subset \sigma(\mathbb{A}_\omega)$$

for almost all ω . To that end, we first notice that there exists $\widehat{\Omega} \subset \Omega$, $\mu(\widehat{\Omega}) = 1$ such that for arbitrary

$$\omega \in \widehat{\Omega}, \quad (b, p, 0) \in \text{supp}(\mu), \quad \text{and } \{m_k\}_{k=1}^\infty \subset \mathbb{N},\tag{5.6}$$

there exists a sequence $\{r_k\}_{k=1}^\infty$ such that for all $k \in \mathbb{N}$ one has

$$b_\omega(r_k + i) = b_i \text{ for all } i \in \{0, \dots, m_k + 1\},\tag{5.7}$$

$$\max_{0 \leq i \leq m_k + 1} |p_\omega(i + r_k) - p_i| = o(1),\tag{5.8}$$

see, for example, [47, Proposition 3.8]. Pick an arbitrary periodic sequence $(b, p, 0) \in \text{supp}(\mu)$ and an arbitrary $E \in \sigma(\mathbb{A}(b, p, 0))$. Then there exist $\{\varphi_k\}_{k=1}^\infty \subset \ell^2(\Gamma)$ and $\{m_k\}_{k=1}^\infty \subset \mathbb{N}$ such that

$$\begin{aligned}\|\varphi_k\|_{\ell^2(\Gamma)} &= 1, \quad \text{supp}(\varphi_k) \subset B(o; m_k), \quad k \in \mathbb{N}, \\ \|(\mathbb{A}(b, p, q) - E)\varphi_k\|_{\ell^2(\Gamma)} &\rightarrow 0, \quad k \rightarrow \infty.\end{aligned}\tag{5.9}$$

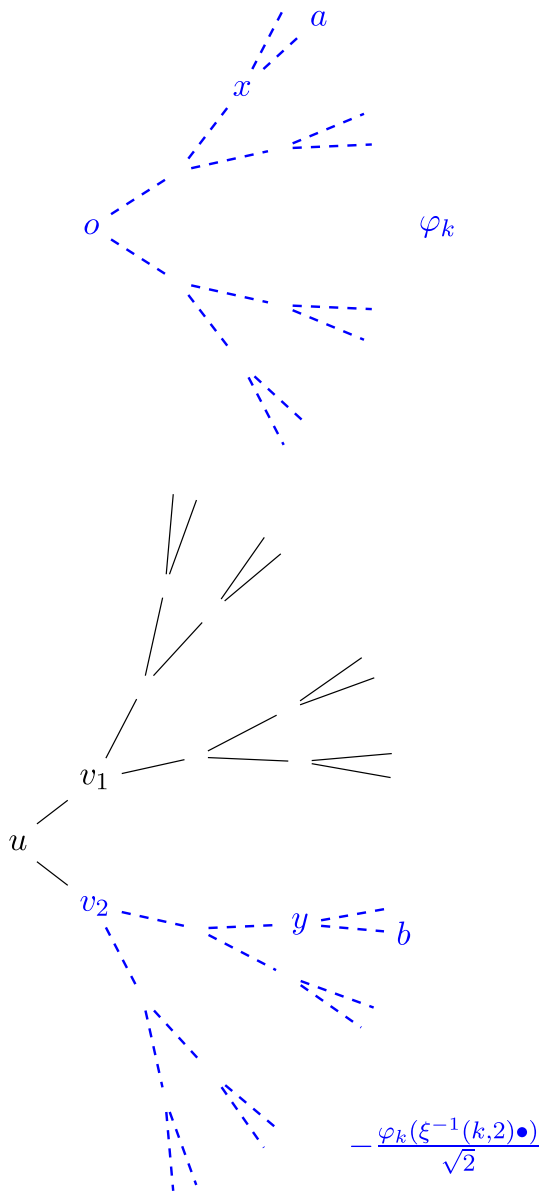
Given (5.6)–(5.9) we are ready to produce a Weyl sequence for \mathbb{A}_ω .

For a fixed $k \in \mathbb{N}$, pick two distinct vertices v_1, v_2 in generation r_k with common backward neighbor $u \in \mathcal{V}$ (in generation $r_k - 1$), see Fig. 2. Then by (5.7) there exists a pair of graph isomorphisms

$$\xi(k, i) : B(o; m_k + 1) \rightarrow T_{v_i} \cap B(v_i; m_k + 1).$$

We notice that

Fig. 2 Top panel: $T(k)$. Bottom panel: vertices v_1, v_2 in generation r_k with common backward neighbor $u, b_{r_k} = 2$. Subtree in blue (dashed) is $T(k, 2)$. The isomorphism $\xi(k, 2)$ maps $T(k)$ onto $T(k, 2)$, in particular $o \mapsto v_2, x \mapsto y, a \mapsto b$, blue (dashed) tree in the top panel gets mapped into the blue (dashed) subtree in the bottom panel



$$\xi(k, i)(o) = v_i, \quad i = 1, 2, \quad k \in \mathbb{N}. \quad (5.10)$$

For brevity, we denote

$$T(k) := \Gamma \cap B(o; m_k + 1), \quad T(k, i) := T_{v_i} \cap B(v_i; m_k + 1). \quad (5.11)$$

Let us define

$$(W_k \varphi)(x) := \begin{cases} 2^{-1/2} \varphi(\xi^{-1}(k, 1)x), & x \in T(k, 1), \\ -2^{-1/2} \varphi(\xi^{-1}(k, 2)x), & x \in T(k, 2), \\ 0, & \text{otherwise} \end{cases} \quad (5.12)$$

for $\varphi \in \ell^2(\mathcal{V})$ which is supported on $B(o, m_k + 1)$. We claim that $\{W_k \varphi_k\}_{k \geq 1}$ is a Weyl sequence for \mathbb{A}_ω , $\omega \in \widehat{\Omega}$. To that end, let us first notice

$$\begin{aligned} & \left| \|(\mathbb{A}(b, p, 0) - E)\varphi_k\|_{\ell^2(\Gamma)} - \|(\mathbb{A}_\omega - E)W_k \varphi_k\|_{\ell^2(\Gamma)} \right| \\ &= \left| \|W_k(\mathbb{A}(b, p, 0) - E)\varphi_k\|_{\ell^2(\Gamma)} - \|(\mathbb{A}_\omega - E)W_k \varphi_k\|_{\ell^2(\Gamma)} \right| \\ &\leq \|W_k(\mathbb{A}(b, p, 0) - E)\varphi_k - (\mathbb{A}_\omega - E)W_k \varphi_k\|_{\ell^2(\Gamma)} \\ &= \|W_k \mathbb{A}(b, p, 0)\varphi_k - \mathbb{A}_\omega W_k \varphi_k\|_{\ell^2(\Gamma)}, \end{aligned}$$

where we used $\|(\mathbb{A}(b, p, 0) - E)\varphi_k\|_{\ell^2(\Gamma)} = \|W_k(\mathbb{A}(b, p, 0) - E)\varphi_k\|_{\ell^2(\Gamma)}$ which follows from the definition of W_k . Next, recalling (5.10) and the fact that u is the common backward neighbor of v_1, v_2 we get

$$\begin{aligned} (\mathbb{A}_\omega W_k \varphi_k)(u) &= p_\omega(u, v_1)[W_k \varphi_k](v_1) + p_\omega(u, v_2)[W_k \varphi_k](v_2) \\ &= \frac{p_\omega(u, v_1)\varphi_k(o) - p_\omega(u, v_2)\varphi_k(o)}{\sqrt{2}} = 0, \end{aligned}$$

since $p_\omega(u, v_1) = p_\omega(u, v_2)$. Further, one has

$$W_k(\mathbb{A}(b, p, q)\varphi_k)(u) = 0 = (\mathbb{A}_\omega W_k \varphi_k)(u), \quad (5.13)$$

where the first equality follows from (5.12). Next, let us fix $i = 1, 2$, $k \in \mathbb{N}$ and use the shorthand $\xi_k := \xi(k, i)$. For $y \in T(k, i)$ let $x := \xi_k^{-1}(y)$, see Fig. 2, then one has

$$W_k(\mathbb{A}(b, p, 0)\varphi_k)(y) - [\mathbb{A}_\omega(W_k \varphi_k)](y) \quad (5.14)$$

$$\begin{aligned} &= \frac{1}{\sqrt{2}}(\mathbb{A}(b, p, 0)\varphi_k)(x) - [\mathbb{A}_\omega(W_k \varphi_k)](y) \\ &= -\frac{1}{\sqrt{2}} \sum_{a \sim x} p(x, a)\varphi_k(a) + \sum_{b \sim y} p_\omega(y, b)(W_k \varphi_k)(b) \\ &= -\frac{1}{\sqrt{2}} \left(\sum_{a \sim x} p(x, a)\varphi_k(a) - \sum_{b \sim \xi_k(x)} p_\omega(\xi_k(x), b)\varphi_k(\xi_k^{-1}b) \right). \end{aligned} \quad (5.15)$$

Let us point out that $\xi_k^{-1}(b)$ is not defined if $b \notin T_{v_1} \cup T_{v_2}$. However, one does have $W_k \varphi_k(b) = 0$ and therefore the equality in (5.15) holds with

$$\varphi_k(\xi_k^{-1}(b)) := W_k \varphi_k(b) = 0. \quad (5.16)$$

Moreover, combining this and (5.7) we obtain

$$\begin{aligned} & \sum_{a \sim x} p(x, a) \varphi_k(a) - \sum_{b \sim y} p_\omega(y, b) \varphi_k(\xi_k^{-1}(b)) \\ &= \sum_{a \sim x} (p(x, a) - p_\omega(\xi_k(x), \xi_k(a))) \varphi_k(a). \end{aligned} \quad (5.17)$$

Given (5.16) and (5.17) we are ready to continue (5.14), (5.15). Changing variables via $b = \xi_k(a)$, we get

$$\begin{aligned} & W_k(\mathbb{A}(b, p, 0)\varphi_k)(y) - [\mathbb{A}_\omega(W_k\varphi_k)](y) \\ &= -\frac{1}{\sqrt{2}} \left(\sum_{a \sim x} [p(x, a) - p_\omega(\xi_k(x), \xi_k(a))] \varphi_k(a) \right), \end{aligned} \quad (5.18)$$

where we made a change of variable $b = \xi_k(a)$. Furthermore we note that (5.18) holds for $y \in \Gamma \setminus (T(k, 1) \cup T(k, 2))$ trivially, i.e., both sides are equal to zero. Recalling $T(k)$ from (5.11) and using (5.8) yield

$$c(k) := \max_{x \in T(k), x \sim a} |p(x, a) - p_\omega(\xi_k(x), \xi_k(a))|^2 \underset{k \rightarrow \infty}{=} o(1). \quad (5.19)$$

Then combining (5.13), (5.18), and (5.19), we obtain

$$\begin{aligned} & \|W_k\mathbb{A}(b, p, 0)\varphi_k - \mathbb{A}_\omega W_k\varphi_k\|_{\ell^2(\Gamma)}^2 \\ &= \sum_{y \in \Gamma} |W_k(\mathbb{A}(b, p, 0)\varphi_k)(y) - [\mathbb{A}_\omega(W_k\varphi_k)](y)|^2 \\ &= \sum_{x \in T(k)} \left| \sum_{a \sim x} [p(x, a) - p_\omega(\xi_k(x), \xi_k(a))] \varphi_k(a) \right|^2 \\ &\leq c(k)C(b^+) \|\varphi_k\|_{\ell^2(\Gamma)}^2 \underset{k \rightarrow \infty}{=} o(1), \end{aligned}$$

where $C(b^+) > 0$ is some fixed constant. Therefore, we get

$$\left| \|\mathbb{A}(b, p, 0) - E\|_{\ell^2(\Gamma)} - \|\mathbb{A}_\omega - E\|_{\ell^2(\Gamma)} \right| \underset{k \rightarrow \infty}{=} o(1).$$

Thus $\{W_k\varphi_k\}_{k \geq 1}$ is a Weyl sequence for \mathbb{A}_ω and $E \in \sigma(\mathbb{A}_\omega)$ as asserted. \square

Remark 5.5 (1) We emphasize that the equality in (5.17) requires special attention if $y \in \partial(T(k, i))$, since in this case the inclusion

$$\xi_k(\{a \in \mathcal{V} : a \sim x\}) \subset \{b \in \mathcal{V} : b \sim y\},$$

could be strict. However, by (5.16) the equality (5.17) holds as asserted even in this special case. Due to this nuance the current proof is not applicable to $\mathbb{J} = \mathbb{S}$.

(Informally, if $q \neq 0$ in (5.1) then we “see” extra bits around v_i which are not observed near o).

- (2) The almost-sure spectrum Σ for $\mathbb{A}_\omega = \mathbb{A}(b_\omega, 1, 0)$ can be computed explicitly if $p \equiv 1, q \equiv 0$, i.e. the random branching model for the adjacency matrix. Indeed, in this case, the quadratic form \mathfrak{a} of the \mathbb{A} is given by

$$\mathfrak{a}[\varphi, \varphi] = - \sum_{u \sim v} \varphi(u) \overline{\varphi(v)}, \varphi \in \ell^2(\Gamma).$$

therefore

$$\|\mathbb{A}(b, 1, 0)\|_{\mathcal{B}(\ell^2(\Gamma))} \leq \|\mathbb{A}(\tilde{b}, 1, 0)\|_{\mathcal{B}(\ell^2(\Gamma))},$$

where $\tilde{b} := \max\{P_1 \text{ supp } \mu\}$ and P_1 is the first coordinate function. Combining this and (5.4) we get

$$\begin{aligned} \Sigma &= \overline{\bigcup_{b \text{ periodic}} \sigma(\mathbb{A}(b, 1, 0))} \subset [-\|\mathbb{A}(\tilde{b}, 1, 0)\|_{\mathcal{B}(\ell^2(\Gamma))}, \|\mathbb{A}(\tilde{b}, 1, 0)\|_{\mathcal{B}(\ell^2(\Gamma))}] \\ &= [-2\sqrt{\tilde{b}}, 2\sqrt{\tilde{b}}] \subset \Sigma. \end{aligned}$$

As before, we note that this proof is not applicable to the case $q \neq 0$ or $p \neq \text{const}$.

- (3) Remark 5.3, the proof of Theorem 5.4, the previous remark, and the question of computing the almost-sure spectrum itself illustrate a subtle distinction between adjacency matrices and Schrödinger operators. This issue arises even in the most simple case $\Gamma = \mathbb{Z}_+, p \equiv 1$, and random q , since [in view of (1.7)]

$$\mathbb{S} = \begin{bmatrix} q(1) & -1 & & & \\ -1 & 2q(2) & -1 & & \\ & -1 & 2q(3) & -1 & \\ & & \ddots & \ddots & \ddots \end{bmatrix}.$$

To be more specific, if one considers

$$\tilde{\mathbb{S}} = \begin{bmatrix} 2q(1) & -1 & & & \\ -1 & 2q(2) & -1 & & \\ & -1 & 2q(3) & -1 & \\ & & \ddots & \ddots & \ddots \end{bmatrix}$$

where $\{q(n)\}$ is a sequence of i.i.d. random variables, then it is well-known that the spectrum of $\tilde{\mathbb{S}}$ is almost surely given by $[-2, 2] + 2 \text{ supp}\{q\}$. Since \mathbb{S} is a rank-one perturbation of $\tilde{\mathbb{S}}$, their essential spectra coincide. However, depending on the support of q , it can happen that \mathbb{S} may have discrete eigenvalues outside of $\sigma_{\text{ess}}(\tilde{\mathbb{S}})$, and these eigenvalues may not be constant almost-surely. Thus, one should not expect the analogue of Theorem 5.4 to hold for random Schrödinger operators on graphs (as opposed to adjacency matrices).

5.2 Breuer-type decomposition

Our next objective is to revise the Breuer decomposition [13, Theorem 2.4] which may be viewed as a discrete version of the orthogonal decomposition of *metric* trees. To point out a difference between the two, we note: The invariant subspaces in (2.10) are parametrized by single vertices, while those in Breuer's decomposition are parametrized by entire generations of vertices.

Theorem 5.6 *Assume Hypothesis 5.1. Then there exists a unitary operator*

$$\Phi_b : \ell^2(\mathcal{V}) \rightarrow \bigoplus_{n=0}^{\infty} \bigoplus_{k=1}^{m(n)} \ell^2(\mathbb{Z}_+),$$

such that

$$\Phi_b \mathbb{J}(b, p, q) \Phi_b^{-1} = \bigoplus_{n=0}^{\infty} \bigoplus_{k=1}^{m(n)} J(T^n b, T^n p, T^n q), \quad (5.20)$$

where $m(n) := b_0 \cdot b_1 \cdots b_{n-1} \cdot (b_n - 1)$, $n \in \mathbb{Z}_+$, and $J(b, p, q)$ denotes the Jacobi matrix acting in $\ell^2(\mathbb{Z}_+)$ and given by

$$J(b, p, q) := \begin{pmatrix} (b_0 p_0 + p_{-1}) q_0 & \sqrt{b_0} p_0 & 0 & & \\ \sqrt{b_0} p_0 & (b_1 p_1 + p_0) q_1 & \sqrt{b_1} p_1 & \ddots & \\ 0 & \sqrt{b_1} p_1 & \ddots & \ddots & \\ & \ddots & \ddots & \ddots & \ddots \end{pmatrix}. \quad (5.21)$$

Proof Breuer's inductive procedure [13, Theorem 2.4] yields an orthonormal basis

$$\{\varphi_{n,k,j} : n \in \mathbb{Z}_+, 1 \leq k \leq m(n), j \in \mathbb{Z}_+\} \subset \ell^2(\mathcal{V}).$$

For all admissible triples n, k, j , the basis elements satisfy

$$\text{supp}(\varphi_{n,k,j}) \subset \{v \in \mathcal{V} : \text{gen}(v) = n + j\}, \quad (5.22)$$

$$\varphi_{n,k,j+1}(v) = \begin{cases} \frac{\varphi_{n,k,j}(u)}{\sqrt{b_{n+j}}}, & u \sim v, \text{ gen}(v) = \text{gen}(u) + 1, \\ 0, & \text{otherwise,} \end{cases} \quad (5.23)$$

and

$$\mathbb{J}(b, p, q) \varphi_{n,k,j} = \begin{cases} \sqrt{b_{n+j-1}} p_{n+j-1} \varphi_{n,k,j-1} \\ \quad + (b_{n+j} p_{n+j} + p_{n+j-1}) q_{n+j} \varphi_{n,k,j} \\ \quad + \sqrt{b_{n+j}} p_{n+j} \varphi_{n,k,j+1}, & j \geq 1, \\ (b_n p_n + p_{n-1}) q_n \varphi_{n,k,0} + \sqrt{b_n} p_n \varphi_{n,k,1}, & j = 0. \end{cases} \quad (5.24)$$

The latter shows that the operator $\mathbb{J}(b, p, q)$ leaves the subspaces

$$\mathcal{H}_{n,k} := \overline{\text{span}\{\varphi_{n,k,j} : j \in \mathbb{Z}_+\}} \subset \ell^2(\mathcal{V})$$

invariant. Thus we have

$$\ell^2(\mathcal{V}) = \bigoplus_{n=0}^{\infty} \bigoplus_{k=1}^{m(n)} \mathcal{H}_{n,k}, \quad \mathbb{J}(b, q, w) P_{\mathcal{H}_{n,k}} = J(T^n b, T^n p, T^n q), \quad (5.25)$$

where $P_{\mathcal{H}_{n,k}}$ denotes an orthogonal projection onto $\mathcal{H}_{n,k}$ in $\ell^2(\mathcal{V})$. Let us define unitary operators

$$\begin{aligned} \mathcal{U}_{n,k} : \mathcal{H}_{n,k} &\rightarrow \ell^2(\mathbb{Z}_+), \quad n \in \mathbb{Z}_+, 1 \leq k \leq m(n), \\ \mathcal{U}_{n,k} \varphi_{n,k,j} &:= \delta_j, \quad j \in \mathbb{Z}_+. \end{aligned}$$

and

$$\Phi_b := \bigoplus_{n=0}^{\infty} \bigoplus_{k=1}^{m(n)} \mathcal{U}_{n,k}.$$

Then (5.24) together with (5.25) yield (5.20) and (5.21) as asserted. \square

5.3 Dynamical and exponential localization for discrete random trees

In this section we discuss spectral and dynamical localization for three discrete models: the random branching model (RBM), the random weights (RWM) model, and random Schrödinger operators (RSO).

Let us denote the nonzero entries of $J(b, p, q)$ by

$$\begin{aligned} \beta_j &= \beta_j(b, p, q) = (b_j p_j + p_{j-1}) q_j, \\ \alpha_j &= \alpha_j(b, p) = \sqrt{b_j} p_j, \quad j \in \mathbb{Z}_+. \end{aligned}$$

Then a sequence $u = \{u_j\}_{j=0}^{\infty}$ satisfies $J(b, p, q)u = Eu$, $E \in \mathbb{R}$, that is,

$$\begin{cases} \alpha_{j-1} u_{j-1} + (\beta_j - E) u_j + \alpha_j u_{j+1} = 0, & j \in \mathbb{N}, \\ (\beta_0 - E) u_0 + \alpha_0 u_1 = 0, \end{cases}$$

if and only if

$$\begin{bmatrix} u_{j+1} \\ \alpha_j u_j \end{bmatrix} = M^{E,j}(b, p, q) \begin{bmatrix} u_j \\ \alpha_{j-1} u_{j-1} \end{bmatrix}, \quad \text{for all } j \in \mathbb{N}.$$

where

$$\begin{aligned} M^{E,j}(b, p, q) &:= \frac{1}{\alpha_j} \begin{bmatrix} E - \beta_j & -1 \\ \alpha_j^2 & 0 \end{bmatrix} \\ &= \begin{bmatrix} \frac{E - (b_j p_j + p_{j-1}) q_j}{\sqrt{b_j p_j}} & -\frac{1}{\sqrt{b_j p_j}} \\ \sqrt{b_j p_j} & 0 \end{bmatrix}. \end{aligned} \quad (5.26)$$

The transfer matrix (5.26) gives rise to an $\mathrm{SL}(2, \mathbb{R})$ -cocycle

$$(T, M^E) : \Omega \times \mathbb{R}^2 \rightarrow \Omega \times \mathbb{R}^2, \quad (T, M^E)(\omega, v) = (T\omega, M^E(\omega)v),$$

where $M^E : \Omega \rightarrow \mathrm{SL}(2, \mathbb{R})$ and

$$M^E(\omega) = \begin{bmatrix} \frac{E - (b_\omega(0)p_\omega(0) + p_\omega(-1))q_\omega(0)}{\sqrt{b_\omega(0)p_\omega(0)}} & -\frac{1}{\sqrt{b_\omega(0)p_\omega(0)}} \\ \sqrt{b_\omega(0)p_\omega(0)} & 0 \end{bmatrix}.$$

The n -step transfer matrix $M_n^E(\omega)$ and the Lyapunov exponent are defined as in (3.2) and (3.3) respectively.

Theorem 5.7 *Assume Hypothesis 5.2. Then there is a set $\mathcal{D} \subseteq \mathbb{R}$ of cardinality at most one such that $G = G_{v(E)}$ enjoys the following properties for $E \in \mathbb{R} \setminus \mathcal{D}$.*

- (i) G is noncompact
- (ii) G is strongly irreducible
- (iii) G is contracting (cf. [19, Definition 2.8])
- (iv) $\mathrm{Fix}(G) = \emptyset$

In particular, L is continuous and positive on $\mathbb{R} \setminus \mathcal{D}$.

Proof Following the proof of Theorem 3.5, we choose

$$(b_1, p_1, q_1) \neq (b_2, p_2, q_2) \in \mathrm{supp} \tilde{\mu},$$

let $M_j(E)$ denote the transfer matrix corresponding to (b_j, p_j, q_j) , and form the matrices $A = M_1 M_2^{-1}$ and $g = [M_1, M_2]$. Let us comment briefly on the method of proof. We can immediately apply [20] to deduce that there is an unspecified discrete set of energies away from which (i)–(iv) hold. In fact, the argument of [20] applies away from energies at which $\mathrm{tr} M_j(E) = 0$ or $\det g(E) = 0$, which allows us to refine this to a discrete set with no more than 3 elements. However, we can do better still: conditions (i)–(iv) hold for any E for which the following criterion is met:

$$\nexists \mathcal{F} \subseteq \mathbb{RP}^1 \text{ with } \#\mathcal{F} \in \{1, 2\} \text{ such that } M_j \mathcal{F} = \mathcal{F} \text{ for } j = 1, 2. \quad (5.27)$$

In particular, (5.27) implies (iii) which in turn implies (i) by standard arguments about $\mathrm{SL}(2, \mathbb{R})$. Once (i) holds, then (5.27) immediately yields (iv) and also implies (ii) (cf. [11]).

Case 1: (5.2) holds. We have $p_1 = p_2 = 1$, so

$$M_j = \frac{1}{\sqrt{b_j}} \begin{bmatrix} E - (b_j + 1)q_j & -1 \\ b_j & 0 \end{bmatrix}.$$

We calculate

$$g = \frac{1}{\sqrt{b_1 b_2}} \begin{bmatrix} b_1 - b_2 & (b_1 + 1)q_1 - (b_2 + 1)q_2 \\ (b_1 - b_2)E + b_2(b_1 + 1)q_1 - b_1(b_2 + 1)q_2 & b_2 - b_1 \end{bmatrix}.$$

Case 1a: $b_1 = b_2$. It follows that $q_1 \neq q_2$ and hence $(b_1 + 1)q_1 \neq (b_2 + 1)q_2$. One can confirm that $\det g(E) \neq 0$ for all E , so that M_1 and M_2 have no eigenvectors in common. Thus, there is no \mathcal{F} of cardinality one with $M_j \mathcal{F} = \mathcal{F}$ for $j = 1, 2$. Now, suppose that an invariant $\mathcal{F} \subseteq \mathbb{R}P^1$ of cardinality two exists. We must then have $\mathcal{F} = \{\bar{u}_1, \bar{u}_2\}$ and $M_j \bar{u}_1 = \bar{u}_2$, $M_j \bar{u}_2 = \bar{u}_1$ for some j ; without loss, assume $j = 1$. This forces $\text{tr } M_1 = 0$. However, since $(b_1 + 1)q_1 \neq (b_2 + 1)q_2$, we must have $\text{tr } M_2 \neq 0$, so $M_2 \mathcal{F} = \mathcal{F}$ forces $M_2 \bar{u}_k = \bar{u}_k$ for $k = 1, 2$, that is to say, each \bar{u}_k is an eigendirection of M_2 . Identifying $\mathbb{C}P^1$ with the Riemann sphere in the usual way, write z_k for the image of \bar{u}_k under the identification $\mathbb{C}P^1 \cong \mathbb{C} \cup \{\infty\}$. Since $M_2 z_k = z_k$, we have

$$\frac{E - (b_2 + 1)q_2}{b_2} - \frac{1}{b_2 z_k} = z_k, \quad k = 1, 2.$$

From this, we deduce $z_1 z_2 = 1/b_2$. On the other hand, since $\text{tr } M_1 = 0$, we observe

$$M_1 z_1 = -\frac{1}{b_1 z_1} \neq z_2, \quad M_1 z_2 = -\frac{1}{b_1 z_2} \neq z_1,$$

a contradiction. Thus, when $b_1 = b_2$, (5.27) holds and we have (i)–(iv) for every $E \in \mathbb{R}$.

Case 1b: $b_1 \neq b_2$. There are two further subcases to consider.

Case 1bi: $(b_1 + 1)q_1 = (b_2 + 1)q_2$. Then, $\det g(E) \neq 0$ for every E . Thus, again M_1 and M_2 never share an eigenvector. At energy $E = E_0 := (b_1 + 1)q_1 = (b_2 + 1)q_2$, both M_1 and M_2 preserve $\mathcal{F} = \{\text{span}(\bar{e}_1), \text{span}(\bar{e}_2)\}$. Since E_0 is the only energy at which $\text{tr } M_j$ vanishes for either j , we have (i)–(iv) for $E \in \mathbb{R} \setminus \{E_0\}$.

Case 1bii: $(b_1 + 1)q_1 \neq (b_2 + 1)q_2$. One can check that $\det g(E)$ vanishes for exactly one value of $E_1 \in \mathbb{R}$. Using the same argument as in Case 1a, we see that there is no invariant \mathcal{F} of cardinality one or two away from $E = E_1$. Thus, (i)–(iv) hold away from $\mathcal{D} = \{E_1\}$.

Case 2: (5.3) holds. Then,

$$M_j = \frac{1}{p_j \sqrt{b_j}} \begin{bmatrix} E & -1 \\ p_j^2 b_j & 0 \end{bmatrix}, \quad \text{and} \quad p_1 \sqrt{b_1} \neq p_2 \sqrt{b_2}.$$

Notice that

$$A := M_1 M_2^{-1} = \frac{1}{p_1 p_2 \sqrt{b_1 b_2}} \begin{bmatrix} p_2^2 b_2 & 0 \\ 0 & p_1^2 b_1 \end{bmatrix}.$$

Since $p_1 \sqrt{b_1} \neq p_2 \sqrt{b_2}$, A is hyperbolic⁶ and any finite set of directions left invariant by M_1 , M_2 , and A must be a subset of $\{\text{span}(\vec{e}_1), \text{span}(\vec{e}_2)\}$. It is easy to see that this cannot happen for $E \neq 0$, so we may take $\mathcal{D} = \{0\}$ in this case. \square

Remark 5.8 Let us note that the need to remove a single point is sharp. For example, in Case 1bi above, one can verify that $L(E_0) = 0$. To see this, write $r = -(b_1/b_2)^{1/2}$ and $R = \text{diag}(r, r^{-1})$, and observe that

$$M_j(E_0)M_k(E_0) = \begin{cases} -I & j = k \\ R^{-1} & (j, k) = (1, 2) \\ R & (j, k) = (2, 1). \end{cases}$$

Thus, by passing to blocks of length two and using the strong law of large numbers, we deduce $L(E_0) = 0$.

Proof of Theorem 1.4 Now that we know that L is positive and obeys a uniform LDT away from \mathcal{D} , spectral and dynamical localization for J_ω follows as in Theorem 3.13, see also [27] where spectral localization was proved for the discrete RBM. Let Ω^* be defined as in (4.14) (where $\widehat{\Omega}$ is as in Theorem 5.4, and $\widetilde{\Omega}$ is a full measure set realizing localization for J_ω) and fix $\omega \in \Omega^*$.

For all $n \in \mathbb{Z}_+$, the spectral subspace $\text{ran}(\chi_I(J_{T^n \omega}))$ enjoys an orthonormal basis $\{f_{n,j}\}_{j=0}^\infty$ of eigenfunctions of $J_{T^n \omega}$ corresponding to energies $E \in I$. If we define $\psi_{n,k,j} := \mathcal{U}_{n,k}^{-1} f_{n,j}$, then

$$\{\psi_{n,k,j} : n \in \mathbb{Z}_+, 1 \leq k \leq m(n), j \in \mathbb{Z}_+\}$$

is an orthonormal basis of $\text{ran}(\chi_I(\mathbb{J}_\omega))$.

Proof of (1.8). For an arbitrary admissible triple n, k, j we will prove (1.8) with $f = \psi_{n,k,j}$. First, we note that by spectral localization for J_ω one has

$$|f_{n,j}(p)| \leq C(f_{n,j})e^{-\lambda p}, \quad p \in \mathbb{Z}_+; \quad \lambda := \min_{E \in I} \frac{L(E)}{2} > 0,$$

for some $C(f_{n,j}) > 0$. Then for $|x| > n$ we get

$$\begin{aligned} |\psi_{n,k,j}(x)| &= |\mathcal{U}_{n,k}^{-1} f_{n,j}(x)| = |f_{n,j}(|x| - n) \varphi_{n,k,|x|-n}(x)| \\ &\stackrel{(5.23)}{\leq} \frac{C(\psi_{n,k,j})e^{-\lambda(|x|-n)}}{\sqrt{w_o(|x|)}}, \end{aligned} \tag{5.28}$$

⁶ I.e., $|\text{tr}(A)| > 2$.

which implies (1.8).

Proof of (1.9). Due to dynamical localization for J_ω one has

$$\sum_{j \in \mathbb{Z}_+} |\langle f_{n,j}(p), f_{n,j}(q) \rangle_{\ell^2(\mathbb{Z}_+)}| \leq C_n e^q e^{-\theta(p-q)}, \quad (5.29)$$

for all $p \geq q$, $\theta < \min_{E \in I} L(E)$, and a constant $C_n = C(n, \omega, \theta) > 0$ (cf., e.g., [19, Proof of Theorem 6.4] where this step is discussed for the standard Anderson Hamiltonian). Next, we have

$$\begin{aligned} & \sup_{t>0} |\langle \delta_x, \chi_I(\mathbb{J}_\omega) e^{-it\mathbb{J}_\omega} \delta_y \rangle_{\ell^2(\mathcal{V})}| \\ & \leq \sum_{\substack{n \in \mathbb{Z}_+ \\ 1 \leq k \leq m(n)}} \sum_{j=0}^{\infty} |\psi_{n,k,j}(x) \psi_{n,k,j}(y)| \\ & \stackrel{(5.22)}{\leq} \sum_{\substack{0 \leq n \leq |y| \\ 1 \leq k \leq m(n)}} \sum_{j=0}^{\infty} |\psi_{n,k,j}(x) \psi_{n,k,j}(y)| \\ & \stackrel{(5.28)}{=} \sum_{\substack{0 \leq n \leq |y| \\ 1 \leq k \leq m(n)}} \sum_{j=0}^{\infty} |f_{n,j}(|x| - n) \varphi_{n,k,|x|-n}(x) f_{n,j}(|y| - n) \varphi_{n,k,|y|-n}(y)| \\ & \stackrel{(5.23)}{\leq} \sum_{\substack{0 \leq n \leq |y| \\ 1 \leq k \leq m(n)}} \sum_{j=0}^{\infty} \frac{|f_{n,j}(|x| - n) f_{n,j}(|y| - n)|}{\sqrt{w_y(|x| - |y| - 1)}} \\ & \stackrel{(5.29)}{\leq} \sum_{\substack{0 \leq n \leq \text{gen}(y) \\ 1 \leq k \leq m(n)}} \frac{C_n e^{|y|} e^{-\theta(|x| - |y|)}}{\sqrt{w_y(|x| - |y| - 1)}} \leq \frac{C_y e^{-\theta(\text{dist}(x,y))}}{\sqrt{w_y(|x| - |y|)}}. \end{aligned}$$

Finally, (1.10) follows from (1.9) by summation in x . \square

Acknowledgements We thank G. Berkolaiko, M. Lukic, and G. Stolz for helpful discussions, and P. Hislop for bringing our attention to this subject and for motivating discussions.

References

1. Aizenman, M., Sims, R., Warzel, S.: Absolutely continuous spectra of quantum tree graphs with weak disorder. *Commun. Math. Phys.* **264**, 371–389 (2006)
2. Aizenman, M., Sims, R., Warzel, S.: Stability of the absolutely continuous spectrum of random Schrödinger operators on tree graphs. *Probab. Theory Relat. Fields* **136**, 363–394 (2006)
3. Aizenman, M., Sims, R., Warzel, S.: Fluctuation Based Proof of the Stability of AC Spectra of Random Operators on Tree Graphs. *Recent Advances in Differential Equations and Mathematical Physics*, *Contemp. Math.*, vol. 412, pp. 1–14. American Mathematical Society, Providence (2006)
4. Aizenman, M., Warzel, S.: Absence of mobility edge for the Anderson random potential on tree graphs at weak disorder. *EPL (Europhysics Letters)* **96**, 37004 (2011)

5. Aizenman, M., Warzel, S.: Resonant delocalization for random Schrödinger operators on tree graphs. *J. Eur. Math. Soc.* **15**, 1167–1222 (2013)
6. Aizenman, M., Warzel, S.: *Random Operators: Disorder Effects on Quantum Spectra and Dynamics*, Graduate Studies in Mathematics, vol. 168. American Mathematical Society, Providence (2015)
7. Albeverio, S., Gesztesy, F., Hoegh-Krohn, R., Holden, H.: with app. by P. Exner, *Solvable Models in Quantum Mechanics*, 2nd edn. AMS-Chelsea Series, Amer. Math. Soc. Providence, RI (2005)
8. Birman, BSh, Solomyak, M.Z.: *Spectral Theory of Self-Adjoint Operators in Hilbert Space*. D. Reidel Publishing Co., Dordrecht (1987)
9. Berkolaiko, G., Kuchment, P.: *Introduction to Quantum Graphs*, Mathematical Surveys and Monographs, vol. 186. AMS, Providence (2012)
10. Berkolaiko, G., Latushkin, Y., Sukhtaiev, S.: Limits of quantum graph operators with shrinking edges. *Adv. Math.* **352**, 632–669 (2019)
11. Bougerol, P., Lacroix, J.: *Products of Random Matrices with Applications to Schrödinger Operators*. Birkhäuser, Basel (1985)
12. Bourgain, J., Schlag, W.: Anderson localization for Schrödinger operators on \mathbb{Z} with strongly mixing potentials. *Commun. Math. Phys.* **215**, 143–175 (2000)
13. Breuer, J.: Singular continuous spectrum for the Laplacian on certain sparse trees. *Commun. Math. Phys.* **219**, 851–857 (2007)
14. Breuer, J.: Localization for the Anderson model on trees with finite dimensions. *Ann. Henri Poincaré* **8**, 1507–1520 (2007)
15. Breuer, J., Denisov, S., Eliasz, L.: On the essential spectrum of Schrödinger operators on trees. *Math. Phys. Anal. Geom.* **21**(4), Art. 33, 25 pp (2018)
16. Breuer, J., Frank, R.: Singular spectrum for radial trees. *Rev. Math. Phys.* **21**, 929–945 (2009)
17. Breuer, J., Keller, M.: Spectral analysis of certain spherically homogeneous graphs. *Oper. Matrices* **7**, 825–847 (2013)
18. Breuer, J., Levi, N.: On the decomposition of the Laplacian on metric graphs. Preprint [arXiv:1901.00349v1](https://arxiv.org/abs/1901.00349v1)
19. Bucaj, V., Damanik, D., Fillman, J., Gerbuz, V., VandenBoom, T., Wang, F., Zhang, Z.: Localization for the one-dimensional Anderson model via positivity and large deviations for the Lyapunov exponent. *Trans. Am. Math. Soc.* **372**(5), 3619–3667 (2019)
20. Bucaj, V., Damanik, D., Fillman, J., Gerbuz, V., VandenBoom, T., Wang, F., Zhang, Z.: Positive Lyapunov exponents and a large deviation theorem for continuum Anderson models, briefly. *J. Funct. Anal.* **277**(9), 3179–3186 (2019)
21. Burenkov, V.I.: *Sobolev Spaces on Domains*. B.G. Teubner, Stuttgart-Leipzig (1998)
22. Carlson, R.: Nonclassical Sturm–Liouville problems and Schrödinger operators on radial trees. *Electron. J. Differ. Equ.* **71**, 1–24 (2000)
23. Craig, W., Simon, B.: Subharmonicity of the Lyapunov index. *Duke Math. J.* **50**, 551–560 (1983)
24. Damanik, D., Lenz, D., Stolz, G.: Lower transport bounds for one-dimensional continuum Schrödinger operators. *Math. Ann.* **336**, 361–389 (2006)
25. Damanik, D., Sims, R., Stolz, G.: *Lyapunov Exponents in Continuum Bernoulli–Anderson models, Operator methods in ordinary and partial differential equations* (Stockholm, 2000), *Oper. Theory Adv. Appl.*, vol. 132, pp. 121–130. Birkhäuser, Basel (2002)
26. Damanik, D., Sims, R., Stolz, G.: Localization for one-dimensional continuum Bernoulli–Anderson models. *Duke Math. J.* **114**, 59–100 (2002)
27. Damanik, D., Sukhtaiev, S.: Anderson localization for radial tree graphs with random branching numbers. *J. Funct. Anal.* **277**, 418–433 (2019)
28. Ekholm, T., Frank, R., Kovarik, H.: Remarks about Hardy inequalities on metric trees. In: Exner, P., et al. (eds.) *Proceedings of Symposium Pure Mathematical Analysis on Graphs and its Applications*, vol. 77, pp. 369–379. American Mathematical Society, Providence (2008)
29. Ekholm, T., Frank, R., Kovarik, H.: Eigenvalue estimates for Schrödinger operators on metric trees. *Adv. Math.* **226**, 5165–5197 (2011)
30. Frank, R., Kovarik, H.: Heat kernels of metric trees and applications. *SIAM J. Math. Anal.* **45**, 1027–1046 (2013)
31. Evans, W.D., Harris, D.J.: Fractals, trees and the Neumann Laplacian. *Math. Ann.* **296**, 493–527 (1993)
32. Evans, W.D., Harris, D.J., Pick, L.: Weighted Hardy and Poincaré inequalities on trees. *J. Lond. Math. Soc.* **52**, 121–136 (1995)

33. Froese, R., Hasler, D., Spitzer, W.: Transfer matrices, hyperbolic geometry and absolutely continuous spectrum for some discrete Schrödinger operators on graphs. *J. Funct. Anal.* **230**, 184–221 (2006)
34. Froese, R., Hasler, D., Spitzer, W.: Absolutely continuous spectrum for the Anderson model on a tree: a geometric proof of Klein's theorem. *Commun. Math. Phys.* **269**, 239–257 (2007)
35. Froese, R., Lee, D., Sadel, C., Spitzer, W., Stolz, G.: Localization for transversally periodic random potentials on binary trees. *J. Spectr. Theory* **6**, 557–600 (2016)
36. Fürstenberg, H.: Noncommuting random products. *Trans. Am. Math. Soc.* **108**, 377–428 (1963)
37. Fürstenberg, H., Kifer, Y.: Random matrix products and measures on projective spaces. *Isr. J. Math.* **46**, 12–32 (1983)
38. Germinet, F., De Bièvre, S.: Dynamical localization for discrete and continuous random Schrödinger operators. *Commun. Math. Phys.* **194**, 323–341 (1998)
39. Goldstein, M., Schlag, W.: Hölder continuity of the integrated density of states for quasi-periodic Schrödinger equations and averages of shifts of subharmonic functions. *Ann. Math. (2)* **154**, 155–203 (2001)
40. Gorodetski, A., Kleptsyn, V.: Parametric Fürstenberg theorem on random products of $SL(2, \mathbb{R})$ matrices. Preprint ([arXiv:1809.00416](https://arxiv.org/abs/1809.00416))
41. Grubb, G.: *Distributions and Operators*. Springer, New York (2009)
42. Grigorchuk, R., Zuk, A.: The lamplighter group as a group generated by a 2-state automaton, and its spectrum. *Geom. Dedicata* **87**, 209–244 (2001)
43. Harrell, E.M., Maltsev, A.V.: On Agmon metrics and exponential localization for quantum graphs. *Commun. Math. Phys.* **359**, 429–448 (2018)
44. Hislop, P., Post, O.: Anderson localization for radial tree-like quantum graphs. *Waves Random Complex Media* **19**, 216–261 (2009)
45. Jitomirskaya, S., Zhu, X.: Large deviations of the Lyapunov exponent and localization for the 1D Anderson model. *Comm. Math. Phys.* **370**, 311–324 (2019)
46. Kato, T.: *Perturbation Theory for Linear Operators*. Springer, Berlin (1980)
47. Kirsch, W.: An invitation to random Schrödinger operators, *Panor. Synthèses, Random Schrödinger Operators*, vol. 25, pp. 1–119. Society of Mathematics, Paris (2008)
48. Kirsch, W., Martinelli, F.: On the spectrum of Schrödinger operators with a random potential. *Commun. Math. Phys.* **85**, 329–350 (1982)
49. Klein, A.: Absolutely continuous spectrum in the Anderson model on the Bethe lattice. *Math. Res. Lett.* **1**, 399–407 (1994)
50. Klein, A.: Spreading of wave packets in the Anderson model on the Bethe lattice. *Commun. Math. Phys.* **177**, 755–773 (1996)
51. Klein, A.: Extended states in the Anderson model on the Bethe lattice. *Adv. Math.* **133**, 163–184 (1998)
52. Naimark, K., Solomyak, M.: Geometry of Sobolev spaces on regular trees and the Hardy inequalities. *Russ. J. Phys.* **8**, 322–335 (2001)
53. Naimark, K., Solomyak, M.: Eigenvalue estimates for the weighted Laplacian on metric trees. *Proc. Lond. Math. Soc.* **80**, 690–724 (2000)
54. Oseledec, V.I.: A multiplicative ergodic theorem. Characteristic Ljapunov, exponents of dynamical systems (Russian). *Trudy Moskov. Mat. Obšč* **19**, 179–210 (1968)
55. Pastur, L., Figotin, A.: *Spectra of Random and Almost-Periodic Operators*. Springer, Berlin (1992)
56. Ruelle, D.: Ergodic theory of differentiable dynamical systems. *Inst. Hautes Études Sci. Publ. Math.* **50**, 27–58 (1979)
57. Schmied, M., Sims, R., Teschl, G.: On the absolutely continuous spectrum of Sturm-Liouville operators with applications to radial quantum trees. *Oper. Matrices* **2**, 417–434 (2008)
58. Sobolev, A., Solomyak, M.: Schrödinger operators on homogeneous metric trees: spectrum in gaps. *Rev. Math. Phys.* **14**, 421–467 (2002)
59. Solomyak, M.: On the spectrum of the Laplacian on regular metric trees. Special section on quantum graphs. *Waves Random Media* **14**, S155–S171 (2004)
60. Stollmann, P.: *Caught by Disorder, Bound States in Random Media*, Progress in Mathematical Physics, vol. 20. Birkhäuser, Boston (2001)
61. Stollmann, P.: Scattering by obstacles of finite capacity. *J. Funct. Anal.* **121**, 416–425 (1994)