Subordinacy theory on star-like graphs

Netanel Levi

Abstract. We study Jacobi matrices on star-like graphs, which are graphs that are given by the pasting of a finite number of half-lines to a compact graph. Specifically, we extend subordinacy theory to this type of graphs, that is, we find a connection between asymptotic properties of solutions to the eigenvalue equations and continuity properties of the spectral measure with respect to the Lebesgue measure. We also use this theory in order to derive results regarding the multiplicity of the singular spectrum.

1. Introduction

In this work, we are interested in studying spectral properties of Jacobi matrices on certain types of graphs. A graph G is given by a pair $\langle V, E \rangle$, where V is the set of vertices, which we assume throughout to be infinite, and $E \subseteq V \times V$ is the set of edges. Given $u, v \in V$, we will denote $u \sim v$ if $(u, v) \in E$. We will also denote by δ_v the function on V which is defined by

$$
\delta_v(w) = \begin{cases} 1 & w = v, \\ 0 & \text{otherwise.} \end{cases}
$$

A Jacobi matrix J on G is given by a set of real numbers, $\{b_v\}_{v \in V}$, and a set of positive numbers, $\{a_e\}_{e \in E}$. J acts on $\ell^2(G)$ by

$$
J(\varphi)(u) = \sum_{w \sim u} a_{(u,w)} \varphi_w + b_u \varphi_u, \tag{1.1}
$$

For simplicity, we assume throughout that $\{b_v\}_{v \in V}$ and $\{a_e\}_{e \in E}$ are bounded, and in that case J is a bounded self-adjoint operator on $\ell^2(V)$. By the spectral theorem (see, e.g., [3, Chapter 6]), J gives rise to a projection-valued measure P on $\sigma(H)$, which is called *the spectral measure* of J. Our goal is to study continuity properties of P

via the asymptotics of solutions to the eigenvalue equation. Namely, given a Jacobi matrix J and $E \in \mathbb{R}$, we study solutions to the formal difference equation

$$J\varphi = E\varphi, \tag{1.2}$$

where by formal we mean that φ need not be in $\ell^2(G)$. We will distinguish between supports of the absolutely continuous and singular parts of P (with respect to the Lebesgue measure) by studying the asymptotic properties of these solutions.

The method extended in this work is known as *subordinacy theory*. It has been developed in [9] for the case of continuum Schrödinger operators on a half-line. In that case, continuity properties of the spectral measure are determined by comparing the growth of solutions which satisfy different boundary conditions to the differential equation associated with the operator. The discrete analogue of subordinacy theory was later developed to Jacobi matrices on $\mathbb{N}$ in [14]. In the latter case, the operator is simply given by a tridiagonal matrix whose entries are bounded and real-valued. Given such an operator J and $\theta \in [0, \pi)$, the operator J_θ is defined by setting $J_\theta = J - \tan(\theta)\langle\delta_1, \cdot\rangle\delta_1$. The method of subordinacy enables one to examine continuity properties of the spectral measure of these operators by comparing formal solutions to the equation

$$J_\theta\varphi = E\varphi, \quad \theta \in [0, \pi), \ E \in \mathbb{R}. \tag{1.3}$$

Given $\theta \in [0, \pi)$, a solution φ to (1.3) can also be regarded as a solution to the same equation with $\theta = 0$, along with the boundary condition

$$\varphi(0)\cos\theta + \varphi(1)\sin\theta = 0.$$

Throughout, the solution to (1.3) with $\theta = 0$ will be referred to as the solution which satisfies a Dirichlet boundary condition.

We now turn to briefly describe the method of subordinacy. Let $\theta \in [0, \pi)$. Note that δ_1 is a cyclic vector for J_θ, namely $\ell^2(\mathbb{N}) = \overline{\text{sp}\{\delta_1, J_\theta\delta_1, J_\theta^2\delta_1, \ldots\}}$, and so the spectral measure of J_θ is equivalent to that of δ_1, in the sense that they have the same null sets. Given $u: \mathbb{N} \to \mathbb{R}$ and $L > 0$, denote

$$\|u\|_L := \Big[\sum_{n=1}^{[L]} |u(n)|^2 + (L - [L])|u([L] + 1)|^2 \Big]^{1/2}.$$

Definition 1.1. A solution ψ to (1.3) will be called *subordinate* if for any other linearly independent solution φ,

$$\lim_{L\to\infty} \frac{\|\psi\|_L}{\|\varphi\|_L} = 0.$$

Denote by μ_θ the spectral measure of δ_1, and by $(\mu_\theta)_s$, $(\mu_\theta)_{ac}$ its singular and absolutely continuous parts (with respect to the Lebesgue measure) respectively. In [14], it is proved that $(\mu_\theta)_s$ is supported on the set of energies for which the solution to (1.3) is subordinate, and $(\mu_\theta)_{ac}$ is supported on the set of energies for which no subordinate solution exists. The theory was further developed in various directions. In [7], it was extended to continuum Schrödinger operators on $\mathbb{R}$. In [11], Jitomirskaya and Last present a strengthening of the theory in the discrete half-line case (i.e. Jacobi matrices on $\mathbb{N}$), and use it to relate the asymptotic properties of the solutions to continuity properties of the spectral measure with respect to various Hausdorff measures. Subordinacy is a very powerful tool as in many cases, the study of asymptotic properties of solutions to (1.2) is more accessible, compared to classical methods such as the direct study of the Borel transform of μ. It has many applications and generalizations, [4, 5, 8, 12, 15, 17, 22] is a very partial list.

In this work, we extend the theory of subordinacy to a certain kind of graphs which we call *star-like*. Roughly speaking, a star-like graph G is a graph which consists of a compact component $C = \langle V_C, E_C \rangle$ along with a finite collection of half-lines attached to it. In Figure 1 we give a few examples of star-like graphs. Although the compact component is not unique, our results do not depend on its choice (see Remark 2.8). Thus, throughout we fix a compact component C and refer to it as *the compact component* of G. For every $v \in V_C$, we denote by G_v the half-line attached to v. If v has no half-line attached to it, then $G_v = \{v\}$.

Definition 1.2. A solution ψ to (1.2) will be called *subordinate* if and only if it is not identically zero, and it is subordinate as a solution on G_v for every $v \in C$ such that G_v is a half-line.

Remark 1.3. Note that, although ψ must not be identically zero, it may vanish on one or more of the half-lines which are attached to C. For example, In graph (a) of Figure 1, there may be eigenvectors which vanish on one of the half-lines. This type of example is discussed in Section 5.2.

Let us illustrate this definition with an example. Suppose $G = \mathbb{Z}$ and C is the subgraph of $\mathbb{Z}$ which consists of the vertices $V_C = \{-1, 0, 1\}$ (in that case, the half-lines are attached to -1 and to 1). Let $E \in \mathbb{R}$ and let $\psi : \mathbb{Z} \to \mathbb{R}$ satisfy $J\psi = E\psi$. Let $J_{\pm 1}$ be the Jacobi matrices restricted to the half-lines $G_{\pm 1}$. Then $\psi|_{G_1}$ is a solution to (1.3) (with J_1 instead of J) with $\theta_1 = \arctan(-\frac{\psi(0)}{\psi(1)})$, and $\psi|_{G_{-1}}$ is a solution to (1.3) (with J_{-1} instead of J) with $\theta_{-1} = \arctan(-\frac{\psi(0)}{\psi(-1)})$. In that case, the solution ψ will be called *subordinate* if both $\psi|_{G_1}$ and $\psi|_{G_{-1}}$ are subordinate as solutions to the half-line problem.

Our main result is the following.
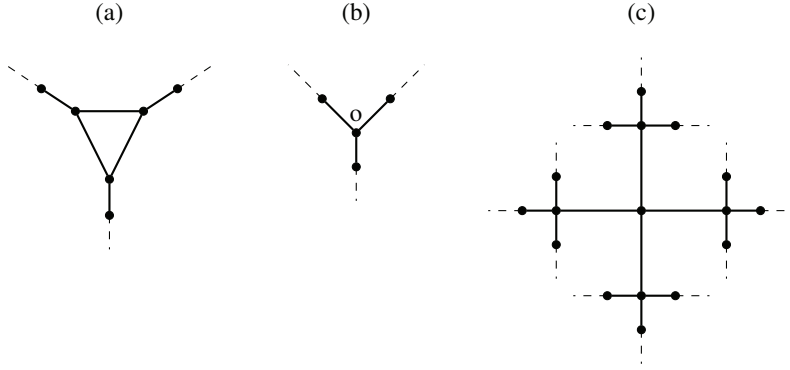
Figure 1. Three examples of star-like graphs. The dashed lines stand for copies of $\mathbb{N}$. The graph in (b) is also called a star graph, while (c) is a "trimming" of a 4-regular tree. In the graph (a), the compact component can be taken to be any finite subgraph which contains the inner triangle. In (b), the compact component can be taken to be any finite subgraph which contains the vertex o and two of its neighbors.

Theorem 1.4. *Let G be a star-like graph, and let $J:\ell^2(G) \to \ell^2(G)$ be as in (1.1). Denote the spectral measure of J by P, and let P_{s}, P_{ac} be its singular and absolutely continuous parts (with respect to the Lebesgue measure) respectively. Then P_{s} is supported on the set*

$$S = \{E \in \mathbb{R}: \text{there exists a subordinate solution to (1.2) on } G\}$$

and P_{ac} is supported on the set

$$N = \bigcup_{v \in C} N_v,$$

where

$$N_v = \{E \in \mathbb{R}: \text{there exists no subordinate solution to (1.2) on } G_v\}.$$

As noted before, when $G = \mathbb{N}$, δ_1 is a cyclic vector for J. In particular, this implies that the spectrum is simple. This is no longer the case for when $G = \mathbb{Z}$, where it can be seen that the absolutely continuous spectrum may have multiplicity 2. Nevertheless, Kac ([13]) showed that the singular spectrum of Schrödinger operators on the real line is simple. Later on, Gilbert ([8]) found a proof of this result using subordinacy theory, and Simon ([19]) found a proof of this fact using the theory of rank one perturbations. The local spectral multiplicity of J can be described via a multiplicity function $N_J:\sigma(J) \to \mathbb{N} \cup \{\infty\}$ (see Section 4 for a precise definition). The results of [8, 13] essentially say that for P_{s}-almost every $E \in \mathbb{R}$, $N_J(E) = 1$.

A generalization of the above result in the continuous setting is given in [21]. They show that for a star-graph with n branches, the local multiplicity is bounded by $n - 1$, and give an explicit formula for N_J. Our second result is the following generalization:

Theorem 1.5. *Let G, J, P be as in Theorem 1.4. Given $E \in \mathbb{R}$, let*

$$S(E) := \{u : G \to \mathbb{R} : u \text{ is a subordinate solution to (1.2) on } G\}.$$

Then, for P_s-almost every $E \in \mathbb{R}$, $N_J(E) \leq \dim S(E)$.

Remark 1.6. It is not hard to show that on each half-line, if a subordinate solution exists then it is unique, which implies that $S(E)$ is a finite-dimensional space and so $\dim S(E)$ is well defined.

When $G = \mathbb{Z}$, it can be seen that $\dim S(E) \leq 1$ for any $E \in \mathbb{R}$, and so simplicity of the singular spectrum in that case is given immediately by Theorem 1.5. With very little additional work, some of the results in [21] can also be obtained in the discrete setting. In general, the inequality cannot be replaced with an equality. We will present an example in which the inequality is strict. Note that by our definition, a solution which is only supported on the compact component is also subordinate and so in general, the multiplicity may exceed the number of half-lines attached to C. However, we will show that in the purely singular continuous part of the spectrum this is not possible. Specifically, we show that for P_s-almost every $E \in \sigma_s(J) \setminus \sigma_{pp}(J)$, $N_J(E)$ is bounded by k, where k is the number of half-lines attached to G. We believe that the bound can be improved to $k - 1$, as this is the case for star-graphs, as shown in [21] for the continuous setting, and in this work for the discrete one. However, our attempts to prove this bound did not succeed.

The rest of the paper is structured as follows. In Section 2, we give some basic measure-theoretic background, present the one-dimensional theory, and introduce the notion of star-like graphs. In Sections 3 and 4, we give proofs of Theorems 1.4 and 1.5 respectively. Section 5 includes some remarks, examples and applications.

2. Preliminaries

We begin by introducing relevant definitions and results regarding the boundary values of Borel transforms of measures.

2.1. Boundary behavior of Borel transforms

Throughout this work we deal with finite complex-valued Borel measures on $\mathbb{R}$. Given such a measure μ, its Borel transform is defined by

$$m_\mu(z) = \int_{\mathbb{R}} \frac{d\mu(x)}{x - z}.$$

It is an analytic function defined on $\mathbb{C}_+ := \{z \in \mathbb{C} : \operatorname{Im} z > 0\}$. As we are interested in the boundary behavior of such functions, given any analytic function $F : \mathbb{C}_+ \to \mathbb{C}_+$ and $E \in \mathbb{R}$, we denote

$$F(E + i0) := \lim_{\varepsilon \to 0} F(E + i\varepsilon)$$

whenever the limit exists. In the case that μ is a positive measure which satisfies

$$\int_{\mathbb{R}} \frac{d\mu(x)}{|x| + 1} < \infty, \tag{2.1}$$

various continuity properties of μ with respect to the Lebesgue measure, which we denote throughout by λ, are related to the boundary behavior of its Borel transform. In particular, we will use the following well-known theorem (see, for example, [18])

Theorem 2.1. *Let μ be a positive measure satisfying (2.1). Denote by $\mu_{\mathrm{ac}}, \mu_{\mathrm{s}}$ the absolutely continuous and singular parts of μ (with respect to the Lebesgue measure) respectively. Then*

(1) *μ_{ac} is supported on the set $\{E \in \mathbb{R} : 0 < \operatorname{Im} m_\mu(E + i0) < \infty\}$;*

(2) *μ_{s} is supported on the set $\{E \in \mathbb{R} : \operatorname{Im} m_\mu(E + i0) = \infty\}$.*

The second type of results that we will need concerns the boundary behavior of ratios of Borel transforms. Specifically, given two Borel measures μ, σ, let $\frac{d\mu}{d\sigma}(E) := \lim_{\varepsilon \to 0} \frac{\mu(E - \varepsilon, E + \varepsilon)}{\sigma(E - \varepsilon, E + \varepsilon)}$ whenever it exists in $\mathbb{C} \cup \{\infty\}$. Note that if $\mu \ll \sigma$, then $\frac{d\mu}{d\sigma}$ coincides with the Radon–Nikodym derivative of μ with respect to σ (as $L^1(\sigma)$ functions). We will use the following version of Poltoratskii's theorem.

Theorem 2.2. *([10, 16]) Let σ, μ be complex-valued Borel measures on $\mathbb{R}$ such that $\mu \ll \sigma$. Denote by σ_{s} the part of σ which is singular with respect to the Lebesgue measure. Then, for σ_{s}-almost every $E \in \mathbb{R}$, the limit $\lim_{\varepsilon \to 0} \frac{m_\mu(E + i\varepsilon)}{m_\sigma(E + i\varepsilon)}$ exists and is equal to $\frac{d\mu}{d\sigma}(E)$.*

We will also need the following result.

Theorem 2.3. ([13]) *Let μ be a real measure and σ be a probability measure, and let $E \in \mathbb{R}$. Assume that $\frac{d\mu}{d\sigma}(E)$ exists in $\mathbb{R}$, and that $\frac{d\sigma}{d\lambda}(E)$ exists, possibly equal to infinity. Then*

$$\lim_{\varepsilon \to 0} \frac{\operatorname{Im}(m_\mu(E + i\varepsilon))}{\operatorname{Im}(m_\sigma(E + i\varepsilon))} = \frac{d\mu}{d\sigma}(E).$$

If μ is a probability measure on $\mathbb{R}$, then as mentioned before, its Borel transform m_μ maps $\mathbb{C}_+ := \{z \in \mathbb{C} : \operatorname{Im} z > 0\}$ to itself analytically. Such functions are called *Herglotz* (sometimes *Nevanlinna* or *Pick*) *functions*. This correspondence also works in the other way (see e.g., [20]), in the following sense: if $m(z) : \mathbb{C}_+ \to \mathbb{C}_+$ is analytic and satisfies

$$m(z) = z^{-1} + O(z^{-2}),$$

then m arises as the Borel transform of some probability measure μ. In this work, we will use the following generalization of the connections mentioned above.

Theorem 2.4. ([6]) *Let $M : \mathbb{C}_+ \to M_n(\mathbb{C})$ be an analytic matrix-valued function such that $\operatorname{Im} M(z)$ is positive semi-definite for all $z \in \mathbb{C}_+$. Then*

(1) *there exists a matrix-valued measure Ω such that*

$$M(z) = C + Dz + \int_{\mathbb{R}} \frac{1 + xz}{x - z} d\Omega(x)$$

where C is a self-adjoint matrix, and D is positive definite;

(2) *the singular part of Ω with respect to λ is supported on the set*

$$\{E \in \mathbb{R} : \operatorname{Im} \operatorname{tr} M(E + i0) = \infty\}.$$

Remark 2.5. For $M \in M_n(\mathbb{C})$, one has $\operatorname{Im} A = \frac{1}{2i}(M - M^*)$.

2.2. Subordinacy theory in the half line case

Let J be a Jacobi matrix on $\mathbb{N}$, and for $\theta \in [0, \pi)$ define the operator J_θ by

$$J_\theta = J - \tan(\theta)\langle \delta_1, \cdot \rangle \delta_1.$$

It can be verified that δ_1 is a cyclic vector for J_θ, namely

$$\ell^2(\mathbb{N}) = \overline{\operatorname{sp}\{J_\theta^k \delta_1 : k \in \mathbb{N} \cup \{0\}\}}. \tag{2.2}$$

Recall that, given a self-adjoint operator H defined on a Hilbert space $\mathcal{H}$ and $\psi, \varphi \in \mathcal{H}$, the spectral measure of ψ and φ with respect to H is defined to be the unique Borel measure which satisfies

$$\langle \psi, f(H)\varphi \rangle = \int_{\sigma(H)} f(t) d\mu_{\psi,\varphi}(t)$$

for any continuous function $f : \sigma(H) \to \mathbb{C}$. If $\psi = \varphi$, we refer to this measure as the spectral measure of ψ with respect to H. Denote by μ_θ the spectral measure of δ_1 with respect to J_θ. It is not hard to see that (2.2) implies that for any $v \in \ell^2(\mathbb{N})$, the spectral measure of v is absolutely continuous with respect to μ_θ. Thus, our task is to study of continuity properties of μ_θ. As noted before, these properties are related to its Borel transform which we denote by m_θ. Note that by the definitions of μ_θ and m_θ, we have

$$m_\theta(z) = \int_{\mathbb{R}} \frac{d\mu_\theta(x)}{x - z} = \langle \delta_1, (J_\theta - z)^{-1} \delta_1 \rangle. \tag{2.3}$$

The main result in half-line subordinacy theory is the following.

Theorem 2.6. *For any $E \in \mathbb{R}$ and $\theta \in [0, \pi)$, $\operatorname{Im} m_\theta(E + i0) = \infty$ if and only if $u_{\theta,E}$ is subordinate.*

Theorem 2.6 was originally proved in the continuous setting by Gilbert and Pearson in [9]. Its discrete analogue is a direct consequence of the results presented in [11].

Remark 2.7. In particular, combining Theorems 2.1 and 2.6, we get that the singular part of μ_θ is supported on the set of energies for which $u_{\theta,E}$ is subordinate.

2.3. Jacobi matrices on star-like graphs

We begin by introducing the notion of a *star-like* graph. Let $C = \langle V_C, E_C \rangle$ be a finite connected graph. A star-like graph is given by selecting a subset of vertices $V_0 \subseteq V_C$, and attaching a copy of $\mathbb{N}$ to each vertex $v \in V_0$. Formally, $G = \langle V, E \rangle$ is given by

$$V = V_C \sqcup \{v_i^u\}_{i \in \mathbb{N}, u \in V_0},$$
$$E = E_C \sqcup \{(u, v_1^u) : i \in V_0\} \sqcup \{(v_i^u, v_{i+1}^u) : i \in \mathbb{N}, u \in V_0\},$$

where $\sqcup$ denotes a disjoint union.

Remark 2.8. A graph G may be constructed by the above procedure in more than one way. For example, $\mathbb{Z}$ can be constructed by selecting C to be either its subgraph which consists of $V_C = \{0, 1\}$, and then $V_0 = V_C$, or by selecting the subgraph which consists of $V_C = \{-1, 0, 1\}$, and then $V_0 = \{-1, 1\}$. However, it is not hard to verify that the sets S and N from Theorem 1.4 do not depend on the choice of C. Thus, throughout we fix C and refer to it as *the compact component* of G.

Let J be a Jacobi matrix on G which acts on $\ell^2(G)$ as in (1.1). Unlike the half-line case, there need not be a cyclic vector for J. However, we have the following.

Claim 2.9. $\ell^2(G) = \overline{\mathrm{sp}\{J^k \delta_v : v \in V_C, k \in \mathbb{N} \cup \{0\}\}}$.

Namely, the set $\{\delta_v : v \in V_C\}$ is cyclic for J.

Proof. We give a sketch of the proof. It suffices to show that

$$\{\delta_v : v \in V\} \subseteq \mathrm{sp}\{J^k \delta_v : v \in V_C, k \in \mathbb{N} \cup \{0\}\} := A.$$

Let $v \in V$. If $v \in V_C$, then clearly $\delta_v \in A$. Otherwise, denote by $d(v, C)$ the length of a minimal path from v to some vertex in V_C, and let $B_k = \{v \in V : d(v, C) = k\}$. If $v \in B_1$, then by the definition of G, there exists some unique $u \in V_C$ such that $v \sim u$. By the definition of J,

$$J \delta_u = b_u \delta_u + a_{(u,v)} \delta_v + \sum_{u \sim w \in C} a_{(u,w)} \delta_w.$$

Now, since $J \delta_u$, $b_u \delta_u + \sum_{u \sim w \in C} a_{(u,w)} \delta_w \in A$, we get that $a_{(u,v)} \delta_v \in A$ as the difference of elements in A, and since $a_{(u,v)} \neq 0$, we get that $\delta_v \in A$. Thus, we get that $B_1 \subseteq A$. Now, noting that by connectedness of G and by the fact that any connected component of the induced graph on $V \setminus V_C$ is isomorphic to $\mathbb{N}$, we get that every $v \in B_2$ has a unique vertex in $u \in B_1$ such that $v \sim u$. From here, one can proceed by induction. ∎

Denote $V_C = \{v_1, \ldots, v_n\}$, and let $\delta_k := \delta_{v_k}$. Claim 2.9 implies that every spectral measure of J is absolutely continuous with respect to the measure

$$\mu := \sum_{k=1}^{n} \mu_k,$$

where μ_k is the spectral measure of δ_k, and so our purpose is to study the continuity properties of μ. Given $z \in \mathbb{C}_+$, define $M(z) \in M_n(\mathbb{C})$ by

$$(M(z))_{kj} = \langle \delta_k, (J - z)^{-1} \delta_j \rangle. \tag{2.4}$$

$M : \mathbb{C}_+ \to M_n(\mathbb{C})$ is analytic and for any $z \in \mathbb{C}_+$, $\mathrm{Im}\, M(z)$ is positive semi-definite. Let Ω be the matrix-valued measure given by Theorem 2.4. By the definition of the spectral measure we get that for any $1 \leq k \leq n$, $\mu_k = \Omega_{kk}$, and so $\mu = \mathrm{tr}\, \Omega$. For $1 \leq k \leq n$, define G_k in the following way.

- If $v_k \in V_0$, then G_k is the graph induced on $\{v_k\} \cup \{v_i^{v_k} : i \in \mathbb{N}\}$. Note that $G_k \cong \mathbb{N}$ with v_k as the origin.

- Otherwise, G_k consists of the singleton $\{v_k\}$.

Finally, denote by J_k the operator $P_k J P_k$ on $\ell^2(G_k)$, where P_k is the projection of $\ell^2(G)$ onto $\ell^2(G_k)$. Note that in the first case, J_k is a Jacobi matrix on the half-line G_k, and in the latter case J_k acts on $\mathbb{C}$ by multiplication by b_{v_k}. In any case, J_k is a self-adjoint operator.

3. Proof of Theorem 1.4

Let $1 \le k \le n$ and let $z \in \mathbb{C}_+$. Denote by $\tilde{u}_k \in \ell^2(G_k)$ the unique ℓ^2 solution to $(J_k - z)u = \delta_k$, and by $\varphi_k \in \ell^2(G)$ the unique ℓ^2 solution to $(J - z)\varphi = \delta_k$ (these solutions exist and are unique since J_k, J are self-adjoint). Let $u_k \in \ell^2(G)$ be defined by

$$u_k(v) = \begin{cases} \tilde{u}_k(v), & v \in G_k, \\ 0, & \text{otherwise.} \end{cases}$$

For $1 \le j < k \le n$, the supports of u_j and u_k are disjoint, and so the collection $\{u_k : 1 \le k \le n\}$ is linearly independent. Denote by m_k the Borel transform of the spectral measure of δ_k with respect to J_k. By (2.3), $m_k(z) = u_k(1)$, and so we have

$$((J - z)u_k)(v) = \begin{cases} 1, & v = v_k, \\ a_{(v_k, v)} m_k(z), & v \sim v_k, v \in V_C, \\ 0, & \text{otherwise.} \end{cases} \tag{3.1}$$

For $1 \le k \le n$, define $w_k \in \mathbb{C}^n$ by $w_k^j = ((J - z)u_k)(v_j)$.

Claim 3.1. *The collection $\{w_k : 1 \le k \le n\}$ is linearly independent.*

Proof. Assume by contradiction that

$$\sum_{k=1}^{n} \alpha_k w_k = 0,$$

and there exists $1 \le j \le n$ such that $\alpha_j \ne 0$. Let

$$u = \sum_{k=1}^{n} \alpha_k u_k.$$

The fact that $\alpha_j \ne 0$ along with the disjointness of the supports of u_j and u_k for any $k \ne j$ implies that $u|_{G_j} \ne 0$. Now, note that for any $1 \le k \le n$, $(J - z)u_k$ vanishes outside V_C which implies that $(J - z)u$ also vanishes outside V_C. In addition,

$$((J - z)u_k)|_{V_C} = w_k,$$

and so

$$((J - z)u)|_{V_C} = \sum_{k=1}^{n} \alpha_k w_k = 0.$$

This implies that z is an eigenvalue of J, which is a contradiction to J being self-adjoint. ∎

The following lemma is crucial to the proof of Theorem 1.4.

Lemma 3.2. *Let $M(z)$ be as in (2.4). Then for any $z \in \mathbb{C}_+$, $M(z)$ is invertible and*

$$M(z)^{-1} = A + \operatorname{diag}\left(\frac{1}{m_1(z)}, \dots, \frac{1}{m_n(z)}\right), \tag{3.2}$$

where A is defined by

$$A_{ij} = \begin{cases} a_{(v_i, v_j)}, & v_i \sim v_j, \\ 0, & \text{otherwise.} \end{cases}$$

Proof. Let $1 \le k \le n$. Note that if $\alpha_1, \dots, \alpha_n \in \mathbb{C}$ satisfy

$$\sum_{j=1}^{n} \alpha_j w_j = \delta_k,$$

then

$$\varphi_k = \sum_{j=1}^{n} \alpha_j u_j.$$

Denote by $F(z)$ the matrix whose j'th column is w_j. By claim 3.1, $F(z)$ is invertible for every $z \in \mathbb{C}_+$, and $\alpha_1, \dots, \alpha_n$ can be retrieved by the equation

$$F(z) \begin{pmatrix} x_1 \\ \vdots \\ x_n \end{pmatrix} = e_k,$$

i.e.,

$$\begin{pmatrix} \alpha_1 \\ \vdots \\ \alpha_n \end{pmatrix} = F(z)^{-1} e_k,$$

where e_k is the k'th element in the standard basis of $\mathbb{C}^n$. Thus, we have that $\alpha_j = (F(z))_{jk}^{-1}$. Note that

$$(M(z))_{kj} = \langle \delta_k, (J - z)^{-1} \delta_j \rangle = \varphi_j(v_k)$$

and since $u_k(v_j) = 0$ whenever $k \ne j$ and $u_k(v_k) = m_k(z)$, we get that

$$(M(z))_{kj} = (F(z))_{jk}^{-1} \cdot m_j(z),$$

i.e.,

$$M(z) = \operatorname{diag}(m_1(z), \dots, m_n(z)) \cdot (F(z))^{-1}.$$

Since for any $1 \le k \le n$, $m_k(z) \in \mathbb{C}_+$, and since $F(z)$ is invertible, we get that $M(z)$ is invertible and

$$M(z)^{-1} = F(z) \cdot \mathrm{diag}\Big(\frac{1}{m_1(z)}, \ldots, \frac{1}{m_n(z)}\Big).$$

and so $M(z)^{-1}$ is given by multiplying the k'th column of $F(z)$ by $\frac{1}{m_k(z)}$. By (3.1), the k'th column of $F(z)$ is given by

$$(F(z)_k)_j = \begin{cases} 1, & v_j = v_k, \\ a_{(v_k, v_j)} m_k(z), & v_j \sim v_k, \\ 0, & \text{otherwise.} \end{cases}$$

Now, (3.2) follows from a straightforward calculation. ∎

Remark 3.3. This is a special case of an identity known as the *Krein, Feshbach,* or *Schur formula* (see, e.g., [1, Chapter 5]), which says that given a self-adjoint operator H acting on a Hilbert space $\mathcal{H}$ and P the orthogonal projection onto a finite-dimensional subspace $\mathcal{H}_0 \subseteq \mathcal{H}$, if we write

$$H = PHP + \hat{H},$$

then, denoting $H_P := PHP$, the following identity holds:

$$P(H - z)^{-1} P = [H_P + [P(\hat{H} - z)^{-1} P]_P^{-1}]_P^{-1}, \tag{3.3}$$

where $[\cdot]_P^{-1}$ indicates that we take the inverse only on $\mathcal{H}_0$. Note that, in our case, $M(z) = P_C(J - z)^{-1} P_C$, and so (3.2) is given by taking the inverse (inside $\mathcal{H}_0$) on both sides of (3.3), and explicitly computing the RHS. The proof of Lemma 3.2 is a more direct way of obtaining this result.

We will also need the following fact.

Claim 3.4. *Let $E \in \mathbb{R}$ and denote $J|_C = P_C J P_C$, where P_C is the projection of $\ell^2(G)$ onto $\ell^2(C)$. Suppose that $\tilde{\varphi} \in \ell^2(C)$ satisfies*

$$(J|_C \tilde{\varphi})(v_k) = E\tilde{\varphi}(v_k)$$

for any $1 \le k \le n$ such that $G_k = \{v_k\}$. Then $\tilde{\varphi}$ can be extended to a function $\varphi \colon G \to \mathbb{C}$ such that $J\varphi = E\varphi$.

Remark 3.5. Note that φ need not be in $\ell^2(G)$, and so the expression $J\varphi = E\varphi$ should be thought of as a difference equation.

Proof. First, note that whenever G_k consists of a single vertex, v_k has no neighbors outside C and so for any $\psi : G \to \mathbb{C}$ which extends $\tilde\varphi$,

$$(J|_C \tilde\varphi)(v_k) = (J\psi)(v_k) = E\psi(v_k).$$

For any $1 \le k \le n$ such that $G_k \cong \mathbb{N}$, the values of φ on G_k can be determined via the difference equation. ∎

We are now ready to prove Theorem 1.4.

Proof. Using the notations from the previous section, for $1 \le j, k \le n$ denote by μ_{jk} the spectral measure of δ_j and δ_k, i.e. μ_{jk} is the unique Borel measure that satisfies

$$\langle \delta_j, f(H)\delta_k \rangle = \int f(x) d\mu_{jk}(x)$$

for any continuous function on the spectrum of J. Also, denote

$$\mu_{kk} := \mu_k$$

and

$$\mu := \sum_{k=1}^{n} \mu_k.$$

Note that

$$(M(z))_{jk} = \int \frac{d\mu_{jk}(x)}{x - z}$$

i.e., $(M(z))_{jk}$ is the Borel transform of μ_{jk}, and $\operatorname{tr} M$ is the Borel transform of μ. For $1 \le k \le n$, let

$$A_k = \left\{ E \in \mathbb{R} : \frac{d\mu_k}{d\mu}(E) > 0 \right\}.$$

Since

$$1 = \frac{d\mu}{d\mu} = \sum_{k=1}^{n} \frac{d\mu_k}{d\mu},$$

we get that

$$A := \bigcup_{k=1}^{n} A_k$$

supports μ. Let $1 \le k \le n$ and let $E \in A_k$. By Theorem 2.1, we may assume that

$$\lim_{\varepsilon \to 0} \operatorname{Im}(\operatorname{tr} M(E + i\varepsilon)) = \infty.$$

By Theorem 2.3, we get that

$$\lim_{\varepsilon \to 0} \frac{\mathrm{Im}(M_{kk}(E + i\varepsilon))}{\mathrm{Im}(\mathrm{tr}\, M(E + i\varepsilon))} > 0$$

and so

$$\lim_{\varepsilon \to 0} \mathrm{Im}(M_{kk}(E + i\varepsilon)) = \infty.$$

Thus, we may assume that $\mu_k|_{A_k}$ is singular with respect to the Lebesgue measure. In addition, $\mu|_{A_k} \ll \mu_k|_{A_k}$, and so for any $1 \leq j \leq n$, $\mu_{jk}|_{A_k} \ll \mu_k|_{A_k}$. By Theorem 2.2, this implies that for μ_k-almost every $E \in A_k$ the limit

$$\lim_{\varepsilon \to 0} \frac{M_{jk}(E + i\varepsilon)}{M_{kk}(E + i\varepsilon)} := \alpha_j$$

exists. In addition, it is not hard to show that we may assume that α_j is real. Now, let e_k be the k'th element in the standard basis of $\mathbb{C}^n$. We have

$$\begin{aligned}
0 &= \lim_{\varepsilon \to 0} \frac{1}{|M_{kk}(E + i\varepsilon)|} \|e_k\| \\
&= \lim_{\varepsilon \to 0} \left\| \frac{e_k}{M_{kk}(E + i\varepsilon)} \right\| \\
&= \lim_{\varepsilon \to 0} \left\| M(E + i\varepsilon)^{-1} \left(M(E + i\varepsilon) \frac{e_k}{M_{kk}(E + i\varepsilon)} \right) \right\|.
\end{aligned}$$

By the fact that

$$M(E + i\varepsilon)(e_k) = \begin{pmatrix} M_{1k}(E + i\varepsilon) \\ \vdots \\ M_{nk}(E + i\varepsilon) \end{pmatrix},$$

we get

$$0 = \lim_{\varepsilon \to 0} M(E + i\varepsilon)^{-1} \begin{pmatrix} \frac{M_{1k}(E+i\varepsilon)}{M_{kk}(E+i\varepsilon)} \\ \vdots \\ \frac{M_{nk}(E+i\varepsilon)}{M_{kk}(E+i\varepsilon)} \end{pmatrix}.$$

Taking into account (3.2), we get that for every $1 \leq j \leq n$,

$$\lim_{\varepsilon \to 0} \sum_{\substack{1 \leq l \leq n \\ v_l \sim v_j}} a_{(v_l, v_j)} \frac{M_{lk}(E + i\varepsilon)}{M_{kk}(E + i\varepsilon)} + \frac{M_{jk}(E + i\varepsilon)}{M_{kk}(E + i\varepsilon)} \cdot \frac{1}{m_j(E + i\varepsilon)} = 0. \tag{3.4}$$

Define $\tilde{\varphi}$ on C by $\tilde{\varphi}(v_j) = \alpha_j$ for every $1 \leq j \leq n$. We claim that $(J|_C \tilde{\varphi})(v_j) = E\tilde{\varphi}(v_j)$ whenever G_j consists of a single vertex, and so by Claim 3.4, it can be extended to a solution φ to (1.2) on G. So, fix such a j. Clearly, we have

$$\lim_{\varepsilon \to 0} \frac{1}{m_j(E + i\varepsilon)} = b_j - E$$

and by (3.4), we get that

$$\sum_{\substack{1 \le l \le n \\ v_l \sim v_j}} a_{(v_l, v_j)} \alpha_l = -(b_j - E)\alpha_j,$$

and so

$$(J|_C \tilde{\varphi})(v_j) = \sum_{\substack{1 \le l \le n \\ v_l \sim v_j}} a_{(v_l, v_j)} \alpha_l + b_j \alpha_j = -(b_j - E)\alpha_j + b_j \alpha_j = E\alpha_j = E\tilde{\varphi}(v_j)$$

as required. Let φ be the completion of $\tilde{\varphi}$ to a solution of (1.2) on G. Note that φ is non-trivial since $\varphi(v_k) = \alpha_k = 1$. We claim that φ is subordinate. Let $1 \le j \le n$ such that G_j is a half-line.

If $\alpha_j = 0$ and there exists a sequence $(\varepsilon_n)_{n=1}^{\infty}$ such that

$$\lim_{n \to \infty} m_j(E + i\varepsilon_n) \ne 0,$$

then we get that

$$\sum_{\substack{1 \le l \le n \\ v_l \sim v_j}} a_{(v_l, v_j)} \alpha_l = 0$$

and so by the difference equation which φ satisfies, we get that φ vanishes on G_j. Otherwise,

$$\lim_{\varepsilon \to 0} m_j(E + i\varepsilon) = 0,$$

and, by Theorem 2.6, any solution which satisfies $\varphi_{v_j} = 0$ is subordinate on G_j, since, in that case,

$$\lim_{\varepsilon \to 0} m_j(E + i\varepsilon) = \cot \theta \quad \text{for } \theta = \frac{\pi}{2}.$$

Assume now that $\alpha_j \ne 0$. By (3.4), this implies that the limit

$$\lim_{\varepsilon \to 0} \frac{1}{m_j(E + i\varepsilon)}$$

exists and is real, and so there exists $\theta \in [0, \pi)$ such that

$$\lim_{\varepsilon \to 0} m(E + i\varepsilon) = \cot(\theta).$$

By Theorem 2.6, this implies that if there exists $\lambda \ne 0$ such that

$$\left(\sum_{\substack{1 \le l \le n \\ v_l \sim v_j}} a_{(v_l, v_j)} \alpha_l, \alpha_j \right) = \lambda(-\sin(\theta), \cos(\theta)) \tag{3.5}$$

then φ is subordinate on G_j. Note that $\alpha_j \neq 0$ implies that $\theta \neq \frac{\pi}{2}$. By (3.4), we get that

$$\sum_{\substack{1 \leq l \leq n \\ v_l \sim v_j}} a_{(v_l, v_j)} \alpha_l = -\frac{\sin(\theta)}{\cos(\theta)} \alpha_j$$

and so

$$\Big(\sum_{\substack{1 \leq l \leq n \\ v_l \sim v_j}} a_{(v_l, v_j)} \alpha_l, \alpha_j \Big) = \Big(-\frac{\sin(\theta)}{\cos(\theta)} \alpha_j, \alpha_j \Big).$$

Multiplying by $\frac{\cos(\theta)}{\alpha_j}$, we get (3.5) as required. ∎

Remark 3.6. Note that the choice of the k'th column, i.e., setting

$$\alpha_j = \lim_{\varepsilon \to 0} \frac{M_{jk}(E + i\varepsilon)}{M_{kk}(E + i\varepsilon)}$$

is only done so that the resulting solution will be non-trivial. Namely, for any $1 \leq l \leq n$, if there exists some $1 \leq p \leq n$ such that

$$\lim_{\varepsilon \to 0} \frac{M_{pl}(E + i\varepsilon)}{M_{kk}(E + i\varepsilon)} \neq 0,$$

then setting $\alpha_j = \lim_{\varepsilon \to 0} \frac{M_{jl}(E+i\varepsilon)}{M_{kk}(E+i\varepsilon)}$ for any $1 \leq j \leq n$ will also generate a subordinate solution.

4. Proof of Theorem 1.5

Let H be a self-adjoint operator acting on a Hilbert space $\mathcal{H}$ and let P be the projection-valued measure associated with H by the spectral theorem. It is well known (see, e.g., [3, Chapter 7]), that there exist a collection of Hilbert spaces $\{\mathcal{H}_E : E \in \mathbb{R}\}$ so that H is unitarily equivalent to multiplication by the free variable E on the space

$$\widetilde{\mathcal{H}} := \int_{\mathbb{R}}^{\oplus} \mathcal{H}_E \, d\mu(E)$$

whenever μ is a Borel measure on $\sigma(H)$ for which $\mu(A) = 0 \iff P(A) = 0$ for any Borel set A. The measure μ is not unique, but it is determined (up to unitary equivalence) by its null sets. In particular, if $v_1, \ldots, v_n$ is cyclic for H, then μ can be taken to be the sum $\mu_{v_1} + \cdots + \mu_{v_n}$. The spectral multiplicity function is then given by $N_H(E) = \dim \mathcal{H}_E$ and is defined μ-almost everywhere.

Let J be a Jacobi matrix on a star-like graph G, and let M, Ω, μ be as in the previous section. We will use the following fact.

Proposition 4.1. *For μ_s-almost every $E \in \mathbb{R}$,*

$$N_J(E) = \operatorname{rank} \omega(E), \tag{4.1}$$

where $\omega(E)$ is defined by $(\omega(E))_{ij} = \frac{d\mu_{ij}}{d\mu}(E)$.

Remark 4.2. As already mentioned, $\omega(E)$ is also defined μ_s-almost everywhere.

Proof. Let $U: \ell^2(G) \to \widetilde{\mathcal{H}}$ be the unitary transformation which satisfies $UJ\varphi = \widetilde{J}U\varphi$ for any $\varphi \in \ell^2(G)$, where $\widetilde{J}$ is the operator of multiplication by the free variable. Note that for any $\varphi \in \ell^2(G)$, $U\varphi$ is a vector-valued function $E \to U\varphi(E)$ such that $U\varphi(E) \in \mathcal{H}_E$ for μ-almost every $E \in \mathbb{R}$, and $\int_{\mathbb{R}} \|U\varphi(E)\|^2 d\mu(E) < \infty$. Recall that $\delta_1, \ldots, \delta_n$ is a cyclic set for J, and denote $\psi_j = U\delta_j$ for $1 \le j \le n$. Let μ_{ψ_i, ψ_j} be the spectral measure of ψ_i and ψ_j with respect to $\widetilde{J}$. For any Borel set $A \subseteq \mathbb{R}$, we have

$$\mu_{\psi_i, \psi_j}(A) = \int \mathbb{1}_A d\mu_{\psi_i, \psi_j} = \langle \psi_i, \mathbb{1}_A(\widetilde{J})\psi_j \rangle = \int_A \langle \psi_i(E), \psi_j(E) \rangle_{\mathcal{H}_E} d\mu(E).$$

Thus, if we denote $f_{ij}(E) = \langle \psi_i(E), \psi_j(E) \rangle_{\mathcal{H}_E}$, then

$$d\mu_{ij} = d\mu_{\psi_i, \psi_j} = f_{ij} d\mu.$$

and so we have

$$(\omega(E))_{ij} = \frac{d\mu_{ij}}{d\mu}(E) = \langle \psi_i(E), \psi_j(E) \rangle_{\mathcal{H}_E}.$$

Now, we claim that for μ-almost every $E \in \mathbb{R}$, $\operatorname{sp}\{\psi_1(E), \ldots, \psi_n(E)\} = \mathcal{H}_E$. Assume not. Then there exists $0 \ne \psi \in \int_{\mathbb{R}}^{\oplus}$ and a Borel set $A \subseteq \mathbb{R}$ such that $\mu(A) > 0$ and $\psi(E) \perp \operatorname{sp}\{\psi_1(E), \ldots, \psi_n(E)\}$ for μ-almost every $E \in A$. Thus, for every $k \in \mathbb{N}$ and for every $1 \le i \le n$, we have

$$\begin{aligned}
0 &= \langle \widetilde{J}^k \psi_j, \psi \rangle_1 = \langle \widetilde{J}^k U U^{-1} \psi_j, U U^{-1} \psi \rangle_1 \\
&= \langle U J^k U^{-1} \psi_j, U U^{-1} \psi \rangle_1 = \langle J^k \delta_j, U^{-1} \psi \rangle_2
\end{aligned}$$

where $\langle \cdot, \cdot \rangle_1$ is the inner product in the direct integral, and $\langle \cdot, \cdot \rangle_2$ is the inner product in $\ell^2(G)$. This implies that $U^{-1}\psi = 0$ and so $\psi \equiv 0$ which contradicts our assumption. Finally, (4.1) follows from the fact that if $\operatorname{sp}\{v_1, \ldots, v_k\} = \mathbb{C}^n$, then the rank of the matrix $A_{ij} = \langle v_i, v_j \rangle$ is n. ∎

We are now ready to prove Theorem 1.5.

Proof. As before, let $A_k = \{E \in \mathbb{R}: \frac{d\mu_k}{d\mu}(E) > 0\}$ for $1 \le k \le n$. By Theorem 2.2, for every $1 \le j, l \le n$ and for μ_s-almost every $E \in A_k$, we have

$$(\omega(E))_{lj} = \frac{d\mu_k}{d\mu}(E) \cdot \lim_{\varepsilon \to 0} \frac{M_{lj}(E + i\varepsilon)}{M_{kk}(E + i\varepsilon)}.$$

Since $0 < \frac{d\mu_k}{d\mu}(E) < \infty$ for μ-almost every $E \in A_k$, we get that if rank $\omega(E) = m$, then, by Remark 3.6, there are at least m independent subordinate solutions to (1.2), as required. ∎

In the next section, we will construct an example in which the inequality in Theorem 1.5 is strict. Also, in the case where $G = \mathbb{Z}$, the inequality becomes an equality, and so in that sense, this result is optimal. However, the next proposition shows that in the purely singular continuous part of the spectrum, Theorem 1.5 can be improved.

Proposition 4.3. *Denote by k the number of half-lines emanating from C, i.e., $k = \#\{j : G_j \cong \mathbb{N}\}$ and assume that $k > 1$. Then for μ-almost every $E \in \sigma_{\mathrm{sc}}(J) \setminus \sigma_{\mathrm{pp}}(J)$, $N_J(E) \leq k$.*

Proof. Assume not. Let $E \in \sigma_{\mathrm{sc}}(J) \setminus \sigma_{\mathrm{pp}}(J)$ and let $\varphi_1, \ldots, \varphi_{k+1}$ be linearly independent solutions of $H\varphi = E\varphi$. For every $1 \leq j \leq k$, let $v_j \in G_j$ be a vertex on which every non-zero subordinate solution on G_j does not vanish. Such a vertex exists due to the uniqueness of the subordinate solution on a half-line. For any $1 \leq i \leq k + 1$, define $u_i = (\varphi_i(v_1), \ldots, \varphi_i(v_k))$. Then $\{u_1, \ldots, u_{k+1}\}$ is a subset of $\mathbb{C}^k$ with $k + 1$ vectors, and so it is linearly dependent. Then there exist $\alpha_1, \ldots, \alpha_{k+1} \in \mathbb{C}$ such that

$$\sum_{i=1}^{k+1} \alpha_i u_i = 0.$$

By the uniqueness of subordinate solutions on half-lines, this implies that the solution

$$\varphi = \sum_{i=1}^{k+1} \alpha_i \varphi_i$$

vanishes on each half-line. But $\varphi \neq 0$ since the set $\varphi_1, \ldots, \varphi_{k+1}$ is linearly independent, and so φ is supported on a finite set, and in particular $\varphi \in \ell^2(G)$. This implies that $E \in \sigma_{\mathrm{pp}}(J)$, which contradicts our assumption. ∎

5. Remarks

5.1. The multiplicity of Schrödinger operators on star-graphs

In [21], the multiplicity of Schrödinger operators on star-graphs is studied in the continuous setting. The graph Γ is given by the gluing of a finite number of half-lines, and a Schrödinger operator H on Γ is given in the following way. On each half-line ℓ, H acts as a Schrödinger operator on ℓ, i.e., there exists $q_\ell : \mathbb{R}_{\geq 0} \to \mathbb{R}$ such that, for every φ in the domain of H,

$$H(\varphi|_\ell) = -(\varphi|_\ell)'' + q_\ell \varphi|_\ell. \tag{5.1}$$

The domain of H consists of functions on Γ which satisfy some natural properties on each half-line in order for (5.1) to make sense, along with a boundary condition at the gluing point. We give a general description of the analysis done in [21]. Assume that Γ consists of the gluing of $\ell_1, \ldots, \ell_n$. For any $1 \leq i \leq n$, define H_i as the Schrödinger operator on ℓ_i which acts by $H\varphi = -\varphi'' + q_{\ell_i}\varphi$, along with a Dirichlet boundary condition at the origin. Denote by μ_i a scalar spectral measure for H_i, and by $(\mu_i)_s$ the singular part of μ_i with respect to the Lebesgue measure. Denote also by P the projection-valued spectral measure of H, and by P_s its singular part with respect to the Lebesgue measure. The main result of [21] essentially says the following. A support for P_s can be given by $S = S_1 \sqcup S_2$, where S_1 consists of energies for which at least two of the singular parts $(\mu_1)_s, \ldots, (\mu_n)_s$ overlap, and $P|_{S_2}$ is mutually singular with respect to each of $(\mu_1)_s, \ldots, (\mu_n)_s$. In addition, the multiplicity on S_1 is equal to the number of overlaps, and the multiplicity on S_2 is 1.

We now describe the discrete analogue of this result, which can be obtained with our approach. In this subsection we restrict our attention to discrete star-graphs (see Figure 1 (b)). Let G be a discrete star-graph, and let J be a Jacobi matrix on G. Assume that G consists of n half-lines. If $G \neq \mathbb{Z}$ (i.e. if $n \neq 2$), then we denote by o the only vertex of G which satisfies $\deg(o) = n$ (for example, the vertex o in Figure 1 (b) is the only vertex in G with degree 3). If $G = \mathbb{Z}$, then we denote $o{=}0$. Denote the neighbors of o by $v_1, \ldots, v_n$. For every $1 \leq k \leq n$, denote by ℓ_k the connected component of v_k in the graph $G \setminus (o, v_k)$, and by J_k the operator $P_{\ell_k} J P_{\ell_k}$. J_k is simply J restricted to ℓ_k along with a Dirichlet boundary condition at the origin. Denote by ρ_k the spectral measure of δ_{v_k} with respect to J_k, and by A_k the support of $(\rho_k)_s$ given by Theorem 2.6. Finally, let μ be the scalar spectral measure of J as defined in Section 2.3, and by S the support of μ_s given by Theorem 1.4.

Remark 5.1. Note that in order to define μ, one needs to choose a compact component for G. Here, we choose C to be the induced graph on $V_C = \{o, v_1, \ldots, v_n\}$. Also note that under this choice of C, the graphs $G_1, \ldots, G_n$ defined in Section 2.3 are not the same as the graphs $\ell_1, \ldots, \ell_n$ defined above, since $C \neq \{v_1, \ldots, v_n\}$.

Theorem 5.2. *Denote by S_1 the set of real numbers for which there exist at least two indices $1 \leq i < j \leq n$ for which $A_i \cap A_j \neq \emptyset$, and denote $S_2 = (S_1)^c$. Then*

(1) $S_1 \subseteq S$ and $N_J|_{S_1} \leq n - 1$ for μ_s-almost every $E \in \mathbb{R}$.

(2) $N_J|_{S_2 \cap S} = 1$ for μ_s-almost every $E \in \mathbb{R}$.

Proof. (1) Let $E \in S_1$. Without loss of generality, assume that $E \in A_1 \cap A_2$. Let φ_1 be the subordinate solution to $J_1\varphi = E\varphi$ which satisfies $\varphi_1(v_1) = 1$, and let φ_2 be the subordinate solution to $J_2\varphi = E\varphi$ which satisfies $\varphi_2(v_2) = -1$. Define

$$\varphi : G \to \mathbb{C}$$

by

$$\varphi(u) = \begin{cases} \varphi_j(u) & \text{if } u \in \ell_j, \ 1 \le j \le 2, \\ 0 & \text{otherwise.} \end{cases}$$

It is not hard to verify that φ satisfies $J\varphi = E\varphi$, and that it is a subordinate solution, and so $E \in S$. For the second part, note that each subordinate solution ψ must satisfy $\psi(o) = 0$ due to the Dirichlet boundary condition which it satisfies on ℓ_1. Thus, the eigenvalue equation at o is given by

$$0 = E\psi(o) = \sum_{j=1}^{n} \psi(v_j)$$

and so there are at most $n - 1$ subordinate solutions. Now, the result follows from Theorem 1.5.

(2) Let $E \in S \cap S_2$, and let ψ be a non-trivial subordinate solution of $J\varphi = E\varphi$. Assume that ψ does not vanish on ℓ_1. We claim that $\psi(o) \neq 0$. Indeed, if $\psi(o) = 0$, then there must be at least one more index $j \neq 1$ such that ψ does not vanish on ℓ_j. Otherwise, we have

$$0 = E\psi(o) = \sum_{j=1}^{n} \psi(v_j) = \psi(v_1)$$

and it follows that $\psi \equiv 0$. Thus, assume, without loss of generality, that $\psi(v_2) \neq 0$. This implies that $E \in A_1 \cap A_2$, which contradicts our assumption that $E \in S_2$. Now, by uniqueness of the subordinate solution on each half-line, we get that $\psi(o)$ determines $\psi(v_j)$ for any $1 \le j \le n$, and so $\dim S(E) = 1$, which implies that $N_J(E) = 1$ for μ_s-almost every $E \in S_2 \cap S$, as required. ∎

5.2. Strict inequality in Theorem 1.5

Consider a star-like graph G for which C is a triangle graph, and there is a half-line attached to every vertex of C (see Figure 1 (a)). As before, denote the vertices of C by v_1, v_2, v_3 and the half-line with v_i as an origin by G_i. It is well known (see, e.g. [2, Chapter 7]) that there is one-to-one correspondence between probability measures with infinite and bounded support and bounded Jacobi matrices on $\mathbb{N}$, and so in order to construct an example, we start by constructing the appropriate measures. Let μ_1 be a probability measure on $[0, 1]$ such that $\mu_1(\{0\}) > 0$, and let μ_2 be defined by

$$d\mu_2 = \frac{1}{c} f d\lambda,$$

where

$$f(x) = \begin{cases} \dfrac{1}{\sqrt{x}} & 0 < x < 1, \\ 0 & \text{otherwise,} \end{cases}$$

and

$$c = \int\limits_0^1 \frac{dx}{\sqrt{x}}.$$

Clearly, μ_2 is a probability measure on $[0, 1]$ which is absolutely continuous with respect to λ. It is not hard to show that the Borel transform of μ_2 satisfies

$$|m_{\mu_2}(0 + i0)| = \infty.$$

For $\theta = \frac{\pi}{4}$ and for $k = 1, 2$, let $J^{(k)}$ be a Jacobi matrix on $\mathbb{N}$ such that the spectral measure of $J_\theta^{(k)}$ is μ_k. We define a Jacobi matrix J on G in the following way: $J|_{G_1} = J|_{G_3} = J^{(1)}$, $J|_{G_2} = J^{(2)}$, and for every $e \in E_C$, $a_e = 1$. The fact that $\mu_1(\{0\}) > 0$ implies that there exists $0 \neq \varphi_1 \in \ell^2(\mathbb{N})$ such that $J_\theta^{(1)} \varphi_1 = 0$ and $\varphi_1(1) = 1$. By Theorem 2.6, any solution to $J^{(2)} \varphi = 0$ which satisfies $\varphi(0) = -\varphi(1)$ is subordinate. We denote by φ_2 the solution which satis111fies $\varphi_2(1) = -1$. Now, it is can be verified that $\psi_1 : G \to \mathbb{C}$ which is defined by

$$\psi|_{G_1} = \varphi_1, \quad \psi|_{G_2} = -\varphi_1, \quad \psi|_{G_3} = 0$$

satisfies $J\psi = 0$. In addition, any other non-trivial eigenvector must be a multiple of ψ, and so 0 is a simple eigenvalue. On the other hand, the solution $\tilde{\psi}$ which is defined by

$$\tilde{\psi}|_{G_1} = \varphi_1, \quad \tilde{\psi}|_{G_2} = 0, \quad \tilde{\psi}|_{G_3} = \varphi_2$$

is also a subordinate solution which is linearly independent of ψ, and so $N_J(0) = 1 < 2 = \dim S(0)$.

References

[1] M. Aizenman and S. Warzel, *Random operators*. Grad. Stud. Math. 168, American Mathematical Society, Providence, RI, 2015 Zbl 1333.82001 MR 3364516

[2] J. M. Berezanskiĭ, *Expansions in eigenfunctions of selfadjoint operators*. Transl. Math. Monogr. 17, American Mathematical Society, Providence, RI, 1968 Zbl 0157.16601 MR 0222718

[3] M. S. Birman and M. Z. Solomjak, *Spectral theory of selfadjoint operators in Hilbert space*. Math. Appl., Sov. Ser., D. Reidel Publishing Co., Dordrecht, 1987 Zbl 0744.47017 MR 1192782

[4] D. Buschmann, A proof of the Ishii–Pastur theorem by the method of subordinacy. *Univ. Iagel. Acta Math.* (1997), no. 34, 29–34 Zbl 0944.34071 MR 1458029

[5] D. Damanik, R. Killip, and D. Lenz, Uniform spectral properties of one-dimensional quasicrystals. III. α-continuity. *Comm. Math. Phys.* **212** (2000), no. 1, 191–204 Zbl 1045.81024 MR 1764367

[6] F. Gesztesy and E. Tsekanovskii, On matrix-valued Herglotz functions. *Math. Nachr.* **218** (2000), 61–138 Zbl 0961.30027 MR 1784638

[7] D. J. Gilbert, On subordinacy and analysis of the spectrum of Schrödinger operators with two singular endpoints. *Proc. Roy. Soc. Edinburgh Sect. A* **112** (1989), no. 3-4, 213–229 Zbl 0678.34024 MR 1014651

[8] D. J. Gilbert, On subordinacy and spectral multiplicity for a class of singular differential operators. *Proc. Roy. Soc. Edinburgh Sect. A* **128** (1998), no. 3, 549–584 Zbl 0909.34022 MR 1632827

[9] D. J. Gilbert and D. B. Pearson, On subordinacy and analysis of the spectrum of one-dimensional Schrödinger operators. *J. Math. Anal. Appl.* **128** (1987), no. 1, 30–56 Zbl 0666.34023 MR 915965

[10] V. Jakšić and Y. Last, A new proof of Poltoratskii's theorem. *J. Funct. Anal.* **215** (2004), no. 1, 103–110 Zbl 1070.47012 MR 2085111

[11] S. Jitomirskaya and Y. Last, Power-law subordinacy and singular spectra. I. Half-line operators. *Acta Math.* **183** (1999), no. 2, 171–189 Zbl 0991.81021 MR 1738043

[12] S. Y. Jitomirskaya and Y. Last, Power law subordinacy and singular spectra. II. Line operators. *Comm. Math. Phys.* **211** (2000), no. 3, 643–658 Zbl 1053.81031 MR 1773812

[13] I. S. Kac, Spectral multiplicity of a second-order differential operator and expansion in eigenfunction (in Russian). *Izv. Akad. Nauk SSSR Ser. Mat.* **27** (1963), 1081–1112 Zbl 0238.34033 MR 0159982

[14] S. Khan and D. B. Pearson, Subordinacy and spectral theory for infinite matrices. *Helv. Phys. Acta* **65** (1992), no. 4, 505–527 MR 1179528

[15] A. Kiselev, Y. Last, and B. Simon, Modified Prüfer and EFGP transforms and the spectral analysis of one-dimensional Schrödinger operators. *Comm. Math. Phys.* **194** (1998), no. 1, 1–45 Zbl 0912.34074 MR 1628290

[16] A. G. Poltoratskiĭ, Boundary behavior of pseudocontinuable functions. *Algebra i Analiz* **5** (1993), no. 2, 189–210, in Russian; English translation, *St. Petersburg Math. J.* **5** (1994), no. 2, 389–406. MR 1223178

[17] C. Remling, Embedded singular continuous spectrum for one-dimensional Schrödinger operators. *Trans. Amer. Math. Soc.* **351** (1999), no. 6, 2479–2497 Zbl 0918.34074 MR 1665336

[18] B. Simon, Spectral analysis of rank one perturbations and applications. In *Mathematical quantum theory. II. Schrödinger operators (Vancouver, BC, 1993)*, pp. 109–149, CRM Proc. Lecture Notes 8, American Mathematical Society, Providence, RI, 1995 Zbl 0824.47019 MR 1332038

[19] B. Simon, On a theorem of Kac and Gilbert. *J. Funct. Anal.* **223** (2005), no. 1, 109–115 Zbl 1071.47027 MR 2139882

[20] B. Simon, *Szegő's theorem and its descendants*. M. B. Porter Lectures, Princeton University Press, Princeton, NJ, 2011 Zbl 1230.33001 MR 2743058

[21] S. Simonov and H. Woracek, Spectral multiplicity of selfadjoint Schrödinger operators on star-graphs with standard interface conditions. *Integral Equations Operator Theory* **78** (2014), no. 4, 523–575 Zbl 1312.47026 MR 3180877

[22] A. Zlatoš, Sparse potentials with fractional Hausdorff dimension. *J. Funct. Anal.* **207** (2004), no. 1, 216–252 Zbl 1038.47026 MR 2027640

Netanel Levi
Einstein Institute of Mathematics, Edmond J. Safra Campus,
The Hebrew University in Jerusalem, 91904 Jerusalem, Israel; netanel.levi@mail.huji.ac.il