# EQUILIBRIUM MEASURES AND CAPACITIES IN SPECTRAL THEORY 

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#### Abstract

This is a comprehensive review of the uses of potential theory in studying the spectral theory of orthogonal polynomials. Much of the article focuses on the Stahl-Totik theory of regular measures, especially the case of OPRL and OPUC. Links are made to the study of ergodic Schrödinger operators where one of our new results implies that, in complete generality, the spectral measure is supported on a set of zero Hausdorff dimension (indeed, of capacity zero) in the region of strictly positive Lyapunov exponent. There are many examples and some new conjectures and indications of new research directions. Included are appendices on potential theory and on Fekete-Szegő theory.


## 1. Introduction

This paper deals with applications of potential theory to spectral and inverse spectral theory, mainly to orthogonal polynomials especially on the real line (OPRL) and unit circle (OPUC). This is an area that has traditionally impacted both the orthogonal polynomial community and the spectral theory community with insufficient interrelation. The OP approach emphasizes the procedure of going from measure to recursion parameters, that is, the inverse spectral problem, while spectral theorists tend to start with recursion parameters and so work with direct spectral theory.

Potential theory ideas in the orthogonal polynomial community go back at least to a deep 1919 paper of Faber [35] and a seminal 1924 paper of Szegő [107] with critical later developments of Kalmár [63] and

[^0]Erdös-Turán [34]. The modern theory was initiated by Ullman [114] (see also $[115,116,117,118,112,119,123]$ and earlier related work of Korovkin [68] and Widom [122]), followed by an often overlooked paper of Van Assche [120], and culminating in the comprehensive and deep monograph of Stahl-Totik [105]. (We are ignoring the important developments connected to variable weights and external potentials, which are marginal to the themes we study; see [91] and references therein.)

On the spectral theory community side, theoretical physicists rediscovered Szegö's potential theory connection of the growth of polynomials and the density of zeros - this is called the Thouless formula after [110], although discovered slightly earlier by Herbert and Jones [51]. The new elements involve ergodic classes of potentials, especially Kotani theory (see [69, 96, 70, 71, 30, 27]).

One purpose of this paper is to make propaganda on both sides: to explain some of the main aspects of the Stahl-Totik results to spectral theorists and the relevant parts of Kotani theory to the OP community. But this article looks forward even more than it looks back. In thinking through the issues, I realized there were many interesting questions to examine. Motivated in part by the remark that one can often learn from wrong conjectures [56], I make several conjectures which, depending on your point of view, can be regarded as either bold or foolhardy (I especially have Conjectures 8.7 and 8.11 in mind).

The potential of a measure $\mu$ on $\mathbb{C}$ is defined by

$$
\begin{equation*}
\Phi_{\mu}(x)=\int \log |x-y|^{-1} d \mu(y) \tag{1.1}
\end{equation*}
$$

which, for each $x$, is well defined (although perhaps $\infty$ ) if $\mu$ has compact support. The relevance of this to polynomials comes from noting that if $P_{n}$ is a monic polynomial,

$$
\begin{equation*}
P_{n}(z)=\prod_{j=1}^{n}\left(z-z_{j}\right) \tag{1.2}
\end{equation*}
$$

and $d \nu_{n}$ its zero counting measure, that is,

$$
\begin{equation*}
\nu_{n}=\frac{1}{n} \sum_{j=1}^{n} \delta_{z_{j}} \tag{1.3}
\end{equation*}
$$

the point measure with $n \nu_{n}(\{w\})=$ multiplicity of $w$ as a root of $P_{n}$, then

$$
\begin{equation*}
\left|P_{n}(z)\right|^{1 / n}=\exp \left(-\Phi_{\nu_{n}}(z)\right) \tag{1.4}
\end{equation*}
$$

If now $d \mu$ is a measure of compact support on $\mathbb{C}$, let $X_{n}(z)$ and $x_{n}(z)$ be the monic orthogonal and orthonormal polynomials for $d \mu$, that is,

$$
\begin{equation*}
X_{n}(z)=z^{n}+\text { lower order } \tag{1.5}
\end{equation*}
$$

with

$$
\begin{equation*}
\int \overline{X_{n}(z)} X_{m}(z) d \mu(z)=\left\|X_{n}\right\|_{L^{2}}^{2} \delta_{n m} \tag{1.6}
\end{equation*}
$$

and

$$
\begin{equation*}
x_{n}(z)=\frac{X_{n}(z)}{\left\|X_{n}\right\|_{L^{2}}} \tag{1.7}
\end{equation*}
$$

Here and elsewhere $\|\cdot\|$ without a subscript means the $L^{2}$ norm for the measure currently under consideration.

When $\operatorname{supp}(d \mu) \subset \mathbb{R}$, we use $P_{n}, p_{n}$ and note (see $\left.[108,39]\right)$ there are Jacobi parameters $\left\{a_{n}, b_{n}\right\}_{n=1}^{\infty} \in[(0, \infty) \times \mathbb{R}]^{\infty}$, so

$$
\begin{gather*}
x p_{n}(x)=a_{n+1} p_{n+1}(x)+b_{n+1} p_{n}(x)+a_{n} p_{n-1}(x)  \tag{1.8}\\
\left\|P_{n}\right\|=a_{1} \ldots a_{n} \mu(\mathbb{R})^{1 / 2}
\end{gather*}
$$

and if $\operatorname{supp}(d \mu) \subset \partial \mathbb{D}$, the unit circle, we use $\Phi_{n}, \varphi_{n}$ and note (see [108, 44, 98, 99]) there are Verblunsky coefficients $\left\{\alpha_{n}\right\}_{n=0}^{\infty} \in \mathbb{D}^{\infty}$, so

$$
\begin{gather*}
\Phi_{n+1}(z)=z \Phi_{n}-\bar{\alpha}_{n} \Phi_{n}^{*}(z)  \tag{1.9}\\
\left\|\Phi_{n}(z)\right\|=\rho_{0} \ldots \rho_{n-1} \mu(\partial \mathbb{D})^{1 / 2} \tag{1.10}
\end{gather*}
$$

where

$$
\begin{equation*}
\Phi_{n}^{*}(z)=z^{n} \overline{\Phi_{n}(1 / \bar{z})} \quad \rho_{j}=\left(1-\left|\alpha_{j}\right|^{2}\right)^{1 / 2} \tag{1.11}
\end{equation*}
$$

As usual, we will use $J$ for the Jacobi matrix formed from the parameters in the OPRL case, that is, $J$ is tridiagonal with $b_{j}$ on diagonal and $a_{j}$ off-diagonal.

The $X_{n}$ minimize $L^{2}$ norms, that is,

$$
\begin{equation*}
\left\|X_{n}\right\|_{L^{2}(d \mu)}=\min \left\{\left\|Q_{n}\right\|_{L^{2}} \mid Q_{n}(z)=z^{n}+\text { lower order }\right\} \tag{1.12}
\end{equation*}
$$

Given a compact $E \subset \mathbb{C}$, the Chebyshev polynomials are defined by ( $L^{\infty}$ is the sup norm over $E$ )

$$
\begin{equation*}
\left\|T_{n}\right\|_{L^{\infty}(E)}=\min \left\{\left\|Q_{n}\right\|_{L^{\infty}} \mid Q_{n}(z)=z^{n}+\text { lower order }\right\} \tag{1.13}
\end{equation*}
$$

These minimum conditions suggest that extreme objects in potential theory, namely, the capacity, $C(E)$, and equilibrium measure, $d \rho_{E}$, discussed in [50, 73, 81, 88, 91] and Appendix A will play a role (terminology from Appendix A is used extensively below). In fact, going back to Szegő [107] (we sketch a proof in Appendix B), one knows

Theorem 1.1 (Szegő [107]). For any compact $E \subset \mathbb{C}$ with Chebyshev polynomials, $T_{n}$, one has

$$
\begin{equation*}
\lim _{n \rightarrow \infty}\left\|T_{n}\right\|_{L^{\infty}(E)}^{1 / n}=C(E) \tag{1.14}
\end{equation*}
$$

This has an immediate corollary (it appears with this argument in Widom [122] and may well be earlier):

Corollary 1.2. Let $\mu$ be a measure of compact support, $E$, in $\mathbb{C}$. Let $X_{n}(z ; d \mu)$ be its monic OPs. Then

$$
\begin{equation*}
\limsup _{n \rightarrow \infty}\left\|X_{n}\right\|_{L^{2}(\mathbb{C}, d \mu)}^{1 / n} \leq C(E) \tag{1.15}
\end{equation*}
$$

Proof. By (1.12),

$$
\begin{equation*}
\left\|X_{n}\right\|_{L^{2}(\mathbb{C}, d \mu)}^{1 / n} \leq\left\|T_{n}\right\|_{L^{2}(\mathbb{C}, d \mu)}^{1 / n} \tag{1.16}
\end{equation*}
$$

where $T_{n}$ are the Chebyshev polynomials for $E$. On $E,\left|T_{n}(z)\right| \leq$ $\left\|T_{n}\right\|_{L^{\infty}(E)}$ so, $\operatorname{since} \operatorname{supp}(d \mu)=E$,

$$
\begin{equation*}
\left\|X_{n}\right\|_{L^{2}(\mathbb{C}, d \mu)}^{1 / n} \leq\left\|T_{n}\right\|_{L^{\infty}(E)}^{1 / n} \mu(E)^{1 / 2 n} \tag{1.17}
\end{equation*}
$$

and (1.15) follows from (1.14).
For OPRL and OPUC, the relation (1.15) says

$$
\begin{align*}
\limsup \left(a_{1} \ldots a_{n}\right)^{1 / n} \leq C(E) & (\mathrm{OPRL})  \tag{1.18}\\
\limsup \left(\rho_{1} \ldots \rho_{n}\right)^{1 / n} \leq C(E) & (\mathrm{OPUC}) \tag{1.19}
\end{align*}
$$

(1.18) is a kind of thickness indication of the spectrum of discrete Schrödinger operators (with $a_{j} \equiv 1$ ) where it is not widely appreciated that $C(E) \geq 1$.

In many cases that occur in spectral theory, one considers discrete and essential spectrum. In this context, $\sigma_{\text {ess }}(d \mu)$ is the nonisolated points of $\operatorname{supp}(d \mu) . \sigma_{\mathrm{d}}(d \mu)=\operatorname{supp}(d \mu) \backslash \sigma_{\text {ess }}(d \mu)$ is a countable discrete set. If $d \nu$ is any measure with finite Coulomb energy $\nu\left(\sigma_{\mathrm{d}}(d \mu)\right)=0$, thus $C(\operatorname{supp}(d \mu))=C\left(\sigma_{\text {ess }}(d \mu)\right)$; so we will often consider $E=\sigma_{\text {ess }}(d \mu)$ in (1.15). In fact, as discussed in Appendix A after Theorem A.13, we should take $E=\sigma_{\text {cap }}(d \mu)$.

The inequality (1.16) suggests singling out a special case. A measure $d \mu$ of compact support, $E$, on $\mathbb{C}$ is called regular if and only if

$$
\begin{equation*}
\lim _{n \rightarrow \infty}\left\|X_{n}\right\|_{L^{2}(d \mu)}^{1 / n}=C(E) \tag{1.20}
\end{equation*}
$$

For $E=[-1,1]$, this class was singled out by Ullman [114]; the general case is due to Stahl-Totik [105].

Example 1.3. The Nevai class, $N(a, b)$, with $a>0, b \in \mathbb{R}$, is the set of probability measures on $\mathbb{R}$ whose Jacobi parameters obey

$$
\begin{equation*}
a_{n} \rightarrow a \quad b_{n} \rightarrow b \tag{1.21}
\end{equation*}
$$

The Jacobi matrix with $a_{n} \equiv a, b_{n} \equiv b$ is easily seen to have spectrum

$$
\begin{equation*}
E(a, b)=[b-2 a, b+2 a] \tag{1.22}
\end{equation*}
$$

so, by (A.8),

$$
\begin{equation*}
C(E(a, b))=a \tag{1.23}
\end{equation*}
$$

By Weyl's theorem, if $\mu \in N(a, b), \sigma_{\text {ess }}(\mu)=E(a, b)$, so $\mu$ is regular.
Example 1.4. Here is an example of a regular measure on $\mathbb{R}$ not in a Nevai class. Let $d \mu$ be the measure with Jacobi parameters $b_{n} \equiv 0$ and

$$
a_{n}= \begin{cases}1 & n \neq k^{2} \text { for all } k  \tag{1.24}\\ \frac{1}{2} & n=k^{2} \text { for some } k\end{cases}
$$

Clearly, $\lim \left(a_{1} \ldots a_{n}\right)^{1 / n}=1$, so if $\operatorname{supp}(d \mu)=[-2,2]$ (which has capacity 1 by (1.23)), we will have a regular measure not in a Nevai class. Since $\left(\begin{array}{ll}0 & c \\ c & 0\end{array}\right) \geq\left(\begin{array}{cc}-c & 0 \\ 0 & -c\end{array}\right)$ for any $c>0$, the Jacobi matrix $J$ associated to $d \mu$ is bounded below by a diagonal matrix with elements either $-\frac{1}{2},-1,-\frac{3}{2}$, or -2 (for $n=1,2, k^{2}$ or $k^{2}+1$ and otherwise), so $J \geq-2$. Similarly, $J \leq 2$. Thus, $\sigma(J)=\operatorname{supp}(d \mu) \subset[-2,2]$. On the other hand, since $\lim \left(a_{1} \ldots a_{n}\right)^{1 / n}=1$, (1.18) implies $C(E) \geq 1$ where $E=\operatorname{supp}(d \mu)$. If $E \varsubsetneqq[-2,2]$, it is missing an open subset and so $|E|<4$ and $C(E)<1$ (by (A.57)). Thus $C(E) \geq 1$ implies $E=[-2,2]$. Alternatively, using plane wave trial functions cut off to live in $\left[k^{2}+2,(k+1)^{2}-1\right]$, we easily see directly that $[-2,2] \subset \sigma(J)$.

This example has no a.c. spectrum by results of Remling [90]. In Section 8 (see Example 8.12), we have models which are regular, not in Nevai class with nonempty a.c. spectrum.

Example 1.5. The CN (for Cesàro-Nevai) class was introduced by Golinskii-Khrushchev [46] for OPUC and it has an OPRL analog. For OPRL, it says

$$
\begin{equation*}
\frac{1}{n}\left[\sum_{j=1}^{n}\left|a_{j}-1\right|+\left|b_{j}\right|\right] \rightarrow 0 \tag{1.25}
\end{equation*}
$$

Example 1.4 is in this class and is regular, but it is not true that every element of CN for OPRL is regular; for example, if $a_{n} \equiv 1$ and each $b_{j}$ is 0 or 1 so (1.25) holds but with arbitrarily long strings of only 0 's and also of only 1's, e.g., $b_{j}=1$ if $n^{2} \leq j \leq n^{2}+n$ and $b_{j}=0$ otherwise. Then $\sigma(J)=[-2,3]$ with $C(\sigma(J))=\frac{5}{4}$, but $\left(a_{1} \ldots a_{n}\right)^{1 / n} \rightarrow 1$. (To see that $\sigma(J)=[-2,3]$, let $J_{+}, J_{-}$be the matrices with $a_{n} \equiv 1$ and
$b_{n} \equiv 1$ (for $J_{+}$) and $b_{n} \equiv 0$ (for $J_{-}$). Then $-2 \leq J_{-} \leq J \leq J_{+} \leq 3$, so $\sigma(J) \subset[-2,3]$. On the other hand, because of the long strings, a variational argument shows $\sigma\left(J_{+}\right) \cup \sigma\left(J_{-}\right) \subset \sigma(J)$.) However, for OPUC, the subset of CN class with $\sup _{n}\left|\alpha_{n}\right|<1$ consists of only regular measures. For by Theorem 4.3.17 of [98], $\sigma(d \mu)=\partial \mathbb{D}$ with capacity 1 . On the other hand, if $A=\sup _{n}\left|\alpha_{n}\right|$ and

$$
L(A)=-\frac{\log (1-A)}{2 A}
$$

then (since $\log |1-x|$ is convex)

$$
-\log \rho_{j} \leq L(A)\left|\alpha_{j}\right|^{2}
$$

so

$$
\frac{1}{n} \sum_{j=0}^{n-1} \log \rho_{j} \leq L(A) \frac{1}{n} \sum_{j=0}^{n-1}\left|\alpha_{j}\right|^{2} \leq L(A) \frac{1}{n} \sum_{j=0}^{n-1}\left|\alpha_{j}\right|
$$

goes to zero so $\left(\rho_{0} \ldots \rho_{n-1}\right)^{1 / n} \rightarrow 1$. It is easy to see that if $\sup _{n}\left|\alpha_{n}\right|<1$ is dropped, regularity can be lost.

Example 1.6. Let $a_{n} \equiv \frac{1}{2}$ and let $b_{n}= \pm 1$ chosen as identically distributed random variables. As above, all these random J's have $-2 \leq J \leq 2$, so $\operatorname{supp}(d \mu) \subset[-2,2]$. Since there will be, with probability 1 , long stretches of $b_{n} \equiv 1$ or $b_{n} \equiv-1$, it is easy to see $\operatorname{supp}(d \mu) \supset([-1,1]+1) \cup([-1,1]-1)=[-2,2]$. Thus, a typical random $d \mu$ has support $[-2,2]$ with capacity 1 , but obviously $\lim \left(a_{1} \ldots a_{n}\right)^{1 / n}=\frac{1}{2}$. This shows there are measures which are not regular. By Example 1.3, random slowly decaying Jacobi matrices are regular, so neither randomness nor pure point measures necessarily destroy regularity. We return to pure point measures in Theorem 5.5 and Corollary 5.6.

Section 8 has many more examples of regular measures. Regularity is important because of its connections to zero distributions and to root asymptotics. Let $d \nu_{n}$ denote the zero distribution for $X_{n}(z)$ defined by (1.3). Then

Theorem 1.7. Let $d \mu$ be a measure on $\mathbb{R}$ with $\sigma_{\text {ess }}(d \mu)=E$ compact and $C(E)>0$. If $\mu$ is regular, then $d \nu_{n}$ converges weakly to $d \rho_{E}$, the equilibrium measure for $E$.

Remarks. 1. For $E=[0,1]$, ideas close to this occur in Erdös-Turán [34]. The full result is in Stahl-Totik [105] who prove a stronger result. Rather than $E \subset \mathbb{R}$, they need that the unbounded component, $\Omega$, of $\mathbb{C} \backslash E$ is dense in $\mathbb{C}$. We will prove Theorem 1.7 in Section 2.
2. The result is false for measures on $\partial \mathbb{D}$. Indeed, it fails for $d \mu=$ $\frac{d \theta}{2 \pi}$, that is, $\alpha_{n} \equiv 0$ and $\Phi_{n}(z)=z^{n}$. However, there is a result for paraorthogonal polynomials and for the balayage of $d \nu_{n}$. Theorem 1.7 is true if $\operatorname{supp}(d \mu) \subset \partial \mathbb{D}$ but is not all of $\partial \mathbb{D}$. We will discuss this further in Section 3 where we also prove a version of Theorem 1.7 for OPUC.

For later purposes, we note
Proposition 1.8. Let $d \mu$ be a measure on $\mathbb{R}$ with $\sigma_{\mathrm{ess}}(d \mu)=E$ compact. Then any limit of $d \nu_{n}$ is supported on $E$.

Proof. It is known (see [98, Sect. 1.2]) that if $(a, b) \cap \operatorname{supp}(d \mu)=\emptyset$, then $P_{n}(x)$ has at most one zero in $(a, b)$. It follows that if $e$ is an isolated point of $\operatorname{supp}(d \mu)$, then $(e-\delta, e+\delta)$ has at most three zeros for $\delta$ small (with more argument, one can get two). Thus, points not in $E$ have neighborhoods, $N$, with $\nu_{n}(N) \leq \frac{3}{n}$.

Stahl-Totik [105] also have the following almost converse (their Sect. 2.2)-for $E \subset \mathbb{R}$, we prove a slightly stronger result; see Theorem 2.5.

Theorem 1.9. Let $d \mu$ be a measure on $\mathbb{R}$ with $\sigma_{\text {ess }}(d \mu)=E$ compact and $C(E)>0$. Suppose that $d \nu_{n} \rightarrow d \rho_{E}$, the equilibrium measure. Then either $d \mu$ is regular or there exists a Borel set, $X$, with $d \mu(\mathbb{R} \backslash$ $X)=0$ and $C(X)=0$.

Remarks. 1. As an example where such an $X$ exists even though $C(\operatorname{supp}(\mu))>0$, consider a $\mu$ which is dense pure point on $[-2,2]$.
2. In Section 8, we will see explicit examples where $d \nu_{n} \rightarrow d \rho_{E}$ but $d \mu$ is not regular.

The other connection is to root asymptotics of the OPs. Recall the Green's function, $G_{E}(z)$, is defined by (A.40); it vanishes q.e. (quasieverywhere, defined in Appendix A) on $E$, is harmonic on $\Omega$, and asymptotic to $\log |z|-\log C(E)+o(1)$ as $|z| \rightarrow \infty$. The main theorem on root asymptotics is:

Theorem 1.10. Let $E \subset \mathbb{C}$ be compact and let $\mu$ be a measure of compact support with $\sigma_{\mathrm{ess}}(\mu)=E$. Then the following are equivalent:
(i) $\mu$ is regular, that is, $\lim _{n \rightarrow \infty}\left\|X_{n}\right\|_{L^{2}(d \mu)}^{1 / n}=C(E)$.
(ii) For all $z$ in $\mathbb{C}$, uniformly on compacts,

$$
\begin{equation*}
\limsup \left|x_{n}(z)\right|^{1 / n} \leq e^{G_{E}(z)} \tag{1.26}
\end{equation*}
$$

(iii) For q.e. $z$ in $\partial \Omega$ (with $\Omega$ the unbounded connected component of $\mathbb{C} \backslash E)$, we have

$$
\begin{equation*}
\lim \sup \left|x_{n}(z)\right|^{1 / n} \leq 1 \tag{1.27}
\end{equation*}
$$

Moreover, if (i)-(iii) hold, then (here cvh $=$ closed convex hull)
(iv) For every $z \in \mathbb{C} \backslash \operatorname{cvh}(\operatorname{supp}(d \mu))$, we have

$$
\begin{equation*}
\lim _{n \rightarrow \infty}\left|x_{n}(z)\right|^{1 / n}=e^{G_{E}(z)} \tag{1.28}
\end{equation*}
$$

(v) For q.e. $z \in \partial \Omega$,

$$
\begin{equation*}
\limsup _{n \rightarrow \infty}\left|x_{n}(z)\right|^{1 / n}=1 \tag{1.29}
\end{equation*}
$$

(vi) For any sequence $Q_{n}(z)$ of polynomials of degree $n$, for all $z \in \mathbb{C}$,

$$
\begin{equation*}
\lim \sup \left|\frac{Q_{n}(z)}{\left\|Q_{n}\right\|_{L^{2}(d \mu)}}\right|^{1 / n} \leq e^{G_{E}(z)} \tag{1.30}
\end{equation*}
$$

Remark. It is easy to see that (iv), (v) or (vi) are equivalent to (i)(iii).

For $E \subset \mathbb{R}$, we will prove this in Section 2. For $E \subset \partial \mathbb{D}$, we prove it in Section 3.

The original result asserting cases where regularity holds was proven in 1940!

Theorem 1.11 (Erdös-Turán [34]). Let $d \mu$ be supported on [-2, 2] and suppose

$$
\begin{equation*}
d \mu(x)=w(x) d x+d \mu_{\mathrm{s}} \tag{1.31}
\end{equation*}
$$

with $d \mu_{\mathrm{s}}$ singular. Suppose $w(x)>0$ for a.e. $x$ in $[-2,2]$. Then $\mu$ is regular.

Remarks. 1. Erdös-Turán [34] worked on $[-1,1]$ and had $d \mu_{\mathrm{s}}=0$.
2. We now have a stronger result than this-namely, Rakhmanov's theorem (see [99, Ch. 9]). If $w(x)>0$ for a.e. $x$, one knows $a_{n} \rightarrow 1$ (and $b_{n} \rightarrow 0$ ) much more than $\left(a_{1} \ldots a_{n}\right)^{1 / n} \rightarrow 1$ (equivalently, we have ratio asymptotics on the $p$ 's and not just root asymptotics). Regularity is a "poor man's" Rakhmanov's theorem. But unlike Rakhmanov's theorem which is only known for a few other $E$ 's (see [29, 90] and the discussion in Section 8), this weaker version holds very generally.
3. In this case, $d \rho_{E}$ is equivalent to $d x$, so (1.31) and $W(x)>0$ for a.e. $x$ is equivalent to saying that $\rho_{E}$ is $\mu$-a.c.

In Section 4, we will prove the following vast generalization of the Erdös-Turán result:

Theorem 1.12 (Widom [122]). Let $\mu$ be a measure on $\mathbb{R}$ with compact support and $E=\sigma_{\text {ess }}(d \mu)$ and $C(E)>0$. Suppose $d \rho_{E}$ is the equilibrium measure for $E$ and

$$
\begin{equation*}
d \mu=w(x) d \rho_{E}(x)+d \mu_{\mathrm{s}} \tag{1.32}
\end{equation*}
$$

where $d \mu_{\mathrm{s}}$ is $d \rho_{E}$-singular. Suppose $w(x)>0$ for $d \rho_{E^{-}}$a.e. $x$. Then $\mu$ is regular.

Remarks. 1. As above, (1.32) $+w(x)>0$ is equivalent to saying that $d \rho_{E}$ is absolutely continuous with respect to $d \mu$.
2. Widom's result is much more general than what we have in this theorem. His $E$ is a general compact set in $\mathbb{C}$. His polynomials are defined by general families of minimum conditions, for example, $L^{p}$ minimizers. Most importantly, he has a general family of support conditions that, as he notes in a one-sentence remark, include the case where $d \rho_{E}$ is a.c. with respect to $d \mu$. Because of its spectral theory connection, we have focused on the $L^{2}$ minimizers, although it is not hard to accommodate more general ones. We focus on the $w>0$ case because if one goes beyond that, it is better to look at conditions that depend on weights and not just supports of the measure as Widom (and Ullman [114]) do (see Theorem 1.13 below).
3. In Section 4, we will give a proof of this theorem due to Van Assche [120] who mentions Widom's paper but says it is not clear his hypotheses apply despite an explicit (albeit terse) aside in Widom's paper. Stahl-Totik [105] state Theorem 1.12 explicitly. They seem to be unaware of Van Assche's paper or Widom's aside.

It was Geronimus [43] who seems to have first noted that there are non-a.c. measures which are regular (and later Widom [122] and Ullman [114]). Of course, with the discovery of Nevai class measures which are not a.c. $[95,32,111,78,79,80,97,67]$, there are many examples, but given a measure, one would like to know effective criteria. Stahl-Totik [105, Ch. 4] have many, of which we single out:

Theorem 1.13 (Stahl-Totik [105]). Let E be a finite union of disjoint closed intervals in $\mathbb{R}$. Suppose $\mu$ is a measure on $\mathbb{R}$ with $\sigma_{\text {ess }}(d \mu)=E$, and for any $\eta>0(|\cdot|$ is Lebesgue measure $)$,

$$
\begin{equation*}
\lim _{m \rightarrow \infty}\left|\left\{x \left\lvert\, \mu\left(\left[x-\frac{1}{m}, x+\frac{1}{m}\right]\right) \leq e^{-m \eta}\right.\right\}\right|=0 \tag{1.33}
\end{equation*}
$$

Then $\mu$ is regular.
Theorem 1.14 (Stahl-Totik [105]). Let E be a finite union of disjoint closed intervals in $\mathbb{R}$. Suppose $\mu$ is a measure on $\mathbb{R}$ and that $\mu$ is
regular. Then for any $\eta>0$,

$$
\begin{equation*}
\lim _{m \rightarrow \infty} C\left(\left\{x \in E \left\lvert\, \mu\left(\left[x-\frac{1}{m}, x+\frac{1}{m}\right]\right) \leq e^{-m \eta}\right.\right\}\right)=0 \tag{1.34}
\end{equation*}
$$

Remarks. 1. We will prove Theorem 1.13 in Section 5.
2. Stahl-Totik [105] state these results for $E=[-1,1]$, but it is easy to accommodate finite unions of disjoint closed intervals; see Corollary 6.6.

In Section 6, we turn to structural results (all due to Stahl-Totik [105]) connected to inheritance of regularity when measures have a relation, for example, when restrictions of regular measures are regular.

Section 7 discusses relations of potential theory and ergodic Jacobi matrices. This theory concerns OPRL (or OPUC) whose recursion coefficients are samples of an ergodic process-as examples, totally random or almost periodic cases. In that case, various ergodic theorems guarantee the existence of $\lim \left(a_{1} \ldots a_{n}\right)^{1 / n}$, of $d \nu_{\infty} \equiv \lim d \nu_{n}$, and of a natural Lyapunov exponent, $\gamma(z)$, which off of $\operatorname{supp}(d \mu)$ is $\lim \left|p_{n}(z ; d \mu)\right|^{1 / n}$ and is subharmonic on $\mathbb{C}$. In that section, we will prove some of the few new results of this paper:

Theorem 1.15. Let $d \mu_{\omega}$ be the measures associated to an ergodic family of $O P R L, d \nu_{\infty}$ and $\gamma$ its density of states and Lyapunov exponent. Let $E=\operatorname{supp}\left(d \nu_{\infty}\right)$. Then the following are equivalent:
(a) $\quad \gamma(x)=0$ for $d \rho_{E}$-a.e. $x$
(b) $\quad \lim _{n \rightarrow \infty}\left(a_{1}(\omega) \ldots a_{n}(\omega)\right)^{1 / n}=C(E)$
for a.e. $\omega$. Moreover, if (a) and (b) hold, then

$$
\begin{equation*}
d \nu_{\infty}(x)=d \rho_{E}(x) \tag{1.36}
\end{equation*}
$$

with $d \rho_{E}$ the equilibrium measure for $E$. Conversely, if (1.36) holds, either (a) and (b) hold, or else for a.e. $\omega, d \mu_{\omega}$ is supported on a set of capacity zero.
Remarks. 1. We will prove that for a.e. $\omega$,

$$
\begin{equation*}
\lim _{n \rightarrow \infty}\left(a_{1}(\omega) \ldots a_{n}(\omega)\right)^{1 / n}=C(E) \exp \left(-\int \gamma(x) d \rho_{E}(x)\right) \tag{1.37}
\end{equation*}
$$

2. In Section 8, we will see examples where $d \nu_{\infty}=d \rho_{E}$ but (1.35) fails. Of course, (a) also fails.

The following is an ultimate version of what is sometimes called the Pastur-Ishii theorem (see Section 7).

Theorem 1.16. Let $d \mu_{\omega}$ be a family of measures associated to an ergodic family of OPRL and let $\gamma$ be its Lyapunov exponent. Let $S \subset \mathbb{R}$ be the Borel set of $x \in \mathbb{R}$ with $\gamma(x)>0$. Then for a.e. $\omega$, there exists $Q_{\omega}$ of capacity zero so $d \mu_{\omega}\left(S \backslash Q_{\omega}\right)=0$. In particular, $d \mu_{\omega} \upharpoonright S$ is of local Hausdorff dimension zero.

We should explain what is really new in this theorem. It has been known since Pastur [86] and Ishii [53] that for ergodic Schrödinger operators, the spectral measures are supported on the eigenvalues union the bad set where Lyapunov behavior fails (this bad set actually occurs, e.g., $[12,60])$. The classic result is that the bad set has Lebesgue measure zero. The new result here (elementary given a potential theoretic point of view!) is that the bad set has capacity zero.

Section 8 describes examples, open questions, and conjectures. Section 9 has some remarks on the possible extensions of these ideas to continuum Schrödinger operators. Appendix A is a primer of potential theory and Appendix B proves Theorem 1.1 on Chebyshev polynomials.

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## 2. Regular Measures for OPRL

In this section, our main goal is to prove Theorems 1.7 and 1.10 for OPRL. The key will be a series of arguments familiar to spectral theorists as the Thouless formula, albeit in a different (nonergodic) guise. The key will be an analog of positivity of the Lyapunov exponent off the spectrum.

Lemma 2.1. (a) Let $J$ be a bounded Jacobi matrix and let $H$ be the convex hull of the spectrum of $J$. For any $\varphi \in L^{2}(\mathbb{R}, d \mu)$,

$$
\begin{equation*}
|\langle\varphi,(J-z) \varphi\rangle| \geq \operatorname{dist}(z, H)\|\varphi\|^{2} \tag{2.1}
\end{equation*}
$$

(b) The Jacobi parameters $a_{n}$ obey (recall $a_{n}>0$ )

$$
\begin{equation*}
a_{n} \leq \frac{1}{2} \operatorname{diam}(H) \tag{2.2}
\end{equation*}
$$

Proof. In a spectral representation, $J$ is multiplication by $x$. If $d=$ $\operatorname{dist}(z, H)$, there is $\omega \in \partial \mathbb{D}$ with $\operatorname{Re}[(x-z) \omega] \geq d$ for all $x \in H$. Thus $\operatorname{Re}(\langle\varphi,(J-z) \varphi\rangle \omega) \geq d\|\varphi\|^{2}$, which yields (2.1).

Let $D=\frac{1}{2} \operatorname{diam}(H)$ and $c=$ center of $H$ so $H=[c-D, c+D]$. Then

$$
\begin{aligned}
a_{n} & =\int x p_{n} p_{n-1} d \mu \\
& =\int(x-c) p_{n} p_{n-1} d \mu \\
& \leq \sup _{\sigma(J)}|x-c|=D
\end{aligned}
$$

proving (2.2).
The following is related to the proof of Theorem 4.3.15 in [98]:
Proposition 2.2. Let $d \mu$ be a measure on $\mathbb{R}$ of compact support, $p_{n}(x, d \mu)$ the normalized OPRL, and $H$ the convex hull of the support of $d \mu$. For $z \notin H$, let $d(z)=\operatorname{dist}(z, H)$. Let $D=\frac{1}{2} \operatorname{diam}(H)$. Then for such $z$,

$$
\begin{equation*}
\left|p_{n}(z, d \mu)\right|^{2} \geq\left(\frac{d}{D}\right)^{2}\left(1+\left(\frac{d}{D}\right)^{2}\right)^{n-1} \tag{2.3}
\end{equation*}
$$

In particular, $p_{n}(z) \neq 0$ for all $n$ and

$$
\begin{equation*}
\liminf \left|p_{n}(z, d \mu)\right|^{1 / n} \geq\left(1+\left(\frac{d}{D}\right)^{2}\right)^{1 / 2}>1 \tag{2.4}
\end{equation*}
$$

Remark. Of course, it is well known that $p_{n}$ has all its zeros on $H$.
Proof. Let $\varphi_{n}(x)$ be the function

$$
\begin{equation*}
\varphi_{n}(x)=\sum_{j=0}^{n} p_{n}(z) p_{n}(x) \tag{2.5}
\end{equation*}
$$

which has components $\varphi_{n}=\left\langle p_{0}(z), \ldots, p_{n}(z), 0,0, \ldots\right\rangle$ in $p_{n}(x)$ basis. Then, by the recursion relation,

$$
\left[(J-z) \varphi_{n}\right]_{j}= \begin{cases}0 & j \neq n, n+1  \tag{2.6}\\ -a_{n+1} p_{n+1}(z) & j=n \\ a_{n+1} p_{n}(z) & j=n+1\end{cases}
$$

(a version of the CD formula!). Thus,

$$
\begin{equation*}
\left\langle\varphi_{n},(J-z) \varphi_{n}\right\rangle=-a_{n+1} p_{n+1}(z) \overline{p_{n}(z)} \tag{2.7}
\end{equation*}
$$

and (2.1) becomes

$$
\begin{equation*}
a_{n+1}\left|p_{n+1}(z) p_{n}(z)\right| \geq d \sum_{j=0}^{n}\left|p_{j}(z)\right|^{2} \tag{2.8}
\end{equation*}
$$

By (2.2),

$$
\begin{equation*}
\left|p_{n+1}(z) p_{n}(z)\right| \geq \frac{d}{D} \sum_{j=0}^{n}\left|p_{j}(z)\right|^{2} \tag{2.9}
\end{equation*}
$$

Next use $2|x y| \leq \alpha x^{2}+\alpha^{-1} y^{2}$ for any $\alpha$ to see

$$
\begin{equation*}
\left|p_{n+1}(z) p_{n}(z)\right| \leq \frac{1}{2} \frac{d}{D}\left|p_{n}(z)\right|^{2}+\frac{1}{2} \frac{D}{d}\left|p_{n+1}(z)\right|^{2} \tag{2.10}
\end{equation*}
$$

which, with (2.9), implies

$$
\begin{equation*}
\left|p_{n+1}(z)\right|^{2} \geq\left(\frac{d}{D}\right)^{2} \sum_{j=0}^{n}\left|p_{j}(z)\right|^{2} \tag{2.11}
\end{equation*}
$$

This implies that

$$
\begin{equation*}
\sum_{j=0}^{n+1}\left|p_{j}(z)\right|^{2} \geq\left[1+\left(\frac{d}{D}\right)^{2}\right] \sum_{j=0}^{n}\left|p_{j}(z)\right|^{2} \tag{2.12}
\end{equation*}
$$

so, since $p_{0}(z)=1$, we obtain

$$
\begin{equation*}
\sum_{j=0}^{n}\left|p_{j}(z)\right|^{2} \geq\left[1+\left(\frac{d}{D}\right)^{2}\right]^{n} \tag{2.13}
\end{equation*}
$$

(2.13) plus (2.11) imply (2.3), and that implies (2.4).

Remark. (2.4) is also related to Schnol's theorem (see [92, 93] and [98, Lemma 4.3.13]) and to Combes-Thomas estimates [23, 1].

This yields the key estimate, given the following equality:
Theorem 2.3. Let $d \mu$ be a measure of compact support on $\mathbb{R}$ with $H$ the convex hull of $\operatorname{supp}(d \mu)$. Let $n(j)$ be a subsequence (i.e., $n(1)<$ $n(2)<n(3)<\ldots$ in $\{0,1,2, \ldots\})$ so that the zero counting measures $d \nu_{n(j)}$ have a weak limit $d \nu_{\infty}$ and so that $\left(a_{1} \ldots a_{n(j)}\right)^{1 / n(j)}$ has a nonzero limit $A$. Then, for any $z \notin H$,

$$
\begin{equation*}
\lim _{j \rightarrow \infty}\left|p_{n(j)}(z)\right|^{1 / n(j)}=A^{-1} \exp \left(-\Phi_{\nu_{\infty}}(z)\right) \tag{2.14}
\end{equation*}
$$

where $\Phi_{\nu}$ is the potential of $\nu$. In particular,

$$
\begin{equation*}
\exp \left(-\Phi_{\nu_{\infty}}(z)\right)>A \tag{2.15}
\end{equation*}
$$

Proof. (1.4) says that

$$
\begin{equation*}
\left|p_{n(j)}(z)\right|^{1 / n(j)}\left(a_{1} \ldots a_{n(j)}\right)^{1 / n(j)}=\exp \left(-\Phi_{\nu_{n(j)}}(z)\right) \tag{2.16}
\end{equation*}
$$

For $z \notin H, \log |z-y|^{-1}$ is continuous on $H$ so since $\nu_{n}$ and so $\nu_{\infty}$ are supported on $H$ (indeed, $\nu_{\infty}$ is supported on $\left.\sigma_{\text {ess }}(d \mu)\right), \Phi_{\nu_{n(j)}}(z) \rightarrow$ $\Phi_{\nu_{\infty}}(z)$ and (2.16) implies (2.14). By (2.4), LHS of (2.14) $>1$, which implies (2.15).

Note: (2.15) implies that $\Phi_{\nu_{\infty}}$ is bounded above which, by arguments of Craig-Simon [25], implies $\nu_{\infty}((-\infty, E])$ is log-Hölder continuous. This is a new result, although the fact for ergodic Jacobi matrices is due to Craig-Simon [25].

This yields an independent proof of Corollary 1.2 for OPRL, and more:

Theorem 2.4. Under the hypotheses of Theorem 2.3, if $E=\sigma_{\text {ess }}(d \mu)$, then $A \leq C(E)$, and if $A=C(E)$, then $d \nu_{\infty}=d \rho_{E}$, the equilibrium measure for $E$. In particular, if $\mu$ is regular $\left(\right.$ i.e., $\lim \left(a_{1} \ldots a_{n}\right)^{1 / n}=$ $C(E)$ ), then $d \nu_{n} \rightarrow d \rho_{E}$ and (1.28) holds for $z \notin H$.

Remark. Thus, we have proven Corollary 1.2 again, Theorem 1.7, and one part of Theorem 1.10.

Proof. By (2.15) for $z \notin H$,

$$
\begin{equation*}
\Phi_{\nu_{\infty}}(z) \leq \log \left(A^{-1}\right) \tag{2.17}
\end{equation*}
$$

By lower semicontinuity, this also holds on $H$. Integrating $d \nu_{\infty}$ using (A.6), we obtain

$$
\begin{equation*}
\mathcal{E}\left(\nu_{\infty}\right) \leq \log \left(A^{-1}\right) \tag{2.18}
\end{equation*}
$$

Since $\inf _{\nu}(\mathcal{E}(\nu))=\log \left(C(E)^{-1}\right)$, we obtain $\log \left(C(E)^{-1}\right) \leq \log \left(A^{-1}\right)$, that is, $A \leq C(E)$. By uniqueness of minimizers, if $A=C(E), d \nu_{\infty}=$ $d \rho_{E}$ and regularity implies $d \rho_{E}$ is the only limit point. By compactness, $d \nu_{n} \rightarrow d \rho_{E}$, and then, by (2.14), we obtain (1.28) for $z \notin H$.

Completion of the Proof of Theorem 1.10 for OPRL. We proved above that (i) $\Rightarrow$ (iv) and so, by the submean property of $|f(z)|^{1 / n}$ (alternatively, by the subharmonicity of $\log |f(z)|$ ) for analytic functions, we get (ii) also on $H$ and thus have (i) $\Rightarrow$ (ii) in full. (ii) $\Rightarrow$ (iii) is trivial.
$\left(\right.$ iii $\Rightarrow$ (i). Pick a subsequence $n(j)$ so that $\left(a_{1} \ldots a_{n(j)}\right)^{1 / n(j)} \rightarrow$ $\liminf \left(a_{1} \ldots a_{n(j)}\right)^{1 / n}=A$, and so $\nu_{n(j)} \rightarrow \nu_{\infty}$. By (1.4) and Theorem A.7, (iii) implies for q.e. $x$ in $E$ we have that $\Phi_{\nu_{\infty}}(x) \geq \log \left(A^{-1}\right)$. Thus, since $d \rho_{E}$ gives zero weight to zero capacity sets (see Proposition A.6) and (A.2),

$$
\begin{align*}
\log \left(A^{-1}\right) & \leq \int \log \Phi_{\nu_{\infty}}(x) d \rho_{E}(x) \\
& =\int \Phi_{\rho_{E}}(x) d \nu_{\infty}(x) \\
& \leq \log \left(C(E)^{-1}\right) \tag{2.19}
\end{align*}
$$

by (A.23). Thus $A^{-1} \leq C(E)^{-1}$, so $C(E) \leq A$. By Theorem 2.4 (or Corollary 1.2), we see $A=C(E)$, that is, $\mu$ is regular.

The reader may be concerned about this argument if $A=0$. But in that case, Theorem A. 7 and (1.27) imply that q.e. on $E, \Phi_{\nu}(x)=\infty$ which is inconsistent with Theorem A. 11 (or with Corollary A.5). Thus (1.27) implies $A>0$.
$(\mathrm{i}) \Rightarrow(\mathrm{v})$. This is immediate from Theorem 2.4, the fact that equality holds q.e. in (A.15) and in (A.23).
(ii) $\Rightarrow$ (vi). Without loss, we can redefine $Q_{n}$ so $\left\|Q_{n}\right\|_{L^{2}(d \mu)}=1$. Then

$$
\begin{equation*}
Q_{n}(z)=\sum_{j=0}^{n} c_{j, n} p_{j}(z, d \mu) \tag{2.20}
\end{equation*}
$$

where

$$
\begin{equation*}
\sum_{j=0}^{n}\left|c_{j, n}\right|^{2}=1 \Rightarrow\left|c_{j, n}\right| \leq 1 \tag{2.21}
\end{equation*}
$$

Thus

$$
\begin{equation*}
\left|Q_{n}(z)\right| \leq n \sup _{0 \leq j \leq n}\left|p_{j}(z, d \mu)\right| \tag{2.22}
\end{equation*}
$$

and $(1.26) \Rightarrow(1.30)$.
This completes the proof of Theorem 1.10 for OPRL and our presentation of the key properties of regular measures for OPRL. We turn to relations between the support of $d \mu$ and regularity of the density of zeros that will include Theorem 1.9.

Theorem 2.5. Let $d \mu$ be a measure of compact support, E, with Jacobi parameters, $\left\{a_{n}, b_{n}\right\}_{n=1}^{\infty}$. Let $n(j)$ be a subsequence so that $d \nu_{n(j)}$ has a limit, $d \rho_{E}$, the equilibrium measure for $E$. Then either
(a)

$$
\begin{equation*}
\lim _{j \rightarrow \infty}\left(a_{1} \ldots a_{n(j)}\right)^{1 / n(j)}=C(E) \tag{2.23}
\end{equation*}
$$

or
(b) $\mu$ is carried by a set of capacity zero, that is, there is $X \subset E$ of capacity zero so $\mu(\mathbb{R} \backslash X)=0$.

Proof. Let $A$ be a limit point of $\left(a_{1} \ldots a_{n(j)}\right)^{1 / n}$. If $A=0$, interpret $A^{-1}$ as $\infty$. By (2.16) and the upper envelope theorem (Theorem A.7), we see for some subsubsequence $\tilde{n}(j)$,

$$
\begin{equation*}
\lim _{j \rightarrow \infty}\left|p_{\tilde{n}(j)}(x)\right|^{1 / \tilde{n}(j)}=A^{-1} \exp \left(-\Phi_{\rho_{E}}(x)\right) \tag{2.24}
\end{equation*}
$$

for q.e. $x$. By Theorem A. $10, \Phi_{\rho_{E}}(x)=\log \left(C(E)^{-1}\right)$ for q.e. $x$. So for q.e. $x \in E$,

$$
\begin{equation*}
\lim _{j \rightarrow \infty}\left|p_{\tilde{n}(j)}(x)\right|^{1 / \tilde{n}(j)}=A^{-1} C(E) \tag{2.25}
\end{equation*}
$$

On the other hand (see (4.14) below), for $\mu$-a.e. $x$, we have

$$
\begin{equation*}
\left|p_{n}(x)\right| \leq C(x)(n+1) \tag{2.26}
\end{equation*}
$$

so for such $x$,

$$
\begin{equation*}
\lim \sup \left|p_{n}(x)\right|^{1 / n} \leq 1 \tag{2.27}
\end{equation*}
$$

If $A<C(E)$, then $\frac{C(E)}{A}>1$, so (2.27) can only hold on the set of capacity zero where (2.25) fails, that is, either $A=C(E)$ (since it is always true that $A \leq C(E)$ ) or $\mu$ is carried by a set of capacity zero.

Before leaving the subject of OPRL, we want to say something about nonregular situations:

Theorem 2.6. Let $\mu$ be a fixed measure of compact support on $\mathbb{R}$.
(a) The set of limit points of $\left(a_{1} \ldots a_{n}\right)^{1 / n}$ is always a closed interval.
(b) The set of limits of zero counting measures $d \nu_{n}$ is always a closed compact set.

Remarks. 1. As quoted in [119], where the first proof of (a) appeared, (a) is a theorem of Freud and Ziegler.
2. Part (b) was conjectured in Ullman [117] and is proven in StahlTotik [105] (see Theorem 2.1.4 of [105]).
3. Stahl-Totik [105] also prove (their Theorem 2.2.1) that so long as no carrier of $\mu$ has capacity zero, the existence of a limit for $d \nu_{n(j)}$ implies the existence of a limit for $\left(a_{1} \ldots a_{n(j)}\right)^{1 / n(j)}$. However, as we will see (Example 2.7), the converse is false.

Proof. We sketch the proof of (a); the proof of (b) can be found in [105] and is similar in spirit. The set of limit points is a closed subset of $[0, C(E)]$. If it is not connected, we can find limit points $A<B$ and $c \in(A, B)$ which is not a limit point.

Thus, there are $N$ and $\varepsilon$ so for $n>N$,

$$
\begin{equation*}
\Gamma_{n} \equiv\left(a_{1} \ldots a_{n}\right)^{1 / n} \notin(c-\varepsilon, c+\varepsilon) \tag{2.28}
\end{equation*}
$$

Suppose $\Gamma_{n}<c-\varepsilon$ and let $D=\frac{1}{2} \operatorname{diam}(\operatorname{cvh}(\operatorname{supp}(\mu))) \geq a_{n}$ by (2.2). Then

$$
\begin{equation*}
\Gamma_{n+1}=a_{n+1}^{1 / n+1} \Gamma_{n}^{n / n+1} \leq D^{1 / n+1}(c-\varepsilon)^{n / n+1} \tag{2.29}
\end{equation*}
$$

Since RHS of (2.29) converges to $c-\varepsilon$, we can find $N_{1}$ so

$$
\begin{equation*}
n \geq N_{1} \Rightarrow \text { RHS of }(2.29) \leq c \tag{2.30}
\end{equation*}
$$

Thus $n \geq N, n \geq N_{1}$, and $\Gamma_{n} \leq c-\varepsilon$ implies $\Gamma_{n+1} \leq c-\varepsilon$ (by (2.28)).
It follows that $\Gamma_{n}$ cannot have both $A$ and $B$ as limit points.
This contradiction proves the set of limit points is an interval.
Example 2.7. This example shows that $\left(a_{1} \ldots a_{n}\right)^{1 / n}$ may have a limit (necessarily strictly less than $C(E)$ ) but $d \nu_{n}$ does not. A more complicated example appears as Example 2.2.7 in [105]. Let $a_{n} \equiv 1$ (so $\left.\left(a_{1} \ldots a_{n}\right)^{1 / n} \rightarrow 1\right)$ and

$$
b_{n}= \begin{cases}1 & N_{2 \ell} \leq n<N_{2 \ell+1} \\ -1 & N_{2 \ell+1} \leq n<N_{2 \ell+2}\end{cases}
$$

where $N_{\ell}=2^{3^{\ell}}$. It is easy to see by looking at traces of powers of the cutoff Jacobi matrix that $d \nu_{N_{2 \ell}^{2}} \rightarrow d \rho_{[-1,3]}$ and $d \nu_{N_{2 \ell+1}^{2}} \rightarrow d \rho_{[-3,1]}$.

There is another result about the set of limit points that should be mentioned in connection with work of Ullman and collaborators. Define $c_{\mu}$ to be inf of the capacity of Borel sets, $S$, which are carriers of $\mu$ in the sense that $\mu(\mathbb{R} \backslash S)=0$. For example, if $\mu$ is a dense pure point measure with support $E=[-2,2], \mu$ is supported on a countable set, so $c_{\mu}=0$ even though $C(E)=1$. Then, in general, Ullman shows that any limit point of $\left(a_{1} \ldots a_{n}\right)^{1 / n}$ lies in $\left[c_{\mu}, C(\operatorname{supp}(d \mu))\right]$, and Wyneken [123] proved that given any $\mu$ and any $[A, B] \subset\left[c_{\mu}, C(\operatorname{supp}(d \mu))\right]$, there is $\eta$ mutually equivalent to $\mu$ so the set of limit points of $\Gamma_{n}(\eta)$ is $[A, B]$ (see also Theorem 5.4 below).

In particular, these results show that if $c_{\mu}=C(\operatorname{supp}(d \mu))$, then $\mu$ is regular-a theorem of Ullman [114], although Widom [122] essentially had the same theorem (this oversimplifies the relation between Widom [122] and Ullman [114]; see [105, Ch. 4]). We have not discussed this result in detail because the Stahl-Totik criterion of Theorem 1.13 essentially subsumes these earlier works (at least for $E$ a finite union of closed intervals) and we will prove that in Section 5.

## 3. Regular Measures for OPUC

In this section, we will prove Theorem 1.10 for OPUC and an analog of Theorem 1.7. Here one issue will be that if $E=\partial \mathbb{D}$, the zero density may not converge to a measure on $\partial \mathbb{D}$. The key step concerns Proposition 2.2, which essentially depended on the CD formula which is only known for OPRL and OPUC, and where the OPUC version is not obviously relevant. Instead, we will see, using operator theoretic methods [101], that there is a kind of "half CD formula" that suffices. We begin with an analog of Lemma 2.1:

Lemma 3.1. (a) Let $\mu$ be a measure of compact support on $\mathbb{C}$ and $H$ the convex hull of the support of $\mu$. Let $M_{z}$ be multiplication by $z$ on $L^{2}(\mathbb{C}, d \mu)$. Then for any $z_{0} \in \mathbb{C}$ and $\varphi \in L^{2}(\mathbb{C}, d \mu)$, we have

$$
\begin{equation*}
\left|\left\langle\varphi,\left(M_{z}-z_{0}\right) \varphi\right\rangle\right| \geq \operatorname{dist}\left(z_{0}, H\right)\|\varphi\|^{2} \tag{3.1}
\end{equation*}
$$

(b) Let $D$ be defined by

$$
\begin{equation*}
D=\min _{w}\left[\max _{z \in H}|z-w|\right] \tag{3.2}
\end{equation*}
$$

(which lies between $\frac{1}{2} \operatorname{diam}(H)$ and $\operatorname{diam}(H)$ ). Then

$$
\begin{equation*}
\left\|X_{n+1}\right\|_{L^{2}(d \mu)} \leq D\left\|X_{n}\right\|_{L^{2}(d \mu)} \tag{3.3}
\end{equation*}
$$

Proof. (a) Let $\omega \in \partial \mathbb{D}$. Then

$$
\begin{aligned}
\left|\left\langle\varphi,\left(M_{z}-z_{0}\right) \varphi\right\rangle\right| & \geq \int \operatorname{Re}\left(\left(z-z_{0}\right) \bar{\omega}\right)|\varphi(z)|^{2} d \mu(z) \\
& \geq \min _{z \in H} \operatorname{Re}\left(\left(z-z_{0}\right) \bar{\omega}\right)\|\varphi\|^{2}
\end{aligned}
$$

Maximizing over $\omega$ yields (3.1).
(b) Since $(z-w) X_{n}$ is a monic polynomial of degree $n+1$,

$$
\left\|X_{n+1}\right\| \leq\left\|(z-w) X_{n}\right\| \leq \max _{z \in H}|z-w|\left\|X_{n}\right\|
$$

Minimizing over $w$ yields (3.3).
To get the analog of (2.7), we need
Proposition 3.2. Let $d \mu$ be a measure of compact support on $\mathbb{C}$ and let $M_{z}$ be multiplication by $z$ on $L^{2}(\mathbb{C}, d \mu)$. Let $K$ be the orthogonal projection in $L^{2}(\mathbb{C}, d \mu)$ onto the $n+1$-dimensional subspace polynomials of degree at most $n$. Then

$$
\begin{equation*}
\left[M_{z}, K\right] K=\frac{\left\|X_{n+1}\right\|}{\left\|X_{n}\right\|}\left[\left\langle x_{n}, \cdot\right\rangle x_{n+1}\right] \tag{3.4}
\end{equation*}
$$

Remark. This is essentially "half" the CD formula; operator theoretic approaches to the CD formula are discussed in [101].

Proof. For any $\varphi$,

$$
\begin{equation*}
\left[M_{z}, K\right] K \varphi=(1-K) z(K \varphi) \tag{3.5}
\end{equation*}
$$

This clearly vanishes if $K \varphi=0$ or if $\varphi \in \operatorname{ran} K_{n-1}$. Thus, it is a rank one operator. Moreover, since $(1-K) z X_{n}=X_{n+1}$, we see

$$
\left[M_{z}, K\right] K X_{n}=X_{n+1}
$$

Since $X_{n+1}=\left\|X_{n+1}\right\| x_{n+1}$ and $X_{n}=\left\|X_{n}\right\| x_{n}$, we see that (3.4) holds.

Proposition 3.3. Let $d \mu$ be a measure of compact support on $\mathbb{C}$, $x_{n}(z ; d \mu)$ the normalized OPs, and $H$ the convex hull of the support of $d \mu$. For $z_{0} \notin H$, let $d\left(z_{0}\right)=\operatorname{dist}\left(z_{0}, H\right)$ and let $D$ be given by (3.2). Then for such $z_{0}$,

$$
\begin{equation*}
\left|x_{n}\left(z_{0} ; d \mu\right)\right|^{2} \geq\left(\frac{d}{D}\right)^{2}\left(1+\left(\frac{d}{D}\right)^{2}\right)^{n-1} \tag{3.6}
\end{equation*}
$$

In particular, $x_{n}\left(z_{0}\right) \neq 0$ for all $n$ and

$$
\begin{equation*}
\liminf \left|x_{n}\left(z_{0} ; d \mu\right)\right|^{1 / n} \geq\left(1+\left(\frac{d}{D}\right)^{2}\right)^{1 / 2}>1 \tag{3.7}
\end{equation*}
$$

Remark. Again, it is well known (a theorem of Fejér) that zeros of $x_{n}$ lie in $H$.

Proof. Define

$$
\begin{equation*}
\varphi_{n}(w)=\sum_{j=0}^{n} \overline{x_{j}\left(z_{0}\right)} x_{j}(w) \tag{3.8}
\end{equation*}
$$

We claim

$$
\begin{equation*}
\left\langle\varphi_{n},\left(M_{z}-z_{0}\right) \varphi_{n}\right\rangle=-\frac{\left\|X_{n+1}\right\|}{\left\|X_{n}\right\|} \overline{x_{n}\left(z_{0}\right)} x_{n+1}\left(z_{0}\right) \tag{3.9}
\end{equation*}
$$

This is precisely an analog of (2.7). Given this and Lemma 3.1, the proof is identical to that of Proposition 2.2.

To prove (3.9), we note the integral kernel of $K_{n}$ is

$$
\begin{equation*}
K_{n}(s, t)=\sum_{j=0}^{n} p_{j}(s) \overline{p_{j}(t)} \tag{3.10}
\end{equation*}
$$

and that (3.4) says

$$
\begin{equation*}
\int(s-w) K_{n}(s, w) K_{n}(w, t) d \mu(w)=\frac{\left\|X_{n+1}\right\|}{\left\|X_{n}\right\|} x_{n+1}(s) \overline{x_{n}\left(z_{0}\right)} \tag{3.11}
\end{equation*}
$$

(3.11) originally holds for a.e. $s, t$ in $\operatorname{supp}(d \mu)$, but since both sides are polynomials in $s$ and $\bar{t}$, for all $s, t$. Setting $s=t=z_{0}$, (3.11) is just (3.9).

Now we want to specialize to OPUC. The zeros in that case lie in $\mathbb{D}$. One defines the balayage of the zeros measure, $d \nu_{n}$, on $\partial \mathbb{D}$ by

$$
\begin{equation*}
\mathcal{P}(d \nu)=F(\theta) \frac{d \theta}{2 \pi} \tag{3.12}
\end{equation*}
$$

where

$$
\begin{equation*}
F(\theta)=\int \frac{1-|z|^{2}}{\left|e^{i \theta}-z\right|^{2}} d \nu_{n}(z) \tag{3.13}
\end{equation*}
$$

It is the unique measure on $\partial \mathbb{D}$ with

$$
\begin{equation*}
\int z^{k} \mathcal{P}\left(d \nu_{n}\right)=\int z^{k} d \nu_{n}(z) \tag{3.14}
\end{equation*}
$$

for $k \geq 0$ (see [98, Prop. 8.2.2]).
Since $|z|>1 \geq|w|$ implies

$$
\begin{equation*}
\log |z-w|^{-1}=-\log |z|+\operatorname{Re}\left(\sum_{j=1}^{\infty} \frac{1}{j}\left(\frac{w}{z}\right)^{j}\right) \tag{3.15}
\end{equation*}
$$

by (3.14), we have

$$
\begin{equation*}
|z|>1 \Rightarrow \Phi_{\nu_{n}}(z)=\Phi_{\mathcal{P}\left(d \nu_{n}\right)}(z) \tag{3.16}
\end{equation*}
$$

If $d \nu_{n} \rightarrow d \nu_{\infty}$, then $\mathcal{P}\left(d \nu_{n}\right) \rightarrow \mathcal{P}\left(d \nu_{\infty}\right)$, and this equals $d \nu_{\infty}$ if $d \nu_{\infty}$ is a measure on $\partial \mathbb{D}$. If $\operatorname{supp}(d \mu) \varsubsetneqq \partial \mathbb{D}$, then it is known that the bulk of the zeros goes to $\partial \mathbb{D}$ (Widom's zero theorem; see [98, Thm. 8.1.8]), so $d \nu_{\infty}$ is a measure on $\partial \mathbb{D}$. It is also known (see [98, Thm. 8.2.7]) that the zero counting measures for the paraorthogonal polynomials (POPUC) have the same weak limits as $\mathcal{P}\left(d \nu_{n}\right)$. The analogs of Theorems 2.3 and 2.4 are thus:

Theorem 3.4. Let $d \mu$ be a measure on $\partial \mathbb{D}$, the unit circle. Let $n(j)$ be a subsequence with $n(1)<n(2)<\ldots$ so that $\left(\rho_{1} \ldots \rho_{n(j)}\right)^{1 / n(j)}$ has a nonzero limit $A$ and so that there is a measure $d \nu_{\infty}$ on $\partial \mathbb{D}$ which is the weak limit of $\mathcal{P}\left(d \nu_{n(j)}\right)$ (equivalently, of $d \nu_{n(j)}$ if $\operatorname{supp}(d \mu) \neq \partial \mathbb{D}$; equivalently, of the zero counting measures of POPUC). Then for any $|z|>1$ or $z \notin \partial \mathbb{D} \backslash \operatorname{supp}(d \mu)$,

$$
\begin{equation*}
\lim _{j \rightarrow \infty}\left|\varphi_{n(j)}(z)\right|^{1 / n(j)}=A^{-1} \exp \left(-\Phi_{\nu_{\infty}}(z)\right) \tag{3.17}
\end{equation*}
$$

In particular,

$$
\begin{equation*}
\exp \left(-\Phi_{\nu_{\infty}}(z)\right) \geq A \tag{3.18}
\end{equation*}
$$

It follows that if $E=\sigma_{\text {ess }}(d \mu)$, then $A \leq C(E)$, and if $\mu$ is regular (i.e., $\left.\left(\rho_{1} \ldots \rho_{n}\right)^{1 / n} \rightarrow C(E)\right)$, then every limit point of $\mathcal{P}\left(d \nu_{n(j)}\right)$ is the equilibrium measure $d \nu_{E}$. So $\mathcal{P}\left(d \nu_{n}\right) \rightarrow d \nu_{E}$ (and if $E \neq \partial \mathbb{D}$, $\left.d \nu_{n} \rightarrow d \nu_{E}\right)$.

Proof. Given the above discussion and results, this is identical to the proofs of Theorems 2.3 and 2.4.

By mimicking the proof we give for Theorem 1.10 for OPRL, we obtain the same result for OPUC.

## 4. Van Assche's Proof of Widom's Theorem

In this section, we will prove Theorem 1.12 using part of Van Assche's approach [120]. The basic idea is simple: By a combination of Chebyshev's inequality and the Borel-Cantelli lemma, if $\left\|P_{n(j)}\right\|_{L^{2}(d \mu)}^{1 / n(j)} \rightarrow A$, then for $d \mu$-a.e. $x$, we have $\lim \sup _{j \rightarrow \infty}\left|P_{n(j)}(x)\right|^{1 / n(j)} \leq A$. By using some potential theory, we will find that the density of zeros measure, $d \nu$, supported on $E$ obeys for q.e. $x, \Phi_{\nu}(x) \geq \log \left(A^{-1}\right)$ a.e. $d \mu$. Since $d \rho_{E}$ is a.c. with respect to $d \mu$, this will imply $\int \Phi_{\rho_{E}}(x) d \nu(x) \geq$ $\log \left(A^{-1}\right)$. But by potential theory again, $\Phi_{\rho_{E}}(x) \leq \log \left(C(E)^{-1}\right)$, so we will have $A^{-1} \leq C(E)^{-1}$, that is, $C(E) \leq A$.

Lemma 4.1. Let $d \mu$ be a probability measure on a measure space $X$. Let $f_{n(j)}$ be a sequence of functions indexed by integers $1 \leq n(1)<$ $n(2)<\ldots$. Suppose for some $1 \leq p<\infty$,

$$
\begin{equation*}
\limsup _{j \rightarrow \infty}\left\|f_{n(j)}\right\|_{L^{p}}^{1 / n(j)}=A \tag{4.1}
\end{equation*}
$$

Then for $d \mu$-a.e. $x$,

$$
\begin{equation*}
\limsup \left|f_{n(j)}(x)\right|^{1 / n(j)} \leq A \tag{4.2}
\end{equation*}
$$

Proof. Fix $B>A$. Then

$$
\begin{equation*}
\mu\left(S_{j}(B)\right) \equiv \mu\left(\left\{x\left|\left|f_{n(j)}(x)\right|>B^{n(j)}\right\}\right) \leq \frac{\left\|f_{n(j)}\right\|_{L^{p}}^{p}}{B^{n(j) p}}\right. \tag{4.3}
\end{equation*}
$$

By (4.1) and $B>A$, we see

$$
\sum_{j} \mu\left(S_{j}(B)\right)<\infty
$$

so for $\mu$-a.e. $x$, there is $J(x)$ with $x \notin S_{j}$ for all $j>J(x)$. Thus, for $\mu$-a.e. $x$,

$$
\limsup \left|f_{n(j)}(x)\right|^{1 / n(j)} \leq B
$$

Since $B$ is arbitrary, we have (4.2).
Proof of Theorem 1.12. Let $A$ be a limit point of $\left\|P_{n(j)}\right\|_{L^{2}(d \mu)}^{1 / n(j)}$. By passing to a subsequence, we can suppose the zero counting measure $d \nu_{n(j)}$ has a limit $d \nu_{\infty}$ which, by Proposition 1.8, is supported on $E$.

By Lemma 4.1 for a.e. $x(d \mu)$,

$$
\begin{equation*}
\limsup \left|P_{n(j)}(x)\right|^{1 / n(j)} \leq A \tag{4.4}
\end{equation*}
$$

By (1.4) for such $x$,

$$
\begin{equation*}
\lim \sup \exp \left(-\Phi_{\nu_{n(j)}}(x)\right) \leq A \tag{4.5}
\end{equation*}
$$

By the upper envelope theorem (Theorem A.7) for q.e. $x \in \mathbb{C}$,

$$
\begin{equation*}
\Phi_{\nu_{\infty}}(x)=\liminf \Phi_{\nu_{n(j)}}(x) \tag{4.6}
\end{equation*}
$$

Thus, there exist sets $S_{1}$ and $S_{2}$ so that $\mu\left(S_{1}\right)=0$ and $C\left(S_{2}\right)=0$, so that for $x \in \mathbb{C} \backslash\left(S_{1} \cup S_{2}\right)$,

$$
\begin{equation*}
\Phi_{\nu_{\infty}}(x) \geq \log \left(A^{-1}\right) \tag{4.7}
\end{equation*}
$$

We can now repeat the argument that led to (2.19). By hypothesis, $d \rho_{E}$ is $d \mu$-a.c. So $\rho_{E}\left(S_{1}\right)=0$ and, of course, since $\mathcal{E}\left(\rho_{E}\right)<\infty, \rho_{E}\left(S_{2}\right)=$ 0 . Thus, (4.7) holds a.e. $d \rho_{E}$.

Therefore, by (A.2),

$$
\begin{aligned}
\log \left(A^{-1}\right) & \leq \int \Phi_{\nu_{\infty}}(x) d \rho_{E}(x) \\
& =\int \Phi_{\rho_{E}}(x) d \nu_{\infty}(x) \\
& \leq \log \left(C(E)^{-1}\right)
\end{aligned}
$$

by (A.23).
Thus, $A^{-1} \leq C(E)^{-1}$ or $C(E) \leq A$. Thus, $\lim \inf \left\|P_{n}\right\|^{1 / n} \geq C(E)$. Since (see (1.15)), limsup $\left\|P_{n}\right\|^{1 / n} \leq C(E)$, we have regularity.

The above proof is basically a part of Van Assche's argument [120] which can be simplified since he proves that $d \nu_{\infty}=d \rho_{E}$ by a direct argument using similar ideas, and we can avoid that because of the general argument in Section 2.

This argument can also prove a related result-we will see examples of this phenomenon at the end of the next section.

Theorem 4.2. Suppose $\mu$ is a measure of compact support on $\mathbb{R}$ so $E \subset \operatorname{supp}(d \mu)$ for an essentially perfect compact set $E$ with $C(E)>0$. Suppose $d \rho_{E}$ is a.c. with respect to $d \mu$, and for some $n(1)<n(2)<\ldots$, we have

$$
\begin{equation*}
\left\|P_{n(j)}\right\|_{L^{2}(\mathbb{R}, d \mu)}^{1 / n(j)} \rightarrow C(E) \tag{4.8}
\end{equation*}
$$

for the monic $P_{n}(x, d \mu)$. Let $d \nu_{n(j)}$ be the corresponding zero counting measure. Then $d \nu_{n(j)} \xrightarrow{w} d \rho_{E}$.

Remarks. 1. We have in mind cases where $E$ is a proper subset of $\operatorname{supp}(d \mu)$. There will be many subsets with the same capacity, but there can only be one that has $d \rho_{E}$ a.e. with respect to $d \mu$.
2. Since $d \mu \upharpoonright E$ is regular (by Theorem 1.12) and $\left\|P_{n(j)}(\cdot, \mu)\right\|_{L^{2}(d \mu)} \geq\left\|P_{n(j)}(\cdot, \mu \upharpoonright E)\right\|_{L^{2}(d \mu \backslash E)}$, we see that

$$
\lim \inf \left\|P_{n}\right\|_{L^{2}(\mathbb{R}, d \mu)}^{1 / n} \geq C(E)
$$

so (4.8) is equivalent to a lim sup assumption.
Proof. Let $d \nu_{\infty}$ be a limit point of $d \nu_{n(j)}$. As in the proof of Theorem 1.12, there exist sets $S_{1}$ with $\mu\left(S_{1}\right)=0$ and $S_{2}$ with $C\left(S_{2}\right)=0$, so for $x \in \mathbb{C} \backslash\left(S_{1} \cup S_{2}\right)$, we have

$$
\begin{equation*}
\Phi_{\nu_{\infty}}(x) \geq \log \left(C(E)^{-1}\right) \tag{4.9}
\end{equation*}
$$

Since $\rho_{E}\left(S_{1}\right)=0$ by the assumption and $\rho_{E}\left(S_{2}\right)=0$ since $C(E)>0$, (4.9) holds for $\rho_{E}$-a.e. $x$. Moreover, since

$$
\begin{equation*}
\Phi_{\rho_{E}}(z) \leq \log \left(C(E)^{-1}\right) \tag{4.10}
\end{equation*}
$$

for all $z$,

$$
\begin{align*}
\log \left(C(E)^{-1}\right) & \leq \int \Phi_{\nu_{\infty}}(x) d \rho_{E}(x) \\
& =\int \Phi_{\rho_{E}}(x) d \nu_{\infty}(x) \\
& \leq \log \left(C(E)^{-1}\right) \tag{4.11}
\end{align*}
$$

Thus, using (4.10), we see

$$
\Phi_{\rho_{E}}(x)=\log \left(C(E)^{-1}\right)
$$

for $\nu_{\infty}$-a.e. $x$. But $\Phi_{\rho_{E}}(z)<\log \left(C(E)^{-1}\right)$ for all $z \notin E$, so $\nu_{\infty}$ is supported on $E$. By (4.11) and (4.9), $\Phi_{\nu_{\infty}}(x)=\log \left(C(E)^{-1}\right)$ for $d \rho_{E^{-}}$ a.e. $x$. By Theorem A.14, $\nu_{\infty}=\rho_{E}$.

There is an alternate way to prove (4.4) without Lemma 4.1 that links it to ideas more familiar to spectral theorists. It is well known that for elliptic PDEs, there are polynomially bounded eigenfunctions for a.e. energy with respect to spectral measures. This is called the BGK expansion in [94] after Berezanskiǐ[14], Browder [20], Gårding [40], Gel'fand [41], and Kac [62]. The translation to OPRL is discussed in Last-Simon [75]. Since $\int\left|p_{n}(x)\right|^{2} d \mu=1$, we have

$$
\begin{equation*}
\sum_{n=0}^{\infty}(n+1)^{-2} \int\left|p_{n}(x)\right|^{2} d \mu<\infty \tag{4.12}
\end{equation*}
$$

and thus, for $d \mu$-a.e. $x$,

$$
\begin{equation*}
\sum_{n=0}^{\infty}(n+1)^{-2}\left|p_{n}(x)\right|^{2}<\infty \tag{4.13}
\end{equation*}
$$

so

$$
\begin{equation*}
\left|P_{n}(x)\right| \leq C(x)(n+1)\left\|P_{n}\right\|_{L^{2}} \tag{4.14}
\end{equation*}
$$

which implies (4.4).

It is interesting to note that if $E$ is such that it is regular and $d \rho_{E}$ is purely absolutely continuous on $E=\operatorname{supp}\left(d \rho_{E}\right)$, one can use these ideas to provide an alternate proof (see Simon [100] for still another alternate proof in this case). For in that case, the measure associated to the second kind polynomials, $q_{n}(x)$, also has a.c. weight $\widetilde{w}(x)>0$ for a.e. $x$ in $E$, and thus

$$
\begin{equation*}
\left|q_{n}(x)\right| \leq C(x)(n+1) \tag{4.15}
\end{equation*}
$$

which, by constancy of the Wronskian, implies

$$
\begin{equation*}
\left|p_{n}(x)\right|^{2}+\left|p_{n+1}(x)\right|^{2} \geq \widetilde{C}(x)(n+1)^{-1} \tag{4.16}
\end{equation*}
$$

If $d \nu_{n(j)} \rightarrow d \nu_{\infty}$, so does $d \nu_{n(j)+1}$ (by interlacing of zeros), and thus, by (4.14) and (4.16), if $\lim \left(a_{1} \ldots a_{n(j)}\right)^{1 / n(j)} \rightarrow A$, then

$$
\begin{equation*}
-\log (A)+\int \log \left(|x-y|^{-1}\right) d \nu_{\infty}(y)=0 \tag{4.17}
\end{equation*}
$$

for $x \in E$ but with a set of Lebesgue measure zero and of capacity zero removed. By Theorem A.14, we conclude that $A=C(E)$ and $d \nu_{\infty}=d \rho_{E}$.
Remark. We note that (4.17) holds a.e. on the a.c. spectrum and by the above arguments, a.e. on that spectrum, $\frac{1}{n} \log \left\|T_{n}(x)\right\| \rightarrow 0$, a deterministic analog of the Pastur-Ishii theorem.

## 5. The Stahl-Totik Criterion

In this section, we will present an exposition of Stahl-Totik's proof [105] of their result, our Theorem 1.13. As a warmup, we prove

Theorem 5.1. Let $d \mu$ be a measure on $\partial \mathbb{D}$ obeying

$$
\begin{equation*}
\inf _{\theta_{0}} \mu\left(\left\{e^{i \theta}| | \theta-\theta_{0} \left\lvert\, \leq \frac{1}{m}\right.\right\}\right) \geq C_{\delta} e^{-\delta m} \tag{5.1}
\end{equation*}
$$

for all $\delta>0$. Then $\mu$ is regular.
Proof. We will use Bernstein's inequality that for any polynomial, $P_{n}$, of degree $n$,

$$
\begin{equation*}
\sup _{z \in \mathbb{D}}\left|P_{n}^{\prime}(z)\right| \leq n \sup _{z \in \partial \mathbb{D}}\left|P_{n}(z)\right| \tag{5.2}
\end{equation*}
$$

Szegő's simple half-page proof of this can be found, for example, in Theorem 2.2.5 of [98].

Applying this to the monic polynomials $\Phi_{n}(z ; d \mu)$, we see that if $\theta_{n}$ is chosen with $\left|\Phi_{n}\left(e^{i \theta_{n}} ; d \mu\right)\right|=\left\|\Phi_{n}\right\|_{\partial \mathbb{D}}$, the sup norm, and $\left|\theta-\theta_{n}\right| \leq \frac{1}{2 n}$, then

$$
\begin{equation*}
\left|\Phi_{n}\left(e^{i \theta} ; d \mu\right)\right| \geq \frac{1}{2}\left\|\Phi_{n}\right\|_{\partial \mathbb{D}} \tag{5.3}
\end{equation*}
$$

Thus, by (5.1) with $m=2 n$,

$$
\begin{equation*}
\left\|\Phi_{n}(\cdot ; d \mu)\right\|_{L^{2}(d \mu)}^{2} \geq\left(\frac{1}{4}\left\|\Phi_{n}\right\|_{\partial \mathbb{D}}^{2}\right) C_{\delta} e^{-2 \delta n} \tag{5.4}
\end{equation*}
$$

Since $\Phi_{n}(\cdot ; d \mu)$ is monic,

$$
\begin{equation*}
\int \Phi_{n}\left(e^{i \theta} ; d \mu\right) e^{-i n \theta} \frac{d \theta}{2 \pi}=1 \tag{5.5}
\end{equation*}
$$

so the sup norm obeys

$$
\begin{equation*}
\left\|\Phi_{n}\right\|_{\partial \mathbb{D}} \geq 1 \tag{5.6}
\end{equation*}
$$

and so (5.4) implies

$$
\liminf \left\|\Phi_{n}(\cdot ; d \mu)\right\|_{L^{2}}^{1 / n} \geq e^{-2 \delta}
$$

Since $\delta$ is arbitrary, the lim inf is larger than or equal to 1 . Since $C(E)=1, \mu$ is regular.

There are two issues with just using these ideas to prove Theorem 1.13. While (5.5) is special for $\partial \mathbb{D}$, its consequence, (5.6), is really only an expression of $\left\|T_{n}\right\|_{E} \geq C(E)^{n}$ (see (B.8)), so it is not an issue.

However, (5.2) only holds because a circle has no ends. The analog for, say, $[-1,1]$ is Bernstein's inequality

$$
\begin{equation*}
\left|p^{\prime}(x)\right| \leq \frac{n}{\sqrt{1-x^{2}}}\|p\|_{[-1,1]} \tag{5.7}
\end{equation*}
$$

or (Markov's inequality)

$$
\begin{equation*}
\left|p^{\prime}(x)\right| \leq n^{2}\|p\|_{[-1,1]} \tag{5.8}
\end{equation*}
$$

Either one can be used to obtain a theorem like Theorem 5.1 on $[-1,1]$ but $e^{-\delta m}$ needs to be replaced by $e^{-\delta \sqrt{m}}$ - interesting, but weaker than Theorem 1.13.

The other difficulty is that (5.1) is global, requiring a result uniform in $\theta_{0}$, and (1.33) needs only a result for most $\theta_{0}$. The problem with using bounds on derivatives is that they only get information on a single set of size $O\left(\frac{1}{n}\right)$ at best. They get $\left|p_{n}(x)\right| \geq \frac{1}{2}\left\|p_{n}\right\|_{E}$ there, but that is overkill-we only need $\left|p_{n}(x)\right| \geq e^{-\delta^{\prime} n}\left\|p_{n}\right\|_{E}$, and that actually holds on a set of size $O(1)$ ! The key will thus be a variant of the Remez inequality in the following form:
Proposition 5.2. Fix $E$ a finite union of closed bounded intervals in $\mathbb{R}$. Then there is $c(\delta)>0$ with $c(\delta) \rightarrow 0$ as $\delta \downarrow 0$, so that for any $F \subset E$ with $|E \backslash F|<\delta$, we have

$$
\begin{equation*}
\left\|Q_{n}\right\|_{E} \leq e^{c(\delta) n}\left\|Q_{n}\right\|_{F} \tag{5.9}
\end{equation*}
$$

for any polynomial, $Q_{n}$, of degree $n$.

Remarks. 1. This is a variant of an inequality of Remez [89]; see the proof for his precise result.
2. The relevance of Remez's inequality to regularity appeared already in Erdös-Turán [34] and was the key to the proof in Freud [39] of the Erdös-Turán theorem, Theorem 1.11. Its use here is due to Stahl-Totik [105].

Proof. If $E=I_{1} \cup \cdots \cup I_{\ell}$ disjoint intervals and $|E \backslash F| \leq \delta$, then $\left|I_{j} \backslash I_{j} \cap F\right| \leq \delta$ for all $j$, so it suffices to prove this result for each single interval and then, by scaling, for $E=[-1,1]$.

In that case, Remez's inequality (due to Remez [89]; see BorweinErdélyi [15] for a proof and further discussion) says that if $F \subset[-1,1]$ and $|[-1,1] \backslash F| \leq \delta$, then with $T_{n}$ the classical first kind Chebyshev polynomials,

$$
\begin{equation*}
\left\|Q_{n}\right\|_{E} \leq T_{n}\left(\frac{2+\delta}{2-\delta}\right)\left\|Q_{n}\right\|_{F} \tag{5.10}
\end{equation*}
$$

(This can be proven by showing the worst case occurs when $F=$ $[-1,1-\delta]$ and $\left.Q_{n}(x)=T_{n}\left(\frac{2 x+\delta}{2-\delta}\right).\right)$

Since

$$
\begin{equation*}
T_{n}(\cosh (x))=\cosh (n x) \leq e^{n x} \tag{5.11}
\end{equation*}
$$

and $\cosh (\varepsilon)=1+\frac{\varepsilon^{2}}{2}+O\left(\varepsilon^{4}\right)$, we have

$$
\begin{equation*}
T_{n}\left(\frac{2+\delta}{2-\delta}\right) \leq \exp \left(n\left[\sqrt{2 \delta}+O\left(\delta^{3 / 2}\right)\right]\right) \tag{5.12}
\end{equation*}
$$

so for $E=[-1,1],(5.9)$ holds with $c(\delta)=\sqrt{2 \delta}+O\left(\delta^{3 / 2}\right)$.
Lemma 5.3. If $P_{n}$ is a real polynomial of degree $n$, and $a>0, S \equiv$ $\left\{\lambda \in \mathbb{R}\left|\left|P_{n}(x)\right|>a\right\}\right.$ is a union of most $(n+1)$ intervals.

Proof. $\partial S$ is the finite set of points where $P_{n}(x)= \pm a$. If all the zeros of $P_{n} \pm a$ are simple, these boundary points are distinct. Including $\pm \infty$ so each interval has two "endpoints," these intervals have at most $2 n+2$ distinct endpoints (and exactly that number if all roots of $P_{n} \pm a$ are real). If some root of $P_{n} \pm a$ is double, two intervals can share an endpoint but that endpoint counts twice in the zeros.

Proof of Theorem 1.13. By Proposition B.3, if $P_{n}(x)=P_{n}(x ; d \mu)$, then

$$
\begin{equation*}
\sup _{x \in E}\left|P_{n}(x)\right| \geq c(E)^{n} \tag{5.13}
\end{equation*}
$$

Fix $\delta_{1}$ and let

$$
\begin{equation*}
F=\left\{x| | P_{n}(x) \mid \leq c(E)^{n} e^{-2 c\left(\delta_{1}\right) n}\right\} \tag{5.14}
\end{equation*}
$$

If $|E \backslash F|<\delta_{1}$, then (5.9) would imply $\left\|P_{n}\right\|_{E} \leq c(E)^{n} e^{-c\left(\delta_{1}\right) n}$, violating (5.13). So

$$
\begin{equation*}
|E \backslash F| \geq \delta_{1} \tag{5.15}
\end{equation*}
$$

By Lemma $5.3, \mathbb{R} \backslash F$ is a union of at most $n+1$ intervals, so if $E$ is a union of $\ell$ intervals, $E \backslash F$ consists of at most $\ell(n+1)$ intervals (a very crude overestimate that suffices for us!).

Some of these intervals may have size less than $\frac{\delta_{1}}{4 n \ell}$, but the total size of those is at most $\frac{\delta_{1}}{2}$, so we can find disjoint intervals $I_{1}^{(n)}, \ldots, I_{k(n)}^{(n)}$ in $E \backslash F$, so

$$
\begin{equation*}
\left|I_{j}^{(n)}\right| \geq \frac{\delta_{1}}{4 n \ell} \quad\left|\bigcup_{j=1}^{k(n)} I_{j}^{(n)}\right| \geq \frac{\delta_{1}}{2} \tag{5.16}
\end{equation*}
$$

Let $\tilde{I}_{j}^{(n)}$ be the interval of size $\frac{1}{2}\left|I_{j}^{(n)}\right|$ and the same center. Then with $L^{(n)}\left(\delta_{1}\right)=\cup_{j=1}^{k(n)} \tilde{I}_{j}^{(n)}$, we have

$$
\begin{gather*}
\left|L^{(n)}\left(\delta_{1}\right)\right| \geq \frac{\delta_{1}}{4}  \tag{5.17}\\
\left|P_{n}(y)\right| \geq c(E)^{n} e^{-2 c\left(\delta_{1}\right) n} \quad \text { if } \quad \operatorname{dist}\left(y, L^{(n)}\left(\delta_{1}\right)\right) \leq \frac{\delta_{1}}{16 n \ell} \tag{5.18}
\end{gather*}
$$

Now define for any $\delta_{2}>0$ and $m$,

$$
J\left(m, \delta_{2}\right)=\left\{x \left\lvert\, \mu\left(x-\frac{1}{m}, x+\frac{1}{m}\right) \geq e^{-\delta_{2} m}\right.\right\}
$$

By hypothesis, for any fixed $\delta_{2}$,

$$
\lim _{m \rightarrow \infty}\left|E \backslash J\left(m, \delta_{2}\right)\right|=0
$$

and, in particular, for any fixed integer $M$, for all large $n$,

$$
\begin{equation*}
\left|E \backslash J\left(M n, \delta_{2}\right)\right|<\frac{\delta_{1}}{4} \tag{5.19}
\end{equation*}
$$

so, in particular,

$$
\begin{equation*}
J\left(M n, \delta_{2}\right) \cap L^{(n)}\left(\delta_{1}\right) \neq \emptyset \tag{5.20}
\end{equation*}
$$

Given $\delta_{1}$, pick $M$ so large that

$$
\begin{equation*}
M^{-1} \leq \frac{\delta_{1}}{16 \ell} \tag{5.21}
\end{equation*}
$$

If $x$ lies in the set on the left side of (5.20), let $I=\left\{y| | x-y \left\lvert\, \leq \frac{1}{M n}\right.\right\}$. Then since $\frac{1}{M n} \leq \frac{\delta_{1}}{16 n \ell}$, for $y \in I$,

$$
\begin{equation*}
\left|P_{n}(y)\right| \geq c(E)^{n} e^{-2 c\left(\delta_{1}\right) n} \tag{5.22}
\end{equation*}
$$

since $x \in L^{(n)}\left(\delta_{1}\right)$ and (5.18) holds. By $x \in J\left(M n, \delta_{2}\right)$,

$$
\begin{equation*}
\mu(I) \geq e^{-M \delta_{2} n} \tag{5.23}
\end{equation*}
$$

Thus,

$$
\begin{equation*}
\left\|P_{n}\right\|_{L^{2}} \geq c(E)^{n} e^{-2 c\left(\delta_{1}\right) n} e^{-M N \delta_{2} / 2} \tag{5.24}
\end{equation*}
$$

so

$$
\begin{equation*}
\lim \inf \left\|P_{n}\right\|_{L^{2}}^{1 / n} \geq c(E) e^{-2 c\left(\delta_{1}\right)} e^{-M \delta_{2} / 2} \tag{5.25}
\end{equation*}
$$

First pick $\delta_{1}$, then fix $M$ by (5.21) (recall $\ell$ is fixed as the number of intervals in $E$ ) and let $\delta_{2}=\frac{\delta_{1}}{M}$. Then take $\delta_{1} \downarrow 0$ and get $\liminf \left\|P_{n}\right\|_{L^{2}}^{1 / n} \geq c(E)$, proving regularity.

Here is a typical application of the Stahl-Totik criterion. It illustrates the limitations of regularity criteria like those of $[114,122]$ that only depend on what sets are carriers for $\mu$. This result is a special case of a theorem of Wyneken [123].

Theorem 5.4. Let $\mu$ be a measure whose support is $E$, a finite union of closed intervals. Then there exists a measure $\eta$ equivalent to $\mu$ which is regular.

Proof. For any $n$, define

$$
\begin{equation*}
\mu_{n}=\sum_{\left\{j \left\lvert\, \mu\left(\left(\frac{j}{n}, \frac{j+1}{n}\right]\right)>0\right.\right\}} \mu\left(\left(\frac{j}{n}, \frac{j+1}{n}\right]\right)^{-1} \mu \upharpoonright\left(\frac{j}{n}, \frac{j+1}{n}\right] \tag{5.26}
\end{equation*}
$$

Then $\mu_{n}$ has total mass at most $n(|E|+\ell)$ where $\ell$ is the number of intervals. Let

$$
\begin{equation*}
\eta=\sum_{n=1}^{\infty} n^{-3} \mu_{n} \tag{5.27}
\end{equation*}
$$

which is easily seen to be equivalent to $\mu$.
Notice that if $\operatorname{dist}(x, \mathbb{R} \backslash E)>\frac{1}{n}$, then $\left[x-\frac{1}{n}, x+\frac{1}{n}\right]$ contains an interval of the form $\left(\frac{j}{2 n}, \frac{j+1}{2 n}\right]$, so $\mu_{2 n}\left(\left[x-\frac{1}{n}, x+\frac{1}{n}\right]\right) \geq 1$. Thus

$$
\begin{equation*}
\left|\left\{x \left\lvert\, \eta\left(\left[x-\frac{1}{n}, x+\frac{1}{n}\right]\right) \leq 8 n^{-3}\right.\right\}\right| \leq \frac{2 \ell}{n} \tag{5.28}
\end{equation*}
$$

and (1.33) holds.
By using point measures, it is easy to construct nonregular measures, including ones that illustrate how close (1.34) is to being ideal. The key is

Theorem 5.5. Let $\left\{x_{j}\right\}_{j=1}^{\infty}$ be a bounded sequence in $\mathbb{R}$ and $\left\{a_{j}\right\}_{j=1}^{\infty}$ an $\ell^{1}$ sequence of positive numbers. Let

$$
\begin{equation*}
\mu=\sum_{j=1}^{\infty} a_{j} \delta_{x_{j}} \tag{5.29}
\end{equation*}
$$

Let

$$
\begin{equation*}
d=\max _{j, k}\left|x_{j}-x_{k}\right| \tag{5.30}
\end{equation*}
$$

Then

$$
\begin{equation*}
\left\|P_{n}(x, d \mu)\right\|_{L^{2}(\mathbb{R}, d \mu)} \leq d^{n}\left(\sum_{j=n+1}^{\infty} a_{j}\right)^{1 / 2} \tag{5.31}
\end{equation*}
$$

Proof. Let $Q_{n}(x)=\prod_{j=1}^{n}\left(x-x_{j}\right)$, which kills the contributions of the pure points at $\left\{x_{j}\right\}_{j=1}^{n}$, so

$$
\begin{equation*}
\left\|Q_{n}\right\|^{2} \leq \sum_{j=n+1}^{\infty} d^{2 n} a_{j} \tag{5.32}
\end{equation*}
$$

by (5.30). Since $\left\|P_{n}\right\| \leq\left\|Q_{n}\right\|$, (5.31) is immediate.
Corollary 5.6. Let $\left\{x_{j}\right\}_{j=1}^{\infty}$ be an arbitrary bounded subset of $\mathbb{R}$. Then there exists a pure point measure $d \mu$ with precisely this set as its set of pure points, so that $\left\|P_{n}\right\|^{1 / n} \rightarrow 0$. In particular, if $E$ is any compact set with $C(E)>0$, there is a measure $\mu$ with $\operatorname{supp}(\mu)=E$ and $\mu$ not regular.
Proof. Pick $a_{j}=e^{-j^{2}}$ so $\left(\sum_{j=n+1}^{\infty} a_{j}\right)^{1 / 2 n} \rightarrow 0$.
Example 5.7. The following illuminates (1.33). For $2^{n} \leq k<2^{n+1}$, let $x_{k}=\frac{k-2^{n}}{2^{n}}$ and let $0<y<1$. Define

$$
\begin{equation*}
d \mu=\sum_{k=1}^{\infty} y^{k} \delta_{x_{k}} \tag{5.33}
\end{equation*}
$$

The $x_{k}$ are not distinct, but that does not change the bound (5.31). Thus

$$
\begin{equation*}
\lim \sup \left\|P_{n}\right\|^{1 / n} \leq y \tag{5.34}
\end{equation*}
$$

Since $C([0,1])=\frac{1}{4}$, the measure is not regular if $y<\frac{1}{4}$. On the other hand, if $2^{n} \leq m \leq 2^{n+1}$ and $x_{0} \in[0,1]$, there is an $x_{k}$ with $\left|x_{k}-x_{0}\right| \leq \frac{1}{m}$ and $2^{n} \leq k<2^{n+1}$. Thus

$$
\mu\left(\left[x_{0}-\frac{1}{m}, x_{0}+\frac{1}{m}\right]\right) \geq y^{k} \geq y^{2^{n+1}} \geq y^{2 m}
$$

so (1.33) holds for $\eta=-\log y^{2}$, that is, for some but not all $\eta$. This shows the exponential rate in Theorem 1.13 cannot be improved.

Example 5.8. We will give an example of a measure $d \mu$ on $[-2,2]$ which is a.c. on $[-2,0]$ and so that among the limit points of the zero counting measures, $d \nu_{n}$ are both $d \rho_{[-2,2]}$ and $d \rho_{[-2,0]}$, the equilibrium measure for $[-2,2]$ and for $[-2,0]$. This will answer a question asked me by Yoram Last, in reaction to Remling [90], whether a.c. spectra
force the existence of a density of states and also show that bounds on limit points of $d \nu_{n}$ of Totik-Ullman [112] and Simon [100] cannot be improved.

We define $d \mu$ by

$$
\begin{equation*}
d \mu \upharpoonright[-2,0]=(-x(x+2))^{-1 / 2} d x \tag{5.35}
\end{equation*}
$$

picked so the OPRL for the restriction are multiples of the Chebyshev polynomials for $[-2,0]$

$$
\begin{equation*}
d \mu \upharpoonright[0,2]=\sum_{n=1}^{\infty} a_{n} d \eta_{n} \tag{5.36}
\end{equation*}
$$

where $d \eta_{n}$ is concentrated uniformly at the dyadic rationals of the form $k / 2^{n}$ not previously "captured," that is,

$$
\begin{equation*}
d \eta_{n}=\sum_{j=0}^{2^{n}-1} \frac{1}{2^{n}} \delta_{(2 j+1) / 2^{n-1}} \tag{5.37}
\end{equation*}
$$

The $a_{n}$ 's are carefully picked as follows. Define $N_{j}$ inductively by

$$
\begin{equation*}
N_{1}=1 \quad N_{j+1}=2^{N_{j}^{3}} \tag{5.38}
\end{equation*}
$$

and

$$
a_{n}= \begin{cases}\frac{1}{n^{2}} & N_{2 k-1}<N \leq N_{2 k} \\ 2^{-n^{4}} & N_{2 k}<n \leq N_{2 k+1}\end{cases}
$$

Our goal will be to prove that

$$
\begin{equation*}
d \nu_{2^{N_{2 k}^{2}}} \rightarrow d \rho_{[-2,0]} \quad d \nu_{2^{N_{2 k+1}^{2}}} \rightarrow d \rho_{[-2,2]} \tag{5.39}
\end{equation*}
$$

Intuitively, for $m=2^{N_{2 k+1}^{2}}$, the measures at level $1 / m$ will be uniformly spaced out (on an exponential scale), so by the Stahl-Totik theorem, the zeros will want to look like the equilibrium measure for $[-2,2]$. But for $m=2^{N_{2 k}^{2}}$, most intervals of size $1 / m$ in $[0,2]$ will have tiny measure, so the zeros will want to almost all lie on $[-2,0]$, where the best strategy for these ( to minimize $\int P_{m}^{2} d \mu$ ) will be to approximate the equilibrium measure for $[-2,0]$.

As a preliminary, we will show

$$
\begin{array}{ll}
\lim \sup \left\|P_{n}\right\|^{1 / n}=1 & \lim \inf \left\|P_{n}\right\|^{1 / n}=\frac{1}{2} \\
\lim _{k \rightarrow \infty}\left\|P_{2^{N_{2 k}^{2}}}\right\|^{1 / 2^{N_{2 k}^{2}}}=\frac{1}{2} & \lim _{k \rightarrow \infty}\left\|P_{2^{N_{2 k+1}^{2}}}\right\|_{1 / 2^{N_{2 k+1}^{2}}}=1 \tag{5.41}
\end{array}
$$

We begin with

$$
\begin{equation*}
\limsup \left\|P_{n}\right\|^{1 / n} \leq 1 \quad \liminf \left\|P_{n}\right\|^{1 / n} \geq \frac{1}{2} \tag{5.42}
\end{equation*}
$$

The first is immediate from (1.15) and $C([-2,2])=1$; the second from $\left\|P_{n}(x, d \mu)\right\| \geq\left\|P_{n}(x, d \mu \upharpoonright[-2,0])\right\|$ (by (1.12)), regularity of $d \mu \upharpoonright[-2,0]$, and $C([-2,0])=\frac{1}{2}$.

Next, we turn to

$$
\begin{equation*}
\limsup _{k \rightarrow \infty}\left\|P_{2^{N_{2 k}^{2}}}\right\|^{1 / 2^{N_{2 k}^{2}}} \leq \frac{1}{2} \tag{5.43}
\end{equation*}
$$

Let $T_{n}(x ;[-2,0])$ be the Chebyshev polynomials for $[-2,0]$ (which are just affinely related to the classic Chebyshev polynomials of the first kind) and let

$$
\begin{equation*}
Q_{2^{N_{2 k}^{2}}}(x)=T_{2^{N_{2 k}^{2}-N_{2 k}}}(x ;[-2,0]) \prod_{\ell=1}^{2^{N_{2 k}}}\left(x-\frac{\ell}{2^{N_{2 k}-1}}\right) \tag{5.44}
\end{equation*}
$$

so by (1.12),

$$
\begin{equation*}
\left\|P_{2^{N_{2 k}^{2}}}\right\| \leq\left\|Q_{2^{N_{2 k}^{2}}}\right\| \leq 1+2 \tag{5.45}
\end{equation*}
$$

where 1 is the contribution of the integral from $[-2,0]$ and 2 from $(0,2)$.

Since $\cos \ell x=2^{\ell-1}(\cos x)^{\ell}+$ lower order and the average of $\cos ^{2} x$ is $\frac{1}{2}$, for any $[a, b]$,

$$
\begin{equation*}
\left\|T_{m}(x ;[a, b])\right\|=\sqrt{2} C([a, b])^{m} \tag{5.46}
\end{equation*}
$$

where the norm is over $L^{2}\left(\mathbb{R}, d \rho_{[a, b]}\right)$. Since the product in (5.44) is bounded by $4^{2^{N_{2 k}}}$ on $[-2,2]$, we have

$$
\begin{equation*}
1 \leq \sqrt{2}\left(\frac{1}{2}\right)^{2^{N_{2 k}^{2}-N_{2 k}}} 4^{2_{2 k}^{N}} \tag{5.47}
\end{equation*}
$$

On the other hand, there is a constant $K$ so

$$
\begin{equation*}
\left\|T_{m}(x ;[-2,0])\right\|_{L^{\infty}([-2,2])} \leq K^{m} \tag{5.48}
\end{equation*}
$$

and the product in (5.44) kills all the pure points up to level $N_{2 k}$ :

$$
\begin{aligned}
\boxed{2} & \leq K^{2^{N_{2 k}}} \sum_{n=N_{2 k}}^{\infty} a_{n} \\
& \leq K^{2^{N_{2 k}}}\left[N_{2 k+1} 2^{-N_{2 k}^{4}}+\left(N_{2 k+1}\right)^{-1}\right]
\end{aligned}
$$

is much smaller than the right side of (5.47) for $k$ large. Thus, by (5.47),

$$
\begin{equation*}
\lim \sup (\sqrt[1]{1}+\sqrt[2]{ })^{1 / 2^{N_{2 k}^{2}}} \leq \frac{1}{2} \tag{5.49}
\end{equation*}
$$

proving (5.43).

In verifying (5.40), we finally prove that

$$
\begin{equation*}
\liminf _{k \rightarrow \infty}\left\|P_{2^{N_{2 k+1}^{2}}}\right\|^{1 / 2^{N_{2 k+1}^{2}}} \geq 1 \tag{5.50}
\end{equation*}
$$

By the fact that the Chebyshev polynomials for $[-2,2]$ obey

$$
\begin{equation*}
T_{n}(2 \cos x ;[-2,2])=2 \cos n x \tag{5.51}
\end{equation*}
$$

has $\left\|T_{n}\right\|_{L^{\infty}([-2,2])}=2$, we see

$$
\begin{equation*}
\left\|P_{m}\right\|_{L^{\infty}([-2,2])} \geq 2 \tag{5.52}
\end{equation*}
$$

By Markov's inequality (5.8), we have

$$
\begin{equation*}
\left\|P_{m}^{\prime}\right\|_{L^{\infty}([-2,2])} \leq \frac{m^{2}}{2}\left\|P_{m}\right\|_{L^{\infty}([-2,2])} \tag{5.53}
\end{equation*}
$$

so there is an interval of size $4 / m^{2}$ where $P_{m}(x) \geq 1$, that is,

$$
\begin{equation*}
\left\|P_{m}(x)\right\|_{L^{2}(d \mu)}^{2} \geq \inf _{y \in[-2,2]} \mu\left(\left[y-\frac{2}{m^{2}}, y+\frac{2}{m^{2}}\right]\right) \tag{5.54}
\end{equation*}
$$

which implies that

$$
\begin{equation*}
\left\|P_{2^{N_{2 k+1}^{2}}}\right\|_{L^{2}(\mu)} \geq a_{2^{N_{2 k+1}^{2}}} 2^{-N_{2 k+1}^{2}} \tag{5.55}
\end{equation*}
$$

so it is bounded from below by a power of $2^{-N_{2 k+1}^{2}}$. Since $m^{-\ell / m} \rightarrow 1$ for any fixed $\ell$, we obtain (5.50).

Clearly, (5.42), (5.43), and (5.50) imply (5.40) and (5.41). We now only need to go from there to results on limits of $d \nu_{n}$. By Theorem 2.4, the second equality in (5.41) implies the second limit result in (5.39). By Theorem 4.2, the first equality in (5.41) implies the first limit result in (5.39).

Example 5.9. Here is an example of a measure $d \mu$ on $[0,1]$ where the density of zeros has a limit singular relative to the equilibrium measure for $[0,1]$. Such examples are discussed in [105] and go back to work of Ullman. Let $\Sigma$ be the classical Cantor set and $d \rho_{\Sigma}$ its equilibrium measure. Let

$$
\begin{equation*}
d \mu=d \rho_{\Sigma}+\sum_{n=1}^{\infty} 2^{-n^{4}}\left(\sum_{j=0}^{2^{n-1}} \frac{1}{2^{n}} \delta_{(2 j+1) / 2^{n}}\right) \tag{5.56}
\end{equation*}
$$

As in the above construction, one shows $\left\|P_{n}\right\|^{1 / n} \rightarrow C(\Sigma)$ and then Theorem 4.2 implies that $d \nu_{n} \rightarrow d \rho_{\Sigma}$ which is singular with respect to Lebesgue measure, and so relative to $d \rho_{[0,1]} \equiv d \rho_{\operatorname{supp}(d \mu)}$.

## 6. Structural Results

In this section, we will focus on the mutual regularity of related measures. There are three main theorems, all from Stahl-Totik [105]:

Theorem 6.1. Let $\mu, \eta$ be two measures of compact support whose supports are equal up to sets of capacity zero. If $\mu \geq \eta$ and $\eta$ is regular, then so is $\mu$.

Theorem 6.2. Let $\left\{E_{n}\right\}_{n=1}^{\infty}$ and $E_{\infty}$ be compact subsets of $\mathbb{C}$ so that $E_{\infty}$ and $\cup_{n=1}^{\infty} E_{n}$ agree up to sets of capacity zero and $C\left(E_{\infty}\right)>0$. Let $\mu$ be a measure with $\operatorname{supp}(d \mu)=E_{\infty}$ so that each $\mu \upharpoonright E_{j}$ which is nonzero is regular. Then $\mu$ is regular.

Remark. By $\mu \upharpoonright K$, we mean the measure

$$
\begin{equation*}
(\mu \upharpoonright K)(S)=\mu(K \cap S) \tag{6.1}
\end{equation*}
$$

To understand why the next theorem is so restrictive compared to Theorem 6.2, consider

Example 6.3. Let $E$ be the standard Cantor set in $[0,1]$. Let $\eta$ be a measure on $E$ which is not regular (see Corollary 5.6) and let

$$
\begin{equation*}
d \mu=d \eta+d x \upharpoonright[0,1] \tag{6.2}
\end{equation*}
$$

By Theorem 6.1, $d \mu$ is regular. But $d \mu \upharpoonright E=d \eta$ is not regular.
Theorem 6.4. Let $I=[a, b]$ be a closed interval with $I \subset E \subset \mathbb{R}$ and $E$ compact. Let $\mu$ be a regular measure with support in $E$ so $C(\operatorname{supp}(\mu \upharpoonright I))>0$. Then $\mu \upharpoonright I$ is regular.

Remarks. 1. We do not require that $\operatorname{supp}(d \mu)=E$ (nor that $I \subset$ $\operatorname{supp}(d \mu))$ but only that $\operatorname{supp}(d \mu) \subset E$ and that $\mu$ is regular in the sense that $C(\operatorname{supp}(d \mu))>0$ and $\left\|P_{n}(\cdot, d \mu)\right\|_{L^{2}(d \mu)}^{1 / n} \rightarrow C(\operatorname{supp}(d \mu))$.
2. The analog of the sets $I$ in [105] must have nonempty twodimensional interior. Our $I$ obviously has empty two-dimensional interior, but if $I=[a, b] \subset E \subset \mathbb{R}$ and if $D$ is the disk $\left\{z\left|\left|z-\frac{1}{2}(a+b)\right| \leq\right.\right.$ $\left.\frac{1}{2}|b-a|\right\}$, then $\mu \upharpoonright D=\mu \upharpoonright I$.

The proofs of Theorems 6.1 and 6.2 will be easy, but Theorem 6.4 will be nontrivial. Here are some consequences of these results:

Corollary 6.5. Let $\mu, \nu$ be two regular measures (with different supports allowed). Then their max, $\mu \vee \nu$, and sum, $\mu+\nu$, are regular.

Remark. See Doob [33] for the definition of $\mu \vee \nu$.

Proof. $\mu+\nu$ and $\mu \vee \nu$ have the same support and $\mu+\nu \geq \mu \vee \nu$ so, by Theorem 6.1, we only need the result for $\mu \vee \nu$. Let $E_{1}=\operatorname{supp}(\mu)$ and $E_{2}=\operatorname{supp}(\nu), \operatorname{sosupp}(\mu \vee \nu)=E_{1} \cup E_{2}$. By definition, $(\mu \vee \nu) \upharpoonright E_{1} \geq \mu$ and they have the same supports. So, by Theorem 6.1, $\mu \vee \nu \upharpoonright E_{1}$ is regular. Similarly, $\mu \vee \nu \upharpoonright E_{2}$ is regular. By Theorem $6.2, \mu \vee \nu$ is regular.

Corollary 6.6. Let $E=I_{1} \cup \cdots \cup I_{\ell}$ be a union of finitely many disjoint closed intervals. Let $\mu$ be a measure on $E$. Then $\mu$ is regular if and only if each $\mu \upharpoonright I_{j}$ is regular.
Proof. Immediate from Theorems 6.2 and 6.4.
Proof of Theorem 6.1. Since (1.12) holds,

$$
\begin{equation*}
\left\|X_{n}\right\|_{L^{2}(d \eta)} \leq\left\|X_{n}\right\|_{L^{2}(d \mu)} \tag{6.3}
\end{equation*}
$$

Given (1.15), we have (with $E=\operatorname{supp}(d \mu))$

$$
\lim \left\|X_{n}\right\|_{L^{2}(d \eta)}^{1 / n}=C(E) \Rightarrow \lim \left\|X_{n}\right\|_{L^{2}(d \mu)}^{1 / n}=C(E)
$$

Proof of Theorem 6.2. Let $\mu_{j}=\mu \upharpoonright E_{j}$ and let $x_{n}(z)$ be the $x_{n}$ 's for $d \mu$. Then

$$
\begin{equation*}
\left\|x_{n}\right\|_{L^{2}\left(d \mu_{j}\right)} \leq\left\|x_{n}\right\|_{L^{2}(d \mu)}=1 \tag{6.4}
\end{equation*}
$$

so

$$
\begin{equation*}
\left|x_{n}(z)\right| \leq \frac{\left|x_{n}(z)\right|}{\left\|x_{n}\right\|_{L^{2}\left(d \mu_{j}\right)}} \tag{6.5}
\end{equation*}
$$

By regularity and Theorem $1.10(\mathrm{vi})$, for q.e. $z \in E_{j}$ (using $G_{E_{j}}(z)=0$ for q.e. $z \in E_{j}$ by Theorem A.10(b) and (A.40)), we see for q.e. $x \in E_{j}$,

$$
\begin{equation*}
\limsup _{n \rightarrow \infty}\left|x_{n}(z)\right|^{1 / n} \leq 1 \tag{6.6}
\end{equation*}
$$

Since $\cup_{j=1}^{\infty} E_{j}$ is q.e. $E$, we have (6.6) q.e. on all of $E$. By (iii) $\Rightarrow$ (i) in Theorem 1.10, $\mu$ is regular.

To prove Theorem 6.4, we first make a reduction:
Proposition 6.7. Suppose there is $I=[a, b] \subset E \subset \mathbb{R}($ with $a<b)$, $\mu$ regular, $C(\operatorname{supp}(\mu \upharpoonright I))>0$ but

$$
\begin{equation*}
\liminf \left\|P_{n}(\cdot, \mu \upharpoonright I)\right\|_{L^{2}(\mu\lceil I)}^{1 / n}<C(\operatorname{supp}(\mu \upharpoonright I)) \tag{6.7}
\end{equation*}
$$

Then there exists a $\mu$, perhaps distinct but also regular and supported on $E$, so that (6.7) holds and

$$
\begin{equation*}
\lim \inf \left\|P_{n}(\cdot, \mu \upharpoonright I)\right\|_{L^{2}(\mu\lceil I)}^{1 / n}>0 \tag{6.8}
\end{equation*}
$$

Proof. If (6.8) holds for the initial $\mu$, we can stop. Otherwise, we will take

$$
\begin{equation*}
\tilde{\mu}=\mu+\rho_{F} \tag{6.9}
\end{equation*}
$$

where $F=\left[x_{0}-\delta, x_{0}+\delta\right] \subset I$ with $x_{0}=\frac{1}{2}(a+b)$ and $\delta$ sufficiently small chosen later.

By Corollary 6.5, $\tilde{\mu}$ is regular and, by (6.3),

$$
\left\|P_{n}(\cdot, \tilde{\mu} \upharpoonright I)\right\|_{L^{2}(\tilde{\mu} \mid I)}^{1 / n} \geq\left\|P_{n}\left(\cdot, \rho_{F}\right)\right\|_{L^{2}\left(\rho_{F}\right)}^{1 / n}
$$

so (6.8) holds since

$$
\lim _{n \rightarrow \infty}\left\|P_{n}\left(\cdot, \rho_{F}\right)\right\|_{L^{2}\left(\rho_{F}\right)}^{1 / n}=C(F)>0
$$

Thus, we need only prove that (6.7) holds for suitable $\delta$. Since we are supposing (6.8) fails for $\mu$, pick $n(j) \rightarrow \infty$ so

$$
\begin{equation*}
\left\|P_{n(j)}(\cdot, \mu \upharpoonright I)\right\|_{L^{2}(\mu\lceil I)}^{1 / n(j)} \rightarrow 0 \tag{6.10}
\end{equation*}
$$

Define

$$
\begin{equation*}
Q_{2 n(j)}(x)=P_{n(j)}(x, \mu \upharpoonright I)\left(x-x_{0}\right)^{n(j)} \tag{6.11}
\end{equation*}
$$

Let $d=\operatorname{diam}(E)$ and note that on $E$, since $x_{0} \in I \subset E$,

$$
\begin{equation*}
\left|x-x_{0}\right|^{n(j)} \leq d^{n(j)} \tag{6.12}
\end{equation*}
$$

and since $P_{n(j)}$ has all its zeros in $\operatorname{cvh}(E)$, for $x \in E$,

$$
\begin{equation*}
\left|P_{n(j)}(x)\right| \leq d^{n(j)} \tag{6.13}
\end{equation*}
$$

Thus,

$$
\begin{align*}
\left\|Q_{2 n(j)}\right\|_{L^{2}(\tilde{\mu} \mid I)}^{2} & =\left\|P_{n(j)}\left(\cdot-x_{0}\right)^{n(j)}\right\|_{L^{2}(\mu \mid I)}^{2}+\left\|P_{n(j)}\left(\cdot-x_{0}\right)^{n(j)}\right\|_{L^{2}\left(d \rho_{E}\right)}^{2} \\
& \leq d^{2 n(j)}\left\|P_{n(j)}\right\|_{L^{2}(\mu \mid I)}^{2}+d^{2 n(j)} \delta^{2 n(j)} \tag{6.14}
\end{align*}
$$

Using (6.10), we see

$$
\limsup \left\|Q_{2 n(j)}\right\|_{L^{2}(\tilde{\mu} \mid I)}^{1 / 2 n(j)} \leq d^{1 / 2} \delta^{1 / 2}<C(\operatorname{supp}(\mu \upharpoonright I)) \leq C(\operatorname{supp}(\tilde{\mu} \upharpoonright I))
$$

if we take $\delta$ small. Since (1.12),

$$
\begin{equation*}
\left\|P_{2 n(j)}(\cdot, \tilde{\mu} \upharpoonright I)\right\|_{L^{2}(\tilde{\mu} \mid I)} \leq\left\|Q_{2 n(j)}\right\|_{L^{2}(\tilde{\mu} \mid I)} \tag{6.15}
\end{equation*}
$$

we see $\tilde{\mu} \upharpoonright I$ is not regular.
Proof of Theorem 6.4. By Proposition 6.7, we can find $\mu$ so $\mu$ is regular, $\mu \upharpoonright I$ is not regular but for some $a>0\left(\right.$ with $\left.\mu_{I}=\mu \upharpoonright I\right)$,

$$
\begin{equation*}
\int_{I}\left|P_{n}\left(x, \mu_{I}\right)\right|^{2} d \mu \geq a^{n} \tag{6.16}
\end{equation*}
$$

Fix $x_{0} \in I_{\text {int }}$. Let $d=\operatorname{diam}(E)$ and for $\ell$ to be picked shortly, let

$$
\begin{equation*}
Q_{n(2 \ell+1)}(x)=P_{n}\left(x, \mu_{I}\right)\left(1-\frac{\left(x-x_{0}\right)^{2}}{d^{2}}\right)^{\ell n} \tag{6.17}
\end{equation*}
$$

Since $P$ obeys (6.13), if we define

$$
\begin{equation*}
\eta=\max _{x \in E \backslash I}\left(1-\frac{\left(x-x_{0}\right)^{2}}{d^{2}}\right)<1 \tag{6.18}
\end{equation*}
$$

we have for $x \notin I$,

$$
\begin{equation*}
\left|Q_{n(2 \ell+2)}(x)\right| \leq \eta^{\ell n} d^{n} \tag{6.19}
\end{equation*}
$$

Choose $\ell$ so

$$
\begin{equation*}
\left(\eta^{\ell} d\right)^{2}<a \tag{6.20}
\end{equation*}
$$

Then, by (6.17), (6.19), and (6.20),

$$
\begin{equation*}
\int_{K}\left|Q_{n(2 \ell+1)}(x)\right|^{2} d \mu \leq 2 \int\left|P_{n}\left(x, \mu_{I}\right)\right|^{2} d \mu \tag{6.21}
\end{equation*}
$$

so

$$
\begin{equation*}
\left|p_{n}\left(x_{0} ; d \mu_{I}\right)\right|^{1 / n} \leq \frac{2^{1 / 2 n}\left|Q_{n(2 \ell+1)}\left(x_{0}\right)\right|^{1 / n}}{\left\|Q_{n(2 \ell+1)}\right\|_{L^{2}(K, d \mu)}^{1 / n}} \tag{6.22}
\end{equation*}
$$

so, by (1.30) and regularity of $\mu$, for $\mu$-a.e. $x_{0}$ in $I^{\text {int }} \cap \operatorname{supp}(d \mu)$,

$$
\lim \sup \left|p_{n}\left(x_{0} ; d \mu_{I}\right)\right|^{1 / n} \leq 1
$$

But then Theorem 1.10 implies $d \mu_{I}$ is regular. This contradiction proves the theorem.

## 7. Ergodic Jacobi Matrices and Potential Theory

In this section, we will explore regularity ideas for ergodic half- and whole-line Jacobi operators and see this is connected to Kotani theory (see [99, Sect. 10.11] and [27] as well as the original papers [69, 96, 42, 70]). A main goal is to prove Theorems 1.15 and 1.16.

Let $(\Omega, d \sigma)$ be a probability measure space. Let $T: \Omega \rightarrow \Omega$ be an invertible ergodic transformation. Let $\tilde{A}, \tilde{B}$ be measurable functions from $\Omega$ to $\mathbb{R}$ with $\tilde{B}$ bounded, $\tilde{A}$ positive, and both $\tilde{A}$ and $\tilde{A}^{-1}$ bounded. For $\omega \in \Omega$ and $n \in \mathbb{Z}$, define $a_{n}(\omega), b_{n}(\omega)$ by

$$
\begin{equation*}
a_{n}(\omega)=\tilde{A}\left(T^{n} \omega\right) \quad b_{n}(\omega)=\tilde{B}\left(T^{n} \omega\right) \tag{7.1}
\end{equation*}
$$

By $J(\omega)$, we mean the Jacobi matrix with parameters $\left\{a_{n}(\omega), b_{n}(\omega)\right\}_{n=1}^{\infty}$. By $\tilde{J}(\omega)$, we mean the two-sided Jacobi matrix with parameters $\left\{a_{n}(\omega), b_{n}(\omega)\right\}_{n=-\infty}^{\infty}$. Occasionally we will use $J_{k}^{+}(\omega)$ for the one-sided matrix with parameters $\left\{a_{k+n}(\omega), b_{k+n}(\omega)\right\}_{n=1}^{\infty}$ and $J_{k}^{-}(\omega)$ for the one-sided matrix with parameters $\left\{a_{k-n}(\omega), b_{k+1-n}(\omega)\right\}_{n=1}^{\infty}$.

Spectral measures for one-sided matrices (and vector $\delta_{1}$ ) are $d \mu_{\omega}, d \mu_{\omega}^{ \pm k}$ and for $\tilde{J}(\omega)$, we use $d \tilde{\mu}_{\omega ; k}$ for vector $\delta_{k}$.

For spectral theory, the transfer matrix is basic. Define for $n \in \mathbb{Z}$,

$$
A_{n}(x, \omega)=\frac{1}{a_{n+1}(\omega)}\left(\begin{array}{cc}
x-b_{n+1}(\omega) & -1  \tag{7.2}\\
a_{n+1}(\omega)^{2} & 0
\end{array}\right)
$$

then

$$
\begin{equation*}
a_{n+1} u_{n+1}+\left(b_{n+1}-x\right) u_{n}+a_{n} u_{n-1}=0 \tag{7.3}
\end{equation*}
$$

is equivalent to

$$
\begin{equation*}
\binom{u_{n+1}}{a_{n+1} u_{n}}=A_{n}\binom{u_{n}}{a_{n} u_{n-1}} \tag{7.4}
\end{equation*}
$$

We define for $n<m$,

$$
\begin{equation*}
T(m, n ; x, \omega)=A_{m}(x, \omega) A_{m-1}(x, \omega) \ldots A_{n+1}(x, \omega) \tag{7.5}
\end{equation*}
$$

and $T(n, n ; x, \omega)=1$ and, for $m<n, T(m, n ; x, \omega)=T(n, m ; x, \omega)^{-1}$. Thus, solutions of (7.3) obey

$$
\begin{equation*}
\binom{u_{m+1}}{a_{m+1} u_{m}}=T(m, n ; x, \omega)\binom{u_{n+1}}{a_{n+1} u_{n}} \tag{7.6}
\end{equation*}
$$

In particular, for $n \geq 1$,

$$
\begin{equation*}
\binom{p_{n+1}(x, \omega)}{a_{n+1} p_{n}(x, \omega)}=T(n,-1 ; x, \omega)\binom{1}{0} \tag{7.7}
\end{equation*}
$$

The ergodic and subadditive ergodic theorems produce the following well-known facts:

Theorem 7.1. There exists $\Omega_{0} \subset \Omega$ of full $\sigma$ measure so that for $\omega \in \Omega_{0}$,
(a) $\sigma(\tilde{J}(\omega))=E$, a fixed perfect subset of $\mathbb{R}$ independent of $\omega\left(\right.$ in $\left.\Omega_{0}\right)$. Moreover, for any $\omega \in \Omega_{0}$, each $J_{k}^{ \pm}$obeys

$$
\begin{equation*}
\sigma_{\mathrm{ess}}\left(J_{k}^{ \pm}(\omega)\right)=E \tag{7.8}
\end{equation*}
$$

(b) There is a measure $d \nu_{\infty}$ with

$$
\begin{equation*}
\operatorname{supp}\left(d \nu_{\infty}\right)=E \tag{7.9}
\end{equation*}
$$

If $d \nu_{n}^{k, \pm, \omega}$ is the zero counting measure for $J_{k}^{ \pm}(\omega)$, then for any $\omega \in \Omega_{0}$, as $n \rightarrow \infty$,

$$
\begin{equation*}
d \nu_{n}^{k, \pm, \omega} \xrightarrow{\omega} d \nu_{\infty} \tag{7.10}
\end{equation*}
$$

(c) Define the Lyapunov exponent $\gamma(z)$ for $z$ by

$$
\begin{equation*}
\gamma(z)=\lim _{n \rightarrow \infty} \mathbb{E}\left(\frac{1}{n} \log \|T(n-1,-1 ; z, \omega)\|\right) \tag{7.11}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathbb{E}(f)=\int f(\omega) d \sigma(\omega) \tag{7.12}
\end{equation*}
$$

and (7.11) includes that the limit exists. Moreover, for any $k \in \mathbb{Z}$,

$$
\begin{align*}
& \gamma(z)=\lim _{n \rightarrow \infty} \mathbb{E}\left(\frac{1}{n} \log \|T(n+k, k ; z, \omega)\|\right)  \tag{7.13}\\
& \gamma(z)=\lim _{n \rightarrow \infty} \mathbb{E}\left(\frac{1}{n} \log \|T(k-n, k ; z, \omega)\|\right) \tag{7.14}
\end{align*}
$$

(d) For any $\omega \in \Omega_{0}$ and $z \notin E$ and $k$ fixed,

$$
\begin{align*}
\lim _{n \rightarrow \infty} \frac{1}{n} \log \|T(n+k, k ; z, \omega)\| & =\gamma(z)  \tag{7.15}\\
\lim _{n \rightarrow \infty} \frac{1}{n} \log \|T(k-n, k ; z, \omega)\| & =\gamma(z) \tag{7.16}
\end{align*}
$$

(e) For any $z \in E$ and $\sigma$-a.e. $\omega \in \Omega_{0}$, (7.15) and (7.16) hold.
(f) For $\omega \in \Omega_{0}$, $\lim _{n \rightarrow \infty}\left(a_{1} \ldots a_{n}\right)^{1 / n}=A$ exists and is $\omega$-independent, and one has the Thouless formula,

$$
\begin{equation*}
\gamma(z)=\log \left(A^{-1}\right)+\int \log (|z-x|) d \nu_{\infty}(x) \tag{7.17}
\end{equation*}
$$

Moreover, for all z,

$$
\begin{equation*}
\gamma(z) \geq 0 \tag{7.18}
\end{equation*}
$$

Remarks. 1. For proofs, see [21, 26, 87, 99]. (7.17) is due (in the physics literature) to Herbert-Jones [51] and Thouless [110]. It is, of course, just (1.4) for $z \notin \operatorname{cvh}(\Sigma)$. Almost everything else here is a simple consequence of the Birkhoff ergodic theorem/the Kingman subadditive ergodic theorem and translation invariance which implies, for example, that the expectation in (7.13) is $k$-independent for each $n$.
2. There are two subtleties to OP readers. First, (7.14) comes from $\left\|A^{-1}\right\|=\|A\|$ for $2 \times 2$ matrices $A$ with $\operatorname{det}(A)=1$. It implies that the Lyapunov exponent is the same in both directions. $\operatorname{det}(T)=1$ also implies (7.18).
3. The second subtlety concerns equality in (7.17) for all $z$, including those in $\Sigma$. This was first proven by Avron-Simon [12]; the simplest proof is due to Craig-Simon [25] who were motivated by work of Herman [52]. The point is that, in general, $\lim \sup \frac{1}{n} \log \|T(n+k, k ; z, \omega)\|$ (and $\limsup \frac{1}{n} \log \left|p_{n}(z, \omega)\right|$ ) may not be upper semicontinuous but $\mathbb{E}\left(\frac{1}{n} \log \|T(n+k, k ; z, \omega)\|\right)$ is because of translation invariance, Hölder's inequality, and

$$
\begin{equation*}
T(n+\ell+k, k ; z, \omega)=T\left(n+\ell+k, \ell+k ; z, T^{\ell} \omega\right) T(\ell+k, k ; z, \omega) \tag{7.19}
\end{equation*}
$$

This implies that the expectation is subadditive so the limit is an inf.

Two main examples are the Anderson model and almost periodic functions. For the former, $\left(a_{n}(\omega), b_{n}(\omega)\right)$ are independent $(0, \infty) \times \mathbb{R}$ valued (bounded with $a_{n}^{-1}$ also bounded) identically distributed random variables. In the almost periodic case, $\Omega$ is a finite- or infinitedimensional torus with $d \sigma$ Haar measure and $\tilde{A}, \tilde{B}$ continuous functions. A key observation (of Avron-Simon [12]) is that in this almost periodic case, the density of states exists for all, not only a.e., $\omega \in \Omega$ so we can then take $\Omega_{0}=\Omega$ in Theorem 7.1.

Here is the first consequence of potential theory ideas in this setting:
Theorem 7.2. E has positive capacity; indeed,

$$
\begin{equation*}
C(E) \geq A \tag{7.20}
\end{equation*}
$$

Moreover, $E$ is always potentially perfect (as defined in Appendix A). Each $d \mu_{\omega}\left(\omega \in \Omega_{0}\right)$ is regular if and only if equality holds in (7.20).

Proof. Use $\gamma \geq 0$ in (7.17), integrating $d \nu_{\infty}$, to see that

$$
\begin{equation*}
\mathcal{E}\left(\nu_{\infty}\right) \leq \log \left(A^{-1}\right) \tag{7.21}
\end{equation*}
$$

so $\mathcal{E}\left(\nu_{\infty}\right)<\infty$, implying $C(E)>0$. By (7.21), we get (7.20).
By (7.21), $\nu_{\infty}$ has finite energy and so, by Proposition A.6, $\nu_{\infty}$ gives zero weight to any set of capacity zero. It follows that if $x \in \operatorname{supp}\left(d \nu_{\infty}\right)$, then $C((x-\delta, x+\delta))>0$ for all $\delta$. By (7.9), $E$ is potentially perfect.

By definition of $A$, regularity for all $\omega \in \Omega_{0}$ is equivalent to $C(E)=$ A.

Proof of Theorem 1.15. We will prove that (1.37) holds. By (7.21), $\nu_{\infty}$ has finite Coulomb energy, so $\nu_{\infty}$ gives zero weight to sets of zero capacity. Since equality holds in (A.23), q.e. on $E$, we conclude that

$$
\begin{align*}
\log \left(C(E)^{-1}\right) & =\int d \nu_{\infty}(x) \Phi_{\rho_{E}}(x) \\
& =\int d \rho_{E}(x) \Phi_{\nu_{\infty}}(x)  \tag{A.2}\\
& =\int d \rho_{E}(x)\left[\log \left(A^{-1}\right)-\gamma(x)\right] \tag{7.17}
\end{align*}
$$

This is (1.37).
By (1.37), we have (1.35) $\Leftrightarrow \int \gamma(x) d \rho_{E}(x)=0$ which, given that $\gamma(x) \geq 0$, holds if and only if $\gamma(x)=0$ for $\rho_{E}$-a.e. $x$.

If (1.35) holds, then each $d \mu_{\omega}$ is regular, so by Theorem 1.7, $d \nu_{\infty}=$ $d \rho_{E}$. The converse part follows from Theorem 1.9.

Note: Remling remarked to me that Theorem 1.15 has a deterministic analog with essentially the same proof.

Kotani theory says something about when $\gamma(x)=0$ but we have not succeeded in making a tight connection, so we will postpone the precise details until we discuss conjectures in the next section. As a final topic, we want to prove Theorem 1.16 and a related result.

Proof of Theorem 1.16. By (4.14) for $d \mu_{\omega}$-a.e. $x$, we have

$$
\begin{equation*}
\limsup \left|p_{n}(x)\right|^{1 / n} \leq 1 \tag{7.22}
\end{equation*}
$$

On the other hand, by the upper envelope theorem and (1.4) for q.e. $x$,

$$
\begin{align*}
\lim \left|p_{n}(x)\right|^{1 / n} & =A^{-1} \exp \left(-\Phi_{\nu_{\infty}}(x)\right)  \tag{7.23}\\
& =\exp (\gamma(x)) \tag{7.24}
\end{align*}
$$

by (7.17). Let $Q_{\omega}$ be the capacity zero set where (7.24) fails.
On $S, \exp (\gamma(x))>1$, so since (7.22) holds for a.e. $x$, we have $d \mu_{\omega}(S \backslash$ $\left.Q_{\omega}\right)=0$ as claimed.

Remark. All we used was that $d \nu_{\infty}$ is the limit of $d \nu_{n}$, so this holds for all $\omega \in \Omega_{0}$. In particular, in the almost periodic case, it holds for all $\omega$ in the hull.

One is also interested in the whole-line operator.
Theorem 7.3. Let $\tilde{J}(\omega)$ be the whole-line Jacobi matrix associated with $\left\{a_{n}(\omega), b_{n}(\omega)\right\}_{n=-\infty}^{\infty}$ and $d \mu_{\omega, k}$ its spectral measures. Let $S \subset \mathbb{R}$ be the Borel set of $x$ with $\gamma(x)>0$. Then for each $\omega \in \Omega_{0}$, there exists a set $\tilde{Q}_{\omega}$ of capacity zero so that

$$
\begin{equation*}
\mu_{\omega, k}\left(S \backslash \tilde{Q}_{\omega}\right)=0 \tag{7.25}
\end{equation*}
$$

for all $k$.
Proof. By (7.7), the transfer matrix $T(n,-1 ; x, \omega)$ has matrix elements given by $p_{n+1}, p_{n}$ and the second kind polynomials $q_{n+1}, q_{n}$. As in the last proof, there is a set $\tilde{Q}_{\omega}^{(1)}$ of capacity zero so for $x \notin Q_{\omega}^{(1)}$,

$$
\begin{equation*}
\lim \left|p_{n}(x)\right|^{1 / n}=\exp (\gamma(x)) \tag{7.26}
\end{equation*}
$$

and (zeros of $p_{n}$ and $q_{n}$ interlace, so the zero counting measure for $q_{n}$ also converges to $d \nu_{\infty}$ )

$$
\begin{equation*}
\lim \left|q_{n}(x)\right|^{1 / n}=\exp (\gamma(x)) \tag{7.27}
\end{equation*}
$$

In particular, for $x \notin \tilde{Q}_{\omega}^{(1)}$, we have

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \frac{1}{n} \log \|T(n,-1 ; x, \omega)\|=\gamma(x) \tag{7.28}
\end{equation*}
$$

By the Ruelle-Osceledec theorem (see, e.g., [99, Sect. 10.5]), for any $w \neq 0 \in \mathbb{C}^{2}$, either

$$
\begin{equation*}
\|T(n,-1 ; x, \omega) w\|^{1 / n} \rightarrow e^{\gamma(x)} \tag{7.29}
\end{equation*}
$$

or

$$
\begin{equation*}
\|T(n,-1 ; x, \omega) w\|^{1 / n} \rightarrow e^{-\gamma(x)} \tag{7.30}
\end{equation*}
$$

Similarly, there is a set $\tilde{Q}_{\omega}^{(2)}$ of capacity zero with similar behavior as $n \rightarrow \infty$.

This says that every solution of (7.3) for $x \notin \tilde{Q}_{\omega}^{(1)} \cup \tilde{Q}_{\omega}^{(2)}$ either grows exponentially at $\pm \infty$ or decays exponentially. Thus, polynomial boundedness implies $\ell^{2}$ solutions. If $\tilde{Q}_{\omega}^{(3)}$ is the set of eigenvalues of $\tilde{J}(\omega)$ which is countable and so of capacity zero, and if $\tilde{Q}_{\omega}=\tilde{Q}_{\omega}^{(1)} \cup$ $\tilde{Q}_{\omega}^{(2)} \cup \tilde{Q}_{\omega}^{(3)}$, then

$$
\tilde{J}(\omega) u=x u \text { with } u \text { polynomially bounded } \Rightarrow x \in \tilde{Q}_{\omega}
$$

By the BGK expansion discussed in Section 4, this implies the spectral measures of $\tilde{J}(\omega)$ are supported on $\tilde{Q}_{\omega}$, that is, (7.25) holds.

Remarks. 1. The reader will recognize this proof as a slight variant of the Pastur-Ishii argument $[86,53]$ that proves absence of a.c. spectrum on $S$.
2. As above, in the almost periodic case, this holds for all $\omega$ in the hull.
3. This is the first result on zero Hausdorff dimension in this generality. But for suitable analytic quasi-periodic Jacobi matrices, the result is known; see Jitomirskaya-Last [59] and Jitomirskaya [56].

## 8. Examples, Open Problems, and Conjectures

Here we consider a number of illustrative examples and raise some open questions and conjectures. The conjectures are sometimes mere guesses and could be wrong. Indeed, when I started writing this paper, I had intended to make a conjecture for which a counterexample appears below as Example 8.12. So the reader should regard the conjectures as an attempt to stimulate work with my own guesses. I will try to explain my guesses, but they are not always compelling.

Example 8.1 (Random and Decaying Random OPRL). Let $\Omega=$ $\times_{n=1}^{\infty}[(0, \infty) \times \mathbb{R}]$ with $d \sigma\left(\left\{a_{n}, b_{n}\right\}\right)=\otimes d \eta\left(a_{n}, b_{n}\right)$, where $\eta$ is a measure of compact support on $(0, \infty) \times \mathbb{R}$. For each $\omega \in \Omega$, there is an associated Jacobi matrix, and we want results on $J(\omega)$ that hold for $\sigma$ a.e. $\omega$. The traditional Anderson model is the case where $a_{n} \equiv 1$ and $b_{n}$ is uniformly distributed on $[\alpha, \beta]$, that is, $d \eta(a, b)=\delta_{a 1} \frac{1}{\beta-\alpha} \chi_{(\alpha, \beta)}(b) d b$.

The decaying random model has two extra parameters, $\lambda \in(0, \infty)$ and $\gamma \in(0,1)$, takes $\tilde{a}_{n}(\omega) \equiv 1, \tilde{b}_{n}(\omega)$ the Anderson model with $\beta=-\alpha=1$, and takes

$$
\begin{align*}
b_{n}(\omega) & =\lambda n^{-\gamma} \tilde{b}_{n}(\omega)  \tag{8.1}\\
a_{n}(\omega) & =1 \tag{8.2}
\end{align*}
$$

The Anderson model is ergodic; the decaying random model is not. The Anderson model goes back to his famous work [2] with the first mathematical results by Kunz-Souillard [72] and the decaying model to Simon [95] (see also [67]).

For the Anderson model, it is known for a.e. $\omega$, $\sigma_{\text {ess }}(J(\omega))=[-2+$ $\alpha, 2+\beta]$, while for the decaying random model, $\sigma_{\text {ess }}(J(\omega))=[-2,2]$ by Weyl's theorem (i.e., $J(\omega)$ is in Nevai class). Clearly, $\left(a_{1} \ldots a_{n}\right)^{1 / n}=1$. For the Anderson model,

$$
\begin{equation*}
C\left(\sigma_{\mathrm{ess}}(J(\omega))\right)=\frac{1}{4}(4+(\beta-\alpha))>1 \tag{8.3}
\end{equation*}
$$

while for the decaying Anderson model,

$$
\begin{equation*}
C\left(\sigma_{\text {ess }}(J(\omega))\right)=1 \tag{8.4}
\end{equation*}
$$

so the former is not regular, while the latter is.
Of course, for the regular model, the density of zeros is the equilibrium measure where $\rho_{E}(x)=\frac{d \rho}{d x}=\frac{1}{\pi}\left(4-x^{2}\right)^{-1 / 2}$ by (A.34). For the Anderson model, on the other hand, $\frac{d \nu}{d x}$ is very different. It is $C^{\infty}$ even at the endpoints (by [103]) and decays exponentially fast to zero at the ends of the spectrum (Lifshitz tails; see [66]).

The Anderson model is known to have dense pure point spectrum and so is the decaying model if $\gamma<\frac{1}{2}$. It is known for the Anderson model (see [31]) that for some $\omega$-dependent labeling of the eigenvalues,

$$
\begin{equation*}
d \mu_{\omega}=\sum w_{n}(\omega) \delta_{e_{n}(\omega)} \tag{8.5}
\end{equation*}
$$

where for some $c>0$,

$$
\begin{equation*}
\left|w_{n}(\omega)\right| \leq e^{-c|n|} \tag{8.6}
\end{equation*}
$$

The same methods should allow one to prove for the decaying model on each $[-A, A] \subset(-2,2)$ that there is a labeling so that

$$
\begin{equation*}
\left|w_{n}(\omega)\right| \leq e^{-c|n|^{1-2 \gamma}} \tag{8.7}
\end{equation*}
$$

One expects that there are lower bounds of the same form and that the labels are such that the $e_{n}(\omega)$ are quasi-uniformly distributed (i.e., for $n \gg m$, the first $n e_{j}(\omega)$ are at least within $\frac{1}{m}$ of each point away from the edge of the spectrum). If these expectations are met, this example nicely illustrates Theorems 1.13 and 1.14.

In the not regular Anderson case, one expects $\mu_{\omega}\left(\left[\frac{j}{m}, \frac{j+1}{m}\right]\right) \sim e^{-c m}$ for fixed $c$, while in the regular decaying random model, one expects $\mu_{\omega}\left(\left[\frac{j}{m}, \frac{j+1}{m}\right]\right) \sim e^{-c m^{1-2 \gamma}}>C_{\eta} e^{-\eta m}$ for any $\eta$.

Example 8.2 (Generic Regular Measures). Fix $a_{n} \equiv 1,0<\gamma<\frac{1}{2}$, and let $\mathcal{B}=\left\{\left\{b_{n}\right\} \mid \lim n^{\gamma} b_{n} \rightarrow 0\right\}$ normed by $\left\|\left||b| \|=\sup _{n}\right| n^{\gamma} b_{n} \mid\right.$. It is known ([97]; see also $[77,22])$ that for a dense $G_{\delta}$ in $\mathcal{B}$, the associated Jacobi matrix has singular continuous measure. We believe there is some suitable sense in which a generic regular measure is singular continuous.

Example 8.3 (Almost Mathieu Equation). Perhaps the most studied model in spectral theory is the whole-line Jacobi matrix with $a_{n} \equiv 1$ and

$$
\begin{equation*}
b_{n}=\lambda \cos (n \alpha+\theta) \tag{8.8}
\end{equation*}
$$

where $\lambda, \alpha, \theta$ are parameters with $\frac{\alpha}{\pi}$ irrational. (See [56] for a review on the state of knowledge.) We will use some of the most refined results and comment on whether they are needed for the main potential theoretic conclusions. We fix $\alpha, \lambda . \theta \in \Omega=[0,2 \pi)$ labels the hull of an almost periodic family.

It is known since Avron-Simon [12] that for $|\lambda|>2$, there is no a.c. spectrum for almost all $\theta$ (and by Kotani [71] and Last-Simon [75], for all $\theta$ ) and by [12] for $\alpha$ which are Liouville numbers (irrational but very well approximated by rationals) only singular continuous spectrum for all $\theta$. Jitomirskaya [55] proved that for $\alpha$ 's with good Diophantine properties and $|\lambda|>2$, there is dense pure point spectrum for a.e. $\theta$ (and there is also singular continuous spectrum for a dense set of $\theta$ 's [60]). On the other hand, Last [74] proved that for $|\lambda|<2$ and all irrational $\alpha$ that the spectrum is a.c. for almost all $\theta$ (now known for all $\theta$ [71, 75]). It is now known the spectrum in this region is purely a.c. (see $[9,8,6]$ ).

At the special point $\lambda=2$, it is known that for all irrational $\alpha$, the spectrum has measure zero $[74,10]$, and therefore for all irrational $\alpha$ and a.e. $\theta$, the spectrum is purely singular continuous [49].

An important special feature for our purposes is Aubry duality (found by Aubry [5]; proven by Avron-Simon [12]) that relates the Lyapunov exponent $\gamma(z)$ and integrated density of states, $k(E)=$ $\int_{-\infty}^{E} d \nu_{\infty}(x)$, for $\alpha$ fixed (they are $\theta$-independent) at $\lambda$ and $\frac{4}{\lambda}$. Making the $\lambda$-dependence explicit,

$$
\begin{equation*}
k\left(E, \frac{4}{\lambda}\right)=k\left(\frac{2 E}{\lambda}, \lambda\right) \quad \gamma\left(z, \frac{4}{\lambda}\right)=\gamma\left(\frac{2 z}{\lambda}, \lambda\right)+\log \left(\frac{\lambda}{2}\right) \tag{8.9}
\end{equation*}
$$

Kotani theory implies in the a.c. region (i.e., $\lambda<2$ ) that $\gamma(E)=0$ for a.e. $E \in \operatorname{spec}(J)$, and Bourgain-Jitomirskaya [19] proved continuity of $\gamma$. So using (8.9),

$$
\begin{equation*}
E \in \operatorname{spec}(J) \Rightarrow \gamma(E)=\max \left(0, \log \left(\frac{\lambda}{2}\right)\right) \tag{8.10}
\end{equation*}
$$

and, in particular, $\gamma(E)=0$ on the spectrum if $|\lambda| \leq 2$.
We thus have:
Theorem 8.4. The density of zeros for the almost Mathieu equation is the equilibrium measure for its spectrum. For $\lambda \leq 2$, the measures are regular (for all $\omega$ ). For $\lambda>2$, they are not regular since

$$
\begin{equation*}
\lim \left(a_{1} \ldots a_{n}\right)^{1 / n}=1<C\left(\operatorname{spec}\left(J_{0}(\omega)\right)\right)=\log \left(\frac{\lambda}{2}\right) \tag{8.11}
\end{equation*}
$$

Proof. By Theorem 1.14, the measures are regular if $\lambda \leq 2$ since $\gamma(E)=$ 0 on the spectrum. That the measure is the equilibrium measure even if $\lambda>2$ follows from (8.9) as does (8.11).

Remarks. 1. Thus we see an example where the density of zeros is the equilibrium measure even though $d \mu_{\omega}$ is not regular. Consistently with Theorem 2.5, $d \mu_{\omega}$ lives on a set of capacity 0 by Theorem 7.3.
2. If we knew a priori that $d \rho_{\sigma_{\text {ess }}\left(J_{\omega}\right)}$ were absolutely continuous, Kotani theory then would suffice for Theorem 8.4. But as it is, we need the continuity result of [19].
3. It should be an exceptional situation that $J(\omega)$ has some singular spectrum but the density of states is still $d \rho_{E}$. In particular, if there are separate regions in $\sigma(J)$ of positive capacity where $\gamma(x)=0$ and where $\gamma(x)>0$, the density of states cannot be $d \rho_{E}$ since, for it, $\gamma(x)$ is constant on $\operatorname{supp}\left(d \rho_{E}\right)$. For examples with such coexistent spectrum (some only worked for the continuum case), see [17, 18, 36, 37].

Example 8.5 (Rotation Invariant Anderson Model OPUC). Let $d \sigma_{0}$ be a rotation invariant measure on the disk, $\mathbb{D}$ (i.e., on $\overline{\mathbb{D}}$ with $\left.\sigma_{0}(\partial \mathbb{D})=0\right)$. Let $\sigma$ on $X_{j=0}^{\infty} \mathbb{D}$ be $\otimes_{j=0}^{\infty} d \sigma_{0}\left(z_{j}\right)$. The ergodic OPUC with Verblunsky coefficients $\alpha_{j}$ distributed by $\sigma$ is called the rotation invariant Anderson model, and it is discussed in [99, Sect. 12.6] and earlier in Teplyaev [109] and Golinskii-Nevai [47].

If

$$
\begin{equation*}
\int-\log (1-|z|) d \sigma_{0}(z)<\infty \tag{8.12}
\end{equation*}
$$

then

$$
\begin{equation*}
\lim _{n \rightarrow \infty}\left(\rho_{0} \ldots \rho_{n-1}\right)^{1 / n}=\exp \left(\int \log \left(1-|z|^{2}\right) d \sigma_{0}(z)\right)>0 \tag{8.13}
\end{equation*}
$$

If also

$$
\begin{equation*}
\int-\log |z| d \sigma_{0}(z)<\infty \tag{8.14}
\end{equation*}
$$

then, by a use of the ergodic theorem,

$$
\lim _{n \rightarrow \infty}\left|\alpha_{n}\right|^{1 / n}=1
$$

with probability 1. By a theorem of Mhaskar-Saff [82] (see [98, Thm. 8.1.1]), any limit point of the zero counting measure lives on $\partial \mathbb{D}$ so, by the ergodic theorem, $\nu_{n}$ has a limit $\nu_{\infty}$ on $\partial \mathbb{D}$.

By the rotation invariance of $\sigma_{0}$, the distribution of $\left\{\alpha_{j}\right\}$ is invariant under $\alpha_{j} \rightarrow e^{i(j+1) \theta} \alpha_{j}$. So the collection of measures is rotation invariant and thus, by ergodicity, $d \nu_{\infty}$ is rotation invariant, that is, it is $\frac{d \theta}{2 \pi}$. By the Thouless formula and (8.13),

$$
\gamma\left(e^{i \theta}\right)=-\int \log \left(1-|z|^{2}\right) d \sigma_{0}(z)>0
$$

so long as $\sigma_{0} \neq \delta_{z=0}$. This is constant on $\partial \mathbb{D}$.
Thus, this family of measures is not regular, but the density of zeros is the equilibrium measure for $\operatorname{supp}\left(d \mu_{\omega}\right)=\partial \mathbb{D}$. This is the simplest example of a nonregular measure for which the density of zeros is the equilibrium measure. As is proven in Theorem 12.6 .1 of [99], the measure is a pure point measure, so $d \mu_{w}$ is for a.e. $\omega$ supported on a countable set, so of zero capacity, consistent with Theorem 2.5.

Example 8.6 (Subshifts). This is a rich class of ergodic Jacobi matrices (with $a_{n} \equiv 1$ ), reviewed in [28] (see also [99, Sect. 12]). For many of them, it is known that $E \equiv \sigma(J)$ is a set of Lebesgue measure zero on which $\gamma(x)$ is everywhere 0 . By Theorem 1.15, $C(E)=1$ and a.e. $\omega$ has regular $d \mu_{\omega}$, so, in particular, $d \nu_{\infty}=d \rho_{E}$.

Notice that, by Craig's argument (see Theorem A.13), if $d \mu$ is any probability measure whose support, $E$, has measure zero, then $G(z)=$ $\int \frac{d \mu(y)}{y-z}$ has the form

$$
\begin{equation*}
G(z)=-\frac{1}{\sqrt{(z-a)(z-b)}} \prod_{j=1}^{\infty} \frac{\left(z-\lambda_{j}\right)}{\sqrt{\left(z-\ell_{j}\right)\left(z-u_{j}\right)}} \tag{8.15}
\end{equation*}
$$

where the gaps in $E$ are $\left(\ell_{j}, u_{j}\right)$ and $a=\inf \ell_{j}, b=\sup \mu_{j}$. This is so regular that we wildly make the following:

Conjecture 8.7. Any ergodic matrix that has a spectrum of measure zero has vanishing Lyapunov exponent on the spectrum; equivalently, $\gamma(x)>0$ for some $x \in \Sigma$ implies $|\Sigma|>0$. Such zero Lyapunov exponent examples would thus be regular.

We note that for analytic functions on the circle with irrational rotation, this result is known to be true [57], following from combining results from Bourgain [16] and Bourgain-Jitomirskaya [19]. Of two experts I consulted, one thought it was false and the other, "likely true but too little support to make it a conjecture." Fools rush in where experts fear to tread.

Open Question 8.8 (The Classical Cantor Set). Of course, one of the simplest of measure zero sets is the classical Cantor set. It would be a good first step to understand its "isospectral tori." Which whole-line Jacobi matrices have $\left\langle\delta_{0},\left(J_{0}-z\right)^{-1} \delta_{0}\right\rangle=(8.15)$ ? Are they regular? As suggested by Deift-Simon [30], are they mainly mutually singular? Are any or all almost periodic?
Conjecture 8.9 (Last's Conjecture). A little more afield from potential theory, but worth mentioning, is the conjecture of Last that any ergodic Jacobi matrix (whole- or half-line) with some a.c. spectrum is almost periodic. Does it help to consider the case where the spectrum is purely a.c.? We note that a result of Kotani [70] implies Last's conjecture if $a_{n}, b_{n}$ take only finitely many values.

And it links up to the next question:
Open Question 8.10 (Denisov-Rakhmanov Theorem). Let $E$ be an essentially perfect set, that is, for every $x \in E$ and $\delta>0, \mid(x-$ $\delta, x+\delta) \cap E \mid>0$. In [29], $E$ was called a DR set if any half-line Jacobi matrix with $\sigma_{\text {ess }}(J)=\Sigma_{\mathrm{ac}}(J)=E$ has a set of right limit points which is uniformly compact (and so the limits are all almost periodic). A classical theorem of Rakhmanov, as extended by Denisov (see [99, Ch. 9]), says that $[-2,2]$ is a DR set. Damanik-Killip-Simon [29] proved a number of $E$ 's, including those associated with periodic problems, are DR sets. Remling [90] recently proved any finite union of closed intervals is a DR set, and he remarks that it is possible to combine his methods with those of Sodin-Yuditskii [104] to prove that any homogeneous set in the sense of Carleson (see [104] for a definition) is a DR set.

Following this section's trend to make (foolhardy?) conjectures:
Conjecture 8.11. Any essentially perfect compact subset of $\mathbb{R}$ is a DR set.

A counterexample would also be very interesting. This is relevant to this paper because, as we have explained, Widom's theorem (Theorem 1.12) is a kind of poor man's DR condition.

Related to this: it would be interesting to find a proof of the almost periodicity of every reflectionless two-sided Jacobi matrix with spectrum a finite union of intervals that did not rely on the theory of meromorphic functions on a Riemann surface.

Example 8.12. Remling [90] has some interesting Jacobi matrices that are regular on $[-2,2]$, not in Nevai class, and have $\Sigma_{\mathrm{ac}}=[0,2]$. Related examples for Schrödinger operators appeared earlier in Molchanov [83].

## 9. Continuum Schrödinger Operators

The theory presented earlier was developed by the OP community dealing with discrete (i.e., difference) equations. The spectral theory community knows there are usually close analogies between difference and differential equations, so it is natural to ask about regularity ideas for continuum Schrödinger operators - a subject that does not seem to have been addressed before. We begin this exploration here. This is more a description of a research project than a final report. We will be discursive without proofs.

The first problem that one needs to address is that there is no natural potential theory for infinite unbounded sets. $\log |x-y|^{-1}$ is unbounded above and below so Coulomb energies can go to $-\infty$. Moreover, the natural measures are no longer probability measures. There is no reasonable notion of capacity, even of renormalized capacity. But at least sometimes there is a natural notion of equilibrium measure and equilibrium potential.

Consider $E=[0, \infty)$. We may not know the precise right question but we know the right answer: For $V=0$, the solutions of $-u^{\prime \prime}+V u=$ $\lambda u$ with $u(0)=0$ are $u(x)=C \sinh (x \sqrt{-\lambda})$, and so

$$
\begin{equation*}
\lim _{x \rightarrow \infty} \frac{\log |u(x)|}{x}=\sqrt{-\lambda} \tag{9.1}
\end{equation*}
$$

which must be the correct analog of the potential theorist Green's function. And there is a huge literature on continuum density of states, which for this case is

$$
\begin{equation*}
d \rho(\lambda)=\chi_{[0, \infty)}(\lambda)(\lambda)^{-1 / 2}(2 \pi)^{-1} d \lambda \tag{9.2}
\end{equation*}
$$

This comes from noting the eigenvalues on [0, 1] with $u(0)=u(L)=0$ boundary conditions are $\left(\frac{\pi n}{L}\right)^{2}, n=1,2, \ldots$. Here is a first attempt to find the right question.

It is the derivative of $\int \log |x-y|^{-1} d \mu(x)$ that is a Herglotz function, so we make

Definition. We say $d \nu$ is an equilibrium measure associated to a set $E \subset[a, \infty)$ for some $a$, if and only if there is a Herglotz function, $F_{E}(z)$, on $\mathbb{C}$ so that
(i) $\operatorname{Im} F(\lambda+i 0)$ is supported on $\lambda \in E$.
(ii) $\operatorname{Re} F(\lambda+i 0)=0$ for a.e. $\lambda \in E$.
(iii) $F(\lambda) \rightarrow 0$ as $\lambda \rightarrow-\infty$.
(iv) $\pi^{-1} F(\lambda+i \varepsilon) d \lambda \xrightarrow{w} d \nu(\lambda)$
(v) For any bounded connected component $(a, b)$ of $\mathbb{R} \backslash E$, we have

$$
\begin{equation*}
\int_{a}^{b} F(\lambda) d \lambda=0 \tag{9.3}
\end{equation*}
$$

We will say $d \nu$ is normalized if

$$
\begin{equation*}
F(\lambda) \sim \frac{1}{2}(-\lambda)^{-1 / 2}(1+o(1)) \tag{9.4}
\end{equation*}
$$

near $-\infty$.
The reason for choosing (9.3) and (9.4) will be made clear shortly. Once we have $F$, we define the equilibrium potential of $E$ by

$$
\begin{equation*}
\Phi_{E}(z)=\operatorname{Re}\left(\int_{x_{0}}^{z} F(\omega) d \omega\right) \tag{9.5}
\end{equation*}
$$

where $x_{0} \in E$ and the integral is in a path in $\mathbb{C} \backslash[a, \infty)$ with $a=$ $\inf \{y \in \mathbb{R} \mid y \in E\}$. That $\operatorname{Re} F=0$ on $E$ and that (9.3) holds show $\Phi_{E}$ is independent of $x_{0}$. (9.3) also implies $\Phi_{E}(z)=0$ on $E$. For this reason, we need to take $E=\sigma_{\text {ess }}\left(-\frac{d^{2}}{d x^{2}}+V\right)$, not $\sigma\left(-\frac{d^{2}}{d x^{2}}+V\right)$.

With (9.4), we have

$$
\begin{equation*}
\Phi_{E}(z)=\operatorname{Re}(\sqrt{-z})(1+o(1)) \tag{9.6}
\end{equation*}
$$

near $-\infty$. We can explain why we normalize as we do. For regular situations, we expect that the absolute value of the eigenfunction, $\psi_{z}(x)$, analogous to OPs (see below) are asymptotic to

$$
\exp \left(x \Phi_{E}(z)\right)
$$

as $x \rightarrow \infty$. This, in turn, is related to integrals of the negative of the real part of

$$
\begin{equation*}
m(z, x)=\frac{\eta_{z}^{\prime}(x)}{\eta_{z}(x)} \tag{9.7}
\end{equation*}
$$

where $\eta$ is the solution of $L^{2}$ at infinity.

It is a result of Atkinson [4] (see also [45]) that in great generality that as $|z| \rightarrow \infty,-\frac{d^{2}}{d x^{2}}+V$ is bounded from below in sectors about $(-\infty, a)$ and, in general, in sectors $|\arg z| \in(c, \pi-\varepsilon)$,

$$
\begin{equation*}
m(z, \lambda)=-\sqrt{-z}+o(1) \tag{9.8}
\end{equation*}
$$

$\psi$ should grow as the inverse of $\eta$, so $\Phi \sim-m$ as $z \rightarrow-\infty$.
This is stronger than (9.6) (if one can interchange limits $x \rightarrow \infty$ and $z \rightarrow \infty$ ) since the error in (9.6) is $o(1) \sqrt{-z}$, while in (9.7) it is $o(1)$. The lack of a constant term is an issue to be understood.

If we take $E=[0, \infty)$ since $F^{\prime}>0$ on $(-\infty, 0)$, we have $F>0$ on $(-\infty, 0)$, and so $\log F(x+i 0)$ has boundary values 0 on $(-\infty, 0)$ and $\frac{1}{2}$ on $(0, \infty)$. This plus $\log F(z)=o(z)$ at $-\infty$ uniquely determine $\log F$, and so $F$, up to an overall constant which is fixed by the normalization yielding

$$
\begin{equation*}
F(z)=\frac{1}{2 \sqrt{-z}} \tag{9.9}
\end{equation*}
$$

so there is a unique "potential" for $[0, \infty)$ that gives the right $\Phi(z)=$ $\sqrt{-z}$.

Similarly, for a finite number of gaps removed from $[0, \infty)$, one gets a unique $F$. Craig's argument yields $F$ up to positions of zeros in the gap, which are then fixed by (9.1).

Open Project 9.1. Develop a formal theory of equilibrium measures and equilibrium potentials for unbounded sets that are "close" to $[0, \infty)$ (e.g., one might require that $E \backslash[0, \infty)$ has finite Lebesgue measure). Can one understand the $o(1)$ in (9.8) from this theory?

With potentials in hand, we can define regularity. We recall first that given any $V$ on $[0, \infty)$ which is locally in $L^{1}$, one can define the regular solution, $\psi(x, z)$, obeying

$$
\begin{gather*}
-\psi^{\prime \prime}(x, z)+V(x) \psi(x, z)=z \psi(x, z)  \tag{9.10}\\
\psi(0, z)=0 \quad \psi^{\prime}(0, z)=1 \tag{9.11}
\end{gather*}
$$

Here $\psi$ is $C^{1}$ (and so, locally bounded), its second distributional derivative is $L^{1}$, and obeys (9.10) as a distribution. For fixed $x, \psi$ is an entire function of $x$ of order $\frac{1}{2}$. If $\psi$ is not $L^{2}$ at infinity for (one and hence all) $z \in \mathbb{C}_{+}, V$ is called limit point at infinity and then there is a unique selfadjoint operator $H$ which is formally $-\frac{d^{2}}{d x^{2}}+V(x)$ with $u(0)=0$ boundary conditions. We only want to consider the case where $H$ is bounded below (which never happens if $V$ is not limit point). $\eta_{z}(x)$ is
then the solution $L^{2}$ at $\infty$ determined up to a constant, so

$$
\begin{equation*}
m(x, z)=\frac{\eta_{z}^{\prime}(x)}{\eta_{z}(x)} \tag{9.12}
\end{equation*}
$$

is determined by $V$.
Definition. Let $E=\sigma_{\text {ess }}(H)$. We say $H$ is regular if and only if for all $z \notin \sigma(H)$,

$$
\begin{equation*}
\limsup _{x \rightarrow \infty} \frac{1}{x} \log |\psi(x, z)|=\Phi_{E}(z) \tag{9.13}
\end{equation*}
$$

Of course, for this to make sense, $E$ has to be a set for which there is a potential. This will eliminate a case like $V(x)=x^{2}$ where $\sigma_{\text {ess }}(H)$ is empty). We expect the following should be easy to prove:

Metatheorem 9.2. (a) If $H$ is regular, for $z \notin \sigma(H)$, limsup in (9.13) can be replaced by lim.
(b) $H$ is regular if and only if (9.13) holds q.e. on $E$ (where $\Phi_{E}(z)$ is q.e. $=0$ ).
(c) If $H$ is regular, the density of states exists and equals the equilibrium measure for $E$.
(d) Conversely, if the density of states exists and equals the equilibrium measure for $E_{0}$, either $H$ is regular or else the spectral measure for $H$ is supported on a set of capacity zero.
(e) If $H$ is regular and

$$
\lim _{n \rightarrow \infty} \int_{n}^{n+1}|(\delta V)(x)| d x=0
$$

then $H+\delta V$ is also regular.
Remarks. 1. Here capacity zero and q.e. are defined in the usual way, that is, any probability measure of compact support contained in $E$ has infinite Coulomb energy.
2. By density of states, we mean the following (see $[11,13,54,61,65$, 84, 85]). Take $H_{L}$ to be the operator $-\frac{d^{2}}{d x^{2}}+V$ with $u(0)=0$ boundary conditions on $L^{2}([0, L], d x)$. This has infinite but discrete spectrum $E_{1, L}<E_{2, L}<E_{3, L}<\ldots$ (the solutions of $\psi(L, z)=0$ ). Let $d \nu_{L}$ be the infinite measure that gives weight $\frac{1}{L}$ to each $E_{j, L}$. If w-lim $d \nu_{L}$ (as functions on continuous functions of compact support) exist, we say the density of states exists and the limit is called the density of states.
3. (e) should follow from a standard use of an iterated DuHamel's formula.

Open Project 9.3. Verify Metatheorem 9.2 and explore, in particular, analogs of
(a) Widom's theorem, Theorem 1.12
(b) The Stahl-Totik criterion, Theorem 1.13
(c) For ergodic continuum Schrödinger operators, the analog of Theorem 1.16.

## Appendix A: A Child's Garden of Potential Theory in the Complex Plane

We summarize the elements of potential theory relevant to this paper. For lucid accounts of the elementary parts of the theory, see the appendix of Stahl-Totik [105], Martinez-Finkelshtein [81], and especially Ransford [88]. More comprehensive are Helms [50], Tsuji [113], and especially Landkof [73]. We will try to sketch some of the most important notions in remarks but refer to the texts, especially for the more technical aspects.

The two-dimensional Coulomb potential is $\log |x-y|^{-1}$ which has two lacks compared to the more familiar $|x-y|^{-1}$ of three dimensions: It is neither positive nor positive definite. We will deal with lack of positivity by only considering measures of compact support, and conditional positive definitiveness can replace positive definitiveness in some situations.

If $\mu$ is a positive measure of compact support on $\mathbb{C}$, its potential is defined by

$$
\begin{equation*}
\Phi_{\mu}(x)=\int \log |x-y|^{-1} d \mu(y) \tag{A.1}
\end{equation*}
$$

Because $\mu$ has compact support, $\log |x-y|^{-1}$ is bounded below for $x$ fixed, so if we allow the value $+\infty, \Phi_{\mu}$ is always well defined and Fubini's theorem is applicable and implies that for another positive measure, $\nu$, also of compact support, we have

$$
\begin{equation*}
\int \Phi_{\mu}(x) d \nu(x)=\int \Phi_{\nu}(x) d \mu(x) \tag{A.2}
\end{equation*}
$$

Sometimes it is useful to fix $M>0$ and define the cutoff

$$
\begin{equation*}
\Phi_{\mu}^{M}(x)=\int \log \left[\min \left(M,|x-y|^{-1}\right)\right] d \mu(y) \tag{A.3}
\end{equation*}
$$

$\Phi_{\mu}^{M}$ is continuous and $\Phi_{\mu}^{M}$ is an increasing sequence in $M$, so
Proposition A.1. $\Phi_{\mu}(x)$ is harmonic on $\mathbb{C} \backslash \operatorname{supp}(d \mu)$, lower semicontinuous on $\mathbb{C}$, and superharmonic there.

One might naively think that $\Phi_{\mu}(x)$ only fails to be continuous because it can go to infinity and that it is continuous in the extended sense - but that is wrong!

Example A.2. Let $x_{n}=-n^{-1}$ and let

$$
\begin{equation*}
d \mu=\sum_{n=1}^{\infty} n^{-2} \delta_{x_{n}} \tag{A.4}
\end{equation*}
$$

Then $\Phi_{\mu}\left(x_{n}\right)=\infty$ and $x_{n} \rightarrow 0$, but

$$
\begin{equation*}
\Phi_{\mu}(0)=\sum_{n=1}^{\infty} n^{-2} \log n<\infty \tag{A.5}
\end{equation*}
$$

Notice that this is consistent with lower semicontinuity, that is, $\Phi_{\mu}\left(\lim x_{n}\right) \leq \lim \inf \Phi_{\mu}\left(x_{n}\right)$. Also notice, given Hydrogen atom spectra, that this example is relevant to spectral theory.

Lest you think this kind of behavior is only consistent with unbounded $\Phi_{\mu}$, one can replace $\delta_{x_{n}}$ by a smeared out probability measure, $\eta_{n}$ (using equilibrium measures on a small interval, $I_{n}$, about $x_{n}$ ), so $\Phi_{\eta_{n}}=\lambda n^{2}$ on $I_{n}$ and have with $\mu=\sum n^{-2} \eta_{n}$, then $\Phi_{\mu}$ is bounded, $\Phi_{\mu}\left(x_{n}\right) \geq \lambda$ while $\Phi_{\mu}(0) \leq 2 \sum_{n=1}^{\infty} n^{-2} \log n$. Hence one loses continuity for $\lambda$ large.

The following is sometimes useful:
Proposition A.3. If $\Phi_{\mu}(x)$ restricted to $\operatorname{supp}(\mu)$ is continuous, then $\Phi_{\mu}$ is continuous on $\mathbb{C}$.

Remarks. 1. The general case can be found in [73, Theorem 1.7]. Here we will sketch the case where $\operatorname{supp}(\mu) \subset \mathbb{R}$ which is most relevant to OPRL.
2. By lower semicontinuity, if $\Phi_{\mu}$ fails to be continuous on $\mathbb{C}$, there exists $z_{n} \rightarrow z_{\infty}$, so $\Phi_{\mu}\left(z_{n}\right) \rightarrow a>\Phi_{\mu}\left(z_{\infty}\right)$. Continuity off $\operatorname{supp}(\mu)$ is easy, so we must have $z_{\infty} \in \mathbb{R}($ since we are $\operatorname{supposing} \operatorname{supp}(\mu) \subset \mathbb{R})$.
3. If $w, x, y \in \mathbb{R}$, then $|w-x-i y|^{-1} \leq|w-x|^{-1}$, so

$$
\Phi_{\mu}(x+i y) \leq \Phi_{\mu}(x)
$$

and thus $\lim \inf \Phi_{\mu}\left(\operatorname{Re} z_{n}\right) \geq a>\Phi_{\mu}\left(\operatorname{Re} z_{\infty}\right)$, so without loss, we can suppose $z_{n}$ are real.
4. If $(\alpha, \beta) \subset \mathbb{R} \backslash \operatorname{supp}(\mu)$ with $\alpha, \beta \in \operatorname{supp}(\mu)$, it is easy to see that $\Phi_{\mu}(x)$ is continuous when restricted to $[\alpha, \beta]$ (using monotone convergence at the endpoints) and convex on $[\alpha, \beta]$ since $\log |x|^{-1}$ is convex. Thus, $\max _{[\alpha, \beta]} \Phi_{\mu}(x)=\max \left(\Phi_{\mu}(\alpha), \Phi_{\mu}(\beta)\right)$. From this, it is easy to see that if such a $z_{n} \in \mathbb{R}$ exists, one can take $z_{n} \in \operatorname{supp}(\mu)$ and so get a contradiction to the assumed continuity of $\Phi_{\mu}$ restricted to $\operatorname{supp}(\mu)$.

The energy or Coulomb energy of $\mu$ is defined by

$$
\begin{equation*}
\mathcal{E}(\mu)=\int \Phi_{\mu}(x) d \mu(x)=\int \log |x-y|^{-1} d \mu(x) d \mu(y) \tag{A.6}
\end{equation*}
$$

where, again, the value $+\infty$ is allowed. If $E \subset \mathbb{C}$ is compact, we say it has capacity zero if $\mathcal{E}(\mu)=\infty$ for all $\mu \in \mathcal{M}_{+, 1}(E)$, the probability measures on $E$. If $E$ does not have capacity zero, then the capacity, $C(E)$, of $E$ is defined by

$$
\begin{equation*}
C(E)=\exp \left(-\min \left(\mathcal{E}(\rho) \mid \rho \in \mathcal{M}_{+, 1}(E)\right)\right) \tag{A.7}
\end{equation*}
$$

One indication that this strange-looking definition is sensible is seen by, as we will show below (see Example A.17),

$$
\begin{equation*}
C([a, b])=\frac{1}{4}(b-a) \tag{A.8}
\end{equation*}
$$

It is useful to define the capacity of any Borel set. For bounded open sets, $U$,

$$
\begin{equation*}
C(U)=\sup (C(K) \mid K \subset U, K \text { compact }) \tag{A.9}
\end{equation*}
$$

and then for arbitrary bounded Borel $X$,

$$
\begin{equation*}
C(X)=\inf (C(U) \mid X \subset U, U \text { open }) \tag{A.10}
\end{equation*}
$$

It can then be proven (see [73, Thm. 2.8]) that

$$
\begin{equation*}
C(X)=\sup (C(K) \mid K \subset X, K \text { compact }) \tag{A.11}
\end{equation*}
$$

for any Borel sets and that (A.10) holds for compact $X$. In particular, $C(X)=0$ if and only if $\mathcal{E}(\mu)=\infty$ for any measure $\mu$ with $\operatorname{supp}(\mu) \subset X$.

The key technical fact behind Theorem 1.16 is the following:
Proposition A.4. If $C(X)>0$ for a Borel set $X$, there exists a probability measure, $\mu$, supported in $X$ so that $\Phi_{\mu}(x)$ is continuous on $\mathbb{C}$.

Remarks. 1. Let $\mu$ have finite energy so $\int \Phi_{\mu}(x) d \mu(x)<\infty$. By Lusin's theorem (see, e.g., the remark after Theorem 6 of Appendix A of Lax [76] for the truly simple proof), we can find compact sets $K \subset$ $\operatorname{supp}(d \mu)$ so $\mu(K)>0$ and $\Phi_{\mu} \upharpoonright K$ is continuous.
2. Let $\nu=\mu \upharpoonright K$, that is, $\nu(S)=\mu(S \cap K)$. Since $\mu(K)>0, \nu$ is a nonzero measure. By general principles, both $\Phi_{\nu}$ and $\Phi_{\mu-\nu}$ are lower semicontinuous on $K$, so since $\Phi_{\mu}$ is continuous,

$$
\begin{equation*}
\Phi_{\nu}=\Phi_{\mu}-\Phi_{\mu-\nu} \tag{A.12}
\end{equation*}
$$

is upper semicontinuous on $K$. Thus, $\Phi_{\nu}$ is continuous on $K$ (since $\Phi_{\mu}$ is continuous on $K$, it is bounded there, so $\Phi_{\nu}$ and $\Phi_{\mu-\nu}$ are both bounded there, so there are no $\infty-\infty$ cancellations in (A.12)).
3. By Proposition A.3, $\Phi_{\nu}$ is continuous on $\mathbb{C}$.

Now suppose $\mu$ is an arbitrary measure of compact support and that $C\left(\left\{x \mid \Phi_{\mu}(x)=\infty\right\}\right)>0$. Then, by the above proposition, there is an $\eta$ supported on that set with $\Phi_{\eta}$ continuous and so bounded above on $\operatorname{supp}(d \mu)$. Thus,

$$
\begin{equation*}
\int \Phi_{\eta}(x) d \mu(x)<\infty \tag{A.13}
\end{equation*}
$$

On the other hand, $\Phi_{\mu}(x)=\infty$ on $\operatorname{supp}(d \eta)$, so

$$
\begin{equation*}
\int \Phi_{\mu}(x) d \eta(x)=\infty \tag{A.14}
\end{equation*}
$$

This contradicts (A.2). We thus see that the last proposition implies:
Corollary A.5. For any measure of compact support, $\mu,\left\{x \mid \Phi_{\mu}(x)=\right.$ $\infty\}$ has capacity zero.

A main reason for defining capacity for any Borel set is that it lets us single out sets of capacity zero (also called polar sets), which are very thin sets (e.g., of Hausdorff dimension zero; see Theorem A.20). We say an event (i.e., a Borel set) occurs quasi-everywhere (q.e.) if and only if it fails on a set of capacity zero. "Nearly everywhere" is also used. A countable union of capacity zero sets is capacity zero. Note that if $\mu$ is any measure of compact support, with $\mathcal{E}(\mu)<\infty$, then $\mathcal{E}(\mu \upharpoonright E)<\infty$ for any compact $E$ (because $\log |x-y|^{-1}$ is bounded below) and thus, $\mu(E)=0$ if $C(E)=0$. It follows (using (A.11)) that

Proposition A.6. If $\mathcal{E}(\mu)<\infty$, then $\mu(X)=0$ for any $X$ with $C(X)=0$.

Here is an important result showing the importance of sets of zero capacity. It is the key to Van Assche's proof in Section 4 and the proof of our new Theorem 1.16 in Section 7.

Theorem A.7. Let $\nu_{n}, \nu$ be measures with supports contained in a fixed compact set $K$ and $\sup _{n} \nu_{n}(K)<\infty$. If $\nu_{n} \rightarrow \nu$ weakly, then

$$
\begin{equation*}
\liminf _{n \rightarrow \infty} \Phi_{\nu_{n}}(x) \geq \Phi_{\nu}(x) \tag{A.15}
\end{equation*}
$$

for all $x \in \mathbb{C}$ and equality holds q.e.
Remarks. 1. (A.15) is called the "Principle of Descent" and the equality q.e. is the "Upper Envelope Theorem."
2. Suppose $\nu_{n}$ has a point mass of weight $\frac{1}{2^{n}}$ at $\left\{\frac{j}{2^{n}}\right\}_{j=0}^{2^{n}-1}$. Then $d \nu_{n} \rightarrow d x \equiv d \nu$, Lebesgue measure. $\Phi_{\nu_{n}}\left(\frac{j}{2^{n}}\right)=\infty$ so $\liminf \Phi_{\nu_{n}}(x)=$ $\infty$ at any dyadic rational, while $\Phi_{\nu}(x)<\infty$ for all $x$. This shows equality may not hold everywhere. This example is very relevant to spectral theory. For the Anderson model, we expect $\limsup \left|p_{n}(x)\right|^{1 / n}=e^{\gamma(x)}$
for almost all $x$ and $\lim \sup \left|p_{n}(x)\right|^{1 / n}=e^{-\gamma(x)}$ at the eigenvalues. Thus, with $\nu_{n}$ the zero counting measure for $p_{n}$, so $\Phi_{\nu_{n}}(x)=-\log \left|p_{n}(x)\right|^{1 / n}$, we have $\lim \inf \Phi_{\nu_{n}}(x)=-\gamma(x)$ for almost all $x$ and $\gamma(x)$ at the eigenvalue consistent with (A.15), and with (A.15) failing on a capacity zero set, including the countable set of eigenvalues.
3. (A.15) is easy. For $\Phi_{\nu}^{M}$ is the convolution with a continuous function so $\lim _{n \rightarrow \infty} \Phi_{\nu_{n}}^{M}(x)=\Phi_{\nu}^{M}(x)$. Since $\Phi_{\nu_{n}}(x) \geq \Phi_{\nu_{n}}^{M}(x)$, we see $\liminf _{n \rightarrow \infty} \Phi_{\nu_{n}}(x) \geq \Phi_{\nu}^{M}(x)$. Taking $M \rightarrow \infty$ yields (A.15).
4. Let $X$ be the set of $x$ for which the inequality in (A.15) is strict. Suppose $C(X)>0$. Then, by Proposition A.2, there is $\eta$ supported on $X$ with $\Phi_{\eta}(x)$ continuous so

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \int \Phi_{\eta}(x) d \nu_{n}=\int \Phi_{\eta}(x) d \nu \tag{A.16}
\end{equation*}
$$

By (A.2) and Fatou's lemma ( $\Phi_{\nu_{n}}(x)$ is uniformly bounded below),

$$
\begin{align*}
\lim _{n \rightarrow \infty} \int \Phi_{\eta}(x) d \nu_{n} & =\lim \int \Phi_{\nu_{n}}(x) d \eta \\
& \geq \int \liminf \Phi_{\nu_{n}}(x) d \eta \\
& >\int \Phi_{\nu}(x) d \eta  \tag{A.17}\\
& =\int \Phi_{\eta}(x) d \nu
\end{align*}
$$

where (A.17) comes from the assumptions $\operatorname{supp}(d \eta) \subset X$ and (A.15) is strict on $X$. This contradiction to (A.16) shows $C(X)=0$, that is, equality holds in (A.15) q.e.

If $\mathcal{E}_{M}(\mu)=\int \Phi_{\mu}^{M} d \mu(x)$, then it is easy to prove $\mathcal{E}_{M}$ is weakly continuous and conditionally positive definite in that

$$
\begin{equation*}
\mu(\mathbb{C})=\nu(\mathbb{C}) \Rightarrow \mathcal{E}_{M}(\mu-\nu) \geq 0 \tag{A.18}
\end{equation*}
$$

where boundedness of $\log \left(\min \left(|x-y|^{-1}, M\right)\right)$ implies $\mathcal{E}_{M}$ makes sense for any signed measure. By taking $M$ to infinity, one obtains

Theorem A.8. The map $\mu \mapsto \mathcal{E}(\mu)$ is weakly lower semicontinuous on $\mathcal{M}_{+, 1}(E)$ for any compact $E \subset \mathbb{C}$. Moreover, it is conditionally positive definite in the sense that for $\mu, \nu \in \mathcal{M}_{+, 1}(E), \mathcal{E}(\mu)<\infty$ and $\mathcal{E}(\nu)<\infty$ imply

$$
\begin{equation*}
\int \Phi_{\nu}(x) d \mu(x) \leq \frac{1}{2} \mathcal{E}(\mu)+\frac{1}{2} \mathcal{E}(\nu) \tag{A.19}
\end{equation*}
$$

with strict inequality if $\mu \neq \nu$.

Remark. The strict inequality requires an extra argument. One can prove that if $\mu, \nu \in \mathcal{M}_{+, 1}(E)$ with finite energy, then $\widehat{\mu}(k)-\widehat{\nu}(k)$ is analytic in $k$ vanishing at $k=0$ and

$$
\begin{equation*}
\mathcal{E}(\mu)+\mathcal{E}(\nu)-2 \int \Phi_{\nu}(x) d \mu(x)=\frac{1}{2 \pi} \int\left|\frac{\widehat{\mu}(k)-\widehat{\nu}(k)}{k}\right|^{2} d^{2} k \tag{A.20}
\end{equation*}
$$

Since the inequality in (A.19) is strict and

$$
\begin{equation*}
\mathcal{E}\left(\frac{1}{2} \mu+\frac{1}{2} \nu\right)=\frac{1}{4} \mathcal{E}(\mu)+\frac{1}{4} \mathcal{E}(\nu)+\frac{1}{2} \int \Phi_{\nu}(x) d \mu(x) \tag{A.21}
\end{equation*}
$$

we see that $\mathcal{E}(\mu)$ is strictly convex on $\mathcal{M}_{+, 1}(E)$, and thus
Theorem A.9. Let $E$ be a compact subset of $\mathbb{C}$ with $C(E)>0$. Then there exists a unique probability measure, $d \rho_{E}$, called the equilibrium measure for $E$, that has

$$
\begin{equation*}
\mathcal{E}\left(\rho_{E}\right)=\log \left(C(E)^{-1}\right) \tag{A.22}
\end{equation*}
$$

The properties of $\rho_{E}$ are summarized in
Theorem A.10. Let $E \subset \mathbb{C}$ be compact. Let $\Omega$ be the unbounded component of $\mathbb{C} \backslash E$ and $\widetilde{\Omega}=\mathbb{C} \backslash(\Omega \cup E)$ the union of the bounded components of $\mathbb{C} \backslash E$. Suppose $C(E)>0$ and $d \rho_{E}$ is its equilibrium measure. Then
(a) For all $x \in \mathbb{C}$,

$$
\begin{equation*}
\Phi_{\rho_{E}}(x) \leq \log \left(C(E)^{-1}\right) \tag{A.23}
\end{equation*}
$$

(b) Equality holds in (A.23) q.e. on $E$ and on $\widetilde{\Omega}$.
(c) Strict inequality holds in (A.23) on $\Omega$.
(d) $\rho_{E}$ is supported on $\partial \Omega$, the boundary viewed as a set in $\mathbb{C}$.
(e) $\Phi_{\rho_{E}}$ is continuous on $\mathbb{C}$ if and only if it continuous when restricted to $\operatorname{supp}\left(d \rho_{E}\right)$ if and only if equality holds in (A.23) on $\operatorname{supp}\left(d \rho_{E}\right)$.
(f) If $I \subset E \subset \mathbb{R}$ with $I=(a, b)$, then $d \rho_{E} \upharpoonright I$ is absolutely continuous with respect to Lebesgue measure, $\frac{d \rho_{E}}{d x} \upharpoonright I$ is real analytic, and equality holds in (A.23) on I.

Remarks. 1. For example, if $E=\partial \mathbb{D}, \Omega=\mathbb{C} \backslash \overline{\mathbb{D}}$ and $\widetilde{\Omega}=\mathbb{D}$.
2. See $[73,88]$ for complete proofs.
3. $(\mathrm{a})+(\mathrm{b})$ is called Frostman's theorem.
4. Equality in (A.23) may not hold everywhere on $E$; for example, if $E=[-1,1] \cup\{2\}$, the equilibrium measure gives zero weight to $\{2\}$, so is the same as the equilibrium measure for $[-1,1]$ and that $d \rho_{E}$ has inequality on $\mathbb{C} \backslash[-1,1]$ by (c).
5. If $f$ is supported on $\operatorname{supp}\left(d \rho_{E}\right)$ and $f$ bounded and Borel, and $\int f d \rho_{E}=0$, then $(1+\varepsilon f) d \rho_{E}$ is a probability measure for $\varepsilon$ small with $\mathcal{E}\left((1+\varepsilon f) d \rho_{E}\right)<\infty$. Since $\frac{d}{d \varepsilon} \mathcal{E}\left((1+\varepsilon f) d \rho_{E}\right)=2 \int f(x) \Phi_{\rho_{E}}(x) d \rho_{E}(x)$, we see $\Phi_{\rho_{E}}(x)$ is a constant for $d \rho_{E^{-}}$-a.e. $x$. Since $\mathcal{E}\left(\rho_{E}\right)=\int d \rho_{E} \Phi_{\rho_{E}}(x)$, the constant must be $\mathcal{E}\left(\rho_{E}\right)=\log \left(C(E)^{-1}\right)$. By lower semicontinuity, (A.23) holds on $\operatorname{supp}\left(d \rho_{E}\right)$. Since $\Phi_{\rho_{E}}$ is harmonic on $\mathbb{C} \backslash \operatorname{supp}\left(d \rho_{E}\right)$ and goes to $-\infty$ as $|x| \rightarrow \infty$, (A.23) holds by the maximum principle.
6. Let $\eta$ be a probability measure on $E$ with $\mathcal{E}(\eta)<\infty$. Then

$$
\begin{equation*}
\mathcal{E}\left((1-t) d \rho_{E}+t d \eta\right)=\mathcal{E}\left(d \rho_{E}\right)+t\left(\int \Phi_{\rho_{E}}(x)\left[d \eta-d \rho_{E}\right]\right)+O\left(t^{2}\right) \tag{A.24}
\end{equation*}
$$

Since $\int d \rho_{E} \Phi_{\rho_{E}}(x)=\mathcal{E}\left(d \rho_{E}\right)=\log \left(C(E)^{-1}\right)$, if $\eta$ is supported on a set where strict inequality holds in (A.23), $\mathcal{E}\left((1-t) d \rho_{E}+t d \eta\right)<\mathcal{E}\left(d \rho_{E}\right)$ for small $t$, violating minimality. Thus the set where (A.23) has inequality cannot support a measure of finite energy, that is, it has zero capacity, proving (b).
7. Since $\Phi_{\rho_{E}}$ is harmonic on $\Omega$ and goes to $-\infty$ at $\infty$, the maximum principle implies $\Phi_{\rho_{E}}(x)$ cannot take its maximum (which is $\log \left(C(E)^{-1}\right)$ ) on $\Omega$. (e) follows from Proposition A.3. (d) is left to the references; see [73, 88].
8. If $I \subset E \subset \mathbb{R}$, one first shows equality holds in (A.23) on $I$ and that $\Phi$ is continuous there. (This uses the theory of "barriers"; see $[73,88]$. One can also prove this using periodic Jacobi matrices and approximations; see [102]). Then one can apply the reflection principle to see that $\Phi_{\rho_{E}}$ has a harmonic continuation across $I$. Indeed, $\Phi_{\rho_{E}}$ is then the real part of a function analytic on $I$ with zero derivative there. That derivative for $\operatorname{Im} z>0$ is the real part of

$$
\begin{equation*}
F(z)=\int \frac{d \rho_{E}(x)}{x-z} \tag{A.25}
\end{equation*}
$$

so, by the standard theory of boundary values of Herglotz functions (see [98, Sect. 1.3]), we have that $d \rho_{E} \upharpoonright I$ is absolutely continuous and

$$
\begin{equation*}
\frac{d \rho_{E}}{d x}=\frac{1}{\pi} \operatorname{Im} F(x+i 0) \tag{A.26}
\end{equation*}
$$

proving real analyticity of this derivative.
9. The same argument as in Remark 8 applies if $I$ is replaced by an analytic arc with a neighborhood $N$ obeying $N \cap E=I$. In particular, if $I$ is an "interval" in $\partial \mathbb{D}$ and $I \subset E \subset \partial \mathbb{D}$, we have absolute continuity and analyticity on $I$.

Here is an interesting consequence of (A.2):

Theorem A.11. Let $\nu$ be a measure of compact support, $E$, so that $C(E)>0$. Then

$$
\begin{equation*}
\Phi_{\nu}(x)<\infty \quad \text { for } d \rho_{E} \text { a.e. } x \tag{A.27}
\end{equation*}
$$

Remarks. 1. This can happen even if $\mathcal{E}(\nu)=\infty$ so $\int \Phi_{\nu}(x) d \nu(x)=$ $\infty$.
2. For (A.23) implies

$$
\int \Phi_{\rho_{E}}(x) d \nu \leq \log \left(C(E)^{-1}\right) \nu(E)<\infty
$$

so (A.2) implies

$$
\begin{equation*}
\int \Phi_{\nu}(x) d \rho_{E}(x)<\infty \tag{A.28}
\end{equation*}
$$

The following illustrates the connection between potential theory and polynomials:

Theorem A. 12 (Bernstein-Walsh Lemma). Let $E$ be a compact set in $\mathbb{C}$ with $C(E)>0$ and let $\Omega$ be the unbounded component of $\mathbb{C} \backslash E$. Let $p_{n}$ be a polynomial of degree $n$ and let

$$
\begin{equation*}
\left\|p_{n}\right\|_{E}=\sup _{z \in E}\left|p_{n}(z)\right| \tag{A.29}
\end{equation*}
$$

Then for all $z \in \Omega$,

$$
\begin{equation*}
\left|p_{n}(z)\right| \leq C(E)^{-n}\left\|p_{n}\right\|_{E}\left[\exp \left(-n \Phi_{\rho_{E}}(z)\right)\right] \tag{A.30}
\end{equation*}
$$

Remarks. 1. This is named after Bernstein and Walsh [121], although the result appears essentially in Szegő [107].
2. Let $\left\{z_{j}\right\}_{j=1}^{n}$ be the zeros of $p_{n}$. Define

$$
\begin{equation*}
g(z)=\log \left|p_{n}(z)\right|+n \Phi_{\rho_{E}}(z)+n \log (C(E)) \tag{A.31}
\end{equation*}
$$

on $\Omega \cup\{\infty\} \backslash\left\{z_{j}\right\}_{j=1}^{n}=\Omega^{\prime} . g$ is harmonic on $\Omega^{\prime}$ including at $\infty$ since both $\log \left|p_{n}(z)\right|$ and $-n \Phi_{\rho_{E}}(z)$ are $n \log |z|$ plus harmonic near $\infty$. Since $g_{n}(z) \rightarrow-\infty$ at the $z_{j} \in \Omega$, we see

$$
\begin{equation*}
\sup _{z \in \Omega^{\prime}}|g(z)| \leq \lim _{\delta \downarrow 0}\left[\sup _{\substack{\operatorname{dist}(w, E)=\delta \\ w \in \Omega}}|g(w)|\right] \tag{A.32}
\end{equation*}
$$

But, by (A.23), $g(z) \leq \log \left|p_{n}(z)\right|$, so

$$
g(z) \leq \log \left\|p_{n}\right\|_{E}
$$

which is (A.30) on $\Omega^{\prime} \backslash\{\infty\}$. (A.30) holds trivially at the $z_{j}$, completing the proof.

Following ideas of Craig [24], one can say much more about $\frac{d \rho_{E}}{d x}$ when $E$ contains an isolated closed interval:

Theorem A.13. Let $E \subset \mathbb{R}$ be compact and $a<c<d<b$ so $E \cap(a, b)=[c, d]$. Then there exists $g$ real, real analytic, and strictly positive on $[c, d]$ so that

$$
\begin{equation*}
d \rho_{E} \upharpoonright[c, d]=g(x)[(d-x)(x-c)]^{-1 / 2} d x \tag{А.33}
\end{equation*}
$$

If $E=\cup_{j=1}^{\ell+1}\left[a_{j}, b_{j}\right]$ with $a_{1}<b_{1}<a_{2}<\cdots<b_{\ell+1}$, there are $x_{j} \in$ $\left(b_{j}, a_{j+1}\right)$ for $j=1,2, \ldots, \ell$ so that

$$
\begin{equation*}
d \rho_{E}(x)=\frac{1}{\pi}\left[\prod_{j=1}^{\ell} \frac{x-x_{j}}{\sqrt{\left(x-b_{j}\right)\left(x-a_{j+1}\right)}}\right] \frac{1}{\sqrt{\left(x-a_{1}\right)\left(b_{\ell+1}-x\right)}} d x \tag{A.34}
\end{equation*}
$$

Remarks. 1. (A.34) is from Craig [24].
2. The idea behind the proof is simple. One lets $F(z)=\int \frac{d \rho_{E}(x)}{x-z}$. By the arguments above, $F$ is pure imaginary on $[c, d]$ as the derivative of $\Phi_{\rho_{E}}(x)$. Thus, $\arg F(x+i 0)$ is $\frac{\pi}{2}$ on $[c, d]$, and by a simple argument, 0 on $[c-\delta, c)$ and $\pi$ on $(d, d+\delta]$. A Herglotz representation for $\log F(x+$ $i 0)$ yields (A.33) and (A.34).
3. The $x_{j}$ 's are uniquely determined by

$$
\begin{equation*}
\int_{b_{j}}^{a_{j+1}} F(x) d x=0 \tag{A.35}
\end{equation*}
$$

Recall a set $S$ is called perfect if it is closed and has no isolated points. A standard argument shows that any compact $E$ has a unique decomposition into disjoint sets, $D \cup S$ where $D$ is a countable set and $S$ is perfect (similarly, any compact $E \subset \mathbb{R}$ can be written $Z \cup F$ where $Z$ has Lebesgue measure zero and $F$ is essentially perfect, that is, $|F \cap(x-\delta, x+\delta)|>0$ for any $x \in F$ and $\delta>0)$.

Similarly, we call a set $P$ potentially perfect (the terminology is new) if $P$ is closed and $C\left(P \cap\left\{x\left|\left|x-x_{0}\right|<\delta\right\}\right)>0\right.$ for all $x_{0} \in P$ and $\delta>0$. It is easy to see that any compact $E \subset \mathbb{C}$ can be uniquely written as a disjoint union $E=Q \cup P$ where $C(Q)=0$ and $P$ is potentially perfect.

These notions are related to equilibrium measures. If $\operatorname{cap}(E)>0$ and $E=Q \cup P$ is this decomposition, then

$$
\begin{equation*}
P=\operatorname{supp}\left(d \rho_{E}\right) \tag{A.36}
\end{equation*}
$$

In particular, $\operatorname{supp}\left(d \rho_{E}\right)=E$ if and only if $E$ is potentially perfect.
Just as one writes $\sigma(d \mu)=\sigma_{\text {disc }}(d \mu) \cup \sigma_{\text {ess }}(d \mu)$, we single out the potentially perfect part of $\sigma(d \mu)$ and call it $\sigma_{\text {cap }}(d \mu)$.

Next, we want to state a kind of converse to Frostman's theorem.

Theorem A.14. Let $E \subset \mathbb{C}$ be compact. Suppose $E$ is potentially perfect. Let $\eta \in \mathcal{M}_{+, 1}(E)$ be a probability measure on $E$ with $\operatorname{supp}(d \eta) \subseteq E$ so that for some constant, $\alpha$,

$$
\begin{equation*}
\Phi_{\eta}(x)=\alpha \quad d \rho_{E} \text {-a.e. } x \tag{A.37}
\end{equation*}
$$

Then $\eta=\rho_{E}$, the equilibrium measure, and $\alpha=\log \left(C(E)^{-1}\right)$.
Remark. By lower semicontinuity, $\Phi_{\eta}(x) \leq \alpha$ on $\operatorname{supp}\left(d \rho_{E}\right)=E$ by hypothesis. Thus,

$$
\begin{equation*}
\mathcal{E}(\eta)=\int \Phi_{\eta}(x) d \eta(x) \leq \alpha<\infty \tag{A.38}
\end{equation*}
$$

so $\eta$ must give zero weight to zero capacity sets. Thus, $\Phi_{\rho_{E}}(x)=$ $\log \left(C(E)^{-1}\right)$ for $d \eta$-a.e. $x$ and thus,

$$
\begin{equation*}
\int \Phi_{\rho_{E}}(x) d \eta(x)=\log \left(C(E)^{-1}\right) \tag{А.39}
\end{equation*}
$$

By (A.2) and (A.37),

$$
\text { LHS of }(\text { A.39 })=\int \Phi_{\eta}(x) d \rho_{E}(x)=\alpha
$$

Thus, $\alpha=\log \left(C(E)^{-1}\right)$, and by (A.38) and uniqueness of minimizers, $\eta=\rho_{E}$.

Next, we note that the Green's function for a compact $E \subset \mathbb{C}$ is defined by

$$
\begin{equation*}
G_{E}(z)=-\Phi_{\rho_{E}}(z)+\log \left(C(E)^{-1}\right) \tag{A.40}
\end{equation*}
$$

It is harmonic on $\mathbb{C} \backslash E, G_{E}(z)-\log |z|$ is harmonic at infinity, and $G_{E}(z)$ has zero boundary values q.e. on $E$. Notice that $G_{E}(z) \geq 0$ on C. If

$$
\begin{equation*}
\lim _{z \rightarrow E} G_{E}(z)=0 \tag{A.41}
\end{equation*}
$$

in the sense that

$$
\begin{equation*}
\lim _{\delta \downarrow 0} \sup _{\operatorname{dist}(z, E)<\delta} G_{E}(z)=0 \tag{A.42}
\end{equation*}
$$

we say $E$ is regular for the Dirichlet problem (just called regular). By Theorem A.10(c), this is true if and only if $\Phi_{\rho_{E}}(x)=\log \left(C(E)^{-1}\right)$ for all $x \in E$. By Theorem A.13, this is true for finite unions of disjoint closed intervals.

Notice that the Bernstein-Walsh lemma (A.30) can be rewritten

$$
\begin{equation*}
\left|p_{n}(z)\right| \leq\left\|p_{n}\right\|_{E} \exp \left(n G_{E}(z)\right) \tag{A.43}
\end{equation*}
$$

Closely related are comparison theorems and limit theorems. We will state them for subsets of $\mathbb{R}$ :

Theorem A.15. Let $E_{1} \subset E_{2} \subset \mathbb{R}$ be compact sets. Then
(i) $C\left(E_{1}\right) \leq C\left(E_{2}\right)$
(ii)

$$
\begin{equation*}
G_{E_{2}}(z) \leq G_{E_{1}}(z) \tag{A.44}
\end{equation*}
$$

for all $z \in \mathbb{C}$
(iii)

$$
\begin{equation*}
d \rho_{E_{2}} \upharpoonright E_{1} \leq d \rho_{E_{1}} \tag{A.45}
\end{equation*}
$$

(iv) If $I=(a, b) \subset E_{1}$, then on $I$,

$$
\begin{equation*}
\frac{d \rho_{E_{2}}}{d x} \leq \frac{d \rho_{E_{1}}}{d x} \tag{A.46}
\end{equation*}
$$

for all $x \in I$.
Theorem A.16. Let $E_{1} \supset E_{2} \supset \ldots$ be compact subsets of $\mathbb{R}$. Let $E_{\infty}=\cap_{j=1}^{\infty} E_{j}$. Then

$$
\begin{equation*}
\lim _{n \rightarrow \infty} C\left(E_{n}\right)=C\left(E_{\infty}\right) \tag{i}
\end{equation*}
$$

(ii) $\rho_{E_{n}} \rightarrow \rho_{E_{\infty}}$ weakly
(iii) For $z \in \mathbb{C} \backslash E_{\infty}$,

$$
\begin{equation*}
\lim _{n \rightarrow \infty} G_{E_{n}}(z)=G_{E_{\infty}}(z) \tag{A.48}
\end{equation*}
$$

and (A.48) holds q.e. on $E_{\infty}$.
(iv) If $I=(a, b) \subset E_{\infty}$, then uniformly on compact subsets of $I$,

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \frac{d \rho_{E_{n}}}{d x}=\frac{d \rho_{E_{\infty}}}{d x} \tag{A.49}
\end{equation*}
$$

Remarks. 1. (A.45) has the pleasing physical interpretation that if one conductor is connected to another, charge leaks out in a way that there is less charge everywhere in the original conductor.
2. Part (i) of each theorem is easy. For Theorem A.15, it follows from the minimum energy definition. For Theorem A.16(i), we note that if $U$ is open with $E_{\infty} \subset U$, then eventually $E_{n} \subset U$, so (A.10) implies (A.47).
3. One proves (ii)-(iv) of Theorem A. 15 first for $E$, a finite union of closed intervals, then proves Theorem A.16, and then for general compact $E_{\infty} \subset \mathbb{R}$ defines $E_{n}=\left\{x \left\lvert\, \operatorname{dist}\left(x, E_{\infty}\right) \leq \frac{1}{n}\right.\right\}$ and proves $\cap_{n} E_{n}=E_{\infty}$ and each $E_{n}$ is a finite union of closed intervals. Theorem A. 16 then yields Theorem A. 15 for general E's (see [102]).
4. For $E_{1}, E_{2}$ finite union of closed intervals and $z \notin E_{2}$, one gets (A.44) by noting the difference $G_{E_{1}}(z)-G_{E_{2}}(z)$ is harmonic on $\mathbb{C} \backslash E_{2}$, zero on $E_{1}$, and positive on $E_{2} \backslash E_{1}$, where $G_{E_{1}}>0$ and $G_{E_{2}}=0$. The
inequality for $z \in E_{2}$ then follows from the fact that any subharmonic function $h$ obeys

$$
\begin{equation*}
h\left(z_{0}\right)=\lim _{r \downharpoonright 0} \frac{1}{2 \pi r} \int_{0}^{2 \pi} h\left(z_{0}+r e^{i \theta}\right) d \theta \tag{A.50}
\end{equation*}
$$

5. For this case, one gets (A.45)/(A.46) by noting that (A.26) can be rewritten

$$
\begin{equation*}
\frac{d \rho_{E}\left(x_{0}\right)}{d x}=\frac{1}{\pi} \lim _{\varepsilon \downarrow 0} \varepsilon^{-1} G_{E}\left(x_{0}+i \varepsilon\right) \tag{A.51}
\end{equation*}
$$

for $x_{0} \in E$ and using (A.44) for $x \in E_{1}$.
6. If $\left\{E_{n}\right\}_{n=1}^{\infty}, E_{\infty}$ are as in Theorem A. 16 and $d \eta$ is a weak limit point of $d \rho_{E_{n}}$, then $\eta$ is supported on $E_{\infty}$, and by lower semicontinuity of the Coulomb energy $\mathcal{E}$,

$$
\begin{aligned}
\mathcal{E}(\eta) & \leq \lim \mathcal{E}\left(\rho_{E_{n}}\right) \\
& =\lim \log \left(C\left(E_{n}\right)^{-1}\right) \\
& =\log \left(C\left(E_{\infty}\right)^{-1}\right)
\end{aligned}
$$

by (A.47), so $\eta=\rho_{E_{\infty}}$, that is, $\rho_{E_{n}} \rightarrow \rho_{E_{\infty}}$ weakly. (A.48) then follows for $z \notin E_{\infty}$ from (A.47) and continuity of $\Phi_{\nu}(z)$ in $\nu$ for $z \notin \operatorname{supp}(d \nu)$. (A.50) implies convergence for $z \in E_{\infty}$.
7. (A.49) follows from $\rho_{E_{n}} \rightarrow \rho_{E_{\infty}}$ and uniform bounds on derivatives of $\frac{d \rho}{d x}$ on $I$, which in turn follow from the proof of (A.33).

Example A.17. Harmonic functions are conformally invariant, which means (since Green's functions are normalized by $G_{E}(z)=\log |z|+$ $O(1)$ near infinity and boundary values of 0 on $E)$, if $Q$ is an analytic bijection of $\mathbb{C} \backslash \overline{\mathbb{D}} \cup\{\infty\}$ to $\Omega \cup\{\infty\}$ with $Q(z)=C z+O(1)$ near infinity, then, since $\log |z|$ is the Green's function for $E=\partial \mathbb{D}, \log \left|Q^{-1}(z)\right|$ is the Green's function for $E$ and $C$ its capacity. In particular, with $Q(z)=z+\frac{1}{z}$, we see

$$
\begin{equation*}
C([-2,2])=1 \tag{A.52}
\end{equation*}
$$

and consistent with (A.34)

$$
\begin{equation*}
d \rho_{[-2,2]}(x)=\frac{1}{\pi} \frac{d x}{\sqrt{4-x^{2}}} \tag{A.53}
\end{equation*}
$$

Notice that, by scaling, if $\lambda>0$ and $\lambda E=\{\lambda z \mid z \in E\}$ and $\mu \in$ $\mathcal{M}_{+, 1}(E)$ is mapped to $\mu_{\lambda}$ in $\mathcal{M}_{+, 1}(\lambda E)$ by scaling, then

$$
\begin{equation*}
\mathcal{E}\left(\mu_{\lambda}\right)=-\log (\lambda)+\mathcal{E}(\mu) \tag{A.54}
\end{equation*}
$$

SO

$$
\begin{equation*}
C(\lambda E)=\lambda C(E) \tag{A.55}
\end{equation*}
$$

This plus translation invariance shows

$$
C([a, b])=\frac{1}{4}(b-a)
$$

Now, let $E \subset \mathbb{R}, E_{-}=E \cap(-\infty, 0), E_{+}=E \backslash E_{-}$, and for $a>0$, let $E^{(a)}=E_{-} \cup\left(E_{+}+a\right)$. Let $d \rho_{E}$ be equilibrium measure for $E$. Since $\log |x+a-y|^{-1}<\log |x-y|^{-1}$ for $x>0, y<0, a>0$, we see

$$
\begin{equation*}
\mathcal{E}\left(\rho \upharpoonright E_{-}+T_{a}\left(\rho \upharpoonright E_{+}\right)\right)<\mathcal{E}(\rho) \tag{A.56}
\end{equation*}
$$

(where $T_{a} \mu$ is the translate of $\mu$ ). Thus

$$
C\left(E^{(a)}\right) \geq C(E)
$$

This is an expression of the repulsive nature of the Coulomb force! Thus, by joining together all the pieces of $E$ (via a limiting argument), one sees that if $|E|$ is the Lebesgue measure of $E \subset \mathbb{R}$, then

$$
\begin{equation*}
C(E) \geq \frac{1}{4}|E| \tag{A.57}
\end{equation*}
$$

and sets of capacity zero have Lebesgue measure zero.
Example A.18. Let $d \mu$ be the conventional Cantor measure on $[0,1]$ which can be thought of as writing $x=\sum \frac{a_{n}(x)}{3^{n}}$ with $a_{n}=0,1$ or 2 and taking $d \mu$ as the infinite product of measures given weight $\frac{1}{2}$ to $a_{n}=0$ or 2 . Looking at $a_{1}$, we get the usual two pieces of mass $\frac{1}{2}$ with minimum distance $\frac{1}{3}$ between them. Look at $a_{1}, \ldots, a_{k}$ and we have $2^{k}$ pieces of mass $2^{-k}$ and minimum distance $3^{-k}$. Given $x, y$ in the Cantor set, dist $|x-y|<3^{-k}$ if and only if they are in the same pieces, that is,

$$
\mu\left(\left\{x\left||x-y|<3^{-k}\right\}\right)=2^{-k}\right.
$$

Thus

$$
\begin{aligned}
\int \log |x-y|^{-1} d \mu(x) d \mu(y) & \leq \sum_{k}(k+1)(\log 3) \mu\left(\left\{x, y| | x-y \mid<3^{-k}\right\}\right) \\
& =\sum_{k}(k+1)(\log 3) 2^{-k}<\infty
\end{aligned}
$$

This shows the Cantor set has positive capacity. Generalizing, we get sets of any Hausdorff dimension $\alpha>0$ with positive capacity. In fact, as we will see shortly, any set of positive Hausdorff dimension has positive capacity.

Example A.19. Fix $a>0$ and let $E=\left(-\frac{a}{2}-\Delta,-\frac{a}{2}\right) \cup\left(\frac{a}{2}, \frac{a}{2}+\Delta\right)$ where $\Delta=\frac{4}{a}$. When $a$ is very large, the equilibrium measure is very
close to the average of the equilibrium measure for the two individual intervals. This measure has energy approximately

$$
2 \frac{1}{4} \log \left(\frac{4}{\Delta}\right)+2 \frac{1}{4} \log (a)=0
$$

so the asymptotic capacity is 1 . This phenomenon of distant pieces of individually small capacity having total capacity bounded away from zero is a two-dimensional phenomenon.

Sets of capacity zero not only have zero Lebesgue measure, but they also have zero $\alpha$-dimensional Hausdorff measure for any $\alpha>0$ :
Theorem A.20. Any compact set $E$ of capacity zero has zero Hausdorff dimension.

Remarks. 1. We will sketch a proof where $E \subset \mathbb{R}$. What one needs to do, for any $\alpha>0, \varepsilon>0$, is to find a cover of $E$ by intervals $I_{1}, \ldots, I_{n} \ldots$ of length $\left|I_{j}\right|$ so that

$$
\begin{equation*}
\sum\left|I_{j}\right|^{\alpha}<\varepsilon \tag{A.58}
\end{equation*}
$$

2. We begin by noting that there is a measure $\mu$ (not necessarily supported by $E$ ) so that $\Phi_{\mu}(x)=\infty$ for all $x \in E$ (we do not care that $E$ is exactly the set where $\Phi_{\mu}=\infty$ but note that by combining the ideas here with Corollary A.5, one can show $E$ is the set where some potential is infinite if and only if $E$ is a $G_{\delta}$-set of zero capacity). Here is how to construct $\mu$. Let $E_{m}=\left\{x \in \mathbb{R} \left\lvert\, \operatorname{dist}(x, E) \leq \frac{1}{m}\right.\right\}$. $E_{m}$ is a finite union of closed intervals and, by (A.10), $C\left(E_{m}\right) \downarrow 0$. Pass to a subsequence $\tilde{E}_{m}$, so $C\left(\tilde{E}_{m}\right) \leq \exp \left(-\frac{1}{m^{2}}\right)$ so $\Phi_{\rho_{\tilde{E}_{m}}}(x) \geq m^{2}$ on $\tilde{E}_{m}$ and so on $E$. Let $\mu=\sum_{m} m^{-2} \rho_{\tilde{E}_{m}} . \mu$ is a finite measure with $\Phi_{\mu}=\infty$ on $E$.
3. Let $x \in E$. Suppose for some $\alpha>0$ and $c>0$, we have with $I_{r}^{x}=(x-r, x+r)$,

$$
\begin{equation*}
\mu\left(I_{r}^{x}\right) \leq c(2 r)^{\alpha} \tag{A.59}
\end{equation*}
$$

Then picking $r=2^{-n}$, we see (with $n_{0}$ large and negative so $\operatorname{supp}(d \mu) \subset$ $I_{2^{-n}}^{x}$ )

$$
\begin{aligned}
\int \log |y-x|^{-1} d \mu(y) & \leq \sum_{n=n_{0}}^{\infty}[(n+1) \log 2] \mu\left(I_{2^{-n}}^{x}\right) \\
& <\infty
\end{aligned}
$$

We thus conclude (A.59) always fails, that is, for any $x \in E$ and any $\alpha$,

$$
\begin{equation*}
\limsup _{r \downarrow 0}(2 r)^{-\alpha} \mu\left(I_{r}^{x}\right)=\infty \tag{A.60}
\end{equation*}
$$

4. Given $\alpha>0, \delta>0$ fixed, by (A.60), we can find for each $x \in E$, so $I_{r_{x}}^{x}$ with

$$
\begin{equation*}
\mu\left(I_{r_{x}}^{x}\right) \geq \delta^{-1}\left(2 r_{x}\right)^{\alpha} \tag{A.61}
\end{equation*}
$$

5. There is standard covering lemma used in the proof of the HardyLittlewood maximal theorem (see [64], p. 74, the proof of the lemma) that we can find a suitable sequence $x_{j}$ with

$$
\begin{equation*}
I_{r_{x_{j}}}^{x_{j}} \cap I_{r_{x_{k}}}^{x_{k}}=\emptyset \tag{A.62}
\end{equation*}
$$

and

$$
\begin{equation*}
\bigcup_{x} I_{r_{x}}^{x} \subset \bigcup_{j=1}^{\infty} I_{4 r_{x_{j}}}^{x_{j}} \tag{А.63}
\end{equation*}
$$

6. Thus, $\left\{I_{4 r_{x_{j}}}^{x_{j}}\right\}$ cover $E$ and, by (A.61),

$$
\begin{aligned}
\sum_{j}\left|I_{4 r_{x_{j}}}^{x_{j}}\right|^{\alpha} & \leq 4^{\alpha} \sum_{j}\left|I_{r_{x_{j}}}^{x_{j}}\right|^{\alpha} \\
& \leq 4^{\alpha} \delta \sum_{j} \mu\left(I_{r_{x_{j}}}^{x_{j}}\right) \\
& \leq 4^{\alpha} \delta \mu(\mathbb{R})
\end{aligned}
$$

by (A.62). Since $\delta$ is arbitrary, we have the required covers to see $\operatorname{dim}(E)=0$.

A final comparison result will be needed in Appendix B:
Theorem A.21. Let $\mu, \nu$ be two probability measures on $\mathbb{R}$ so that for all z near infinity,

$$
\begin{equation*}
\Phi_{\mu}(z) \geq \Phi_{\nu}(z) \tag{A.64}
\end{equation*}
$$

Then $\mu=\nu$. In particular, if either $\Phi_{\mu}(z) \geq \Phi_{\rho_{E}}(z)$ or $\Phi_{\mu}(z) \leq \Phi_{\rho_{E}}(z)$ for all $z$ near infinity, then $\mu=\rho_{E}$.

## Remark.

$$
\begin{aligned}
\Phi_{\mu}(z)+\log |z| & =-\operatorname{Re} \int \log \left(1-\frac{w}{z}\right) d \mu(w) \\
& =\operatorname{Re}\left[\sum_{n=1}^{\infty} z^{-n} \int w^{n} d \mu(w)\right]
\end{aligned}
$$

Thus $\tilde{\Phi}_{\mu}(z) \equiv \Phi_{\mu}(z)+\log |z|$ is harmonic near infinity with $\tilde{\Phi}_{\mu}(\infty)=0$. Thus, $\Phi_{\mu}-\Phi_{\nu}=\Phi_{\mu}-\Phi_{\nu}$ is harmonic and vanishing at $\infty$. The only way it can have a definite sign near infinity is if it is identically 0 . By harmonicity off $\mathbb{R}, \Phi_{\mu}=\Phi_{\nu}$ on $\mathbb{C} \backslash \mathbb{R}$ and then, by (A.50), on $\mathbb{R}$. Thus, $\Phi_{\mu}=\Phi_{\nu}$ as distributions. Since $-\Delta \Phi_{\mu}=2 \pi \mu$, we see $\mu=\nu$.

## Appendix B: Chebyshev Polynomials, Fekete Sets, and Capacity

For further discussion of the issues in this appendix, see AndrievskiiBlatt [3], Goluzin [48], and Saff-Totik [91] whose discussion overlaps ours here. Let $E \subset \mathbb{C}$ be compact and infinite. The Chebyshev polynomials, $T_{n}(x)$, are defined as those monic polynomials of degree $n$ which minimize

$$
\begin{equation*}
\left\|T_{n}\right\|_{E}=\sup _{z \in E}\left|T_{n}(z)\right| \tag{B.1}
\end{equation*}
$$

By $T_{n}^{R}$, the restricted Chebyshev polynomials, we mean the monic polynomials, all of whose zeros lie in $E$, which minimize $\|\cdot\|_{E}$ among all such polynomials. They can be distinct: for example, if $E=\partial \mathbb{D}, T_{n}(z)=z^{n}$ while $T_{n}^{R}(z)=1+z^{n}$ (not unique). It can be proven (see [113, Thm. III.23]) that Chebyshev (but not restricted Chebyshev) polynomials are unique.

Clearly,

$$
\begin{equation*}
\left\|T_{n}\right\|_{E} \leq\left\|T_{n}^{R}\right\|_{E} \tag{B.2}
\end{equation*}
$$

and since $T_{n} T_{m}$ is a monic polynomial of degree $n+m$,

$$
\begin{equation*}
\left\|T_{n+m}\right\|_{E} \leq\left\|T_{n}\right\|_{E}\left\|T_{m}\right\|_{E} \tag{B.3}
\end{equation*}
$$

so $\lim _{n \rightarrow \infty}\left\|T_{n}\right\|_{E}^{1 / n}$ exists, and similarly, so does $\lim _{n \rightarrow \infty}\left\|T_{n}^{R}\right\|_{E}^{1 / n}$.
An $n$ point Fekete set is a set $\left\{z_{j}\right\}_{j=1}^{n} \subset E$ that maximizes

$$
\begin{equation*}
q_{n}\left(z_{1}, \ldots, z_{n}\right)=\prod_{i \neq j}\left|z_{i}-z_{j}\right| \tag{B.4}
\end{equation*}
$$

There are $n(n-1)$ terms in the product and the Fekete constant is defined by

$$
\begin{equation*}
\zeta_{n}(E)=q_{n}\left(z_{1}, \ldots, z_{n}\right)^{1 / n(n-1)} \tag{B.5}
\end{equation*}
$$

for the maximizing set. The set may not be unique: for example, if $E=\partial \mathbb{D}$ and $\omega_{n}$ is an $n$th root of unity, $\left\{z_{k}=z_{0} \omega_{n}^{k}\right\}$ is a minimizer for any $z_{0} \in \partial \mathbb{D}$.

Let $z_{1}, \ldots, z_{n+1}$ be an $n+1$-point Fekete set. For each $j$,

$$
\begin{equation*}
\prod_{\substack{k, \ell \neq j \\ \ell \neq k}}\left|z_{k}-z_{\ell}\right| \leq \zeta_{n}^{n(n-1)} \tag{B.6}
\end{equation*}
$$

Thus, taking the product over the $n+1$ values of $j$ and noting that each $z_{k}-z_{\ell}$ occurs $n-1$ times,

$$
\begin{equation*}
\left[\zeta_{n+1}^{(n+1) n}\right]^{n-1} \leq\left[\zeta_{n}^{n(n-1)}\right]^{n+1} \tag{B.7}
\end{equation*}
$$

so $\zeta_{n}$ is monotone decreasing. Thus $\zeta_{n}$ has a limit, called the transfinite diameter. The main theorem relating these notions and capacity is

Theorem B.1. For any compact set $E \subset \mathbb{C}$, we have

$$
\begin{equation*}
C(E) \leq\left\|T_{n}\right\|_{E}^{1 / n} \leq\left\|T_{n}^{R}\right\|_{E}^{1 / n} \leq \zeta_{n+1}(E) \tag{B.8}
\end{equation*}
$$

Moreover,

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \zeta_{n}(E) \leq C(E) \tag{B.9}
\end{equation*}
$$

(so all limits equal $C(E)$ ). Finally, if $C(E)>0$,
(i) The normalized density of Fekete sets converges to $d \rho_{E}$, the equilibrium measure for $E$.
(ii) If $E \subset \mathbb{R}$, the normalized zero counting measure for $T_{n}$ and for $T_{n}^{R}$ converges to $d \rho_{E}$.

Remarks. 1. Normalized densities and zero counting measure are the point measures that give weight $k / n$ to a point in the set of multiplicity $k$ (for Fekete sets, $k=1$, but for polynomials there can be zeros of multiplicity $k>1$ ).
2. If $E=\partial \mathbb{D}, T_{n}(z)=z^{n}$, so (ii) fails for $T_{n}$. If $E=\mathbb{D}, T_{n}^{R}(z)=z^{n}$ and (ii) fails for $T_{n}^{R}$ also. It can be shown that if $E \subset \partial \mathbb{D}, E \neq \partial \mathbb{D}$, (ii) also holds.
3. Fekete sets have the interpretation of sets minimizing the point Coulomb energy $\sum_{j \neq k} \log \left|z_{j}-z_{k}\right|^{-1}$. Parts of this theorem can be interpreted as saying the point minimizer and associated energy without self-energies converge to the minimizing continuum distribution and energy, which is physically pleasing!
4. The equality of $\lim \zeta_{n}$ and $\lim \left\|T_{n}\right\|^{1 / n}$ is due to Fekete [38]. The rest is due to Szegő [107], whose proof we partly follow.
5. Stieltjes [106] considered what we call Fekete sets for $E=[-1,1]$, proving that, in that case, the set is unique and consists of $1,-1$, and the $n-2$ zeros of a suitable Jacobi polynomial (see [108]). The general set up is due to Fekete [38].
6. When $E \subset \partial \mathbb{D}$, there are two other sets of polynomials related to minimizing $\left\|P_{n}\right\|_{\infty, E}$. We can restrict to either
(a) "Quasi-real" monic polynomials, that is, degree $n$ polynomials, so for some $\varphi, e^{-i \varphi} e^{-i n \theta / 2} P_{n}\left(e^{i \theta}\right)$ is real for $\theta$ real (these are exactly polynomials for which $P_{n}^{*}(z)=e^{-2 i \varphi} P_{n}(z)$ where * is the Szegő dual). Equivalently, zeros are symmetric about $\partial \mathbb{D}$.
(b) Monic $P_{n}$ all of whose zeros lie on $\partial \mathbb{D}$. These Chebyshev-like polynomials are used in [100].
Since there are classes of polynomials between all monic and monic with zeros on $E$, the $n$th roots of the norms also converge to $C(E)$.

$$
\left\|T_{n}\right\|_{E}^{1 / n} \leq\left\|T_{n}^{R}\right\|_{E}^{1 / n} \text { is (B.2). Here is the last inequality in (B.8): }
$$

## Proposition B.2.

$$
\begin{equation*}
\left\|T_{n}^{R}\right\|_{E} \leq \zeta_{n+1}(E)^{n} \tag{B.10}
\end{equation*}
$$

Proof. Let $\left\{z_{j}\right\}_{j=1}^{n+1}$ be an $(n+1)$-Fekete set. Let

$$
\begin{equation*}
P_{k}(z)=\prod_{j \neq k}\left(z-z_{j}\right) \tag{B.11}
\end{equation*}
$$

called a Fekete polynomial. (Note: There is a different set of polynomials occurring in a different context also called Fekete polynomials.) By the maximizing property of Fekete sets,

$$
\begin{equation*}
\left\|P_{k}\right\|_{E}=\prod_{j \neq k}\left|z_{k}-z_{j}\right| \tag{B.12}
\end{equation*}
$$

since if $z_{j}^{\prime}=z_{j}(j \neq k), z_{k}^{\prime}=z$, then $\prod_{\ell \neq k}\left|z_{\ell}^{\prime}-z_{k}^{\prime}\right| \leq \prod_{\ell \neq k}\left|z_{\ell}-z_{k}\right|$. Since $\left\|T_{n}^{R}\right\|_{E} \leq\left\|P_{k}\right\|_{E}$ (by the minimizing property of $\left\|T_{n}^{R}\right\|_{E}$ ), taking the $n+1$ choices of $k$,

$$
\left\|T_{n}^{R}\right\|_{E}^{n+1} \leq \prod_{k=1}^{n+1}\left\|P_{k}\right\|_{E}=\prod_{\text {all }}^{j \neq k}| | z_{k}-z_{j} \mid=\zeta_{n+1}^{n(n+1)}
$$

which is (B.10).
The following completes the proof of (B.8):
Proposition B.3. For any monic polynomial $P_{n}(z)$,

$$
\begin{equation*}
\left\|P_{n}\right\|_{E} \geq C(E)^{n} \tag{B.13}
\end{equation*}
$$

Proof. There is nothing to prove if $C(E)=0$, so suppose $C(E)>0$. By the Bernstein-Walsh lemma (A.30),

$$
\begin{equation*}
\left|P_{n}(z)\right| \leq\left\|P_{n}\right\|_{E} C(E)^{-n} \exp \left(-n \Phi_{\rho_{E}}(z)\right) \tag{B.14}
\end{equation*}
$$

Divide by $|z|^{n}$ and take $z \rightarrow \infty$. The left side of (B.14) goes to 1 . Since $\Phi_{\rho_{E}}(z)=-\log |z|+o(1)$, the right side goes to $\left\|P_{n}\right\|_{E} C(E)^{-n}$.

Next we turn to the convergence of Fekete set counting measures to $d \rho_{E}$.

Proposition B.4. Let $d \nu_{n}$ be finite point probability measures supported at $\left\{z_{j}^{(n)}\right\}_{j=1}^{N_{n}}$ with weight $\nu_{n, j}=\nu\left(\left\{z_{j}^{(n)}\right\}\right)$. Suppose $d \nu_{n} \rightarrow d \eta$ weakly for some measure $\eta$. Suppose there is a compact $K \subset \mathbb{C}$ containing all the $\left\{z_{j}^{(n)}\right\}$ and that as $n \rightarrow \infty$,

$$
\begin{equation*}
\sum_{j} \nu_{n, j}^{2} \rightarrow 0 \tag{B.15}
\end{equation*}
$$

Then

$$
\begin{equation*}
\limsup _{n \rightarrow \infty} \prod_{j \neq k}\left|z_{j}^{(n)}-z_{k}^{(n)}\right|^{\nu_{n, j} \nu_{n, k}} \leq \exp \left(\int d \eta(z) d \eta(w) \log |z-w|\right) \tag{B.16}
\end{equation*}
$$

Remark. Since $\sum_{j} \nu_{n, j}=1$,

$$
\begin{equation*}
\left(\max _{j} \nu_{n, j}\right)^{2} \leq \sum_{j} \nu_{n, j}^{2} \leq \max \nu_{n, j} \tag{B.17}
\end{equation*}
$$

(B.15) is equivalent to

$$
\max _{j} \nu_{n, j} \rightarrow 0
$$

Proof. Fix $m \geq 0$ and let

$$
\begin{equation*}
g_{m}(z, w) \equiv \log (\max (m,|z-w|)) \tag{B.18}
\end{equation*}
$$

Then

$$
\begin{equation*}
m^{\sum_{j} \nu_{n, j}^{2}} \prod_{j \neq k}\left|z_{j}^{(n)}-z_{k}^{(n)}\right|^{\nu_{n, j} \nu_{n, k}} \leq \exp \left(\int d \nu_{n}(z) d \nu_{n}(w) g_{m}(z, w)\right) \tag{B.19}
\end{equation*}
$$

Now take $n \rightarrow \infty$. By (B.15), $m^{\sum \nu_{n, j}^{2}} \rightarrow 1$, and by continuity of $g_{m}(z, w)$ and the weak convergence,

$$
\begin{equation*}
\int d \nu_{n}(z) d \nu_{n}(w) g_{m}(z, w) \rightarrow \int d \eta(z) d \eta(w) g_{m}(z, w) \tag{B.20}
\end{equation*}
$$

we have

$$
\begin{equation*}
\text { LHS of }(\mathrm{B} .16) \leq \exp \left(\int d \eta(x) d \eta(y) g_{m}(z, w)\right) \tag{B.21}
\end{equation*}
$$

Now take $m \rightarrow 0$ using monotone convergence to get (B.16).
Lemma B.5. Let $E \subset \mathbb{R}$. Let $(a, b) \cap E=\emptyset$. Then $T_{n}(z)$ has at most one zero in $(a, b)$ which is simple. If $(a, b) \cap \operatorname{cvh}(E)=\emptyset, T_{n}$ has no zero in $(a, b)$ (where $\operatorname{cvh}(E)$ is the convex hull of $E)$. In particular, if $d \eta$ is a limit point of the normalized zero counting measure for $T_{n}$, then $\operatorname{supp}(d \eta) \subset E$.

Proof. Suppose $x_{1}, x_{2}$ are two zeros in $(a, b)$ with $x_{1}<x_{2}$. Then

$$
\left(z-\left(x_{1}-\delta\right)\right)\left(z-\left(x_{2}+\delta\right)\right)=\left(z-x_{1}\right)\left(z-x_{2}\right)-\delta\left(x_{2}-x_{1}\right)-\delta^{2}
$$

so uniformly on $E$ where $\left(z-x_{1}\right)\left(z-x_{2}\right)>0$,

$$
\left|\left(z-\left(x_{1}-\delta\right)\right)\left(z-\left(x_{2}+\delta\right)\right)\right|<\left|\left(z-x_{1}\right)\left(z-x_{2}\right)\right|
$$

for $\delta$ small. Thus, $\left\|T_{n}(z)\right\|_{E}$ is decreased by changing those two zeros. Similarly, if $x$ is a zero below $\operatorname{cvh}(E), T_{n}$ is decreased by moving the
zero up slightly. If $x_{j}$ is a complex zero, $\left\|T_{n}(z)\right\|_{E}$ is decreased by replacing $x_{j}$ by $\operatorname{Re} x_{j}$.

The final statement is immediate if we note that if $f$ is a continuous function supported in $(a, b)$, then $\int f d \eta=0$.
Proof of Theorem B.1. We have already proved (B.8). The Fekete points are distinct, so $\nu_{n, j}=1 / n$, in the language of Proposition B.4. So if we pass to a subsequence for which $d \nu_{n(j)}$ has a weak limit $\eta$, we see (using $\lim \zeta_{n}$ exists)

$$
\begin{align*}
\lim \zeta_{n}=\lim _{j \rightarrow \infty} \zeta_{n(j)}^{(n(j)-1) / n(j)} & \leq \exp (-\mathcal{E}(d \eta)) \\
& \leq \exp \left(-\inf _{\text {all } d \rho} \mathcal{E}(d \rho)\right)  \tag{B.22}\\
& =C(E)
\end{align*}
$$

By (B.8), $\lim \zeta_{n} \geq C(E)$, so we have equality in (B.22) and $d \eta=d \rho_{E}$. Thus, any limit point is $d \rho_{E}$. By compactness, we have (i).

That leaves the proof of (ii). By the Berstein-Walsh lemma (A.30), for $z \in \mathbb{C} \backslash E$,

$$
\begin{equation*}
\frac{1}{n} \log \left|T_{n}(z)\right| \leq \log \left(\frac{\left\|T_{n}\right\|_{E}^{1 / n}}{C(E)}\right)-\Phi_{\rho_{E}}(z) \tag{B.23}
\end{equation*}
$$

and similarly for $T_{n}^{R}$.
Now let $d \eta$ be a limit point of the normalized density of zeros of $T_{n}(z)$. By the last lemma, $d \eta$ is supported on $E$, so (B.23) plus $\lim \left\|T_{n}\right\|_{\infty}^{1 / n}=$ $C(E)$ implies

$$
\begin{equation*}
\Phi_{\eta}(z) \geq \Phi_{\rho_{E}}(z) \tag{B.24}
\end{equation*}
$$

for $z \in \mathbb{C} \backslash E$. By Theorem A.21, this implies $d \eta=d \rho_{E}$. Thus, $d \rho_{E}$ is the only limit point of the zeros, and so the limit is $d \rho_{E}$.

Note added in proof. Since completion of this manuscript, I have found a result (to appear in "Regularity and the Cesàro-Nevai class", in prep.) relevant to the subject of the current review. In the simplest case, it states that if a measure has $[-2,2]$ as its essential support and is regular, then (1.25) holds.

## References

[1] S. Agmon, Lectures on Exponential Decay of Solutions of Second-Order Elliptic Equations: Bounds on Eigenfunctions of $N$-body Schrödinger Operators, Mathematical Notes, 29, Princeton Univ. Press, Princeton, N.J.; Univ. of Tokyo Press, Tokyo, 1982.
[2] P. W. Anderson, Absence of diffusion in certain random lattices, Phys. Rev. 109 (1958), 1492-1505.
[3] V. V. Andrievskii and H.-P. Blatt, Discrepancy of Signed Measures and Polynomial Approximation, Springer Monographs in Mathematics, SpringerVerlag, New York, 2002.
[4] F. V. Atkinson, On the location of the Weyl circles, Proc. Roy. Soc. Edinburgh Sect. A 88 (1981), 345-356.
[5] S. Aubry, Metal insulator transition in one-dimensional deformable lattices, in "Bifurcation Phenomena in Mathematical Physics and Related Topics," (C. Bardos and D. Bessis, eds.), pp. 163-184, NATO Advanced Study Institute Series, Ser. C: Mathematical and Physical Sciences, 54, D. Reidel Publishing, Dordrecht-Boston, 1980.
[6] A. Avila, in preparation.
[7] A. Avila, in preparation.
[8] A. Avila and D. Damanik, Absolute continuity of the integrated density of states for the almost Mathieu operator with non-critical coupling, preprint.
[9] A. Avila and S. Jitomirskaya, Almost localization and almost reducibility, preprint.
[10] A. Avila and R. Krikorian, Reducibility or nonuniform hyperbolicity for quasiperiodic Schrödinger cocycles, Annals of Math. 164 (2006), 911-940.
[11] J. Avron and B. Simon, Almost periodic Schrödinger operators. I. Limit periodic potentials, Comm. Math. Phys. 82 (1981/82), 101-120.
[12] J. Avron and B. Simon, Almost periodic Schrödinger operators. II. The integrated density of states, Duke Math. J. 50 (1983), 369-391.
[13] M. Benderskii and L. Pastur, The spectrum of the one-dimensional Schrödinger equation with a random potential, Math. USSR Sb. 11 (1970), 245-256; Russian original in Mat. Sb. (N.S.) 82(124) (1970), 273-284.
[14] Ju. M. Berezanskiǐ, Expansions in eigenfunctions of selfadjoint operators, Translations of Mathematical Monographs, 17, American Mathematical Society, Providence, R.I., 1968.
[15] P. Borwein and T. Erdélyi, Polynomials and Polynomial Inequalities, Graduate Texts in Mathematics, 161, Springer-Verlag, New York, 1995.
[16] J. Bourgain, On the spectrum of lattice Schrödinger operators with deterministic potential, J. Anal. Math. 87 (2002), 37-75.
[17] J. Bourgain, On the spectrum of lattice Schrödinger operators with deterministic potential. II., J. Anal. Math. 88 (2002), 221-254.
[18] J. Bourgain, Green's Function Estimates for Lattice Schrödinger Operators and Applications, Annals of Mathematics Studies, 158, Princeton University Press, Princeton, N.J., 2005.
[19] J. Bourgain and S. Jitomirskaya, Continuity of the Lyapunov exponent for quasiperiodic operators with analytic potential, J. Statist. Phys. 108 (2002), 1203-1218.
[20] F. E. Browder, Eigenfunction expansions for formally self-adjoint partial differential operators. I, Proc. Nat. Acad. Sci. U.S.A. 42 (1956), 769-771.
[21] R. Carmona and J. Lacroix, Spectral Theory of Random Schrödinger Operators. Probability and Its Applications, Birkhäuser, Boston, 1990.
[22] J. R. Choksi and M. G. Nadkarni, Genericity of certain classes of unitary and self-adjoint operators, Canad. Math. Bull. 41 (1998), 137-139.
[23] S. Clark and F. Gesztesy, On Povzner-Wienholtz-type self-adjointness results for matrix-valued Sturm-Liouville operators, Proc. Math. Soc. Edinburgh 133A (2003), 747-758.
[24] W. Craig, The trace formula for Schrödinger operators on the line, Comm. Math. Phys. 126 (1989), 379-407.
[25] W. Craig and B. Simon, Subharmonicity of the Lyaponov index, Duke Math. J. 50 (1983), 551-560.
[26] H. L. Cycon, R. G. Froese, W. Kirsch, and B. Simon, Schrödinger Operators With Application to Quantum Mechanics and Global Geometry, Texts and Monographs in Physics, Springer, Berlin, 1987.
[27] D. Damanik, Lyapunov exponents and spectral analysis of ergodic Schrödinger operators: A survey of Kotani theory and its applications, in "Spectral Theory and Mathematical Physics: A Festschrift in Honor of Barry Simon's 60th Birthday" (F. Gesztesy et al., eds.), pp. 539-563, Proc. Symp. Pure Math., 76.2, American Mathematical Society, Providence, R.I., 2007.
[28] D. Damanik, Strictly ergodic subshifts and associated operators, in "Spectral Theory and Mathematical Physics: A Festschrift in Honor of Barry Simon's 60th Birthday" (F. Gesztesy et al., eds.), pp. 505-538, Proc. Symp. Pure Math., 76.2, American Mathematical Society, Providence, R.I., 2007.
[29] D. Damanik, R. Killip, and B. Simon, Perturbations of orthogonal polynomials with periodic recursion coefficients, preprint.
[30] P. A. Deift and B. Simon, Almost periodic Schrödinger operators, III. The absolutely continuous spectrum in one dimension, Comm. Math. Phys. 90 (1983), 389-411.
[31] R. del Rio, S. Jitomirskaya, Y. Last, and B. Simon, Operators with singular continuous spectrum, IV. Hausdorff dimensions, rank one perturbations, and localization, J. Anal. Math. 69 (1996), 153-200.
[32] F. Delyon, B. Simon, and B. Souillard, From power pure point to continuous spectrum in disordered systems, Ann. Inst. H. Poincaré 42 (1985), 283-309.
[33] J. L. Doob, Measure Theory, Graduate Texts in Mathematics, 143, SpringerVerlag, New York, 1994.
[34] P. Erdös and P. Turán, On interpolation. III. Interpolatory theory of polynomials, Annals of Math. (2) 41 (1940), 510-553.
[35] G. Faber, Über Tschebyscheffsche Polynome, J. Reine Angew. Math. 150 (1919), 79-106.
[36] A. Fedotov and F. Klopp, Strong resonant tunneling, level repulsion and spectral type for one-dimensional adiabatic quasi-periodic Schrödinger operators, Ann. Sci. École Norm. Sup. (4) 38 (2005), 889-950.
[37] A. Fedotov and F. Klopp, Weakly resonant tunneling interactions for adiabatic quasi-periodic Schrödinger operators, Mém. Soc. Math. Fr. (N.S.), No. 104 (2006).
[38] M. Fekete, Über die Verteilung der Wurzeln bei gewissen algebraischen Gleichungen mit ganzzahligen Koeffizienten, Math. Z. 17 (1923), 228-249.
[39] G. Freud, Orthogonal Polynomials, Pergamon Press, Oxford-New York, 1971.
[40] L. Gårding, Eigenfunction expansions connected with elliptic differential operators, Tolfte Skandinaviska Matematikerkongressen, Lund, 1953, pp. 44-55, Lunds Universitets Matematiska Institution, Lund, 1954.
[41] I. M. Gel'fand, Expansion in characteristic functions of an equation with periodic coefficients, Dokl. Akad. Nauk SSSR (N.S.) 73, (1950), 1117-1120 (Russian).
[42] J. S. Geronimo, Polynomials orthogonal on the unit circle with random recurrence coefficients, in "Methods of Approximation Theory in Complex Analysis and Mathematical Physics" (Leningrad, 1991), pp. 43-61, Lecture Notes in Math., 1550, Springer, Berlin, 1993.
[43] Ya. L. Geronimus, On certain asymptotic properties of polynomials, Mat. Sb . N. S. 23(65), (1948), 77-88 (Russian).
[44] Ya. L. Geronimus, Orthogonal Polynomials: Estimates, Asymptotic Formulas, and Series of Polynomials Orthogonal on the Unit Circle and on an Interval, Consultants Bureau, New York, 1961.
[45] F. Gesztesy and B. Simon, A new approach to inverse spectral theory, II. General real potentials and the connection to the spectral measure, Annals of Math. 152 (2000), 593-643.
[46] L. Golinskii and S. Khrushchev, Cesàro asymptotics for orthogonal polynomials on the unit circle and classes of measures, J. Approx. Theory 115 (2002), 187-237.
[47] L. Golinskii and P. Nevai, Szegő difference equations, transfer matrices and orthogonal polynomials on the unit circle, Comm. Math. Phys. 223 (2001), 223-259.
[48] G. M. Goluzin, Geometric Theory of Functions of a Complex Variable, Translations of Mathematical Monographs, 26, American Mathematical Society, Providence, R.I., 1969.
[49] A. Gordon, S. Jitomirskaya, Y. Last and B. Simon, Duality and singular continuous spectrum in the almost Mathieu equation, Acta Math. 178 (1997), 169-183.
[50] L. L. Helms, Introduction to Potential Theory, Pure and Applied Mathematics, 22, Wiley-Interscience, New York, 1969.
[51] D. Herbert and R. Jones, Localized states in disordered systems, J. Phys. C: Solid State Phys. 4 (1971), 1145-1161.
[52] M. Herman, Une méthode pour minorer les exposants de Lyapounov et quelques exemples montrant le caractre local d'un théorème d'Arnol'd et de Moser sur le tore de dimension 2, Comment. Math. Helv. 58 (1983), 453-502.
[53] K. Ishii, Localization of eigenstates and transport phenomena in the onedimensional disordered system, Supp. Prog. Theor. Phys. 53 (1973), 77-138.
[54] K. Ishii and H. Matsuda, Localization of normal modes and energy transport in the disordered harmonic chain, Suppl. Prog. Theor. Phys. 45 (1970), 56-86.
[55] S. Jitomirskaya, Metal-insulator transition for the almost Mathieu operator, Annals of Math. 150 (1999), 1159-1175.
[56] S. Jitomirskaya, Ergodic Schrödinger Operators (on one foot), in "Spectral Theory and Mathematical Physics: A Festschrift in Honor of Barry Simon's 60th Birthday" (F. Gesztesy et al., eds.), pp. 613-647, Proc. Symp. Pure Math., 76.2, American Mathematical Society, Providence, R.I., 2007.
[57] S. Jitomirskaya, private communication.
[58] S. Jitomirskaya and M. Landrigan, Zero-dimensional spectral measures for quasi-periodic operators with analytic potential, J. Statist. Phys. 100 (2000), 791-796.
[59] S. Jitomirskaya and Y. Last, Power law subordinacy and singular spectra, II. Line operators, Comm. Math. Phys. 211 (2000), 643-658.
[60] S. Jitomirskaya and B. Simon, Operators with singular continuous spectrum: III. Almost periodic Schrödinger operators, Comm. Math. Phys. 165 (1994), 201-205.
[61] K. Johansson, On Szegő's asymptotic formula for Toeplitz determinants and generalizations, Bull. Sci. Math. (2) 112 (1988), 257-304.
[62] G. I. Kac, Expansion in characteristic functions of self-adjoint operators, Dokl. Akad. Nauk SSSR (N.S.) 119 (1958), 19-22 (Russian).
[63] L. Kalmár, Az interpolatiorol, Mathematikai es Fizikai Lapok, 1927, pp. 120149 (Hungarian).
[64] Y. Katznelson, An Introduction to Harmonic Analysis, 2nd corrected edition, Dover Publications, New York, 1976.
[65] R. Killip, A. Kiselev, and Y. Last, Dynamical upper bounds on wavepacket spreading, Am. J. Math. 125 (2003), 1165-1198.
[66] W. Kirsch and B. Metzger, The integrated density of states and random Schrödinger operators, in "Spectral Theory and Mathematical Physics: A Festschrift in Honor of Barry Simon's 60th Birthday" (F. Gesztesy et al., eds.), pp. 649-696, Proc. Symp. Pure Math., 76.2, American Mathematical Society, Providence, R.I., 2007.
[67] 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), 1-45.
[68] P. P. Korovkin, The asymptotic representation of polynomials orthogonal over a region, Dokl. Akad. Nauk SSSR (N.S.) 58 (1947), 1883-1885 (Russian).
[69] S. Kotani, Ljapunov indices determine absolutely continuous spectra of stationary random one-dimensional Schrödinger operators, in "Stochastic Analysis" (Katata/Kyoto, 1982), pp. 225-247, North-Holland Math. Library, 32, North-Holland, Amsterdam, 1984.
[70] S. Kotani, Jacobi matrices with random potentials taking finitely many values, Rev. Math. Phys. 1 (1989), 129-133.
[71] S. Kotani, Generalized Floquet theory for stationary Schrödinger operators in one dimension, Chaos Solitons Fractals 8 (1997), 1817-1854.
[72] H. Kunz and B. Souillard, Sur le spectre des opérateurs aux différences finies aléatoires, Comm. Math. Phys. 78 (1980/81), 201-246.
[73] N. S. Landkof, Foundations of Modern Potential Theory, Springer-Verlag, Berlin-New York, 1972.
[74] Y. Last, On the measure of gaps and spectra for discrete 1D Schrödinger operators, Comm. Math. Phys. 149 (1992), 347-360.
[75] Y. Last and B. Simon, Eigenfunctions, transfer matrices, and absolutely continuous spectrum of one-dimensional Schrödinger operators, Invent. Math. 135 (1999), 329-367.
[76] P. D. Lax, Functional Analysis, Wiley-Interscience, New York, 2002.
[77] D. Lenz and P. Stollmann, Generic sets in spaces of measures and generic singular continuous spectrum for Delone Hamiltonians, Duke Math. J. 131 (2006), 203-217.
[78] D. S. Lubinsky, Jump distributions on $[-1,1]$ whose orthogonal polynomials have leading coefficients with given asymptotic behavior, Proc. Amer. Math. Soc. 104 (1988), 516-524.
[79] D. S. Lubinsky, Singularly continuous measures in Nevai's class M, Proc. Amer. Math. Soc. 111 (1991), 413-420.
[80] A. P. Magnus and W. Van Assche, Sieved orthogonal polynomials and discrete measures with jumps dense in an interval, Proc. Amer. Math. Soc. 106 (1989), 163-173.
[81] A. Martínez-Finkelshtein, Equilibrium problems of potential theory in the complex plane, in "Orthogonal Polynomials and Special Functions," pp. 79117, Lecture Notes in Math., 1883, Springer, Berlin, 2006.
[82] H. N. Mhaskar and E. B. Saff, On the distribution of zeros of polynomials orthogonal on the unit circle, J. Approx. Theory 63 (1990), 30-38.
[83] S. Molchanov, Multiscale averaging for ordinary differential equations. Applications to the spectral theory of one-dimensional Schrödinger operator with sparse potentials, in "Homogenization," pp. 316-397, Ser. Adv. Math. Appl. Sci., 50, World Scientific Publishing, River Edge, N.J., 1999.
[84] S. N. Naboko, Schrödinger operators with decreasing potential and with dense point spectrum, Soviet Math. Dokl. 29 (1984), 688-691; Russian original in Dokl. Akad. Nauk SSSR 276 (1984), 1312-1315.
[85] L. Pakula, Asymptotic zero distribution of orthogonal polynomials in sinusoidal frequency estimation, IEEE Trans. Inform. Theory 33 (1987), 569-576.
[86] L. A. Pastur, Spectral properties of disordered systems in the one-body approximation Comm. Math. Phys. 75 (1980), 179-196.
[87] L. Pastur and A. Figotin, Spectra of Random and Almost-Periodic Operators, Grundlehren der Mathematischen Wissenschaften [Fundamental Principles of Mathematical Sciences], 297, Springer, Berlin, 1992.
[88] T. Ransford, Potential Theory in the Complex Plane, Press Syndicate of the Univesity of Cambridge, New York, 1995.
[89] E. J. Remez, Sur une propriété des polynômes de Tchebyscheff, Comm. Inst. Sci. Kharkow 13 (1936), 93-95.
[90] C. Remling, The absolutely continuous spectrum of Jacobi matrices, preprint.
[91] E. B. Saff and V. Totik, Logarithmic Potentials With External Fields, Grundlehren der Mathematischen Wissenschaften [Fundamental Principles of Mathematical Sciences], 316, Springer-Verlag, Berlin, 1997.
[92] J. C. Santos-León, A Szegő quadrature formula for a trigonometric polynomial modification of the Lebesgue measure, Rev. Acad. Canaria Cienc. 11 (1999), 183-191.
[93] J. A. Shohat, Théorie Générale des Polinomes Orthogonaux de Tchebichef, Mémorial des Sciences Mathématiques, 66, pp. 1-69, Paris, 1934.
[94] B. Simon, Schrödinger semigroups, Bull. Amer. Math. Soc. 7 (1982), 447-526.
[95] B. Simon, Some Jacobi matrices with decaying potential and dense point spectrum, Comm. Math. Phys. 87 (1982), 253-258.
[96] B. Simon, Kotani theory for one dimensional stochastic Jacobi matrices, Comm. Math. Phys. 89 (1983), 227-234.
[97] B. Simon, Operators with singular continuous spectrum, I. General operators Annals of Math. 141 (1995), 131-145.
[98] B. Simon, Orthogonal Polynomials on the Unit Circle, Part 1: Classical Theory, AMS Colloquium Series, 54.1, American Mathematical Society, Providence, R.I., 2005.
[99] B. Simon, Orthogonal Polynomials on the Unit Circle, Part 2: Spectral Theory, AMS Colloquium Series, 54.2, American Mathematical Society, Providence, R.I., 2005.
[100] B. Simon, Weak convergence of CD kernels and applications, preprint.
[101] B. Simon, The Christoffel-Darboux kernel, in preparation.
[102] B. Simon, Szegő's Theorem and Its Descendants: Spectral Theory for L ${ }^{2}$ Perturbations of Orthogonal Polynomials, in preparation; to be published by Princeton University Press.
[103] B. Simon and M. Taylor, Harmonic analysis on $S L(2, \mathbb{R})$ and smoothness of the density of states in the one-dimensional Anderson model, Comm. Math. Phys. 101 (1985), 1-19.
[104] M. Sodin and P. Yuditskii, Almost periodic Jacobi matrices with homogeneous spectrum, infinite-dimensional Jacobi inversion, and Hardy spaces of character-automorphic functions, J. Geom. Anal. 7 (1997), 387-435.
[105] H. Stahl and V. Totik, General Orthogonal Polynomials, in "Encyclopedia of Mathematics and its Applications," 43, Cambridge University Press, Cambridge, 1992.
[106] T. Stieltjes, Sur les racines de l'équation $X_{n}=0$, Acta Math. 9 (1886), 385-400; Oeuvres Complètes 2, 73-88.
[107] G. Szegő, Bemerkungen zu einer Arbeit von Herrn M. Fekete: Über die Verteilung der Wurzeln bei gewissen algebraischen Gleichungen mit ganzzahligen Koeffizienten, Math. Z. 21 (1924), 203-208.
[108] G. Szegő, Orthogonal Polynomials, Amer. Math. Soc. Colloq. Publ., 23, American Mathematical Society, Providence, R.I., 1939; 3rd edition, 1967.
[109] A. V. Teplyaev, The pure point spectrum of random orthogonal polynomials on the circle, Soviet Math. Dokl. 44 (1992), 407-411; Russian original in Dokl. Akad. Nauk SSSR 320 (1991), 49-53.
[110] D. J. Thouless, Electrons in disordered systems and the theory of localization, Phys. Rep. 13 (1974), 93.
[111] V. Totik, Orthogonal polynomials with ratio asymptotics, Proc. Amer. Math. Soc. 114 (1992), 491-495.
[112] V. Totik and J. L. Ullman, Local asymptotic distribution of zeros of orthogonal polynomials, Trans. Amer. Math. Soc. 341 (1994), 881-894.
[113] M. Tsuji, Potential Theory in Modern Function Theory, reprinting of the 1959 original, Chelsea, New York, 1975.
[114] J. L. Ullman, On the regular behaviour of orthogonal polynomials, Proc. London Math. Soc. (3) 24 (1972), 119-148.
[115] J. L. Ullman, Orthogonal polynomials for general measures. I, in "Rational Approximation and Interpolation" (Tampa, FL, 1983), pp. 524-528, Lecture Notes in Math., 1105, Springer, Berlin, 1984.
[116] J. L. Ullman, Orthogonal polynomials for general measures. II, in "Orthogonal Polynomials and Applications" (Bar-le-Duc, 1984), pp. 247-254, Lecture Notes in Math., 1171, Springer, Berlin, 1985.
[117] J. L. Ullman, Orthogonal polynomials for general measures, in "Orthogonal Polynomials and Their Applications" (Laredo, 1987), pp. 95-99, Lecture Notes in Pure and Appl. Math., 117, Dekker, New York, 1989.
[118] J. L. Ullman and M. F. Wyneken, Weak limits of zeros of orthogonal polynomials, Constr. Approx. 2 (1986), 339-347.
[119] J. L. Ullman, M. F. Wyneken, and L. Ziegler, Norm oscillatory weight measures, J. Approx. Theory 46 (1986), 204-212.
[120] W. Van Assche, Invariant zero behaviour for orthogonal polynomials on compact sets of the real line, Bull. Soc. Math. Belg. Ser. B 38 (1986), 1-13.
[121] J. L. Walsh, Interpolation and Approximation by Rational Functions in the Complex Domain, fourth edition, American Mathematical Society Colloquium Publications, XX, American Mathematical Society, Providence, R.I., 1965.
[122] H. Widom, Polynomials associated with measures in the complex plane, J. Math. Mech. 16 (1967), 997-1013.
[123] M. F. Wyneken, On norm and zero asymptotics of orthogonal polynomials for general measures. $I$, in "Constructive Function Theory" (Edmonton, AB, 1986), Rocky Mountain J. Math. 19 (1989), 405-413.


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