# The Analytic Theory of Matrix Orthogonal Polynomials 

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#### Abstract

We survey the analytic theory of matrix orthogonal polynomials. MSC: 42C05, 47B36, 30C10 keywords: orthogonal polynomials, matrix-valued measures, block Jacobi matrices, block CMV matrices


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## 1 Introduction

### 1.1 Introduction and Overview

Orthogonal polynomials on the real line (OPRL) were developed in the nineteenth century and orthogonal polynomials on the unit circle (OPUC) were initially developed around 1920 by Szegő. Their matrix analogues are of much more recent vintage. They were originally developed in the MOPUC case indirectly in the study of prediction theory [116, 117, 129, 131, 132, 138, 196] in the period 1940-1960. The connection to OPUC in the scalar case was discovered by Krein [131].

Much of the theory since is in the electrical engineering literature [36, 37, 38, 39, 40, 41, 120, 121, $122,123,203]$; see also [84, 86, 87, 88, 142].

The corresponding real line theory (MOPRL) is still more recent: Following early work of Krein [133] and Berezan'ski [9] on block Jacobi matrices, mainly as applied to self-adjoint extensions, there was a seminal paper of Aptekarev-Nikishin [4] and a flurry of papers since the 1990s [10, 11, 12, $14,16,17,19,20,21,22,29,35,43,45,46,47,48,49,50,51,52,53,54,55,56,57,58,59,60,61$, $62,64,65,66,67,68,69,71,73,74,75,76,77,79,83,85,102,103,104,105,106,107,108,109$, $110,111,112,113,137,139,140,143,144,145,148,149,150,155,156,157,154,161,162,179$, 186, 198, 200, 201, 202, 204]; see also [7].

There is very little on the subject in monographs - the more classical ones (e.g., [23, 82, 93, 184]) predate most of the subject; see, however, Atkinson [5, Section 6.6]. Ismail [118] has no discussion and Simon $[167,168]$ has a single section! Because of the use of MOPRL in [33], we became interested in the subject and, in particular, we needed some basic results for that paper which we couldn't find in the literature or which, at least, weren't very accessible. Thus, we decided to produce this comprehensive review that we hope others will find useful.

As with the scalar case, the subject breaks into two parts, conveniently called the analytic theory (general structure results) and the algebraic theory (the set of non-trivial examples). This survey deals entirely with the analytic theory. We note, however, that one of the striking developments in recent years has been the discovery that there are rich classes of genuinely new MOPRL, even at the classical level of Bochner's theorem; see $[20,55,70,72,102,109,110,111,112,113,156,161]$ and the forthcoming monograph [63] for further discussion of this algebraic side.

In this introduction, we will focus mainly on the MOPRL case. For scalar OPRL, a key issue is the passage from measure to monic OPRL, then to normalized OPRL, and finally to Jacobi parameters. There are no choices in going from measure to monic OP, $P_{n}(x)$. They are determined by

$$
\begin{equation*}
P_{n}(x)=x^{n}+\text { lower order }, \quad\left\langle x^{j}, P_{n}\right\rangle=0 \quad j=1, \ldots, n-1 . \tag{1.1}
\end{equation*}
$$

However, the basic condition on the orthonormal polynomials, namely,

$$
\begin{equation*}
\left\langle p_{n}, p_{m}\right\rangle=\delta_{n m} \tag{1.2}
\end{equation*}
$$

does not uniquely determine the $p_{n}(x)$. The standard choice is

$$
p_{n}(x)=\frac{P_{n}(x)}{\left\|P_{n}\right\|} .
$$

However, if $\theta_{0}, \theta_{1}, \ldots$ are arbitrary real numbers, then

$$
\begin{equation*}
\tilde{p}_{n}(x)=\frac{e^{i \theta_{n}} P_{n}(x)}{\left\|P_{n}\right\|} \tag{1.3}
\end{equation*}
$$

also obey (1.2). If the recursion coefficients (aka Jacobi parameters), are defined via

$$
\begin{equation*}
x p_{n}=a_{n+1} p_{n+1}+b_{n+1} p_{n}+a_{n} p_{n-1}, \tag{1.4}
\end{equation*}
$$

then the choice (1.3) leads to

$$
\begin{equation*}
\tilde{b}_{n}=b_{n}, \quad \tilde{a}_{n}=e^{i \theta_{n}} a_{n} e^{-i \theta_{n-1}} . \tag{1.5}
\end{equation*}
$$

The standard choice is, of course, most natural here; for example, if

$$
\begin{equation*}
p_{n}(x)=\kappa_{n} x^{n}+\text { lower order }, \tag{1.6}
\end{equation*}
$$

then $a_{n}>0$ implies $\kappa_{n}>0$. It would be crazy to make any other choice.
For MOPRL, these choices are less clear. As we will explain in Section 1.2, there are now two matrix-valued "inner products" formally written as

$$
\begin{align*}
\langle\langle f, g\rangle\rangle_{R} & =\int f(x)^{\dagger} d \mu(x) g(x),  \tag{1.7}\\
\left\langle\langle f, g\rangle_{L}\right. & =\int g(x) d \mu(x) f(x)^{\dagger} \tag{1.8}
\end{align*}
$$

where now $\mu$ is a matrix-valued measure and ${ }^{\dagger}$ denotes the adjoint, and corresponding two sets of monic OPRL: $P_{n}^{R}(x)$ and $P_{n}^{L}(x)$. The orthonormal polynomials are required to obey

$$
\begin{equation*}
\left\langle\left\langle p_{n}^{R}, p_{m}^{R}\right\rangle\right\rangle_{R}=\delta_{n m} \mathbf{1} \tag{1.9}
\end{equation*}
$$

The analogue of (1.3) is

$$
\begin{equation*}
\tilde{p}_{n}^{R}(x)=P_{n}^{R}(x)\left\langle\left\langle P_{n}^{R}, P_{n}^{R}\right\rangle\right\rangle^{-1 / 2} \sigma_{n} \tag{1.10}
\end{equation*}
$$

for a unitary $\sigma_{n}$. For the immediately following, use $p_{n}^{R}$ to be the choice $\sigma_{n} \equiv 1$. For any such choice, we have a recursion relation,

$$
\begin{equation*}
x p_{n}^{R}(x)=p_{n+1}^{R}(x) A_{n+1}^{\dagger}+p_{n}^{R}(x) B_{n+1}+p_{n-1}^{R}(x) A_{n} \tag{1.11}
\end{equation*}
$$

with the analogue of (1.5) (comparing $\sigma_{n} \equiv \mathbf{1}$ to general $\sigma_{n}$ )

$$
\begin{equation*}
\tilde{B}_{n}=\sigma_{n}^{\dagger} B_{n} \sigma_{n} \quad \tilde{A}_{n}=\sigma_{n-1}^{\dagger} A_{n} \sigma_{n} \tag{1.12}
\end{equation*}
$$

The obvious analogue of the scalar case is to pick $\sigma_{n} \equiv \mathbf{1}$, which makes $\kappa_{n}$ in

$$
\begin{equation*}
p_{n}^{R}(x)=\kappa_{n} x^{n}+\text { lower order } \tag{1.13}
\end{equation*}
$$

obey $\kappa_{n}>0$. Note that (1.11) implies

$$
\begin{equation*}
\kappa_{n}=\kappa_{n+1} A_{n+1}^{\dagger} \tag{1.14}
\end{equation*}
$$

or, inductively,

$$
\begin{equation*}
\kappa_{n}=\left(A_{n}^{\dagger} \ldots A_{1}^{\dagger}\right)^{-1} \tag{1.15}
\end{equation*}
$$

In general, this choice does not lead to $A_{n}$ positive or even Hermitian. Alternatively, one can pick $\sigma_{n}$ so $\tilde{A}_{n}$ is positive. Besides these two "obvious" choices, $\kappa_{n}>0$ or $A_{n}>0$, there is a third that $A_{n}$ be lower triangular that, as we will see in Section 1.4, is natural. Thus, in the study of MOPRL one needs to talk about equivalent sets of $p_{n}^{R}$ and of Jacobi parameters, and this is a major theme of Chapter 2. Interestingly enough for MOPUC, the commonly picked choice equivalent to $A_{n}>0$ (namely, $\rho_{n}>0$ ) seems to suffice for applications. So we do not discuss equivalence classes for MOPUC.

Associated to a set of matrix Jacobi parameters is a block Jacobi matrix, that is, a matrix which when written in $l \times l$ blocks is tridiagonal; see (2.29) below.

In Chapter 2, we discuss the basics of MOPRL while Chapter 3 discusses MOPUC. Chapter 4 discusses the Szegő mapping connection of MOPUC and MOPRL. Finally, Chapter 5 discusses the extension of the theory of regular OPs [180] to MOPRL.

While this is mainly a survey, it does have numerous new results, of which we mention:
(a) The clarification of equivalent Jacobi parameters and several new theorems (Theorems 2.8 and 2.9).
(b) A new result (Theorem 2.28) on the order of poles or zeros of $m(z)$ in terms of eigenvalues of $J$ and the once stripped $J^{(1)}$.
(c) Formulas for the resolvent in the MOPRL (Theorem 2.29) and MOPUC (Theorem 3.24) cases.
(d) A theorem on zeros of $\operatorname{det}\left(\Phi_{n}^{R}\right)$ (Theorem 3.7) and eigenvalues of a cutoff CMV matrix (Theorem 3.10).
(e) A new proof of the Geronimus relations (Theorem 4.2).
(f) Discussion of regular MOPRL (Chapter 5).

There are numerous open questions and conjectures in this paper, of which we mention:
(1) We prove that type 1 and type 3 Jacobi parameters in the Nevai class have $A_{n} \rightarrow \mathbf{1}$ but do not know if this is true for type 2 and, if so, how to prove it.
(2) Determine which monic matrix polynomials, $\Phi$, can occur as monic MOPUC. We know $\operatorname{det}(\Phi(z))$ must have all of its zeros in the unit disk in $\mathbb{C}$, but unlike the scalar case where this is sufficient, we do not know necessary and sufficient conditions.
(3) Generalize Khrushchev theory $[125,126,101]$ to MOPUC; see Section 3.13.
(4) Provide a proof of Geronimus relations for MOPUC that uses the theory of canonical moments [43]; see the discussion at the start of Chapter 4.
(5) Prove Conjecture 5.9 extending a result of Stahl-Totik [180] from OPRL to MOPRL.

### 1.2 Matrix-Valued Measures

Let $\mathcal{M}_{l}$ denote the ring of all $l \times l$ complex-valued matrices; we denote by $\alpha^{\dagger}$ the Hermitian conjugate of $\alpha \in \mathcal{M}_{l}$. (Because of the use of * for Szegő dual in the theory of OPUC, we do not use it for adjoint.) For $\alpha \in \mathcal{M}_{l}$, we denote by $\|\alpha\|$ its Euclidean norm (i.e., the norm of $\alpha$ as a linear operator on $\mathbb{C}^{l}$ with the usual Euclidean norm). Consider the set $\mathcal{P}$ of all polynomials in $z \in \mathbb{C}$ with coefficients from $\mathcal{M}_{l}$. The set $\mathcal{P}$ can be considered either as a right or as a left module over $\mathcal{M}_{l}$; clearly, conjugation makes the left and right structures isomorphic. For $n=0,1, \ldots, \mathcal{P}_{n}$ will denote those polynomials in $\mathcal{P}$ of degree at most $n$. The set $\mathcal{V}$ denotes the set of all polynomials in $z \in \mathbb{C}$ with coefficients from $\mathbb{C}^{l}$. The standard inner product in $\mathbb{C}^{l}$ is denoted by $\langle\cdot, \cdot\rangle_{\mathbb{C}^{l}}$.

A matrix-valued measure, $\mu$, on $\mathbb{R}$ (or $\mathbb{C}$ ) is the assignment of a positive semi-definite $l \times l$ matrix $\mu(X)$ to every Borel set $X$ which is countably additive. We will usually normalize it by requiring

$$
\begin{equation*}
\mu(\mathbb{R})=\mathbf{1} \tag{1.16}
\end{equation*}
$$

(or $\mu(\mathbb{C})=\mathbf{1}$ ) where $\mathbf{1}$ is the $l \times l$ identity matrix. (We use $\mathbf{1}$ in general for an identity operator, whether in $\mathcal{M}_{l}$ or in the operators on some other Hilbert space, and $\mathbf{0}$ for the zero operator or matrix.) Normally, our measures for MOPRL will have compact support and, of course, our measures for MOPUC will be supported on all or part of $\partial \mathbb{D}(\mathbb{D}$ is the unit disk in $\mathbb{C})$.

Associated to any such measures is a scalar measure

$$
\begin{equation*}
\mu_{\mathrm{tr}}(X)=\operatorname{Tr}(\mu(X)) \tag{1.17}
\end{equation*}
$$

(the trace normalized by $\operatorname{Tr}(\mathbf{1})=l$ ). $\mu_{\text {tr }}$ is normalized by $\mu_{\text {tr }}(\mathbb{R})=l$.
Applying the Radon-Nikodym theorem to the matrix elements of $\mu$, we see there is a positive semi-definite matrix function $M_{i j}(x)$ so

$$
\begin{equation*}
d \mu_{i j}(x)=M_{i j}(x) d \mu_{\operatorname{tr}}(x) . \tag{1.18}
\end{equation*}
$$

Clearly, by (1.17),

$$
\begin{equation*}
\operatorname{Tr}(M(x))=1 \tag{1.19}
\end{equation*}
$$

for $d \mu_{\mathrm{tr}}$-a.e. $x$. Conversely, any scalar measure with $\mu_{\mathrm{tr}}(\mathbb{R})=l$ and positive semi-definite matrixvalued function $M$ obeying (1.19) define a matrix-valued measure normalized by (1.16).

Given $l \times l$ matrix-valued functions $f, g$, we define the $l \times l$ matrix $\langle\langle f, g\rangle\rangle_{R}$ by

$$
\begin{equation*}
\langle\langle f, g\rangle\rangle_{R}=\int f(x)^{\dagger} M(x) g(x) d \mu_{\operatorname{tr}}(x), \tag{1.20}
\end{equation*}
$$

that is, its $(j, k)$ entry is

$$
\begin{equation*}
\sum_{n m} \int \overline{f_{n j}(x)} M_{n m}(x) g_{m k}(x) d \mu_{\operatorname{tr}}(x) . \tag{1.21}
\end{equation*}
$$

Since $f^{\dagger} M f \geq 0$, we see that

$$
\begin{equation*}
\langle\langle f, f\rangle\rangle_{R} \geq 0 . \tag{1.22}
\end{equation*}
$$

One might be tempted to think of $\langle\langle f, f\rangle\rangle_{R}^{1 / 2}$ as some kind of norm, but that is doubtful. Even if $\mu$ is supported at a single point, $x_{0}$, with $M=l^{-1} \mathbf{1}$, this "norm" is essentially the absolute value of $A=f\left(x_{0}\right)$, which is known not to obey the triangle inequality! (See [169, Sect. I.1] for an example.)

However, if one looks at

$$
\begin{equation*}
\|f\|_{R}=\left(\operatorname{Tr}\langle\langle f, f\rangle\rangle_{R}\right)^{1 / 2}, \tag{1.23}
\end{equation*}
$$

one does have a norm (or, at least, a semi-norm). Indeed,

$$
\begin{equation*}
\langle f, g\rangle_{R}=\operatorname{Tr}\left\langle\langle f, g\rangle_{R}\right. \tag{1.24}
\end{equation*}
$$

is a sesquilinear form which is positive semi-definite, so (1.23) is the semi-norm corresponding to an inner product and, of course, one has a Cauchy-Schwarz inequality

$$
\begin{equation*}
\mid \operatorname{Tr}\left\langle\langle f, g\rangle_{R}\right| \leq\|f\|_{R}\|g\|_{R} \tag{1.25}
\end{equation*}
$$

We have not specified which $f$ 's and $g$ 's can be used in (1.20). We have in mind mainly polynomials in $x$ in the real case and Laurent polynomials in $z$ in the $\partial \mathbb{D}$ case although, obviously, continuous functions are okay. Indeed, it suffices that $f$ (and $g$ ) be measurable and obey

$$
\begin{equation*}
\int \operatorname{Tr}\left(f^{\dagger}(x) f(x)\right) d \mu_{\operatorname{tr}}(x)<\infty \tag{1.26}
\end{equation*}
$$

for the integrals in (1.21) to converge. The set of equivalence classes under $f \sim g$ if $\|f-g\|_{R}=0$ defines a Hilbert space, $\mathcal{H}$, and $\langle f, g\rangle_{R}$ is the inner product on this space.

Instead of (1.20), we use the suggestive shorthand

$$
\begin{equation*}
\left\langle\langle f, g\rangle_{R}=\int f(x)^{\dagger} d \mu(x) g(x) .\right. \tag{1.27}
\end{equation*}
$$

The use of $R$ here comes from "right" for if $\alpha \in \mathcal{M}_{l}$,

$$
\begin{align*}
& \langle\langle f, g \alpha\rangle\rangle_{R}=\langle\langle f, g\rangle\rangle_{R} \alpha,  \tag{1.28}\\
& \langle\langle f \alpha, g\rangle\rangle_{R}=\alpha^{\dagger}\langle\langle f, g\rangle\rangle_{R}, \tag{1.29}
\end{align*}
$$

but, in general, $\langle\langle f, \alpha g\rangle\rangle_{R}$ is not related to $\langle\langle f, g\rangle\rangle_{R}$.
While $\left(\operatorname{Tr}\langle\langle f, f\rangle\rangle_{R}\right)^{1 / 2}$ is a natural analogue of the norm in the scalar case, it will sometimes be useful to instead consider

$$
\begin{equation*}
\left[\operatorname{det}\langle\langle f, f\rangle\rangle_{R}\right]^{1 / 2} \tag{1.30}
\end{equation*}
$$

Indeed, this is a stronger "norm" in that det $>0 \Rightarrow \operatorname{Tr}>0$ but not vice-versa.
When $d \mu$ is a "direct sum," that is, each $M(x)$ is diagonal, one can appreciate the difference. In that case, $d \mu=d \mu_{1} \oplus \cdots \oplus d \mu_{l}$ and the MOPRL are direct sums (i.e., diagonal matrices) of scalar OPRL

$$
\begin{equation*}
P_{n}^{R}(x, d \mu)=P_{n}\left(x, d \mu_{1}\right) \oplus \cdots \oplus P_{n}\left(x, d \mu_{l}\right) \tag{1.31}
\end{equation*}
$$

Then

$$
\begin{equation*}
\left\|P_{n}^{R}\right\|_{R}=\left(\sum_{j=1}^{l}\left\|P_{n}\left(\cdot, d \mu_{j}\right)\right\|_{L^{2}\left(d \mu_{j}\right)}^{2}\right)^{1 / 2} \tag{1.32}
\end{equation*}
$$

while

$$
\begin{equation*}
\left(\operatorname{det}\left\langle\left\langle P_{n}^{R}, P_{n}^{R}\right\rangle\right\rangle_{R}\right)^{1 / 2}=\prod_{j=1}^{l}\left\|P_{n}\left(\cdot, d \mu_{j}\right)\right\|_{L^{2}\left(d \mu_{j}\right)} . \tag{1.33}
\end{equation*}
$$

In particular, in terms of extending the theory of regular measures [180], $\left\|P_{n}^{R}\right\|_{R}^{1 / n}$ is only sensitive to $\max \left\|P_{n}\left(\cdot, d \mu_{j}\right)\right\|_{L^{2}\left(d \mu_{j}\right)}^{1 / 2}$ while $\left(\operatorname{det}\left\langle\left\langle P_{n}^{R}, P_{n}^{R}\right\rangle\right\rangle_{R}\right)^{1 / 2}$ is sensitive to them all. Thus, det will be needed for that theory (see Chapter 5).

There will also be a left inner product and, correspondingly, two sets of MOPRL and MOPUC. We discuss this further in Sections 2.1 and 3.1.

Occasionally, for $\mathbb{C}^{l}$ vector-valued functions $f$ and $g$, we will want to consider the scalar

$$
\begin{equation*}
\sum_{k, j} \int \overline{f_{k}(x)} M_{k j}(x) g_{j}(x) d \mu_{\operatorname{tr}}(x) \tag{1.34}
\end{equation*}
$$

which we will denote

$$
\begin{equation*}
\int d\langle f(x), \mu(x) g(x)\rangle_{\mathbb{C}^{l}} \tag{1.35}
\end{equation*}
$$

We next turn to module Fourier expansions. A set $\left\{\varphi_{j}\right\}_{j=1}^{N}$ in $\mathcal{H}$ ( $N$ may be infinite) is called orthonormal if and only if

$$
\begin{equation*}
\left\langle\left\langle\varphi_{j}, \varphi_{k}\right\rangle\right\rangle_{R}=\delta_{j k} \mathbf{1} \tag{1.36}
\end{equation*}
$$

This natural terminology is an abuse of notation since (1.36) implies orthogonality in $\langle\cdot, \cdot\rangle_{R}$ but not normalization, and is much stronger than orthogonality in $\langle\cdot, \cdot\rangle_{R}$.

Suppose for a moment that $N<\infty$. For any $a_{1}, \ldots, a_{N} \in \mathcal{M}_{l}$, we can form $\sum_{j=1}^{N} \varphi_{j} a_{j}$ and, by the right multiplication relations (1.28), (1.29), and (1.36), we have

$$
\begin{equation*}
\left\langle\left\langle\sum_{j=1}^{N} \varphi_{j} a_{j}, \sum_{j=1}^{N} \varphi_{j} b_{j}\right\rangle\right\rangle_{R}=\sum_{j=1}^{N} a_{j}^{\dagger} b_{j} . \tag{1.37}
\end{equation*}
$$

We will denote the set of all such $\sum_{j=1}^{N} \varphi_{j} a_{j}$ by $\mathcal{H}_{\left(\varphi_{j}\right)}$-it is a vector subspace of $\mathcal{H}$ of dimension (over $\mathbb{C}) N l^{2}$.

Define for $f \in \mathcal{H}$,

$$
\begin{equation*}
\pi_{\left(\varphi_{j}\right)}(f)=\sum_{j=1}^{N} \varphi_{j}\left\langle\left\langle\varphi_{j}, f\right\rangle\right\rangle_{R} \tag{1.38}
\end{equation*}
$$

It is easy to see it is the orthogonal projection in the scalar inner product $\langle\cdot, \cdot\rangle_{R}$ from $\mathcal{H}$ to $\mathcal{H}_{\left(\varphi_{j}\right)}$.
By the standard Hilbert space calculation, taking care to only multiply on the right, one finds the Pythagorean theorem,

$$
\begin{equation*}
\langle\langle f, f\rangle\rangle_{R}=\left\langle\left\langle f-\pi_{\left(\varphi_{j}\right)} f, f-\pi_{\left(\varphi_{j}\right)} f\right\rangle\right\rangle_{R}+\sum_{j=1}^{N}\left\langle\left\langle\varphi_{j}, f\right\rangle\right\rangle_{R}^{\dagger}\left\langle\left\langle\varphi_{j}, f\right\rangle\right\rangle_{R} . \tag{1.39}
\end{equation*}
$$

As usual, this proves for infinite $N$ that

$$
\begin{equation*}
\sum_{j=1}^{N}\left\langle\left\langle\varphi_{j}, f\right\rangle\right\rangle_{R}^{\dagger}\left\langle\left\langle\varphi_{j}, f\right\rangle\right\rangle_{R} \leq\langle\langle f, f\rangle\rangle_{R} \tag{1.40}
\end{equation*}
$$

and the convergence of

$$
\begin{equation*}
\sum_{j=1}^{N} \varphi_{j}\left\langle\left\langle\varphi_{j}, f\right\rangle\right\rangle_{R} \equiv \pi_{\left(\varphi_{j}\right)}(f) \tag{1.41}
\end{equation*}
$$

allowing the definition of $\pi_{\left(\varphi_{j}\right)}$ and of $\mathcal{H}_{\left(\varphi_{j}\right)} \equiv \operatorname{Ran} \pi_{\left(\varphi_{j}\right)}$ for $N=\infty$.
An orthonormal set is called complete if $\mathcal{H}_{\left(\varphi_{j}\right)}=\mathcal{H}$. In that case, equality holds in (1.40) and $\pi_{\left(\varphi_{j}\right)}(f)=f$.

For orthonormal bases, we have the Parseval relation from (1.39)

$$
\begin{equation*}
\langle\langle f, f\rangle\rangle_{R}=\sum_{j=1}^{\infty}\left\langle\left\langle\varphi_{j}, f\right\rangle\right\rangle_{R}^{\dagger}\left\langle\left\langle\varphi_{j}, f\right\rangle\right\rangle_{R} \tag{1.42}
\end{equation*}
$$

and

$$
\begin{equation*}
\|f\|_{R}^{2}=\sum_{j=1}^{\infty} \operatorname{Tr}\left(\left\langle\left\langle\varphi_{j}, f\right\rangle\right\rangle_{R}^{\dagger}\left\langle\left\langle\varphi_{j}, f\right\rangle\right\rangle_{R}\right) . \tag{1.43}
\end{equation*}
$$

### 1.3 Matrix Möbius Transformations

Without an understanding of matrix Möbius transformations, the form of the MOPUC Geronimus theorem we will prove in Section 3.10 will seem strange-looking. To set the stage, recall that scalar fractional linear transformations (FLT) are associated to matrices $T=\left(\begin{array}{cc}a & b \\ c & d\end{array}\right)$ with $\operatorname{det} T \neq 0$ via

$$
\begin{equation*}
f_{T}(z)=\frac{a z+b}{c z+d} . \tag{1.44}
\end{equation*}
$$

Without loss, one can restrict to

$$
\begin{equation*}
\operatorname{det}(T)=1 \tag{1.45}
\end{equation*}
$$

Indeed, $T \mapsto f_{T}$ is a 2 to 1 map of $\mathbb{S L}(2, \mathbb{C})$ to maps of $\mathbb{C} \cup\{\infty\}$ to itself. One advantage of the matrix formalism is that the map is a matrix homomorphism, that is,

$$
\begin{equation*}
f_{T \circ S}=f_{T} \circ f_{S}, \tag{1.46}
\end{equation*}
$$

which shows that the group of FLTs is $\mathbb{S L}(2, \mathbb{C}) /\{\mathbf{1},-\mathbf{1}\}$.
While (1.46) can be checked by direct calculation, a more instructive way is to look at the complex projective line. $u, v \in \mathbb{C}^{2} \backslash\{0\}$ are called equivalent if there is $\lambda \in \mathbb{C} \backslash\{0\}$ so that $u=\lambda v$. Let [.] denote equivalence classes. Except for $\left.\left[\begin{array}{l}1 \\ 0\end{array}\right)\right]$, every equivalence class contains exactly one point of the form $\binom{z}{1}$ with $z \in \mathbb{C}$. If $\left.\left[\begin{array}{l}1 \\ 0\end{array}\right)\right]$ is associated with $\infty$, the set of equivalence classes is naturally associated with $\mathbb{C} \cup\{\infty\}$. $f_{T}$ then obeys

$$
\begin{equation*}
\left[T\binom{z}{1}\right]=\left[\binom{f_{T}(z)}{1}\right] \tag{1.47}
\end{equation*}
$$

from which (1.46) is immediate.
By Möbius transformations we will mean those FLTs that map $\mathbb{D}$ onto itself. Let

$$
J=\left(\begin{array}{rr}
1 & 0  \tag{1.48}\\
0 & -1
\end{array}\right) .
$$

Then $\left.[u]=\left[\begin{array}{l}z \\ 1 \\ 1\end{array}\right)\right]$ with $|z|=1$ (resp. $|z|<1$ ) if and only if $\langle u, J u\rangle=0$ (resp. $\langle u, J u\rangle<0$ ). From this, it is not hard to show that if $\operatorname{det}(T)=1$, then $f_{T}$ maps $\mathbb{D}$ invertibly onto $\mathbb{D}$ if and only if

$$
\begin{equation*}
T^{\dagger} J T=J \tag{1.49}
\end{equation*}
$$

If $T$ has the form $\left(\begin{array}{ll}a & b \\ c & d\end{array}\right)$, this is equivalent to

$$
\begin{equation*}
|a|^{2}-|c|^{2}=1, \quad|b|^{2}-|d|^{2}=-1, \quad \bar{a} b-\bar{c} d=0 \tag{1.50}
\end{equation*}
$$

The set of $T$ 's obeying $\operatorname{det}(T)=1$ and (1.49) is called $\mathbb{S U}(1,1)$. It is studied extensively in [168, Sect. 10.4].

The self-adjoint elements of $\mathbb{S U}(1,1)$ are parametrized by $\alpha \in \mathbb{D}$ via $\rho=\left(1-|\alpha|^{2}\right)^{1 / 2}$,

$$
T_{\alpha}=\frac{1}{\rho}\left(\begin{array}{ll}
1 & \alpha  \tag{1.51}\\
\bar{\alpha} & 1
\end{array}\right)
$$

associated to

$$
\begin{equation*}
f_{T_{\alpha}}(z)=\frac{z+\alpha}{1+\bar{\alpha} z} . \tag{1.52}
\end{equation*}
$$

Notice that

$$
\begin{equation*}
T_{\alpha}^{-1}=T_{-\alpha} \tag{1.53}
\end{equation*}
$$

and that

$$
\forall z \in \mathbb{D}, \exists!\alpha \text { such that } T_{\alpha}(0)=z,
$$

namely, $\alpha=z$.
It is a basic theorem that every holomorphic bijection of $\mathbb{D}$ to $\mathbb{D}$ is an $f_{T}$ for some $T$ in $\mathbb{S U}(1,1)$ (unique up to $\pm \mathbf{1}$ ).

With this in place, we can turn to the matrix case. Let $\mathcal{M}_{l}$ be the space of $l \times l$ complex matrices with the Euclidean norm induced by the vector norm $\langle\cdot, \cdot\rangle_{\mathbb{C}^{l}}^{1 / 2}$. Let

$$
\begin{equation*}
\mathbb{D}_{l}=\left\{A \in \mathcal{M}_{l}:\|A\|<1\right\} . \tag{1.54}
\end{equation*}
$$

We are interested in holomorphic bijections of $\mathbb{D}_{l}$ to itself, especially via a suitable notion of FLT. There is a huge (and diffuse) literature on the subject, starting with its use in analytic number
theory. It has also been studied in connection with electrical engineering filters and indefinite matrix Hilbert spaces. Among the huge literature, we mention [1, 3, 78, 99, 114, 166]. Especially relevant to MOPUC is the book of Bakonyi-Constantinescu [6].

Consider $\mathcal{M}_{l} \oplus \mathcal{M}_{l}=\mathcal{M}_{l}[2]$ as a right module over $\mathcal{M}_{l}$. The $\mathcal{M}_{l}$-projective line is defined by saying $\left[\begin{array}{c}X \\ Y\end{array}\right] \sim\left[\begin{array}{c}X^{\prime} \\ Y^{\prime}\end{array}\right]$, both in $\mathcal{M}_{l}[2] \backslash\{\mathbf{0}\}$, if and only if there exists $\Lambda \in \mathcal{M}_{l}, \Lambda$ invertible so that

$$
\begin{equation*}
X=X^{\prime} \Lambda, \quad Y=Y^{\prime} \Lambda \tag{1.55}
\end{equation*}
$$

Let $T$ be a map of $\mathcal{M}_{l}[2]$ of the form

$$
T=\left(\begin{array}{ll}
A & B  \tag{1.56}\\
C & D
\end{array}\right)
$$

acting on $\mathcal{M}_{l}[2]$ by

$$
T\left[\begin{array}{l}
X  \tag{1.57}\\
Y
\end{array}\right]=\left[\begin{array}{l}
A X+B Y \\
C X+D Y
\end{array}\right]
$$

Because this acts on the left and $\Lambda$ equivalence on the right, $T$ maps equivalence classes to themselves. In particular, if $C X+D$ is invertible, $T$ maps the equivalence class of $\left[\begin{array}{l}X \\ 1\end{array}\right]$ to the equivalence class of $\left[\begin{array}{c}f_{T}[X] \\ 1\end{array}\right]$, where

$$
\begin{equation*}
f_{T}[X]=(A X+B)(C X+D)^{-1} \tag{1.58}
\end{equation*}
$$

So long as $C X+D$ remains invertible, (1.46) remains true. Let $J$ be the $2 l \times 2 l$ matrix in $l \times l$ block form

$$
J=\left(\begin{array}{rr}
\mathbf{1} & \mathbf{0}  \tag{1.59}\\
\mathbf{0} & -\mathbf{1}
\end{array}\right) .
$$

Note that $\left(\right.$ with $\left.\left[\begin{array}{l}X \\ 1\end{array}\right]^{\dagger}=\left[X^{\dagger} \mathbf{1}\right]\right)$

$$
\left[\begin{array}{c}
X  \tag{1.60}\\
1
\end{array}\right]^{\dagger} J\left[\begin{array}{c}
X \\
1
\end{array}\right] \leq \mathbf{0} \Leftrightarrow X^{\dagger} X \leq \mathbf{1} \Leftrightarrow\|X\| \leq \mathbf{1}
$$

Therefore, if we define $\mathbb{S U}(l, l)$ to be those $T$ 's with $\operatorname{det} T=1$ and

$$
\begin{equation*}
T^{\dagger} J T=J \tag{1.61}
\end{equation*}
$$

then

$$
\begin{equation*}
T \in \mathbb{S U}(l, l) \Rightarrow f_{T}\left[\mathbb{D}_{l}\right]=\mathbb{D}_{l} \text { as a bijection. } \tag{1.62}
\end{equation*}
$$

If $T$ has the form (1.56), then (1.61) is equivalent to

$$
\begin{align*}
& A^{\dagger} A-C^{\dagger} C=D^{\dagger} D-B^{\dagger} B=\mathbf{1}  \tag{1.63}\\
& A^{\dagger} B=C^{\dagger} D \tag{1.64}
\end{align*}
$$

(the fourth relation $B^{\dagger} A=D^{\dagger} C$ is equivalent to (1.64)).
This depends on
Proposition 1.1. If $T=\left(\begin{array}{cc}A & B \\ C & D\end{array}\right)$ obeys (1.61) and $\|X\|<1$, then $C X+D$ is invertible.

Proof. (1.61) implies that

$$
\begin{align*}
T^{-1} & =J T^{\dagger} J  \tag{1.65}\\
& =\left(\begin{array}{rr}
A^{\dagger} & -C^{\dagger} \\
-B^{\dagger} & D^{\dagger}
\end{array}\right) . \tag{1.66}
\end{align*}
$$

Clearly, (1.61) also implies $T^{-1} \in \mathbb{S U}(l, l)$. Thus, by (1.63) for $T^{-1}$,

$$
\begin{equation*}
D D^{\dagger}-C C^{\dagger}=\mathbf{1} \tag{1.67}
\end{equation*}
$$

This implies first that $D D^{\dagger} \geq \mathbf{1}$, so $D$ is invertible, and second that

$$
\begin{equation*}
\left\|D^{-1} C\right\| \leq 1 . \tag{1.68}
\end{equation*}
$$

Thus, $\|X\|<1$ implies $\left\|D^{-1} C X\right\|<1$ so $\mathbf{1}+D^{-1} C X$ is invertible, and thus so is $D\left(\mathbf{1}+D^{-1} C X\right)$.

It is a basic result of Cartan [18] (see Helgason [114] and the discussion therein) that
Theorem 1.2. A holomorphic bijection, $g$, of $\mathbb{D}_{l}$ to itself is either of the form

$$
\begin{equation*}
g(X)=f_{T}(X) \tag{1.69}
\end{equation*}
$$

for some $T \in \mathbb{S U}(l, l)$ or

$$
\begin{equation*}
g(X)=f_{T}\left(X^{t}\right) \tag{1.70}
\end{equation*}
$$

Given $\alpha \in \mathcal{M}_{l}$ with $\|\alpha\|<1$, define

$$
\begin{equation*}
\rho^{L}=\left(\mathbf{1}-\alpha^{\dagger} \alpha\right)^{1 / 2}, \quad \rho^{R}=\left(\mathbf{1}-\alpha \alpha^{\dagger}\right)^{1 / 2} \tag{1.71}
\end{equation*}
$$

Lemma 1.3. We have

$$
\begin{align*}
\alpha \rho^{L} & =\rho^{R} \alpha, & \alpha^{\dagger} \rho^{R} & =\rho^{L} \alpha^{\dagger},  \tag{1.72}\\
\alpha\left(\rho^{L}\right)^{-1} & =\left(\rho^{R}\right)^{-1} \alpha, & \alpha^{\dagger}\left(\rho^{R}\right)^{-1} & =\left(\rho^{L}\right)^{-1} \alpha^{\dagger} . \tag{1.73}
\end{align*}
$$

Proof. Let $f$ be analytic in $\mathbb{D}$ with $f(z)=\sum_{n=0}^{\infty} c_{n} z^{n}$ its Taylor series at $z=0$. Since $\left\|\alpha^{\dagger} \alpha\right\|<1$, we have

$$
\begin{equation*}
f\left(\alpha^{\dagger} \alpha\right)=\sum_{n=0}^{\infty} c_{n}\left(\alpha^{\dagger} \alpha\right)^{n} \tag{1.74}
\end{equation*}
$$

norm convergent, so $\alpha\left(\alpha^{\dagger} \alpha\right)^{n}=\left(\alpha \alpha^{\dagger}\right)^{n} \alpha$ implies

$$
\begin{equation*}
\alpha f\left(\alpha^{\dagger} \alpha\right)=f\left(\alpha \alpha^{\dagger}\right) \alpha \tag{1.75}
\end{equation*}
$$

which implies the first halves of (1.72) and (1.73). The other halves follow by taking adjoints.
Theorem 1.4. There is a one-one correspondence between $\alpha$ 's in $\mathcal{M}_{l}$ obeying $\|\alpha\|<1$ and positive self-adjoint elements of $\mathbb{S U}(l, l)$ via

$$
T_{\alpha}=\left(\begin{array}{cc}
\left(\rho^{R}\right)^{-1} & \left(\rho^{R}\right)^{-1} \alpha  \tag{1.76}\\
\left(\rho^{L}\right)^{-1} \alpha^{\dagger} & \left(\rho^{L}\right)^{-1}
\end{array}\right)
$$

Proof. A straightforward calculation using Lemma 1.3 proves that $T_{\alpha}$ is self-adjoint and $T_{\alpha}^{\dagger} J T_{\alpha}=J$. Conversely, if $T$ is self-adjoint, $T=\left(\begin{array}{c}A \\ C\end{array} \underset{D}{B}\right)$ and in $\mathbb{S U}(l, l)$, then $T^{\dagger}=T \Rightarrow A^{\dagger}=A, B^{\dagger}=C$, so (1.63) becomes

$$
\begin{equation*}
A A^{\dagger}-B B^{\dagger}=\mathbf{1} \tag{1.77}
\end{equation*}
$$

so if

$$
\begin{equation*}
\alpha=A^{-1} B \tag{1.78}
\end{equation*}
$$

then (1.77) becomes

$$
\begin{equation*}
A^{-1}\left(A^{-1}\right)^{\dagger}+\alpha \alpha^{\dagger}=\mathbf{1} \tag{1.79}
\end{equation*}
$$

Since $T \geq 0, A \geq 0$ so (1.79) implies $A=\left(\rho^{R}\right)^{-1}$, and then (1.78) implies $B=\left(\rho^{R}\right)^{-1} \alpha$.
By Lemma 1.3,

$$
\begin{equation*}
C=B^{\dagger}=\alpha^{\dagger}\left(\rho^{R}\right)^{-1}=\left(\rho^{L}\right)^{-1} \alpha^{\dagger} \tag{1.80}
\end{equation*}
$$

and then (by $D=D^{\dagger}, C^{\dagger}=B$, and (1.63)) $D D^{\dagger}-C C^{\dagger}=1$ plus $D>0$ implies $D=\left(\rho^{L}\right)^{-1}$.
Corollary 1.5. For each $\alpha \in \mathbb{D}_{l}$, the map

$$
\begin{equation*}
f_{T_{\alpha}}(X)=\left(\rho^{R}\right)^{-1}(X+\alpha)\left(\mathbf{1}+\alpha^{\dagger} X\right)^{-1}\left(\rho^{L}\right) \tag{1.81}
\end{equation*}
$$

takes $\mathbb{D}_{l}$ to $\mathbb{D}_{l}$. Its inverse is given by

$$
\begin{equation*}
f_{T_{\alpha}}^{-1}(X)=f_{T_{-\alpha}}(X)=\left(\rho^{R}\right)^{-1}(X-\alpha)\left(1-\alpha^{\dagger} X\right)^{-1}\left(\rho^{L}\right) . \tag{1.82}
\end{equation*}
$$

There is an alternate form for the right side of (1.81).
Proposition 1.6. The following identity holds true for any $X,\|X\| \leq 1$ :

$$
\begin{equation*}
\rho^{R}\left(1+X \alpha^{\dagger}\right)^{-1}(X+\alpha)\left(\rho^{L}\right)^{-1}=\left(\rho^{R}\right)^{-1}(X+\alpha)\left(1+\alpha^{\dagger} X\right)^{-1} \rho^{L} . \tag{1.83}
\end{equation*}
$$

Proof. By the definition of $\rho^{L}$ and $\rho^{R}$, we have

$$
X\left(\rho^{L}\right)^{-2}\left(1-\alpha^{\dagger} \alpha\right)=\left(\rho^{R}\right)^{-2}\left(1-\alpha \alpha^{\dagger}\right) X
$$

Expanding, using (1.73) and rearranging, we get

$$
X\left(\rho^{L}\right)^{-2}+\alpha\left(\rho^{L}\right)^{-2} \alpha^{\dagger} X=\left(\rho^{R}\right)^{-2} X+X \alpha^{\dagger}\left(\rho^{R}\right)^{-2} \alpha
$$

Adding $\alpha\left(\rho^{L}\right)^{-2}+X\left(\rho^{L}\right)^{-2} \alpha^{\dagger} X$ to both sides and using (1.73) again, we obtain

$$
\begin{aligned}
X\left(\rho^{L}\right)^{-2}+\alpha\left(\rho^{L}\right)^{-2}+X\left(\rho^{L}\right)^{-2} \alpha^{\dagger} X+ & \alpha\left(\rho^{L}\right)^{-2} \alpha^{\dagger} X \\
& =\left(\rho^{R}\right)^{-2} X+\left(\rho^{R}\right)^{-2} \alpha+X \alpha^{\dagger}\left(\rho^{R}\right)^{-2} X+X \alpha^{\dagger}\left(\rho^{R}\right)^{-2} \alpha
\end{aligned}
$$

which is the same as

$$
(X+\alpha)\left(\rho^{L}\right)^{-2}\left(1+\alpha^{\dagger} X\right)=\left(1+X \alpha^{\dagger}\right)\left(\rho^{R}\right)^{-2}(X+\alpha) .
$$

Multiplying by $\left(1+X \alpha^{\dagger}\right)^{-1}$ and $\left(1+\alpha^{\dagger} X\right)^{-1}$, we get

$$
\left(1+X \alpha^{\dagger}\right)^{-1}(X+\alpha)\left(\rho^{L}\right)^{-2}=\left(\rho^{R}\right)^{-2}(X+\alpha)\left(1+\alpha^{\dagger} X\right)^{-1}
$$

and the statement follows.

### 1.4 Applications and Examples

There are a number of simple examples which show that beyond their intrinsic mathematical interest, MOPRL and MOPUC have wide application.

## (a) Jacobi matrices on a strip

Let $\Lambda \subset \mathbb{Z}^{\nu}$ be a subset (perhaps infinite) of the $\nu$-dimensional lattice $\mathbb{Z}^{\nu}$ and let $\ell^{2}(\Lambda)$ be square summable sequences indexed by $\Lambda$. Suppose a real symmetric matrix $\alpha_{i j}$ is given for all $i, j \in \Lambda$ with $\alpha_{i j}=0$ unless $|i-j|=1$ (nearest neighbors). Let $\beta_{i}$ be a real sequence indexed by $i \in \Lambda$. Suppose

$$
\begin{equation*}
\sup _{i, j}\left|\alpha_{i j}\right|+\sup _{i}\left|\beta_{i}\right|<\infty . \tag{1.84}
\end{equation*}
$$

Define a bounded operator, $J$, on $\ell^{2}(\Lambda)$ by

$$
\begin{equation*}
(J u)_{i}=\sum_{j} \alpha_{i j} u_{j}+\beta_{i} u_{i} . \tag{1.85}
\end{equation*}
$$

The sum is finite with at most $2 \nu$ elements.
The special case $\Lambda=\{1,2, \ldots\}$ with $b_{i}=\beta_{i}, a_{i}=\alpha_{i, i+1}>0$ corresponds precisely to classical semi-infinite tridiagonal Jacobi matrices.

Now consider the situation where $\Lambda^{\prime} \subset \mathbb{Z}^{\nu-1}$ is a finite set with $l$ elements and

$$
\begin{equation*}
\Lambda=\left\{j \in \mathbb{Z}^{\nu}: j_{1} \in\{1,2, \ldots\} ;\left(j_{2}, \ldots j_{\nu}\right) \in \Lambda^{\prime}\right\} \tag{1.86}
\end{equation*}
$$

a "strip" with cross-section $\Lambda^{\prime}$. $J$ then has a block $l \times l$ matrix Jacobi form where $\left(\gamma, \delta \in \Lambda^{\prime}\right)$

$$
\begin{array}{rlrl}
\left(B_{i}\right)_{\gamma \delta} & =b_{(i, \gamma)}, & & (\gamma=\delta), \\
& =a_{(i, \gamma)(i, \delta)}, & & (\gamma \neq \delta), \\
\left(A_{i}\right)_{\gamma \delta} & =a_{(i, \gamma)(i+1, \delta)} . & \tag{1.89}
\end{array}
$$

The nearest neighbor condition says $\left(A_{i}\right)_{\gamma \delta}=0$ if $\gamma \neq \delta$. If

$$
\begin{equation*}
a_{(i, \gamma)(i+1, \gamma)}>0 \tag{1.90}
\end{equation*}
$$

for all $i, \gamma$, then $A_{i}$ is invertible and we have a block Jacobi matrix of the kind described in Section 2.2 below.

By allowing general $A_{i}, B_{i}$, we obtain an obvious generalization of this model-an interpretation of general MOPRL.

Schrödinger operators on strips have been studied in part as approximations to $\mathbb{Z}^{\nu}$; see $[31,95$, $130,134,151,164]$. From this point of view, it is also natural to allow periodic boundary conditions in the vertical directions. Furthermore, there is closely related work on Schrödinger (and other) operators with matrix-valued potentials; see, for example, [8, 24, 25, 26, 27, 28, 30, 96, 97, 165].

## (b) Two-sided Jacobi matrices

This example goes back at least to Nikishin [153]. Consider the case $\nu=2, \Lambda^{\prime}=\{0,1\} \subset \mathbb{Z}$, and $\Lambda$ as above. Suppose (1.90) holds, and in addition,

$$
\begin{align*}
a_{(1,0)(1,1)} & >0,  \tag{1.91}\\
a_{(i, 0)(i, 1)} & =0, \quad i=2,3, \ldots . \tag{1.92}
\end{align*}
$$

Then there are no links between the rungs of the "ladder," $\{1,2, \ldots\} \times\{0,1\}$ except at the end and the ladder can be unfolded to $\mathbb{Z}$ ! Thus, a two-sided Jacobi matrix can be viewed as a special kind of one-sided $2 \times 2$ matrix Jacobi operator.

It is known that for two-sided Jacobi matrices, the spectral theory is determined by the $2 \times 2$ matrix

$$
d \mu=\left(\begin{array}{ll}
d \mu_{00} & d \mu_{01}  \tag{1.93}\\
d \mu_{10} & d \mu_{11}
\end{array}\right)
$$

where $d \mu_{k l}$ is the measure with

$$
\begin{equation*}
\left\langle\delta_{k},(J-\lambda)^{-1} \delta_{l}\right\rangle=\int \frac{d \mu_{k l}(x)}{x-\lambda}, \tag{1.94}
\end{equation*}
$$

but also that it is very difficult to specify exactly which $d \mu$ correspond to two-sided Jacobi matrices.
This difficulty is illuminated by the theory of MOPRL. By Favard's theorem (see Theorem 2.11), every such $d \mu$ (given by (1.93) and positive definite and non-trivial in a sense we will describe in Lemma 2.1) yields a unique block Jacobi matrix with $A_{j}>0$ (positive definite). This $d \mu$ comes from a two-sided Jacobi matrix if and only if
(a) $B_{j}$ is diagonal for $j=2,3, \ldots$.
(b) $A_{j}$ is diagonal for $j=1,2, \ldots$.
(c) $B_{j}$ has strictly positive off-diagonal elements.

These are very complicated indirect conditions on $d \mu$ !

## (c) Banded matrices

Classical Jacobi matrices are semi-infinite symmetric tridiagonal matrices, that is,

$$
\begin{equation*}
J_{k m}=0 \quad \text { if }|k-m|>1 \tag{1.95}
\end{equation*}
$$

with

$$
\begin{equation*}
J_{k m}>0 \quad \text { if }|k-m|=1 . \tag{1.96}
\end{equation*}
$$

A natural generalization are $(2 l+1)$-diagonal symmetric matrices, that is,

$$
\begin{array}{cc}
J_{k m}=0 & \text { if }|k-m|>l, \\
J_{k m}>0 & \text { if }|k-m|=l . \tag{1.98}
\end{array}
$$

Such a matrix can be partitioned into $l \times l$ blocks, which is tridiagonal in block. The conditions (1.97) and (1.98) are equivalent to $A_{k} \in \mathcal{L}$, the set of lower triangular matrices; and conversely, $A_{k} \in \mathcal{L}$, with $A_{k}, B_{k}$ real (and $B_{k}$ symmetric) correspond precisely to such banded matrices. This is why we introduce type 3 MOPRL.

Banded matrices correspond to certain higher-order difference equations. Unlike the secondorder equation (which leads to tridiagonal matrices) where every equation with positive coefficients is equivalent via a variable rescaling to a symmetric matrix, only certain higher-order difference equations correspond to symmetric block Jacobi matrices.

## (d) Magic formula

In [33], Damanik, Killip, and Simon studied perturbations of Jacobi and CMV matrices with periodic Jacobi parameters (or Verblunsky coefficients). They proved that if $\Delta$ is the discriminant of a two-sided periodic $J_{0}$, then a bounded two-sided $J$ has $\Delta(J)=S^{p}+S^{-p}\left((S u)_{n} \equiv u_{n+1}\right)$ if and only if $J$ lies in the isospectral torus of $J_{0}$. They call this the magic formula.

This allows the study of perturbations of the isospectral torus by studying $\Delta(J)$ which is a polynomial in $J$ of degree $p$, and so a $2 p+1$ banded matrix. Thus, the study of perturbations of periodic problems is connected to perturbations of $S^{p}+S^{-p}$ as block Jacobi matrices. Indeed, it was this connection that stimulated our interest in MOPRL, and [33] uses some of our results here.

## (e) Vector-valued prediction theory

As noted in Section 1.1, both prediction theory and filtering theory use OPUC and have natural MOPUC settings that motivated much of the MOPUC literature.

## 2 Matrix Orthogonal Polynomials on the Real Line

### 2.1 Preliminaries

OPRL are the most basic and developed of orthogonal polynomials, and so this chapter on the matrix analogue is the most important of this survey. We present the basic formulas, assuming enough familiarity with the scalar case (see [23, 82, 167, 176, 184, 185]) that we do not need to explain why the objects we define are important.

### 2.1.1 Polynomials, Inner Products, Norms

Let $d \mu$ be an $l \times l$ matrix-valued Hermitian positive semi-definite finite measure on $\mathbb{R}$ normalized by $\mu(\mathbb{R})=\mathbf{1} \in \mathcal{M}_{l}$. We assume for simplicity that $\mu$ has a compact support. However, many of the results below do not need the latter restriction and in fact can be found in the literature for matrix-valued measures with unbounded support.

Define (as in (1.20))

$$
\begin{array}{lll}
\left\langle\langle f, g\rangle_{R}=\int f(x)^{\dagger} d \mu(x) g(x),\right. & \|f\|_{R}=\left(\operatorname{Tr}\langle\langle f, f\rangle\rangle_{R}\right)^{1 / 2}, & f, g \in \mathcal{P} \\
\langle\langle f, g\rangle\rangle_{L}=\int g(x) d \mu(x) f(x)^{\dagger}, & \|f\|_{L}=\left(\operatorname{Tr}\left\langle\langle f, f\rangle_{L}\right)^{1 / 2},\right. & f, g \in \mathcal{P}
\end{array}
$$

Clearly, we have

$$
\begin{align*}
\langle\langle f, g\rangle\rangle_{R}^{\dagger} & =\langle\langle g, f\rangle\rangle_{R}, & \langle\langle f, g\rangle\rangle_{L}^{\dagger} & =\langle\langle g, f\rangle\rangle_{L},  \tag{2.1}\\
\langle\langle f, g\rangle\rangle_{L} & =\left\langle\left\langle g^{\dagger}, f^{\dagger}\right\rangle\right\rangle_{R}, & \|f\|_{L} & =\left\|f^{\dagger}\right\|_{R} . \tag{2.2}
\end{align*}
$$

As noted in Section 1.2, we have the left and right analogues of the Cauchy inequality

$$
\left|\operatorname{Tr}\langle\langle f, g\rangle\rangle_{R}\right| \leq\|f\|_{R}\|g\|_{R}, \quad\left|\operatorname{Tr}\langle\langle f, g\rangle\rangle_{L}\right| \leq\|f\|_{L}\|g\|_{L} .
$$

Thus, $\|\cdot\|_{R}$ and $\|\cdot\|_{L}$ are semi-norms in $\mathcal{P}$. Indeed, as noted in Section 1.2, they are associated to an inner product. The sets $\left\{f:\|f\|_{R}=0\right\}$ and $\left\{f:\|f\|_{L}=0\right\}$ are linear subspaces. Let $\mathcal{P}_{R}$
be the completion of $\mathcal{P} /\left\{f:\|f\|_{R}=0\right\}$ (viewed as a right module over $\mathcal{M}_{l}$ ) with respect to the norm $\|\cdot\|_{R}$. Similarly, let $\mathcal{P}_{L}$ be the completion of $\mathcal{P} /\left\{f:\|f\|_{L}=0\right\}$ (viewed as a left module) with respect to the norm $\|\cdot\|_{L}$.

The set $\mathcal{V}$ defined in Section 1.2 is a linear space. Let us introduce a semi-norm in $\mathcal{V}$ by

$$
\begin{equation*}
|f|=\left\{\int d\langle f(x), \mu(x) f(x)\rangle_{\mathbb{C}^{l}}\right\}^{1 / 2} \tag{2.3}
\end{equation*}
$$

Let $\mathcal{V}_{0} \subset \mathcal{V}$ be the linear subspace of all polynomials such that $|f|=0$ and let $\mathcal{V}_{\infty}$ be the completion of the quotient space $\mathcal{V} / \mathcal{V}_{0}$ with respect to the norm $|\cdot|$.

Lemma 2.1. The following are equivalent:
(1) $\|f\|_{R}>0$ for every non-zero $f \in \mathcal{P}$.
(2) For all $n$, the dimension in $\mathcal{P}_{R}$ of the set of all polynomials of degree at most $n$ is $(n+1) l^{2}$.
(3) $\|f\|_{L}>0$ for every non-zero $f \in \mathcal{P}$.
(4) For all $n$, the dimension in $\mathcal{P}_{L}$ of the set of all polynomials of degree at most $n$ is $(n+1) l^{2}$.
(5) For every non-zero $v \in \mathcal{V}$, we have that $|v| \neq 0$.
(6) For all $n$, the dimension in $\mathcal{V}_{\infty}$ of all vector-valued polynomials of degree at most $n$ is $(n+1) l$. The measure $d \mu$ is called non-trivial if these equivalent conditions hold.

Remark. If $l=1$, these are equivalent to the usual non- triviality condition, that is, $\operatorname{supp}(\mu)$ is infinite. For $l>1$, we cannot define triviality in this simple way, as can be seen by looking at the direct sum of a trivial and non-trivial measure. In that case, the measure is not non-trivial in the above sense but its support is infinite.

Proof. The equivalences $(1) \Leftrightarrow(2),(3) \Leftrightarrow(4)$, and (5) $\Leftrightarrow(6)$ are immediate. The equivalence (1) $\Leftrightarrow(3)$ follows from (2.2). Let us prove the equivalence (1) $\Leftrightarrow$ (5). Assume that (1) holds and let $v \in \mathcal{V}$ be non-zero. Let $f \in \mathcal{M}_{l}$ denote the matrix that has $v$ as its leftmost column and that has zero columns otherwise. Then, $0 \neq\|f\|_{R}^{2}=\operatorname{Tr}\langle\langle f, f\rangle\rangle_{R}=|v|^{2}$ and hence (5) holds. Now assume that (1) fails and let $f \in \mathcal{P}$ be non-zero with $\|f\|_{R}=0$. Then, at least one of the column vectors of $f$ is non-zero. Suppose for simplicity that this is the first column and denote this column vector by $v$. Let $t \in \mathcal{M}_{l}$ be the matrix $t_{i j}=\delta_{i 1} \delta_{j 1}$; then we have

$$
\|f\|_{R}=0 \Rightarrow\langle\langle f, f\rangle\rangle_{R}=0 \Rightarrow 0=\operatorname{Tr}\left(t^{*}\langle\langle f, f\rangle\rangle_{R} t\right)=|v|^{2}
$$

and hence (5) fails.
Throughout the rest of this chapter, we assume the measure $d \mu$ to be non-trivial.

### 2.1.2 Monic Orthogonal Polynomials

Lemma 2.2. Let $d \mu$ be a non-trivial measure.
(i) There exists a unique monic polynomial $P_{n}^{R}$ of degree $n$, which minimizes the norm $\left\|P_{n}^{R}\right\|_{R}$.
(ii) The polynomial $P_{n}^{R}$ can be equivalently defined as the monic polynomial of degree $n$ which satisfies

$$
\begin{equation*}
\left\langle\left\langle P_{n}^{R}, f\right\rangle\right\rangle_{R}=\mathbf{0} \quad \text { for any } f \in \mathcal{P}, \quad \operatorname{deg} f<n . \tag{2.4}
\end{equation*}
$$

(iii) There exists a unique monic polynomial $P_{n}^{L}$ of degree $n$, which minimizes the norm $\left\|P_{n}^{L}\right\|_{L}$.
(iv) The polynomial $P_{n}^{L}$ can be equivalently defined as the monic polynomial of degree $n$ which satisfies

$$
\begin{equation*}
\left\langle\left\langle P_{n}^{L}, f\right\rangle\right\rangle_{L}=\mathbf{0} \quad \text { for any } f \in \mathcal{P}, \quad \operatorname{deg} f<n . \tag{2.5}
\end{equation*}
$$

(v) One has $P_{n}^{L}(x)=P_{n}^{R}(x)^{\dagger}$ for all $x \in \mathbb{R}$ and

$$
\begin{equation*}
\left\langle\left\langle P_{n}^{R}, P_{n}^{R}\right\rangle\right\rangle_{R}=\left\langle\left\langle P_{n}^{L}, P_{n}^{L}\right\rangle\right\rangle_{L} . \tag{2.6}
\end{equation*}
$$

Proof. As noted, $\mathcal{P}$ has an inner product $\langle\cdot, \cdot\rangle_{R}$, so there is an orthogonal projection $\pi_{n}^{(R)}$ onto $\mathcal{P}_{n}$ discussed in Section 1.2. Then

$$
\begin{equation*}
P_{n}^{R}(x)=x^{n}-\pi_{n-1}^{(R)}\left(x^{n}\right) . \tag{2.7}
\end{equation*}
$$

As usual, in inner product spaces, this uniquely minimizes $x^{n}-Q$ over all $Q \in \mathcal{P}_{n-1}$. It clearly obeys

$$
\begin{equation*}
\operatorname{Tr}\left(\left\langle\left\langle P_{n}^{R}, f\right\rangle\right\rangle_{R}\right)=0 \tag{2.8}
\end{equation*}
$$

for all $f \in \mathcal{P}_{n-1}$. But then for any matrix $\alpha$,

$$
\operatorname{Tr}\left(\left\langle\left\langle P_{n}^{R}, f\right\rangle\right\rangle_{R} \alpha\right)=\operatorname{Tr}\left(\left\langle\left\langle P_{n}^{R}, f \alpha\right\rangle\right\rangle_{R}\right)=0
$$

so (2.4) holds.
This proves (i) and (ii). (iii) and (iv) are similar. To prove (v), note that $P_{n}^{L}(x)=P_{n}^{R}(x)^{\dagger}$ follows from the criteria (2.4), (2.5). The identity (2.6) follows from (2.2).

Lemma 2.3. Let $\mu$ be non-trivial. For any monic polynomial $P$, we have $\operatorname{det}\langle\langle P, P\rangle\rangle_{R} \neq 0$ and $\operatorname{det}\langle\langle P, P\rangle\rangle_{L} \neq 0$.

Proof. Let $P$ be a monic polynomial of degree $n$ such that $\langle\langle P, P\rangle\rangle_{R}$ has a non-trivial kernel. Then one can find $\alpha \in \mathcal{M}_{l}, \alpha \neq \mathbf{0}$, such that $\alpha^{\dagger}\langle\langle P, P\rangle\rangle_{R} \alpha=\mathbf{0}$. It follows that $\|P \alpha\|_{R}=0$. But since $P$ is monic, the leading coefficient of $P \alpha$ is $\alpha$, so $P \alpha \neq \mathbf{0}$, which contradicts the non-triviality assumption. A similar argument works for $\langle\langle P, P\rangle\rangle_{L}$.

By the orthogonality of $Q_{n}-P_{n}^{R}$ to $P_{n}^{R}$ for any monic polynomial $Q_{n}$ of degree $n$, we have

$$
\begin{equation*}
\left\langle\left\langle Q_{n}, Q_{n}\right\rangle\right\rangle_{R}=\left\langle\left\langle Q-P_{n}^{R}, Q-P_{n}^{R}\right\rangle\right\rangle_{R}+\left\langle\left\langle P_{n}^{R}, P_{n}^{R}\right\rangle\right\rangle_{R} \tag{2.9}
\end{equation*}
$$

and, in particular,

$$
\begin{equation*}
\left\langle\left\langle P_{n}^{R}, P_{n}^{R}\right\rangle\right\rangle_{R} \leq\left\langle\left\langle Q_{n}, Q_{n}\right\rangle\right\rangle_{R} \tag{2.10}
\end{equation*}
$$

with (by non-triviality) equality if and only if $Q_{n}=P_{n}^{R}$. Since $\operatorname{Tr}$ and det are strictly monotone on strictly positive matrices, we have the following variational principles ((2.11) restates (i) of Lemma 2.2):

Theorem 2.4. For any monic $Q_{n}$ of degree $n$, we have

$$
\begin{align*}
\left\|Q_{n}\right\|_{R} & \geq\left\|P_{n}^{R}\right\|_{R},  \tag{2.11}\\
\operatorname{det}\left\langle\left\langle Q_{n}, Q_{n}\right\rangle\right\rangle_{R} & \geq \operatorname{det}\left\langle\left\langle P_{n}^{R}, P_{n}^{R}\right\rangle_{R}\right. \tag{2.12}
\end{align*}
$$

with equality if and only if $P_{n}^{R}=Q_{n}$.

### 2.1.3 Expansion

Theorem 2.5. Let $d \mu$ be non-trivial.
(i) We have

$$
\begin{equation*}
\left\langle\left\langle P_{k}^{R}, P_{n}^{R}\right\rangle_{R}=\gamma_{n} \delta_{k n}\right. \tag{2.13}
\end{equation*}
$$

for some positive invertible matrices $\gamma_{n}$.
(ii) $\left\{P_{k}^{R}\right\}_{k=0}^{n}$ are a right-module basis for $\mathcal{P}_{n}$; indeed, any $f \in \mathcal{P}_{n}$ has a unique expansion,

$$
\begin{equation*}
f=\sum_{j=0}^{n} P_{j}^{R} f_{j}^{R} \tag{2.14}
\end{equation*}
$$

Indeed, essentially by (1.38),

$$
\begin{equation*}
f_{j}^{R}=\gamma_{j}^{-1}\left\langle\left\langle P_{j}^{R}, f\right\rangle\right\rangle_{R} \tag{2.15}
\end{equation*}
$$

Remark. There are similar formulas for $\langle\langle\cdot, \cdot\rangle\rangle_{L}$. By (2.6),

$$
\begin{equation*}
\left\langle\left\langle P_{k}^{L}, P_{n}^{L}\right\rangle_{L}=\gamma_{n} \delta_{k n}\right. \tag{2.16}
\end{equation*}
$$

(same $\gamma_{n}$, which is why we use $\gamma_{n}$ and not $\gamma_{n}^{R}$ ).
Proof. (i) (2.13) for $n<k$ is immediate from (2.5) and for $n>k$ by symmetry. $\gamma_{n} \geq 0$ follows from (1.22). By Lemma 2.3, $\operatorname{det}\left(\gamma_{n}\right) \neq 0$, so $\gamma_{n}$ is invertible.
(ii) $\operatorname{Map}\left(\mathcal{M}_{l}\right)^{n+1}$ to $\mathcal{P}_{n}$ by

$$
\left\langle\alpha_{0}, \ldots, \alpha_{n}\right\rangle \mapsto \sum_{j=0}^{n} P_{j}^{R} \alpha_{j} \equiv X\left(\alpha_{0}, \ldots, \alpha_{n}\right)
$$

By (2.13),

$$
\alpha_{j}=\gamma_{j}^{-1}\left\langle\left\langle P_{j}^{R}, X\left(\alpha_{0}, \ldots, \alpha_{n}\right)\right\rangle\right\rangle
$$

so that map is one-one. By dimension counting, it is onto.

### 2.1.4 Recurrence Relations for Monic Orthogonal Polynomials

Denote by $\zeta_{n}^{R}$ (resp. $\zeta_{n}^{L}$ ) the coefficient of $x^{n-1}$ in $P_{n}^{R}(x)$ (resp. $P_{n}^{L}(x)$ ), that is,

$$
\begin{aligned}
& P_{n}^{R}(x)=x^{n} \mathbf{1}+\zeta_{n}^{R} x^{n-1}+\text { lower order terms }, \\
& P_{n}^{L}(x)=x^{n} \mathbf{1}+\zeta_{n}^{L} x^{n-1}+\text { lower order terms } .
\end{aligned}
$$

Since $P_{n}^{R}(x)^{\dagger}=P_{n}^{L}(x)$, we have $\left(\zeta_{n}^{R}\right)^{\dagger}=\zeta_{n}^{L}$. Using the parameters $\gamma_{n}$ of (2.13) and $\zeta_{n}^{R}, \zeta_{n}^{L}$ one can write down recurrence relations for $P_{n}^{R}(x), P_{n}^{L}(x)$.

Lemma 2.6. (i) We have a commutation relation

$$
\begin{equation*}
\gamma_{n-1}\left(\zeta_{n}^{R}-\zeta_{n-1}^{R}\right)=\left(\zeta_{n}^{L}-\zeta_{n-1}^{L}\right) \gamma_{n-1} \tag{2.17}
\end{equation*}
$$

(ii) We have the recurrence relations

$$
\begin{align*}
& x P_{n}^{R}(x)=P_{n+1}^{R}(x)+P_{n}^{R}(x)\left(\zeta_{n}^{R}-\zeta_{n+1}^{R}\right)+P_{n-1}(x) \gamma_{n-1}^{-1} \gamma_{n},  \tag{2.18}\\
& x P_{n}^{L}(x)=P_{n+1}^{L}(x)+\left(\zeta_{n}^{L}-\zeta_{n+1}^{L}\right) P_{n}^{L}(x)+\gamma_{n} \gamma_{n-1}^{-1} P_{n-1}^{L}(x) . \tag{2.19}
\end{align*}
$$

Proof. (i) We have

$$
P_{n}^{R}(x)-x P_{n-1}^{R}(x)=\left(\zeta_{n}^{R}-\zeta_{n-1}^{R}\right) x^{n-1}+\text { lower order terms }
$$

and so

$$
\begin{aligned}
\left(\zeta_{n}^{L}-\zeta_{n-1}^{L}\right) \gamma_{n-1} & =\left(\zeta_{n}^{R}-\zeta_{n-1}^{R}\right)^{\dagger}\left\langle\left\langle P_{n-1}^{R}, P_{n-1}^{R}\right\rangle\right\rangle_{R} \\
& =\left(\zeta_{n}^{R}-\zeta_{n-1}^{R}\right)^{\dagger}\left\langle\left\langle x^{n-1}, P_{n-1}^{R}\right\rangle\right\rangle_{R} \\
& =\left\langle\left\langle P_{n}^{R}-x P_{n-1}^{R}, P_{n-1}^{R}\right\rangle\right\rangle_{R} \\
& =\left\langle\left\langle P_{n}^{R}, P_{n-1}^{R}\right\rangle\right\rangle_{R}-\left\langle\left\langle x P_{n-1}^{R}, P_{n-1}^{R}\right\rangle\right\rangle_{R} \\
& =-\left\langle\left\langle x P_{n-1}^{R}, P_{n-1}^{R}\right\rangle\right\rangle_{R} \\
& =-\left\langle\left\langle P_{n-1}^{R}, x P_{n-1}^{R}\right\rangle\right\rangle_{R} \\
& =\left\langle\left\langle P_{n-1}^{R}, P_{n}^{R}-x P_{n-1}^{R}\right\rangle\right\rangle_{R} \\
& =\left\langle\left\langle P_{n-1}^{R}, x^{n-1}\left(\zeta_{n}^{R}-\zeta_{n-1}^{R}\right)\right\rangle\right\rangle_{R} \\
& =\left\langle\left\langle P_{n-1}^{R}, x^{n-1}\right\rangle\right\rangle_{R}\left(\zeta_{n}^{R}-\zeta_{n-1}^{R}\right) \\
& =\gamma_{n-1}\left(\zeta_{n}^{R}-\zeta_{n-1}^{R}\right) .
\end{aligned}
$$

(ii) By Theorem 2.5,

$$
x P_{n}^{R}(x)=P_{n+1}^{R}(x) C_{n+1}+P_{n}^{R}(x) C_{n}+P_{n-1}^{R}(x) C_{n-1}+\cdots+P_{0}^{R} C_{0}
$$

with some matrices $C_{0}, \ldots, C_{n+1}$. It is straightforward that $C_{n+1}=\mathbf{1}$ and $C_{n}=\zeta_{n}^{R}-\zeta_{n+1}^{R}$. By the orthogonality property (2.4), we find $C_{0}=\cdots=C_{n-2}=\mathbf{0}$. Finally, it is easy to calculate $C_{n-1}$ :

$$
\begin{aligned}
\gamma_{n} & =\left\langle\left\langle P_{n}^{R}, x P_{n-1}^{R}\right\rangle_{R}=\left\langle\left\langle x P_{n}^{R}, P_{n-1}^{R}\right\rangle\right\rangle_{R}\right. \\
& =\left\langle\left\langle P_{n+1}^{R}+P_{n}^{R}\left(\zeta_{n}^{R}-\zeta_{n+1}^{R}\right)+P_{n-1}^{R} C_{n-1}, P_{n-1}^{R}\right\rangle\right\rangle_{R} \\
& =C_{n-1}^{\dagger} \gamma_{n-1}
\end{aligned}
$$

and so, taking adjoints and using self-adjointness of $\gamma_{j}, C_{n-1}=\gamma_{n-1}^{-1} \gamma_{n}$. This proves (2.18); the other relation (2.19) is obtained by conjugation.

### 2.1.5 Normalized Orthogonal Polynomials

We call $p_{n}^{R} \in \mathcal{P}$ a right orthonormal polynomial if $\operatorname{deg} p_{n}^{R} \leq n$ and

$$
\begin{gather*}
\left\langle\left\langle p_{n}^{R}, f\right\rangle\right\rangle_{R}=\mathbf{0} \text { for every } f \in \mathcal{P} \text { with } \operatorname{deg} f<n,  \tag{2.20}\\
\left\langle\left\langle p_{n}^{R}, p_{n}^{R}\right\rangle\right\rangle_{R}=\mathbf{1} . \tag{2.21}
\end{gather*}
$$

Similarly, we call $p_{n}^{L} \in \mathcal{P}$ a left orthonormal polynomial if $\operatorname{deg} p_{n}^{L} \leq n$ and

$$
\begin{gather*}
\left\langle\left\langle p_{n}^{L}, f\right\rangle\right\rangle_{L}=\mathbf{0} \text { for every } f \in \mathcal{P} \text { with } \operatorname{deg} f<n,  \tag{2.22}\\
\left\langle\left\langle p_{n}^{L}, p_{n}^{L}\right\rangle_{L}=\mathbf{1} .\right. \tag{2.23}
\end{gather*}
$$

Lemma 2.7. Any orthonormal polynomial has the form

$$
\begin{equation*}
p_{n}^{R}(x)=P_{n}^{R}(x)\left\langle\left\langle P_{n}^{R}, P_{n}^{R}\right\rangle\right\rangle_{R}^{-1 / 2} \sigma_{n}, \quad p_{n}^{L}(x)=\tau_{n}\left\langle\left\langle P_{n}^{L}, P_{n}^{L}\right\rangle\right\rangle_{L}^{-1 / 2} P_{n}^{L}(x) \tag{2.24}
\end{equation*}
$$

where $\sigma_{n}, \tau_{n} \in \mathcal{M}_{l}$ are unitaries. In particular, $\operatorname{deg} p_{n}^{R}=\operatorname{deg} p_{n}^{L}=n$.
Proof. Let $K_{n}$ be the coefficient of $x^{n}$ in $p_{n}^{R}$. Consider the polynomial $q(x)=P_{n}^{R}(x) K_{n}-p_{n}^{R}(x)$, where $P_{n}^{R}$ is the monic orthogonal polynomial from Lemma 2.2. Then $\operatorname{deg} q<n$ and so from (2.4) and (2.20), it follows that $\langle\langle q, q\rangle\rangle_{R}=0$ and so $q(x)$ vanishes identically. Thus, we have

$$
\begin{equation*}
\mathbf{1}=\left\langle\left\langle p_{n}^{R}, p_{n}^{R}\right\rangle\right\rangle_{R}=K_{n}^{\dagger}\left\langle\left\langle P_{n}^{R}, P_{n}^{R}\right\rangle\right\rangle_{R} K_{n} \tag{2.25}
\end{equation*}
$$

and so $\operatorname{det}\left(K_{n}\right) \neq 0$. From (2.25) we get $\left(K_{n}^{\dagger}\right)^{-1} K_{n}^{-1}=\left\langle\left\langle P_{n}^{R}, P_{n}^{R}\right\rangle\right\rangle_{R}$, and so $K_{n} K_{n}^{\dagger}=\left\langle\left\langle P_{n}^{R}, P_{n}^{R}\right\rangle\right\rangle_{R}^{-1}$. From here we get $K_{n}=\left\langle\left\langle P_{n}^{R}, P_{n}^{R}\right\rangle\right\rangle_{R}^{-1 / 2} \sigma_{n}$ with a unitary $\sigma_{n}$. The proof for $p_{n}^{L}$ is similar.

By Theorem 2.5, the polynomials $p_{n}^{R}$ form a right orthonormal module basis in $\mathcal{P}_{R}$. Thus, for any $f \in \mathcal{P}_{R}$, we have

$$
\begin{equation*}
f(x)=\sum_{m=0}^{\infty} p_{m}^{R} f_{m}, \quad f_{m}=\left\langle\left\langle p_{m}^{R}, f\right\rangle\right\rangle_{R} \tag{2.26}
\end{equation*}
$$

and the Parseval identity

$$
\begin{equation*}
\sum_{m=0}^{\infty} \operatorname{Tr}\left(f_{m} f_{m}^{\dagger}\right)=\|f\|_{R}^{2} \tag{2.27}
\end{equation*}
$$

holds true. Obviously, since $f$ is a polynomial, there are only finitely many non-zero terms in (2.26) and (2.27).

### 2.2 Block Jacobi Matrices

The study of block Jacobi matrices goes back at least to Krein [133].

### 2.2.1 Block Jacobi Matrices as Matrix Representations

Suppose that a sequence of unitary matrices $\mathbf{1}=\sigma_{0}, \sigma_{1}, \sigma_{2}, \ldots$ is fixed, and $p_{n}^{R}$ are defined according to (2.24). As noted above, $p_{n}^{R}$ form a right orthonormal basis in $\mathcal{P}_{R}$.

The map $f(x) \mapsto x f(x)$ can be considered as a right homomorphism in $\mathcal{P}_{R}$. Consider the matrix $J_{n m}$ of this homomorphism with respect to the basis $p_{n}^{R}$, that is,

$$
\begin{equation*}
J_{n m}=\left\langle\left\langle p_{n-1}^{R}, x p_{m-1}^{R}\right\rangle_{R} .\right. \tag{2.28}
\end{equation*}
$$

Following Killip-Simon [128] and Simon [167, 168, 176], our Jacobi matrices are indexed with $n=1,2, \ldots$ but, of course, $p_{n}$ has $n=0,1,2, \ldots$. That is why (2.28) has $n-1$ and $m-1$.

As in the scalar case, using the orthogonality properties of $p_{n}^{R}$, we get that $J_{n m}=\mathbf{0}$ if $|n-m|>1$. Denote

$$
B_{n}=J_{n n}=\left\langle\left\langle p_{n-1}^{R}, x p_{n-1}^{R}\right\rangle\right\rangle_{R}
$$

and

$$
A_{n}=J_{n, n+1}=J_{n+1, n}^{\dagger}=\left\langle\left\langle p_{n-1}^{R}, x p_{n}^{R}\right\rangle\right\rangle_{R} .
$$

Then we have

$$
J=\left(\begin{array}{cccc}
B_{1} & A_{1} & \mathbf{0} & \cdots  \tag{2.29}\\
A_{1}^{\dagger} & B_{2} & A_{2} & \cdots \\
\mathbf{0} & A_{2}^{\dagger} & B_{3} & \cdots \\
\vdots & \vdots & \vdots & \ddots
\end{array}\right)
$$

Applying (2.26) to $f(x)=x p_{n}^{R}(x)$, we get the recurrence relation

$$
\begin{equation*}
x p_{n}^{R}(x)=p_{n+1}^{R}(x) A_{n+1}^{\dagger}+p_{n}^{R}(x) B_{n+1}+p_{n-1}^{R}(x) A_{n}, \quad n=1,2, \ldots . \tag{2.30}
\end{equation*}
$$

If we set $p_{-1}^{R}(x)=\mathbf{0}$ and $A_{0}=\mathbf{1}$, the relation (2.30) also holds for $n=0$. By (2.2), we can always pick $p_{n}^{L}$ so that for $x$ real, $p_{n}^{L}(x)=p_{n}^{R}(x)^{\dagger}$, and thus for complex $z$,

$$
\begin{equation*}
p_{n}^{L}(z)=p_{n}^{R}(\bar{z})^{\dagger} \tag{2.31}
\end{equation*}
$$

by analytic continuation. By conjugating (2.30), we get

$$
\begin{equation*}
x p_{n}^{L}(x)=A_{n+1} p_{n+1}^{L}(x)+B_{n+1} p_{n}^{L}(x)+A_{n}^{\dagger} p_{n-1}^{L}(x), \quad n=0,1,2, \ldots . \tag{2.32}
\end{equation*}
$$

Comparing this with the recurrence relations (2.18), (2.19), we get

$$
\begin{equation*}
A_{n}=\sigma_{n-1}^{\dagger} \gamma_{n-1}^{-1 / 2} \gamma_{n}^{1 / 2} \sigma_{n}, \quad B_{n}=\sigma_{n-1}^{\dagger} \gamma_{n-1}^{1 / 2}\left(\zeta_{n-1}^{R}-\zeta_{n}^{R}\right) \gamma_{n-1}^{-1 / 2} \sigma_{n-1} \tag{2.33}
\end{equation*}
$$

In particular, $\operatorname{det} A_{n} \neq 0$ for all $n$.
Notice that since $\sigma_{n}$ is unitary, $\left|\operatorname{det}\left(\sigma_{n}\right)\right|=1$, so (2.33) implies $\operatorname{det}\left(\gamma_{n}^{1 / 2}\right)=\operatorname{det}\left(\gamma_{n-1}^{1 / 2}\right)\left|\operatorname{det}\left(A_{n}\right)\right|$ which, by induction, implies that

$$
\begin{equation*}
\operatorname{det}\left\langle\left\langle P_{n}^{R}, P_{n}^{R}\right\rangle\right\rangle=\left|\operatorname{det}\left(A_{1} \cdots A_{n}\right)\right|^{2} . \tag{2.34}
\end{equation*}
$$

Any block matrix of the form (2.29) with $B_{n}=B_{n}^{\dagger}$ and $\operatorname{det} A_{n} \neq 0$ for all $n$ will be called a block Jacobi matrix corresponding to the Jacobi parameters $A_{n}$ and $B_{n}$.

### 2.2.2 Basic Properties of Block Jacobi Matrices

Suppose we are given a block Jacobi matrix $J$ corresponding to Jacobi parameters $A_{n}$ and $B_{n}$, where $B_{n}=B_{n}^{\dagger}$ and $\operatorname{det} A_{n} \neq 0$ for each $n$.

Consider the Hilbert space $\mathcal{H}_{v}=\ell^{2}\left(\mathbb{Z}_{+}, \mathbb{C}^{l}\right)$ (here $\left.\mathbb{Z}_{+}=\{1,2,3, \ldots\}\right)$ with inner product

$$
\langle f, g\rangle_{\mathcal{H}_{v}}=\sum_{n=1}^{\infty}\left\langle f_{n}, g_{n}\right\rangle_{\mathbb{C}^{l}}
$$

and orthonormal basis $\left\{e_{k, j}\right\}_{k \in \mathbb{Z}_{+}, l \leq j \leq l}$, where

$$
\left(e_{k, j}\right)_{n}=\delta_{k, n} v_{j}
$$

and $\left\{v_{j}\right\}_{1 \leq j \leq l}$ is the standard basis of $\mathbb{C}^{l} . J$ acts on $\mathcal{H}_{v}$ via

$$
\begin{equation*}
(J f)_{n}=A_{n-1}^{\dagger} f_{n-1}+B_{n} f_{n}+A_{n} f_{n+1}, \quad f \in \mathcal{H}_{v} \tag{2.35}
\end{equation*}
$$

(with $f_{0}=0$ ) and defines a symmetric operator on this space. Note that using invertibility of the $A_{n}$ 's, induction shows

$$
\begin{equation*}
\operatorname{span}\left\{e_{k, j}: 1 \leq k \leq n, 1 \leq j \leq l\right\}=\operatorname{span}\left\{J^{k-1} e_{1, j}: 1 \leq k \leq n, 1 \leq j \leq l\right\} \tag{2.36}
\end{equation*}
$$

for every $n \geq 1$. We want to emphasize that elements of $\mathcal{H}_{v}$ and $\mathcal{H}$ are vector-valued and matrixvalued, respectively. For this reason, we will be interested in both matrix- and vector-valued solutions of the basic difference equations.

We will consider only bounded block Jacobi matrices, that is, those corresponding to Jacobi parameters satisfying

$$
\begin{equation*}
\sup _{n} \operatorname{Tr}\left(A_{n}^{\dagger} A_{n}+B_{n}^{\dagger} B_{n}\right)<\infty . \tag{2.37}
\end{equation*}
$$

Equivalently,

$$
\begin{equation*}
\sup _{n}\left(\left\|A_{n}\right\|+\left\|B_{n}\right\|\right)<\infty \tag{2.38}
\end{equation*}
$$

In this case, $J$ is a bounded self-adjoint operator. This is equivalent to $\mu$ having compact support.
We call two Jacobi matrices $J$ and $\tilde{J}$ equivalent if there exists a sequence of unitaries $u_{n} \in \mathcal{M}_{l}$, $n \geq 1$, with $u_{1}=\mathbf{1}$ such that $\tilde{J}_{n m}=u_{n}^{\dagger} J_{n m} u_{m}$. From Lemma 2.7 it is clear that if $p_{n}^{R}, \tilde{p}_{n}^{R}$ are two sequences of normalized orthogonal polynomials, corresponding to the same measure (but having different normalization), then the Jacobi matrices $J_{n m}=\left\langle\left\langle p_{n-1}^{R}, x p_{m-1}^{R}\right\rangle\right\rangle_{R}$ and $\tilde{J}_{n m}=$ $\left\langle\left\langle\tilde{p}_{n-1}^{R}, x \tilde{p}_{m-1}^{R}\right\rangle_{R}\right.$ are equivalent $\left(u_{n}=\sigma_{n-1}^{\dagger} \tilde{\sigma}_{n-1}\right)$. Thus,

$$
\begin{equation*}
\tilde{B}_{n}=u_{n}^{\dagger} B_{n} u_{n}, \quad \tilde{A}_{n}=u_{n}^{\dagger} A_{n} u_{n+1} \tag{2.39}
\end{equation*}
$$

Therefore, we have a map

$$
\begin{align*}
& \Phi: \mu \mapsto\left\{J: J_{m n}=\left\langle\left\langle p_{n-1}^{R}, x p_{m-1}^{R}\right\rangle\right\rangle_{R}, p_{n}^{R}\right.  \tag{2.40}\\
& \quad \text { correspond to } d \mu \text { for some normalization }\}
\end{align*}
$$

from the set of all Hermitian positive semi-definite non-trivial compactly supported measures to the set of all equivalence classes of bounded block Jacobi matrices. Below, we will see how to invert this map.

### 2.2.3 Special Representatives of the Equivalence Classes

Let $J$ be a block Jacobi matrix with the Jacobi parameters $A_{n}, B_{n}$. We say that $J$ is:

- of type 1 , if $A_{n}>\mathbf{0}$ for all $n$;
- of type 2 , if $A_{1} A_{2} \cdots A_{n}>\mathbf{0}$ for all $n$;
- of type 3 , if $A_{n} \in \mathcal{L}$ for all $n$.

Here, $\mathcal{L}$ is the class of all lower triangular matrices with strictly positive elements on the diagonal. Type 3 is of interest because they correspond precisely to bounded Hermitian matrices with $2 l+1$ non-vanishing diagonals with the extreme diagonals strictly positive; see Section 1.4(c). Type 2 is the case where the leading coefficients of $p_{n}^{R}$ are strictly positive definite.

Theorem 2.8. (i) Each equivalence class of block Jacobi matrices contains exactly one element each of type 1, type 2, or type 3.
(ii) Let $J$ be a block Jacobi matrix corresponding to a sequence of polynomials $p_{n}^{R}$ as in (2.24). Then $J$ is of type 2 if and only if $\sigma_{n}=\mathbf{1}$ for all $n$.
Proof. The proof is based on the following two well-known facts:
(a) For any $t \in \mathcal{M}_{l}$ with $\operatorname{det}(t) \neq 0$, there exists a unique unitary $u \in \mathcal{M}_{l}$ such that $t u$ is Hermitian positive semi-definite: $t u \geq 0$.
(b) For any $t \in \mathcal{M}_{l}$ with $\operatorname{det}(t) \neq 0$, there exists a unique unitary $u \in \mathcal{M}_{l}$ such that $t u \in \mathcal{L}$.

We first prove that every equivalence class of block Jacobi matrices contains at least one element of type 1. For a given sequence $A_{n}$, let us construct a sequence $u_{1}=\mathbf{1}, u_{2}, u_{3}, \ldots$ of unitaries such that $u_{n}^{\dagger} A_{n} u_{n+1} \geq \mathbf{0}$. By the existence part of (a), we find $u_{2}$ such that $A_{1} u_{2} \geq \mathbf{0}$, then find $u_{3}$ such that $u_{2}^{\dagger} A_{2} u_{3} \geq \mathbf{0}$, etc. This, together with (2.39), proves the statement. In order to prove the uniqueness part, suppose we have $A_{n} \geq \mathbf{0}$ and $u_{n}^{\dagger} A_{n} u_{n+1} \geq \mathbf{0}$ for all $n$. Then, by the uniqueness part of (a), $A_{1} \geq \mathbf{0}$ and $A_{1} u_{2} \geq \mathbf{0}$ imply $u_{2}=\mathbf{1}$; next, $A_{2} \geq \mathbf{0}$ and $u_{2}^{\dagger} A_{2} u_{3}=A_{2} u_{3} \geq \mathbf{0}$ imply $u_{3}=1$, etc.

The statement (i) concerning type 3 can be proven in the same way, using (b) instead of (a).
The statement (i) concerning type 2 can be proven similarly. Existence: find $u_{2}$ such that $A_{1} u_{2} \geq \mathbf{0}$, then $u_{3}$ such that $\left(A_{1} u_{2}\right)\left(u_{2}^{\dagger} A_{2} u_{3}\right)=A_{1} A_{2} u_{3} \geq \mathbf{0}$, etc. Uniqueness: if $A_{1} \ldots A_{n} \geq \mathbf{0}$ and $A_{1} \cdots A_{n} u_{n+1} \geq \mathbf{0}$, then $u_{n+1}=\mathbf{1}$.

By (2.33), we have $A_{1} A_{2} \cdots A_{n}=\gamma_{n}^{1 / 2} \sigma_{n}$ and the statement (ii) follows from the positivity of $\gamma_{n}$.

We say that a block Jacobi matrix $J$ belongs to the Nevai class if

$$
B_{n} \rightarrow \mathbf{0} \text { and } A_{n}^{\dagger} A_{n} \rightarrow \mathbf{1} \text { as } n \rightarrow \infty .
$$

It is clear that $J$ is in the Nevai class if and only if all equivalent Jacobi matrices belong to the Nevai class.

Theorem 2.9. If $J$ belongs to the Nevai class and is of type 1 or type 3, then $A_{n} \rightarrow 1$ as $n \rightarrow \infty$. Proof. If $J$ is of type 1 , then $A_{n}^{\dagger} A_{n}=A_{n}^{2} \rightarrow \mathbf{1}$ clearly implies $A_{n} \rightarrow \mathbf{1}$ since square root is continuous on positive Hermitian matrices.

Suppose $J$ is of type 3 . We shall prove that $A_{n} \rightarrow \mathbf{1}$ by considering the rows of the matrix $A_{n}$ one by one, starting from the $l$ th row. Denote $\left(A_{n}\right)_{j k}=a_{j, k}^{(n)}$. We have

$$
\left(A_{n}^{\dagger} A_{n}\right)_{l l}=\left(a_{l, l}^{(n)}\right)^{2} \rightarrow 1, \text { and so } a_{l, l}^{(n)} \rightarrow 1
$$

Then, for any $k<l$, we have

$$
\left(A_{n}^{\dagger} A_{n}\right)_{l k}=a_{l, l}^{(n)} a_{l, k}^{(n)} \rightarrow 0, \text { and so } a_{l, k}^{(n)} \rightarrow 0
$$

Next, consider the $(l-1)$ st row. We have

$$
\left(A_{n}^{\dagger} A_{n}\right)_{l-1, l-1}=\left(a_{l-1, l-1}^{(n)}\right)^{2}+\left|a_{l, l-1}^{(n)}\right|^{2} \rightarrow 1
$$

and so, using the previous step, $a_{l-1, l-1}^{(n)} \rightarrow 1$ as $n \rightarrow \infty$. Then for all $k<l-1$, we have

$$
\left(A_{n}^{\dagger} A_{n}\right)_{l-1, k}=\overline{a_{l-1, l-1}^{(n)}} a_{l-1, k}^{(n)}+\overline{a_{l, l-1}^{(n)}} a_{l, k}^{(n)} \rightarrow 0
$$

and so, using the previous steps, $a_{l-1, k} \rightarrow 0$. Continuing this way, we get $a_{j, k}^{(n)} \rightarrow \delta_{j, k}$ as required.
It is an interesting open question if this result also applies to the type 2 case.

### 2.2.4 Favard's Theorem

Here we construct an inverse of the mapping $\Phi$ (defined by (2.40)). Thus, $\Phi$ sets up a bijection between non-trivial measures of compact support and equivalence classes of bounded block Jacobi matrices.

Before proceeding to do this, let us prove:
Lemma 2.10. The mapping $\Phi$ is injective.
Proof. Let $\mu$ and $\tilde{\mu}$ be two Hermitian positive semi-definite non-trivial compactly supported measures. Suppose that $\Phi(\mu)=\Phi(\tilde{\mu})$.

Let $p_{n}^{R}$ and $\tilde{p}_{n}^{R}$ be normalized orthogonal polynomials corresponding to $\mu$ and $\tilde{\mu}$. Suppose that the normalization both for $p_{n}^{R}$ and for $\tilde{p}_{n}^{R}$ has been chosen such that $\sigma_{n}=\mathbf{1}$ (see (2.24)), that is, type 2. From Lemma 2.8 and the assumption $\Phi(\mu)=\Phi(\tilde{\mu})$ it follows that the corresponding Jacobi matrices coincide, that is, $\left\langle\left\langle p_{n}^{R}, x p_{m}^{R}\right\rangle\right\rangle_{R}=\left\langle\left\langle\tilde{p}_{n}^{R}, x \tilde{p}_{m}^{R}\right\rangle\right\rangle_{R}$ for all $n$ and $m$. Together with the recurrence relation (2.30) this yields $p_{n}^{R}=\tilde{p}_{n}^{R}$ for all $n$.

For any $n \geq 0$, we can represent $x^{n}$ as

$$
x^{n}=\sum_{k=0}^{n} p_{k}^{R}(x) C_{k}^{(n)}=\sum_{k=0}^{n} \tilde{p}_{k}^{R}(x) \tilde{C}_{k}^{(n)} .
$$

The coefficients $C_{k}^{(n)}$ and $\tilde{C}_{k}^{(n)}$ are completely determined by the coefficients of the polynomials $p_{n}^{R}$ and $\tilde{p}_{n}^{R}$ and so $C_{k}^{(n)}=\tilde{C}_{k}^{(n)}$ for all $n$ and $k$.

For the moments of the measure $\mu$, we have

$$
\int x^{n} d \mu(x)=\left\langle\left\langle\mathbf{1}, x^{n}\right\rangle\right\rangle_{R}=\sum_{k=0}^{n}\left\langle\left\langle\mathbf{1}, p_{k}^{R} C_{k}^{(n)}\right\rangle\right\rangle_{R}=\langle\langle\mathbf{1}, \mathbf{1}\rangle\rangle_{R} C_{0}^{(n)}=C_{0}^{(n)} .
$$

Since the same calculation is valid for the measure $\tilde{\mu}$, we get

$$
\int x^{n} d \mu(x)=\int x^{n} d \tilde{\mu}(x)
$$

for all $n$. It follows that

$$
\int f(x) d \mu(x) g(x)=\int f(x) d \tilde{\mu}(x) g(x)
$$

for all matrix-valued polynomials $f$ and $g$, and so the measures $\mu$ and $\tilde{\mu}$ coincide.
We can now construct the inverse of the map $\Phi$. Let a block Jacobi matrix $J$ be given. By a version of the spectral theorem for self-adjoint operators with finite multiplicity (see, e.g., [2, Sect. 72]), there exists a matrix-valued measure $d \mu$ with

$$
\begin{equation*}
\left\langle e_{1, j}, f(J) e_{1, k}\right\rangle_{\mathcal{H}_{v}}=\int f(x) d \mu_{j, k}(x) \tag{2.41}
\end{equation*}
$$

and an isometry

$$
R: \mathcal{H}_{v} \rightarrow L^{2}\left(\mathbb{R}, d \mu ; \mathbb{C}^{l}\right)
$$

such that (recall that $\left\{v_{j}\right\}$ is the standard basis in $\mathbb{C}^{l}$ )

$$
\begin{equation*}
\left[R e_{1, j}\right](x)=v_{j}, \quad 1 \leq j \leq l \tag{2.42}
\end{equation*}
$$

and, for any $g \in \mathcal{H}_{v}$, we have

$$
\begin{equation*}
(R J g)(x)=x(R g)(x) \tag{2.43}
\end{equation*}
$$

If the Jacobi matrices $J$ and $\tilde{J}$ are equivalent, then we have $\tilde{J}=U^{*} J U$ for some $U=\oplus_{n=1}^{\infty} u_{n}$, $u_{1}=\mathbf{1}$. Thus,

$$
\left\langle e_{1, j}, f(\tilde{J}) e_{1, k}\right\rangle_{\mathcal{H}_{v}}=\left\langle U e_{1, j}, f(J) U e_{1, k}\right\rangle_{\mathcal{H}_{v}}=\left\langle e_{1, j}, f(J) e_{1, k}\right\rangle_{\mathcal{H}_{v}}
$$

and so the measures corresponding to $J$ and $\tilde{J}$ coincide. Thus, we have a map

$$
\begin{equation*}
\Psi:\{\tilde{J}: \tilde{J} \text { is equivalent to } J\} \mapsto \mu \tag{2.44}
\end{equation*}
$$

from the set of all equivalence classes of bounded block Jacobi matrices to the set of all Hermitian positive semi-definite compactly supported measures.

Theorem 2.11. (i) All measures in the image of the map $\Psi$ are non- degenerate.
(ii) $\Phi \circ \Psi=i d$.
(iii) $\Psi \circ \Phi=i d$.

Proof. (i) To put things in the right context, we first recall that $\|\cdot\|_{\mathcal{H}_{v}}$ is a norm (rather than a semi-norm), whereas $|\cdot|$ on $\mathcal{V}$ (cf. (2.3)) is, in general, a semi-norm. Using the assumption that $\operatorname{det}\left(A_{k}\right) \neq 0$ for all $k$ (which is included in our definition of a Jacobi matrix), we will prove that $|\cdot|$ is in fact a norm. More precisely, we will prove that $|p|>0$ for any non-zero polynomial $p \in \mathcal{V}$; by Lemma 2.1 this will imply that $\mu$ is non-degenerate.

Let $p \in \mathcal{V}$ be a non-zero polynomial, $\operatorname{deg} p=n$. Notice that (2.42) and (2.43) give

$$
\begin{equation*}
\left[R J^{k} e_{1, j}\right](x)=x^{k} v_{j} \tag{2.45}
\end{equation*}
$$

for every $k \geq 0$ and $1 \leq j \leq l$. This shows that $p$ can be represented as $p=R g$, where $g=$ $\sum_{k=0}^{n} J^{k} f_{k}$, and $f_{0}, \ldots, f_{n}$ are vectors in $\mathcal{H}_{v}$ such that $\left\langle f_{i}, e_{j, k}\right\rangle_{\mathcal{H}_{v}}=0$ for all $i=0, \ldots, n, j \geq 2$, $k=1, \ldots, l$ (i.e., the only non-zero components of $f_{j}$ are in the first $\mathbb{C}^{l}$ in $\mathcal{H}_{v}$ ). Assumption $\operatorname{deg} p=n$ means $f_{n} \neq 0$.

Since $R$ is isometric, we have $|p|=\|g\|_{\mathcal{H}_{v}}$, and so we have to prove that $g \neq 0$. Indeed, suppose that $g=0$. Using the assumption $\operatorname{det}\left(A_{k}\right) \neq 0$ and the tri-diagonal nature of $J$, we see that $\sum_{k=0}^{n} J^{k} f_{k}=0$ yields $f_{n}=0$, contrary to our assumption.
(ii) Consider the elements $R e_{n, k} \in L^{2}\left(\mathbb{R}, d \mu ; \mathbb{C}^{l}\right)$. First note that, by (2.36) and (2.45), $R e_{n, k}$ is a polynomial of degree at most $n-1$. Next, by the unitarity of $R$, we have

$$
\begin{equation*}
\left\langle R e_{n, k}, R e_{m, j}\right\rangle_{L^{2}\left(\mathbb{R}, d \mu ; \mathbb{C}^{l}\right)}=\delta_{m, n} \delta_{k, j} \tag{2.46}
\end{equation*}
$$

Let us construct matrix-valued polynomials $q_{n}(x)$, using $R e_{n, 1}, R e_{n, 2}, \ldots, R e_{n, l}$ as columns of $q_{n-1}(x)$ :

$$
\left[q_{n-1}(x)\right]_{j, k}=\left[R e_{n, k}(x)\right]_{j} .
$$

We have $\operatorname{deg} q_{n} \leq n$ and $\left\langle\left\langle q_{m}, q_{n}\right\rangle\right\rangle_{R}=\delta_{m, n} \mathbf{1}$; the last relation is just a reformulation of (2.46). Hence the $q_{n}$ 's are right normalized orthogonal polynomials with respect to the measure $d \mu$. We find

$$
\begin{aligned}
J_{n m} & =\left[\left\langle e_{n, j}, J e_{m, k}\right\rangle_{\mathcal{H}_{v}}\right]_{1 \leq j, k \leq l} \\
& =\left[\left\langle R e_{n, j}, R J e_{m, k}\right\rangle_{L^{2}\left(\mathbb{R}, d \mu ; \mathbb{C}^{l}\right)}\right]_{1 \leq j, k \leq l}
\end{aligned}
$$

$$
\begin{aligned}
& =\left[\left\langle R e_{n, j}, x R e_{m, k}\right\rangle_{L^{2}\left(\mathbb{R}, d \mu ; \mathbb{C}^{l}\right)}\right]_{1 \leq j, k \leq l} \\
& =\left[\left\langle\left[q_{n-1}(x)\right]_{, j}, x\left[q_{m-1}(x)\right]_{\cdot, k}\right\rangle_{L^{2}\left(\mathbb{R}, d \mu ; \mathbb{C}^{l}\right)}\right]_{1 \leq j, k \leq l} \\
& =\left\langle\left\langle q_{n-1}, x q_{m-1}\right\rangle\right\rangle_{R}
\end{aligned}
$$

as required.
(iii) Follows from (ii) and from Lemma 2.10.

### 2.3 The $m$-Function

### 2.3.1 The Definition of the $m$-Function

We denote the Borel transform of $d \mu$ by $m$ :

$$
\begin{equation*}
m(z)=\int \frac{d \mu(x)}{x-z}, \quad \operatorname{Im} z>0 \tag{2.47}
\end{equation*}
$$

It is a matrix-valued Herglotz function, that is, it is analytic and obeys $\operatorname{Im} m(z)>0$. For information on matrix-valued Herglotz functions, see [98] and references therein. Extensions to operator-valued Herglotz functions can be found in [94].
Lemma 2.12. Suppose $d \mu$ is given, $p_{n}^{R}$ are right normalized orthogonal polynomials, and $J$ is the associated block Jacobi matrix. Then,

$$
\begin{equation*}
m(z)=\left\langle\left\langle p_{0}^{R},(x-z)^{-1} p_{0}^{R}\right\rangle\right\rangle_{R} \tag{2.48}
\end{equation*}
$$

and

$$
\begin{equation*}
m(z)=\left\langle e_{1, .},(J-z)^{-1} e_{1, .}\right\rangle_{\mathcal{H}_{v}} . \tag{2.49}
\end{equation*}
$$

Proof. Since $p_{0}^{R}=1,(2.48)$ is just a way of rewriting the definition of $m$. The second identity, (2.49), is a consequence of (2.41) and Theorem 2.11(iii).

### 2.3.2 Coefficient Stripping

If $J$ is a block Jacobi matrix corresponding to invertible $A_{n}$ 's and Hermitian $B_{n}$ 's, we denote the $k$-times stripped block Jacobi matrix, corresponding to $\left\{A_{k+n}, B_{k+n}\right\}_{n \geq 1}$, by $J^{(k)}$. That is,

$$
J^{(k)}=\left(\begin{array}{cccc}
B_{k+1} & A_{k+1} & \mathbf{0} & \cdots \\
A_{k+1}^{\dagger} & B_{k+2} & A_{k+2} & \cdots \\
\mathbf{0} & A_{k+2}^{\dagger} & B_{k+3} & \cdots \\
\vdots & \vdots & \vdots & \ddots
\end{array}\right)
$$

The $m$-function corresponding to $J^{(k)}$ will be denoted by $m^{(k)}$. Note that, in particular, $J^{(0)}=J$ and $m^{(0)}=m$.

Proposition 2.13. Let $J$ be a block Jacobi matrix with $\sigma_{\text {ess }}(J) \subseteq[a, b]$. Then, for every $\varepsilon>0$, there is $k_{0} \geq 0$ such that for $k \geq k_{0}$, we have that $\sigma\left(J^{(k)}\right) \subseteq[a-\varepsilon, b+\varepsilon]$.

Proof. This is an immediate consequence of (the proof of) [42, Lemma 1].

Proposition 2.14 (Due to Aptekarev-Nikishin [4]). We have that

$$
m^{(k)}(z)^{-1}=B_{k+1}-z-A_{k+1} m^{(k+1)}(z) A_{k+1}^{\dagger}
$$

for $\operatorname{Im} z>0$ and $k \geq 0$.
Proof. It suffices to handle the case $k=0$. Given (2.49), this is a special case of a general formula for $2 \times 2$ block operator matrices, due to Schur [163], that reads

$$
\left(\begin{array}{ll}
A & B \\
C & D
\end{array}\right)^{-1}=\left(\begin{array}{cc}
\left(A-B D^{-1} C\right)^{-1} & -A^{-1} B\left(D-C A^{-1} B\right)^{-1} \\
-D^{-1} C\left(A-B D^{-1} C\right)^{-1} & \left(D-C A^{-1} B\right)^{-1}
\end{array}\right)
$$

and which is readily verified. Here $A=B_{1}-z, B=A_{1}, C=A_{1}^{\dagger}$, and $D=J^{(1)}-z$.

### 2.4 Second Kind Polynomials

Define the second kind polynomials by $q_{-1}^{R}(z)=-\mathbf{1}$,

$$
q_{n}^{R}(z)=\int_{\mathbb{R}} d \mu(x) \frac{p_{n}^{R}(z)-p_{n}^{R}(x)}{z-x}, \quad n=0,1,2, \ldots
$$

As is easy to see, for $n \geq 1, q_{n}^{R}$ is a polynomial of degree $n-1$. For future reference, let us display the first several polynomials $p_{n}^{R}$ and $q_{n}^{R}$ :

$$
\begin{array}{rlrl}
p_{-1}^{R}(x)=\mathbf{0}, & & p_{0}^{R}(x)=\mathbf{1}, & p_{1}^{R}(x)=\left(x-B_{1}\right) A_{1}^{-1}, \\
q_{-1}^{R}(x)=-\mathbf{1}, & q_{0}^{R}(x)=\mathbf{0}, & q_{1}^{R}(x)=A_{1}^{-1} . \tag{2.51}
\end{array}
$$

The polynomials $q_{n}^{R}$ satisfy the equation (same form as (2.30))

$$
\begin{equation*}
x q_{n}^{R}(x)=q_{n+1}^{R}(x) A_{n+1}^{\dagger}+q_{n}^{R}(x) B_{n+1}+q_{n-1}^{R}(x) A_{n}, \quad n=0,1,2, \ldots . \tag{2.52}
\end{equation*}
$$

For $n=0$, this can be checked by a direct substitution of (2.51). For $n \geq 1$, as in the scalar case, this can be checked by taking (2.30) for $x$ and for $z$, subtracting, dividing by $x-z$, integrating over $d \mu$, and taking into account the orthogonality relation

$$
\int d \mu(x) p_{n}^{R}(x)=\mathbf{0}, \quad n \geq 1
$$

Finally, let us define

$$
\psi_{n}^{R}(z)=q_{n}^{R}(z)+m(z) p_{n}^{R}(z) .
$$

According to the definition of $q_{n}^{R}$, we have

$$
\psi_{n}^{R}(z)=\left\langle\left\langle f_{z}, p_{n}^{R}\right\rangle\right\rangle_{R}, \quad f_{z}(x)=(x-\bar{z})^{-1}
$$

By the Parseval identity, this shows that for all $\operatorname{Im} z>0$, the sequence $\psi_{n}^{R}(z)$ is in $\ell^{2}$, that is,

$$
\begin{equation*}
\sum_{n=0}^{\infty} \operatorname{Tr}\left(\psi_{n}^{R}(z)^{\dagger} \psi_{n}^{R}(z)\right)<\infty \tag{2.53}
\end{equation*}
$$

In the same way, we define $q_{-1}^{L}(z)=-\mathbf{1}$,

$$
q_{n}^{L}(z)=\int_{\mathbb{R}} \frac{p_{n}^{L}(z)-p_{n}^{L}(x)}{z-x} d \mu(x), \quad n=0,1,2, \ldots
$$

and $\psi_{n}^{L}(z)=q_{n}^{L}(z)+p_{n}^{L}(z) m(z)$.

### 2.5 Solutions to the Difference Equations

For $\operatorname{Im} z>0$, consider the solutions to the equations

$$
\begin{array}{ll}
z u_{n}(z)=\sum_{m=1}^{\infty} u_{m}(z) J_{m n}, & n=2,3, \ldots, \\
z v_{n}(z)=\sum_{m=1}^{\infty} J_{n m} v_{m}(z), & n=2,3, \ldots \tag{2.55}
\end{array}
$$

Clearly, $u_{n}(z)$ solves (2.54) if and only if $v_{n}(z)=\left(u_{n}(\bar{z})\right)^{\dagger}$ solves (2.55). In the above, we normally intend $z$ as a fixed parameter but it then can be used as a variable. That is, $z$ is fixed and $u_{n}(z)$ is a fixed sequence, not a $z$-dependent function. A statement like $v_{n}(z)=\left(u_{n}(\bar{z})\right)^{\dagger}$ means if $u_{n}$ is a sequence solving (2.54) for $z=\bar{z}_{0}$, then $v_{n}$ obeys (2.55) for $z=z_{0}$. Of course, if $u_{n}(z)$ is a function of $z$ in a region, we can apply our estimates to all $z$ in the region. For any solution $\left\{u_{n}(z)\right\}_{n=1}^{\infty}$ of (2.54), let us define

$$
\begin{equation*}
u_{0}(z)=z u_{1}(z)-u_{1}(z) B_{1}-u_{2}(z) A_{1}^{\dagger} . \tag{2.56}
\end{equation*}
$$

With this definition, the equation (2.54) for $n=1$ is equivalent to $u_{0}(z)=0$. In the same way, we set

$$
v_{0}(z)=z v_{1}(z)-B_{1} v_{1}(z)-A_{1} v_{2}(z) .
$$

Lemma 2.15. Let $\operatorname{Im} z>0$ and suppose $\left\{u_{n}(z)\right\}_{n=0}^{\infty}$ solves (2.54) (for $n \geq 2$ ) and (2.56) and belongs to $\ell^{2}$. Then

$$
\begin{equation*}
(\operatorname{Im} z) \sum_{n=1}^{\infty} \operatorname{Tr}\left(u_{n}(z)^{\dagger} u_{n}(z)\right)=-\operatorname{Im} \operatorname{Tr}\left(u_{1}(z) u_{0}(z)^{\dagger}\right) \tag{2.57}
\end{equation*}
$$

In particular, $u_{n}(z)=\alpha p_{n-1}^{R}(z)$ is in $\ell^{2}$ only if $\alpha=\mathbf{0}$.
Proof. Denote $s_{n}=\operatorname{Tr}\left(u_{n}(z) A_{n-1}^{\dagger} u_{n-1}(z)^{\dagger}\right)$. Here $A_{0}=1$. Multiplying (2.54) for $n \geq 2$ and (2.56) for $n=1$ by $u_{n}(z)^{\dagger}$ on the right, taking traces, and summing over $n$, we get

$$
z \sum_{n=1}^{N} \operatorname{Tr}\left(u_{n}(z) u_{n}(z)^{\dagger}\right)=\sum_{n=1}^{N} s_{n+1}+\sum_{n=1}^{N} \operatorname{Tr}\left(u_{n}(z) B_{n+1} u_{n}(z)^{\dagger}\right)+\sum_{n=1}^{N} \overline{s_{n}} .
$$

Taking imaginary parts and letting $N \rightarrow \infty$, we obtain (2.57) since the middle sum is real and the outer sums cancel up to boundary terms. Applying (2.57) to $u_{n}(z)=\alpha p_{n-1}^{R}(z)$, we get zero in the right-hand side:

$$
(\operatorname{Im} z) \sum_{n=1}^{\infty} \operatorname{Tr}\left(\alpha p_{n-1}^{R}(z) p_{n-1}^{R}(z)^{\dagger} \alpha^{\dagger}\right)=0
$$

and hence $\alpha=\mathbf{0}$ since $p_{0}^{R}=\mathbf{1}$.
Theorem 2.16. Let $\operatorname{Im} z>0$.
(i) Any solution $\left\{u_{n}(z)\right\}_{n=0}^{\infty}$ of (2.54) (for $\left.n \geq 2\right)$ can be represented as

$$
\begin{equation*}
u_{n}(z)=a p_{n-1}^{R}(z)+b q_{n-1}^{R}(z) \tag{2.58}
\end{equation*}
$$

for suitable $a, b \in \mathcal{M}_{l}$. In fact, $a=u_{1}(z)$ and $b=-u_{0}(z)$.
(ii) A sequence (2.58) satisfies (2.54) for all $n \geq 1$ if and only if $b=0$.
(iii) A sequence (2.58) belongs to $\ell^{2}$ if and only if $u_{n}(z)=c \psi_{n-1}^{R}(z)$ for some $c \in \mathcal{M}_{l}$. Equivalently, a sequence (2.58) belongs to $\ell^{2}$ if and only if $u_{1}(z)+u_{0}(z) m(z)=0$.

Proof. (i) Let $u_{n}(z)$ be a solution to (2.54). Consider

$$
\tilde{u}_{n}(z)=u_{n}(z)-u_{1}(z) p_{n-1}^{R}(z)+u_{0}(z) q_{n-1}^{R}(z) .
$$

Then $\tilde{u}_{n}(z)$ also solves (2.54) and $\tilde{u}_{0}(z)=\tilde{u}_{1}(z)=0$. It follows that $\tilde{u}_{n}(z)=0$ for all $n$. This proves (i).
(ii) A direct substitution of (2.58) into (2.54) for $n=1$ yields the statement.
(iii) We already know that $c \psi_{n-1}^{R}$ is an $\ell^{2}$ solution. Suppose that $u_{n}(z)$ is an $\ell^{2}$ solution to (2.54). Rewrite (2.58) as

$$
u_{n}(z)=(a-b m(z)) p_{n-1}^{R}(z)+b \psi_{n-1}^{R}(z) .
$$

Since $\psi_{n}^{R}$ is in $\ell^{2}$ and $c p_{n}^{R}$ is not in $\ell^{2}$, we get $a=b m(z)$, which is equivalent to $u_{1}(z)+u_{0}(z) m(z)=0$ or to $u_{n}(z)=b \psi_{n-1}^{R}(z)$.

By conjugation, we obtain:
Theorem 2.17. Let $\operatorname{Im} z>0$.
(i) Any solution $\left\{v_{n}(z)\right\}_{n=0}^{\infty}$ of (2.55) (for $\left.n \geq 2\right)$ can be represented as

$$
\begin{equation*}
v_{n}(z)=p_{n-1}^{L}(z) a+q_{n-1}^{L}(z) b . \tag{2.59}
\end{equation*}
$$

In fact, $a=v_{1}(z)$ and $b=-v_{0}(z)$.
(ii) A sequence (2.59) satisfies (2.55) for all $n \geq 1$ if and only if $b=0$.
(iii) A sequence (2.59) belongs to $\ell^{2}$ if and only if $v_{n}(z)=\psi_{n-1}^{L}(z) c$ for some $c \in \mathcal{M}_{l}$. Equivalently, a sequence (2.59) belongs to $\ell^{2}$ if and only if $v_{1}(z)+m(z) v_{0}(z)=0$.

### 2.6 Wronskians and the Christoffel-Darboux Formula

For any two $\mathcal{M}_{l}$-valued sequences $u_{n}, v_{n}$, define the Wronskian by

$$
\begin{equation*}
W_{n}(u, v)=u_{n} A_{n} v_{n+1}-u_{n+1} A_{n}^{\dagger} v_{n} . \tag{2.60}
\end{equation*}
$$

Note that $W_{n}(u, v)=-W_{n}\left(v^{\dagger}, u^{\dagger}\right)^{\dagger}$. If $u_{n}(z)$ and $v_{n}(z)$ are solutions to (2.54) and (2.55), then by a direct calculation, we see that $W_{n}(u(z), v(z))$ is independent of $n$. Put differently, if both $u_{n}(z)$ and $v_{n}(z)$ are solutions to (2.54), then $W_{n}\left(u(z), v(\bar{z})^{\dagger}\right)$ is independent of $n$. Or, if both $u_{n}(z)$ and $v_{n}(z)$ are solutions to (2.55), then $W_{n}\left(u(\bar{z})^{\dagger}, v(z)\right)$ is independent of $n$. In particular, by a direct evaluation for $n=0$, we get

$$
\begin{aligned}
& W_{n}\left(p_{\cdot-1}^{R}(z), p_{\cdot-1}^{R}(\bar{z})^{\dagger}\right)=W_{n}\left(q_{\cdot-1}^{R}(z), q_{\cdot-1}^{R}(\bar{z})^{\dagger}\right)=\mathbf{0}, \\
& W_{n}\left(p_{\cdot-1}^{L}(\bar{z})^{\dagger}, p_{\cdot-1}^{L}(z)\right)=W_{n}\left(q_{\cdot-1}^{L}(\bar{z})^{\dagger}, q_{\cdot-1}^{L}(z)\right)=\mathbf{0}, \\
& W_{n}\left(p_{\cdot{ }_{-1}}^{R}(z), q_{\cdot{ }_{-1}}^{R}(\bar{z})^{\dagger}\right)=W_{n}\left(p_{\cdot-1}^{L}(\bar{z})^{\dagger}, q_{\cdot-1}^{L}(z)\right)=\mathbf{1} .
\end{aligned}
$$

Let both $u(z)$ and $v(z)$ be solutions to (2.54) of the type (2.58), namely,

$$
u_{n}(z)=a p_{n-1}^{R}(z)+b q_{n-1}^{R}(z), \quad v_{n}(z)=c p_{n-1}^{R}(z)+d q_{n-1}^{R}(z) .
$$

Then the above calculation implies

$$
W_{n}\left(u(z), v(\bar{z})^{\dagger}\right)=a d^{\dagger}-b c^{\dagger} .
$$

Theorem 2.18 (CD Formula). For any $x, y \in \mathbb{C}$ and $n \geq 1$, one has

$$
\begin{equation*}
(x-y) \sum_{k=0}^{n} p_{k}^{R}(x) p_{k}^{L}(y)=-W_{n+1}\left(p_{\cdot-1}^{R}(x), p_{\cdot-1}^{L}(y)\right) . \tag{2.61}
\end{equation*}
$$

Proof. Multiplying (2.30) by $p_{n}^{L}(y)$ on the right and (2.32) (with $y$ in place of $x$ ) by $p_{n}^{R}(x)$ on the left and subtracting, we get

$$
(x-y) p_{n}^{R}(x) p_{n}^{R}(y)=W_{n}\left(p_{\cdot-1}^{R}(x), p_{\cdot-1}^{L}(y)\right)-W_{n+1}\left(p_{\cdot-1}^{R}(x), p_{\cdot-1}^{L}(y)\right)
$$

Summing over $n$ and noting that $W_{0}\left(p^{R}(x), p^{L}(y)\right)=0$, we get the required statement.

### 2.7 The CD Kernel

The CD kernel is defined for $z, w \in \mathbb{C}$ by

$$
\begin{align*}
K_{n}(z, w) & =\sum_{k=0}^{n} p_{k}^{R}(z) p_{k}^{R}(\bar{w})^{\dagger}  \tag{2.62}\\
& =\sum_{k=0}^{n} p_{k}^{L}(\bar{z})^{\dagger} p_{k}^{L}(w) . \tag{2.63}
\end{align*}
$$

(2.63) follows from (2.62) and (2.31). Notice that $K$ is independent of the choices $\sigma_{n}, \tau_{n}$ in (2.24) and that (2.61) can be written

$$
\begin{equation*}
(z-\bar{w}) K_{n}(z, w)=-W_{n+1}\left(p_{\cdot-1}^{R}(z), p_{\cdot-1}^{R}(\bar{w})^{\dagger}\right) \tag{2.64}
\end{equation*}
$$

The independence of $K_{n}$ of $\sigma, \tau$ can be understood by noting that if $f_{m}$ is given by (2.26), then

$$
\begin{equation*}
\int K_{n}(z, w) d \mu(w) f(w)=\sum_{m=0}^{n} p_{m}^{R}(z) f_{m} \tag{2.65}
\end{equation*}
$$

so $K$ is the kernel of the orthogonal projection onto polynomials of degree up to $n$, and so intrinsic. Similarly, if $f_{m}^{(L)}=\left\langle\left\langle f, p_{m}^{L}\right\rangle\right\rangle_{L}$, so

$$
\begin{equation*}
f(x)=\sum_{m=0}^{\infty} f_{m}^{(L)} p_{m}^{L}(x), \tag{2.66}
\end{equation*}
$$

then, by (2.63),

$$
\begin{equation*}
\int f(z) d \mu(z) K_{n}(z, w)=\sum_{m=0}^{n} f_{m}^{(L)} p_{m}^{L}(w) \tag{2.67}
\end{equation*}
$$

One has

$$
\begin{equation*}
\int K_{n}(z, w) d \mu(w) K_{n}(w, \zeta)=K_{n}(z, \zeta) \tag{2.68}
\end{equation*}
$$

as can be checked directly and which is an expression of the fact that the map in (2.65) is a projection, and so its own square.

We will let $\pi_{n}$ be the map of $L^{2}(d \mu)$ to itself given by (2.65) (or by (2.67)). (2.64) can then be viewed (for $z, w \in \mathbb{R}$ ) as an expression of the integral kernel of the commutator $\left[J, \pi_{n}\right]$, which leads to another proof of it [175].

Let $J_{n ; F}$ be the finite $n l \times n l$ matrix obtained from $J$ by taking the top leftmost $n 2$ blocks. It is the matrix of $\pi_{n-1} M_{x} \pi_{n-1}$ where $M_{x}$ is multiplication by $x$ in the $\left\{p_{j}^{R}\right\}_{j=0}^{n-1}$ basis. For $y \in \mathbb{C}$ and $\gamma \in \mathbb{C}^{l}$, let $\varphi_{n, \gamma}(y)$ be the vector whose components are

$$
\begin{equation*}
\left(\varphi_{n, \gamma}(y)\right)_{j}=p_{j-1}^{L}(y) \gamma \tag{2.69}
\end{equation*}
$$

for $j=1,2, \ldots, n$. We claim that

$$
\begin{equation*}
\left[\left(J_{n ; F}-y\right) \varphi_{n, \gamma}(y)\right]_{j}=-\delta_{j n} A_{n} p_{n}^{L}(y) \gamma \tag{2.70}
\end{equation*}
$$

as follows immediately from (2.32).
This is intimately related to (2.61) and (2.64). For recalling $J$ is the matrix in $p^{R}$ basis, $\varphi_{n, \gamma}(y)$ corresponds to the function

$$
\sum_{j=0}^{n-1} p_{j}^{R}(x)\left(\varphi_{n, \gamma}(y)\right)_{j-1}=K_{n}(x, y) \gamma
$$

As we will see in the next two sections, (2.70) has important consequences.

### 2.8 Christoffel Variational Principle

There is a matrix version of the Christoffel variational principle (see Nevai [152] for a discussion of uses in the scalar case; this matrix case is discussed by Duran-Polo [76]):

Theorem 2.19. For any non-trivial $l \times l$ matrix- valued measure, $d \mu$, on $\mathbb{R}$, we have that for any $n$, any $x_{0} \in \mathbb{R}$, and matrix polynomials $Q_{n}(x)$ of degree at most $n$ with

$$
\begin{equation*}
Q_{n}\left(x_{0}\right)=\mathbf{1}, \tag{2.71}
\end{equation*}
$$

we have that

$$
\begin{equation*}
\left\langle\left\langle Q_{n}, Q_{n}\right\rangle\right\rangle_{R} \geq K_{n}\left(x_{0}, x_{0}\right)^{-1} \tag{2.72}
\end{equation*}
$$

with equality if and only if

$$
\begin{equation*}
Q_{n}(x)=K_{n}\left(x, x_{0}\right) K_{n}\left(x_{0}, x_{0}\right)^{-1} \tag{2.73}
\end{equation*}
$$

Remark. (2.72) also holds for $\langle\langle\cdot, \cdot\rangle\rangle_{L}$ but the minimizer is then $K_{n}\left(x_{0}, x_{0}\right)^{-1} K_{n}\left(x, x_{0}\right)$.
Proof. Let $Q_{n}^{(0)}$ denote the right-hand side of (2.73). Then for any polynomial $R_{n}$ of degree at most $n$, we have

$$
\begin{equation*}
\left\langle\left\langle Q_{n}^{(0)}, R_{n}\right\rangle\right\rangle_{R}=K_{n}\left(x_{0}, x_{0}\right)^{-1} R_{n}\left(x_{0}\right) \tag{2.74}
\end{equation*}
$$

because of (2.65). Since $Q_{n}\left(x_{0}\right)=Q_{n}^{(0)}\left(x_{0}\right)=\mathbf{1}$, we conclude

$$
\begin{equation*}
\left\langle\left\langle Q_{n}-Q_{n}^{(0)}, Q_{n}-Q_{n}^{(0)}\right\rangle\right\rangle_{R}=\left\langle\left\langle Q_{n}, Q_{n}\right\rangle\right\rangle_{R}-K_{n}\left(x_{0}, x_{0}\right)^{-1} \tag{2.75}
\end{equation*}
$$

from which (2.72) is immediate and, given the supposed non- triviality, the uniqueness of minimizer.

With this, one easily gets an extension of a result of Máté-Nevai [146] to MOPRL. (They had it for scalar OPUC. For OPUC, it is in Máté-Nevai-Totik [147] on [-1, 1] and in Totik [187] for general OPRL. The result below can be proven using polynomial mappings à la Totik [188] or Jost solutions à la Simon [174].)

Theorem 2.20. Let $d \mu$ be a non-trivial $l \times l$ matrix-valued measure on $\mathbb{R}$ with compact support, $E$. Let $I=(a, b)$ be an open interval with $I \subset E$. Then for Lebesgue a.e. $x \in I$,

$$
\begin{equation*}
\lim \sup (n+1) K_{n}(x, x)^{-1} \leq w(x) . \tag{2.76}
\end{equation*}
$$

Remark. This is intended in terms of expectations in any fixed vector.
We state this explicitly since we will need it in Section 5.4, but we note that the more detailed results of Máté-Nevai-Totik [147], Lubinsky [141], Simon [173], and Totik [189] also extend.

### 2.9 Zeros

We next look at zeros of $\operatorname{det}\left(P_{n}^{L}(z)\right)$, which we will prove soon is also $\operatorname{det}\left(P_{n}^{R}(z)\right)$. Following [66, 177], we will identify these zeros with eigenvalues of $J_{n ; F}$. It is not obvious a priori that these zeros are real and, unlike the scalar situation, where the classical arguments on the zeros rely on orthogonality, we do not know how to get reality just from that (but see the remark after Theorem 2.25).

Lemma 2.21. Let $C(z)$ be an $l \times l$ matrix-valued function analytic near $z=0$. Let

$$
\begin{equation*}
k=\operatorname{dim}(\operatorname{ker}(C(0))) . \tag{2.77}
\end{equation*}
$$

Then $\operatorname{det}(C(z))$ has a zero at $z=0$ of order at least $k$.
Remarks. 1. Even in the $1 \times 1$ case, where $k=1$, the zeros can clearly be of higher order than $k$ since $c_{11}(z)$ can have a zero of any order!
2. The temptation to take products of eigenvalues will lead at best to a complicated proof as the cases $C(z)=\left(\begin{array}{ll}0 & z \\ 1 & 0\end{array}\right)$ and $C(z)=\left(\begin{array}{cc}0 & z^{2} \\ 1 & 0\end{array}\right)$ illustrate.

Proof. Let $e_{1}, \ldots, e_{l}$ be an orthonormal basis with $e_{1}, \ldots, e_{k} \in \operatorname{ker}(C(0))$. By Hadamard's inequality (see Bhatia [13]),

$$
\begin{aligned}
|\operatorname{det}(C(z))| & \leq\left\|C(z) e_{1}\right\| \cdots\left\|C(z) e_{l}\right\| \\
& \leq C|z|^{k}
\end{aligned}
$$

since $\left\|C(z) e_{j}\right\| \leq C|z|$ if $j=1, \ldots, k$ and $\left\|C(z) e_{j}\right\| \leq d$ for $j=k+1, \ldots, l$.
The following goes back at least to [34]; see also [66, 165, 177, 178].
Theorem 2.22. We have that

$$
\begin{equation*}
\operatorname{det}_{\mathbb{C}^{l}}\left(P_{n}^{L}(z)\right)=\operatorname{det}_{\mathbb{C}^{n l}}\left(z-J_{n ; F}\right) . \tag{2.78}
\end{equation*}
$$

Proof. By (2.70), if $\gamma$ is such that $p_{n}^{L}(y) \gamma=0$, then $\varphi_{n, \gamma}(y)$ is an eigenvector for $J_{n ; F}$ with eigenvalue $y$. Conversely, if $\varphi$ is an eigenvector and $\gamma$ is defined as that vector in $\mathbb{C}^{l}$ whose components are the first $l$ components of $\varphi$, then a simple inductive argument shows $\varphi=\varphi_{n, \gamma}(y)$ and then, by (2.70) and the fact that $A_{n}$ is invertible, we see that $p_{n}^{L}(y) \gamma=0$. This shows that for any $y$,

$$
\begin{equation*}
\operatorname{dim} \operatorname{ker}\left(P_{n}^{L}(y)\right)=\operatorname{dim} \operatorname{ker}\left(J_{n ; F}-y\right) \tag{2.79}
\end{equation*}
$$

By Lemma 2.21, if $y$ is any eigenvalue of $J_{n ; F}$ of multiplicity $k$, then $\operatorname{det}\left(P_{n}^{L}(z)\right)$ has a zero of order at least $k$ at $y$. Now let us consider the polynomials in $z$ on the left and right in (2.78). Since $J_{n ; F}$ is Hermitian, the sum of the multiplicities of the zeros on the right is $n l$. Since the polynomial on the left is of degree $n l$, by a counting argument it has the same set of zeros with the same multiplicities as the polynomial on the right. Since both polynomials are monic, they are equal.

Corollary 2.23. All the zeros of $\operatorname{det}\left(P_{n}^{L}(z)\right)$ are real. Moreover,

$$
\begin{equation*}
\operatorname{det}\left(P_{n}^{R}(z)\right)=\operatorname{det}\left(P_{n}^{L}(z)\right) \tag{2.80}
\end{equation*}
$$

Proof. Since $J_{n ; F}$ is Hermitian, all its eigenvalues are real, so (2.78) implies the zeros of $\operatorname{det}\left(P_{n}^{L}(z)\right)$ are real. Thus, since the polynomial is monic,

$$
\begin{equation*}
\overline{\operatorname{det}\left(P_{n}^{L}(\bar{z})\right)}=\operatorname{det}\left(P_{n}^{L}(z)\right) . \tag{2.81}
\end{equation*}
$$

By Lemma 2.2(v), we have

$$
\begin{equation*}
P_{n}^{R}(z)=P_{n}^{L}(\bar{z})^{\dagger} \tag{2.82}
\end{equation*}
$$

since both sides are analytic and agree if $z$ is real. Thus,

$$
\operatorname{det}\left(P_{n}^{R}(z)\right)=\operatorname{det}\left(P_{n}^{L}(\bar{z})^{\dagger}\right)=\overline{\operatorname{det}\left(P_{n}^{L}(\bar{z})\right)}
$$

proving (2.80).
The following appeared before in [178]; see also [165].
Corollary 2.24. Let $\left\{x_{n, j}\right\}_{j=1}^{n l}$ be the zeros of $\operatorname{det}\left(P_{n}^{L}(x)\right)$ counting multiplicity ordered by

$$
\begin{equation*}
x_{n, 1} \leq x_{n, 2} \leq \cdots \leq x_{n, n l} \tag{2.83}
\end{equation*}
$$

Then

$$
\begin{equation*}
x_{n+1, j} \leq x_{n, j} \leq x_{n+1, j+l} \tag{2.84}
\end{equation*}
$$

Remarks. 1. This is interlacing if $l=1$ and replaces it for general $l$.
2. Using $A_{n}$ invertible, one can show the inequalities in (2.84) are strict.

Proof. The min-max principle [159] says that

$$
\begin{equation*}
x_{n, j}=\max _{\substack{L \subset \mathbb{C}^{n l} \\ \operatorname{dim}(L) \leq j-1}} \min _{\substack{f \in L^{L^{\perp}} \\\|f\|=1}}\left\langle f, J_{n ; F} f\right\rangle_{\mathbb{C}^{n l}} . \tag{2.85}
\end{equation*}
$$

If $P: \mathbb{C}^{(n+1) l} \rightarrow \mathbb{C}^{n l}$ is the natural projection, then $\left\langle P f, J_{n+1 ; F} P f\right\rangle_{\mathbb{C}^{(n+1) l}}=\left\langle P f, J_{n ; F} P f\right\rangle_{\mathbb{C}^{n l l}}$ and as $L$ runs through all subspaces of $\mathbb{C}^{(n+1) l}$ dimension at most $j-1, P[L]$ runs through all subspaces of dimension at most $j-1$ in $\mathbb{C}^{n l}$, so (2.85) implies $x_{n+1, j} \leq x_{n, j}$. Using the same argument on $-J_{n ; F}$ and $-J_{n+1 ; F}$ shows $x_{j}\left(-J_{n ; F}\right) \geq x_{j}\left(-J_{n+1 ; F}\right)$. But $x_{j}\left(-J_{n ; F}\right)=-x_{n l+1-j}\left(J_{n ; F}\right)$ and $x_{j}\left(-J_{n+1 ; F}\right)=-x_{(n+1) l+1-j}\left(J_{n+1 ; F}\right)$. That yields the other inequality.

### 2.10 Lower Bounds on $p$ and the Stieltjes-Weyl Formula for $m$

Next, we want to exploit (2.70) to prove uniform lower bounds on $\left\|p_{n}^{L}(y) \gamma\right\|$ when $y \notin \operatorname{cvh}(\sigma(J))$, the convex hull of the support of $J$, and thereby uniform bounds $\left\|p_{n}^{L}(y)^{-1}\right\|$. We will then use that to prove that for $z \notin \sigma(J)$, we have

$$
\begin{equation*}
m(z)=\lim _{n \rightarrow \infty}-p_{n}^{L}(z)^{-1} q_{n}^{L}(z) \tag{2.86}
\end{equation*}
$$

the matrix analogue of a formula that spectral theorists associate with Weyl's definition of the $m$-function [193], although for the discrete case, it goes back at least to Stieltjes [181].

We begin by mimicking an argument from [171]. Let $H=\operatorname{cvh}(\sigma(J))=[c-D, c+D]$ with

$$
\begin{equation*}
D=\frac{1}{2} \operatorname{diam}(H) \tag{2.87}
\end{equation*}
$$

By the definition of $A_{n}$,

$$
\begin{equation*}
\left\|A_{n}\right\|=\left\|\left\langle\left\langle p_{n-1}^{R},(x-c) p_{n}^{R}\right\rangle\right\rangle_{R}\right\| \leq D \tag{2.88}
\end{equation*}
$$

Suppose $y \notin H$ and let

$$
\begin{equation*}
d=\operatorname{dist}(y, H) \tag{2.89}
\end{equation*}
$$

By the spectral theorem, for any vector $\varphi \in \mathcal{H}_{v}$,

$$
\begin{equation*}
\left|\langle\varphi,(J-y) \varphi\rangle_{\mathcal{H}_{v}}\right| \geq d\|\varphi\|^{2} \tag{2.90}
\end{equation*}
$$

By (2.70), with $\varphi=\varphi_{n, \gamma}(y)$,

$$
\begin{equation*}
|\langle\varphi,(J-y) \varphi\rangle| \leq\left\|A_{n}\right\|\left\|p_{n}^{L}(y) \gamma\right\|\left\|p_{n-1}^{L}(y) \gamma\right\| \tag{2.91}
\end{equation*}
$$

while

$$
\begin{equation*}
\|\varphi\|^{2}=\sum_{j=0}^{n-1}\left\|p_{j}^{L}(y) \gamma\right\|^{2} \tag{2.92}
\end{equation*}
$$

As in [171, Prop. 2.2], we get:
Theorem 2.25. If $y \notin H$, for any $\gamma$,

$$
\begin{equation*}
\left\|p_{n}^{L}(y) \gamma\right\| \geq\left(\frac{d}{D}\right)\left(1+\left(\frac{d}{D}\right)^{2}\right)^{(n-1) / 2}\|\gamma\| \tag{2.93}
\end{equation*}
$$

In particular,

$$
\begin{equation*}
\left\|p_{n}^{L}(y)^{-1}\right\| \leq \frac{D}{d} . \tag{2.94}
\end{equation*}
$$

Remark. (2.93) implies $\operatorname{det}\left(p_{n}^{L}(y)\right) \neq 0$ if $\operatorname{Im} y>0$, providing another proof that its zeros are real.
By the discussion after (2.52), if $\operatorname{Im} z>0, q_{n}^{L}(z)+p_{n}^{L}(z) m(z)$ is in $\ell^{2}$, so goes to zero. Since $p_{n}^{L}(z)^{-1}$ is bounded, we obtain:
Corollary 2.26. For any $z \in \mathbb{C}_{+}=\{z: \operatorname{Im} z>0\}$,

$$
\begin{equation*}
m(z)=\lim _{n \rightarrow \infty}-p_{n}^{L}(z)^{-1} q_{n}^{L}(z) \tag{2.95}
\end{equation*}
$$

Remark. This holds for $z \notin H$.
Taking adjoints using (2.82) and $m(z)^{\dagger}=m(\bar{z})$, we see that

$$
\begin{equation*}
m(z)=\lim _{n \rightarrow \infty}-q_{n}^{R}(z) p_{n}^{R}(z)^{-1} \tag{2.96}
\end{equation*}
$$

(2.95) and (2.96) are due to [47], which uses the proof based on the Gauss-Jacobi quadrature formula.

### 2.11 Wronskians of Vector-Valued Solutions

Let $\alpha, \beta$ be two vector-valued solutions $\left(\mathbb{C}^{l}\right)$ of $(2.55)$ for $n=2,3, \ldots$ Define their scalar Wronskian as (Euclidean inner product on $\mathbb{C}^{l}$ )

$$
\begin{equation*}
W_{n}(\alpha, \beta)=\left\langle\alpha_{n}, A_{n} \beta_{n+1}\right\rangle-\left\langle A_{n} \alpha_{n+1}, \beta_{n}\right\rangle \tag{2.97}
\end{equation*}
$$

for $n=2,3, \ldots$. One can obtain two matrix solutions by using $\alpha$ or $\beta$ for one column and 0 for the other columns. The scalar Wronskian is just a matrix element of the resulting matrix Wronskian, so $W_{n}$ is constant (as can also be seen by direct calculation). Here is an application:

Theorem 2.27. Let $z_{0} \in \mathbb{R} \backslash \sigma_{\mathrm{ess}}(J)$. For $k=0,1$, let $q_{k}$ be the multiplicity of $z_{0}$ as an eigenvalue of $J^{(k)}$. Then, $q_{0}+q_{1} \leq l$.
Proof. If $\tilde{\beta}$ is an eigenfunction for $J^{(1)}$ and we define $\beta$ by

$$
\beta_{n}= \begin{cases}0, & n=1,  \tag{2.98}\\ \tilde{\beta}_{n-1}, & n \geq 2\end{cases}
$$

then $\beta$ solves (2.55) for $n \geq 2$. If $\alpha$ is an eigenfunction for $J=J^{(0)}$, it also solves (2.55). Since $\alpha_{n} \rightarrow 0, \beta_{n} \rightarrow 0$, and $A_{n}$ is bounded, $W_{n}(\alpha, \beta) \rightarrow 0$ as $n \rightarrow \infty$ and so it is identically zero. But since $\beta_{1}=0$,

$$
\begin{equation*}
\mathbf{0}=W_{1}(\alpha, \beta)=\left\langle\alpha_{1}, A_{1} \beta_{2}\right\rangle=\left\langle\alpha_{1}, A_{1} \tilde{\beta}_{1}\right\rangle \tag{2.99}
\end{equation*}
$$

Let $V^{(k)}$ be the set of values of eigenfunctions of $J^{(k)}$ at $n=1$. (2.99) says

$$
\begin{equation*}
V^{(0)} \subset\left[A_{1} V^{(1)}\right]^{\perp} . \tag{2.100}
\end{equation*}
$$

Since $q_{k}=\operatorname{dim}\left(V^{(k)}\right)$ and $A_{1}$ is non-singular, (2.100) implies that $q_{0} \leq l-q_{1}$.

### 2.12 The Order of Zeros/Poles of $m(z)$

Theorem 2.28. Let $z_{0} \in \mathbb{R} \backslash \sigma_{\text {ess }}(J)$. For $k=0,1$, let $q_{k}$ be the multiplicity of $z_{0}$ as an eigenvalue of $J^{(k)}$. If $q_{1}-q_{0} \geq 0$, then $\operatorname{det}(m(z))$ has a zero at $z=z_{0}$ of order $q_{1}-q_{0}$. If $q_{1}-q_{0}<0$, then $\operatorname{det}(m(z))$ has a pole at $z=z_{0}$ of order $q_{0}-q_{1}$.
Remarks. 1. To say $\operatorname{det}(m(z))$ has a zero at $z=z_{0}$ of order 0 means it is finite and non-vanishing at $z_{0}$ !
2. Where $d \mu$ is a direct sum of scalar measures, so is $m(z)$, and $\operatorname{det}(m(z))$ is then a product. In the scalar case, $m(z)$ has a pole at $z_{0}$ if $J^{(0)}$ has $z_{0}$ as eigenvalue and a zero at $z_{0}$ if $J^{(1)}$ has $z_{0}$ as eigenvalue. In the direct sum case, we see there can be cancellations, which helps explain why $q_{1}-q_{0}$ occurs.
3. Formally, one can understand this theorem as follows. Cramer's rule suggests $\operatorname{det}(m(z))=$ $\operatorname{det}\left(J^{(1)}-z\right) / \operatorname{det}\left(J^{(0)}-z\right)$. Even though $\operatorname{det}\left(J^{(k)}-z\right)$ is not well-defined in the infinite case, we expect a cancellation of zeros of order $q_{1}$ and $q_{0}$. For $z_{0} \notin H$, the convex hull of $\sigma\left(J^{(0)}\right)$, one can use (2.95) to prove the theorem following this intuition. Spurious zeros in gaps of $\sigma\left(J^{(0)}\right)$ make this strategy difficult in gaps.
4. Unlike in Lemma 2.21, we can write $m$ as a product of eigenvalues and analyze that directly because $m(x)$ is self-adjoint for $x$ real, which sidesteps some of the problems associated with nontrivial Jordan normal forms.
5. This proof gives another demonstration that $q_{0}+q_{1} \leq l$.
$\operatorname{Proof.} m(z)$ has a simple pole at $z=z_{0}$ with a residue which is rank $q_{0}$ and strictly negative definite on its range. Let

$$
f(z)=\left(z-z_{0}\right) m(z) .
$$

$f$ is analytic near $z_{0}$ and self-adjoint for $z$ real and near $z_{0}$. Thus, by eigenvalue perturbation theory $[124,159], f(z)$ has $l$ eigenvalues $\rho_{1}(z), \ldots, \rho_{l}(z)$ analytic near $z_{0}$ with $\rho_{1}, \rho_{2}, \ldots, \rho_{q_{0}}$ non-zero at $z_{0}$ and $\rho_{q_{0}+1}, \ldots, \rho_{l}$ zero at $z_{0}$.

Thus, $m(z)$ has $l$ eigenvalues near $z_{0}, \lambda_{j}(z)=\rho_{j}(z) /\left(z-z_{0}\right)$, where $\lambda_{1}, \ldots, \lambda_{q_{0}}$ have simple poles and the others are regular.

By Proposition 2.14, $m(z)^{-1}$ has a simple pole at $z_{0}$ with residue of rank $q_{1}$ (because $A_{1}$ is non-singular), so $m(z)^{-1}$ has $q_{1}$ eigenvalues with poles. That means $q_{1}$ of $\lambda_{q_{0}+1}, \ldots, \lambda_{l}$ have simple zeros at $z_{0}$ and the others are non-zero. Thus, $\operatorname{det}(m(z))=\prod_{j=1}^{l} \lambda_{j}(z)$ has a pole/zero of order $q_{0}-q_{1}$.

### 2.13 Resolvent of the Jacobi Matrix

Consider the matrix

$$
G_{n m}(z)=\left\langle\left\langle p_{n-1}^{R},(x-z)^{-1} p_{m-1}^{R}\right\rangle_{R} .\right.
$$

Theorem 2.29. One has

$$
G_{n m}(z)= \begin{cases}\psi_{n-1}^{L}(z) p_{m-1}^{R}(z), & \text { if } n \geq m  \tag{2.101}\\ p_{n-1}^{L}(z) \psi_{m-1}^{R}(z), & \text { if } n \leq m\end{cases}
$$

Proof. We have

$$
\begin{array}{ll}
\sum_{m=1}^{\infty} G_{k m}(z) J_{m n}=z G_{k n}(z), & n \neq k \\
\sum_{m=1}^{\infty} J_{n m} G_{m k}(z)=z G_{n k}(z), & n \neq k
\end{array}
$$

Fix $k \geq 0$ and let $u_{m}(z)=G_{k m}(z)$. Then $u_{m}(z)$ satisfies the equation (2.54) for $n \neq k$, and so we can use Theorem 2.16 to describe this solution. First suppose $k>1$. As $u_{m}$ is an $\ell^{2}$ solution and $u_{m}$ satisfies (2.54) for $n=1$, we have

$$
G_{k m}(z)= \begin{cases}a_{k}(z) p_{m-1}^{R}(z), & m \leq k  \tag{2.102}\\ b_{k}(z) \psi_{m-1}^{R}(z), & m \geq k\end{cases}
$$

If $k=1,(2.102)$ also holds true. For $m \geq k$, this follows by the same argument, and for $m=k=1$, this is a trivial statement. Next, similarly, let us consider $v_{m}(z)=G_{m k}(z)$. Then $v_{m}(z)$ solves (2.55) and so, using Theorem 2.17, we obtain

$$
G_{m k}(z)= \begin{cases}p_{m-1}^{L}(z) c_{k}(z), & m \leq k  \tag{2.103}\\ \psi_{m-1}^{L}(z) d_{k}(z), & m \geq k\end{cases}
$$

Comparing (2.102) and (2.103), we find

$$
a_{k}(z) p_{m-1}^{R}(z)=\psi_{k-1}^{L}(z) d_{m}(z)
$$

$$
b_{k}(z) \psi_{m-1}^{R}(z)=p_{k-1}^{L}(z) c_{m}(z)
$$

As $p_{0}^{R}=p_{0}^{L}=\mathbf{1}$, it follows that

$$
a_{k}(z)=\psi_{k-1}^{L}(z) d_{1}(z), \quad c_{m}(z)=b_{1}(z) \psi_{m-1}^{R}(z)
$$

and so we obtain

$$
G_{n m}(z)= \begin{cases}\psi_{n-1}^{L}(z) d_{1}(z) p_{m-1}^{R}(z) & \text { if } n \geq m  \tag{2.104}\\ p_{n-1}^{L}(z) b_{1}(z) \psi_{m-1}^{R}(z) & \text { if } n \leq m\end{cases}
$$

It remains to prove that

$$
\begin{equation*}
b_{1}(z)=d_{1}(z)=\mathbf{1} . \tag{2.105}
\end{equation*}
$$

Consider the case $m=n=1$. By the definition of the resolvent,

$$
G_{11}(z)=\left\langle\left\langle p_{0}^{R},(J-z)^{-1} p_{0}^{R}\right\rangle\right\rangle_{R}=\int \frac{d \mu(x)}{x-z}=m(z) .
$$

On the other hand, by (2.104),

$$
\begin{aligned}
& G_{11}(z)=\psi_{0}^{L}(z) d_{1}(z) p_{0}^{R}(z)=m(z) d_{1}(z) \\
& G_{11}(z)=p_{0}^{L}(z) b_{1}(z) \psi_{0}^{R}(z)=b_{1}(z) m(z)
\end{aligned}
$$

which proves (2.105).

## 3 Matrix Orthogonal Polynomials on the Unit Circle

### 3.1 Definition of MOPUC

In this chapter, $\mu$ is an $l \times l$ matrix-valued measure on $\partial \mathbb{D} .\langle\langle\cdot, \cdot\rangle\rangle_{R}$ and $\langle\langle\cdot, \cdot\rangle\rangle_{L}$ are defined as in the MOPRL case. Non-triviality is defined as for MOPRL. We will always assume $\mu$ is non-trivial. We define monic matrix polynomials $\Phi_{n}^{R}, \Phi_{n}^{L}$ by applying Gram-Schmidt to $\{\mathbf{1}, z \mathbf{1}, \ldots\}$, that is, $\Phi_{n}^{R}$ is the unique matrix polynomial $z^{n} 1+$ lower order with

$$
\begin{equation*}
\left\langle\left\langle z^{k} \mathbf{1}, \Phi_{n}^{R}\right\rangle\right\rangle_{R}=0 \quad k=0,1, \ldots, n-1 . \tag{3.1}
\end{equation*}
$$

We will define the normalized MOPUC shortly. We will only consider the analogue of what we called type 1 for MOPRL because only those appear to be useful. Unlike in the scalar case, the monic polynomials do not appear much because it is for the normalized, but not monic, polynomials that the left and right Verblunsky coefficients are the same.

### 3.2 The Szegő Recursion

Szegő [184] included the scalar Szegő recursion for the first time. It seems likely that Geronimus had it independently shortly after Szegő. Not knowing of the connection with this work, Levinson [138] rederived the recursion but with matrix coefficients! So the results of this section go back to 1947.

For a matrix polynomial $P_{n}$ of degree $n$, we define the reversed polynomial $P_{n}^{*}$ by

$$
\begin{equation*}
P_{n}^{*}(z)=z^{n} P_{n}(1 / \bar{z})^{\dagger} . \tag{3.2}
\end{equation*}
$$

Notice that

$$
\begin{equation*}
\left(P_{n}^{*}\right)^{*}=P_{n} \tag{3.3}
\end{equation*}
$$

and for any $\alpha \in \mathcal{M}_{l}$,

$$
\begin{equation*}
\left(\alpha P_{n}\right)^{*}=P_{n}^{*} \alpha^{\dagger}, \quad\left(P_{n} \alpha\right)^{*}=\alpha^{\dagger} P_{n}^{*} . \tag{3.4}
\end{equation*}
$$

Lemma 3.1. We have

$$
\begin{equation*}
\langle\langle f, g\rangle\rangle_{L}=\langle\langle g, f\rangle\rangle_{L}^{\dagger}, \quad\langle\langle f, g\rangle\rangle_{R}=\langle\langle g, f\rangle\rangle_{R}^{\dagger} \tag{3.5}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\langle\left\langle f^{*}, g^{*}\right\rangle\right\rangle_{L}=\langle\langle f, g\rangle\rangle_{R}^{\dagger}, \quad\left\langle\left\langle f^{*}, g^{*}\right\rangle\right\rangle_{R}=\langle\langle f, g\rangle\rangle_{L}^{\dagger} . \tag{3.6}
\end{equation*}
$$

Proof. The first and second identities follow immediately from the definition. The third identity is derived as follows:

$$
\begin{aligned}
\left\langle\left\langle f^{*}, g^{*}\right\rangle\right\rangle_{L} & =\int e^{i n \theta} g(\theta)^{\dagger} d \mu(\theta)\left(e^{i n \theta} f(\theta)^{\dagger}\right)^{\dagger} \\
& =\int e^{i n \theta} g(\theta)^{\dagger} d \mu(\theta) e^{-i n \theta} f(\theta) \\
& =\int g(\theta)^{\dagger} d \mu(\theta) f(\theta) \\
& =\langle\langle g, f\rangle\rangle_{R} \\
& =\langle\langle f, g\rangle\rangle_{R}^{\dagger}
\end{aligned}
$$

The proof of the last identity is analogous.
Lemma 3.2. If $P_{n}$ has degree $n$ and is left-orthogonal with respect to $z \mathbf{1}, \ldots, z^{n} \mathbf{1}$, then $P_{n}=c\left(\Phi_{n}^{R}\right)^{*}$ for some suitable matrix c.

Proof. By assumption,

$$
\mathbf{0}=\left\langle\left\langle P_{n}, z^{j} \mathbf{1}\right\rangle\right\rangle_{L}=\left\langle\left\langle\left(z^{j} \mathbf{1}\right)^{*}, P_{n}^{*}\right\rangle_{R}=\left\langle\left\langle z^{n-j} \mathbf{1}, P_{n}^{*}\right\rangle\right\rangle_{R} \text { for } 1 \leq j \leq n .\right.
$$

Thus, $P_{n}^{*}$ is right-orthogonal with respect to $\mathbf{1}, z \mathbf{1}, \ldots, z^{n-1} \mathbf{1}$ and hence it is a right-multiple of $\Phi_{n}^{R}$. Consequently, $P_{n}$ is a left- multiple of $\left(\Phi_{n}^{R}\right)^{*}$.

Let us define normalized orthogonal matrix polynomials by

$$
\varphi_{0}^{L}=\varphi_{0}^{R}=\mathbf{1}, \quad \varphi_{n}^{L}=\kappa_{n}^{L} \Phi_{n}^{L} \quad \text { and } \quad \varphi_{n}^{R}=\Phi_{n}^{R} \kappa_{n}^{R}
$$

where the $\kappa$ 's are defined according to the normalization condition

$$
\left\langle\left\langle\varphi_{n}^{R}, \varphi_{m}^{R}\right\rangle\right\rangle_{R}=\delta_{n m} \mathbf{1} \quad\left\langle\left\langle\varphi_{n}^{L}, \varphi_{m}^{L}\right\rangle\right\rangle_{L}=\delta_{n m} \mathbf{1}
$$

along with (a type 1 condition)

$$
\begin{equation*}
\kappa_{n+1}^{L}\left(\kappa_{n}^{L}\right)^{-1}>\mathbf{0} \quad \text { and } \quad\left(\kappa_{n}^{R}\right)^{-1} \kappa_{n+1}^{R}>\mathbf{0} \tag{3.7}
\end{equation*}
$$

Notice that $\kappa_{n}^{L}$ are determined by the normalization condition up to multiplication on the left by unitaries; these unitaries can always be uniquely chosen so as to satisfy (3.7).

Now define

$$
\rho_{n}^{L}=\kappa_{n}^{L}\left(\kappa_{n+1}^{L}\right)^{-1} \quad \text { and } \quad \rho_{n}^{R}=\left(\kappa_{n+1}^{R}\right)^{-1} \kappa_{n}^{R}
$$

Notice that as inverses of positives matrices, $\rho_{n}^{L}>0$ and $\rho_{n}^{R}>0$. In particular, we have that

$$
\kappa_{n}^{L}=\left(\rho_{n-1}^{L} \ldots \rho_{0}^{L}\right)^{-1} \quad \text { and } \quad \kappa_{n}^{R}=\left(\rho_{0}^{R} \ldots \rho_{n-1}^{R}\right)^{-1}
$$

Theorem 3.3 (Szegő Recursion). (a) For suitable matrices $\alpha_{n}^{L, R}$, one has

$$
\begin{align*}
z \varphi_{n}^{L}-\rho_{n}^{L} \varphi_{n+1}^{L} & =\left(\alpha_{n}^{L}\right)^{\dagger} \varphi_{n}^{R, *}  \tag{3.8}\\
z \varphi_{n}^{R}-\varphi_{n+1}^{R} \rho_{n}^{R} & =\varphi_{n}^{L, *}\left(\alpha_{n}^{R}\right)^{\dagger} . \tag{3.9}
\end{align*}
$$

(b) The matrices $\alpha_{n}^{L}$ and $\alpha_{n}^{R}$ are equal and will henceforth be denoted by $\alpha_{n}$.
(c) $\rho_{n}^{L}=\left(\mathbf{1}-\alpha_{n}^{\dagger} \alpha_{n}\right)^{1 / 2}$ and $\rho_{n}^{R}=\left(\mathbf{1}-\alpha_{n} \alpha_{n}^{\dagger}\right)^{1 / 2}$.

Proof. (a) The matrix polynomial $z \varphi_{n}^{L}$ has leading term $z^{n+1} \kappa_{n}^{L}$. On the other hand, the matrix polynomial $\rho_{n}^{L} \varphi_{n+1}^{L}$ has leading term $z^{n+1} \rho_{n}^{L} \kappa_{n+1}^{L}$. By definition of $\rho_{n}^{L}$, these terms are equal. Consequently, $z \varphi_{n}^{L}-\rho_{n}^{L} \varphi_{n+1}^{L}$ is a matrix polynomial of degree at most $n$. Notice that it is leftorthogonal with respect to $z \mathbf{1}, \ldots, z^{n} \mathbf{1}$ since

$$
\left\langle\left\langle z \varphi_{n}^{L}-\rho_{n}^{L} \varphi_{n+1}^{L}, z^{j} \mathbf{1}\right\rangle\right\rangle_{L}=\left\langle\left\langle\varphi_{n}^{L}, z^{j-1} \mathbf{1}\right\rangle\right\rangle_{L}-\left\langle\left\langle\rho_{n}^{L} \varphi_{n+1}^{L}, z^{j} \mathbf{1}\right\rangle\right\rangle_{L}=\mathbf{0}-\mathbf{0}=\mathbf{0} .
$$

Now apply Lemma 3.2. The other claim is proved in the same way.
(b) By part (a) and identities established earlier,

$$
\begin{aligned}
\left(\alpha_{n}^{L}\right)^{\dagger} & =\mathbf{0}+\left(\alpha_{n}^{L}\right)^{\dagger} \mathbf{1} \\
& =\left\langle\left\langle\varphi_{n}^{R, *}, \rho_{n}^{L} \varphi_{n+1}^{L}\right\rangle\right\rangle_{L}+\left(\alpha_{n}^{L}\right)^{\dagger}\left\langle\left\langle\varphi_{n}^{R}, \varphi_{n}^{R}\right\rangle\right\rangle_{R} \\
& =\left\langle\left\langle\varphi_{n}^{R, *}, \rho_{n}^{L} \varphi_{n+1}^{L}\right\rangle\right\rangle_{L}+\left(\alpha_{n}^{L}\right)^{\dagger}\left\langle\left\langle\varphi_{n}^{R, *}, \varphi_{n}^{R, *}\right\rangle_{L} \quad \text { by }(3.6)\right. \\
& =\left\langle\left\langle\varphi_{n}^{R, *}, \rho_{n}^{L} \varphi_{n+1}^{L}+\left(\alpha_{n}^{L}\right)^{\dagger} \varphi_{n}^{R, *}\right\rangle\right\rangle_{L} \\
& =\left\langle\left\langle\varphi_{n}^{R, *}, z \varphi_{n}^{L}\right\rangle\right\rangle_{L} \\
& =\left\langle\left\langle z \varphi_{n}^{R}, \varphi_{n}^{L, *}\right\rangle\right\rangle_{R}^{\dagger} \quad(\text { using the }(n+1) \text {-degree } *) \\
& =\left\langle\left\langle\varphi_{n+1}^{R} \rho_{n}^{R}+\varphi_{n}^{L, *}\left(\alpha_{n}^{R}\right)^{\dagger}, \varphi_{n}^{L, *}\right\rangle\right\rangle_{R}^{\dagger} \\
& =\left\langle\left\langle\varphi_{n+1}^{R} \rho_{n}^{R}, \varphi_{n}^{L, *}\right\rangle\right\rangle_{R}^{\dagger}+\left\langle\left\langle\varphi_{n}^{L, *}\left(\alpha_{n}^{R}\right)^{\dagger}, \varphi_{n}^{L, *}\right\rangle\right\rangle_{R}^{\dagger} \\
& =\mathbf{0}+\left\langle\left\langle\varphi_{n}^{L, *}, \varphi_{n}^{L, *}\left(\alpha_{n}^{R}\right)^{\dagger}\right\rangle\right\rangle_{R} \\
& =\left\langle\left\langle\varphi_{n}^{L, *}, \varphi_{n}^{L, *}\right\rangle\right\rangle_{R}\left(\alpha_{n}^{R}\right)^{\dagger} \\
& =\left\langle\left\langle\varphi_{n}^{L}, \varphi_{n}^{L}\right\rangle\right\rangle_{L}\left(\alpha_{n}^{R}\right)^{\dagger} \\
& =\left(\alpha_{n}^{R}\right)^{\dagger} .
\end{aligned}
$$

(c) Using parts (a) and (b), we see that

$$
\begin{aligned}
\mathbf{1} & =\left\langle\left\langle z \varphi_{n}^{L}, z \varphi_{n}^{L}\right\rangle\right\rangle_{L} \\
& =\left\langle\left\langle\rho_{n}^{L} \varphi_{n+1}^{L}+\alpha_{n}^{\dagger} \varphi_{n}^{R, *}, \rho_{n}^{L} \varphi_{n+1}^{L}+\alpha_{n}^{\dagger} \varphi_{n}^{R, *}\right\rangle\right\rangle_{L}
\end{aligned}
$$

$$
\begin{aligned}
& =\rho_{n}^{L}\left\langle\left\langle\varphi_{n+1}^{L}, \varphi_{n+1}^{L}\right\rangle\right\rangle_{L}\left(\rho_{n}^{L}\right)^{\dagger}+\alpha_{n}^{\dagger}\left\langle\left\langle\varphi_{n}^{R, *}, \varphi_{n}^{R, *}\right\rangle\right\rangle_{L} \alpha_{n} \\
& =\left(\rho_{n}^{L}\right) 2+\alpha_{n}^{\dagger}\left\langle\left\langle\varphi_{n}^{R}, \varphi_{n}^{R}\right\rangle\right\rangle_{R} \alpha_{n} \\
& =\left(\rho_{n}^{L}\right) 2+\alpha_{n}^{\dagger} \alpha_{n} .
\end{aligned}
$$

A similar calculation yields the other claim.
The matrices $\alpha_{n}$ will henceforth be called the Verblunsky coefficients associated with the measure $d \mu$. Since $\rho_{n}^{L}$ is invertible, we have

$$
\begin{equation*}
\left\|\alpha_{n}\right\|<1 \tag{3.10}
\end{equation*}
$$

We will eventually see (Theorem 3.12) that any set of $\alpha_{n}$ 's obeying (3.10) occurs as the set of Verblunsky coefficients for a unique non-trivial measure.

Note that the Szegő recursion for the monic orthogonal polynomials is

$$
\begin{align*}
& z \Phi_{n}^{L}-\Phi_{n+1}^{L}=\left(\kappa_{n}^{L}\right)^{-1} \alpha_{n}^{\dagger}\left(\kappa_{n}^{R}\right)^{\dagger} \Phi_{n}^{R, *} \\
& z \Phi_{n}^{R}-\Phi_{n+1}^{R}=\Phi_{n}^{L, *}\left(\kappa_{n}^{L}\right)^{\dagger} \alpha_{n}^{\dagger}\left(\kappa_{n}^{R}\right)^{-1} \tag{3.11}
\end{align*}
$$

so the coefficients in the $L$ and $R$ equations are not equal and depend on all the $\alpha_{j}, j=1, \ldots, n$.
Let us write the Szegő recursion in matrix form, starting with left-orthogonal polynomials. By Theorem 3.3,

$$
\begin{aligned}
\varphi_{n+1}^{L} & =\left(\rho_{n}^{L}\right)^{-1}\left[z \varphi_{n}^{L}-\alpha_{n}^{\dagger} \varphi_{n}^{R, *}\right] \\
\varphi_{n+1}^{R} & =\left[z \varphi_{n}^{R}-\varphi_{n}^{L, *} \alpha_{n}^{\dagger}\right]\left(\rho_{n}^{R}\right)^{-1}
\end{aligned}
$$

which implies that

$$
\begin{align*}
\varphi_{n+1}^{L} & =z\left(\rho_{n}^{L}\right)^{-1} \varphi_{n}^{L}-\left(\rho_{n}^{L}\right)^{-1} \alpha_{n}^{\dagger} \varphi_{n}^{R, *},  \tag{3.12}\\
\varphi_{n+1}^{R, *} & =\left(\rho_{n}^{R}\right)^{-1} \varphi_{n}^{R, *}-z\left(\rho_{n}^{R}\right)^{-1} \alpha_{n} \varphi_{n}^{L} . \tag{3.13}
\end{align*}
$$

In other words,

$$
\begin{equation*}
\binom{\varphi_{n+1}^{L}}{\varphi_{n+1}^{R, *}}=A^{L}\left(\alpha_{n}, z\right)\binom{\varphi_{n}^{L}}{\varphi_{n}^{R, *}} \tag{3.14}
\end{equation*}
$$

where

$$
A^{L}(\alpha, z)=\left(\begin{array}{cc}
z\left(\rho^{L}\right)^{-1} & -\left(\rho^{L}\right)^{-1} \alpha^{\dagger} \\
-z\left(\rho^{R}\right)^{-1} \alpha & \left(\rho^{R}\right)^{-1}
\end{array}\right)
$$

and $\rho^{L}=\left(\mathbf{1}-\alpha^{\dagger} \alpha\right)^{1 / 2}, \rho^{R}=\left(\mathbf{1}-\alpha \alpha^{\dagger}\right)^{1 / 2}$. Note that, for $z \neq 0$, the inverse of $A^{L}(\alpha, z)$ is given by

$$
A^{L}(\alpha, z)^{-1}=\left(\begin{array}{cc}
z^{-1}\left(\rho^{L}\right)^{-1} & z^{-1}\left(\rho^{L}\right)^{-1} \alpha^{\dagger} \\
\left(\rho^{R}\right)^{-1} \alpha & \left(\rho^{R}\right)^{-1}
\end{array}\right)
$$

which gives rise to the inverse Szegő recursion (first emphasized in the scalar and matrix cases by Delsarte el al. [37])

$$
\begin{aligned}
\varphi_{n}^{L} & =z^{-1}\left(\rho_{n}^{L}\right)^{-1} \varphi_{n+1}^{L}+z^{-1}\left(\rho_{n}^{L}\right)^{-1} \alpha_{n}^{\dagger} \varphi_{n+1}^{R, *}, \\
\varphi_{n}^{R, *} & =\left(\rho_{n}^{R}\right)^{-1} \alpha_{n} \varphi_{n+1}^{L}+\left(\rho_{n}^{R}\right)^{-1} \varphi_{n+1}^{R, *} .
\end{aligned}
$$

For right-orthogonal polynomials, we find the following matrix formulas. By Theorem 3.3,

$$
\begin{align*}
\varphi_{n+1}^{R} & =z \varphi_{n}^{R}\left(\rho_{n}^{R}\right)^{-1}-\varphi_{n}^{L, *} \alpha_{n}^{\dagger}\left(\rho_{n}^{R}\right)^{-1},  \tag{3.15}\\
\varphi_{n+1}^{L, *} & =\varphi_{n}^{L, *}\left(\rho_{n}^{L}\right)^{-1}-z \varphi_{n}^{R} \alpha_{n}\left(\rho_{n}^{L}\right)^{-1} . \tag{3.16}
\end{align*}
$$

In other words,

$$
\left(\begin{array}{ll}
\varphi_{n+1}^{R} & \varphi_{n+1}^{L, *}
\end{array}\right)=\left(\begin{array}{ll}
\varphi_{n}^{R} & \varphi_{n}^{L, *}
\end{array}\right) A^{R}\left(\alpha_{n}, z\right)
$$

where

$$
A^{R}(\alpha, z)=\left(\begin{array}{cc}
z\left(\rho^{R}\right)^{-1} & -z \alpha\left(\rho^{L}\right)^{-1} \\
-\alpha^{\dagger}\left(\rho^{R}\right)^{-1} & \left(\rho^{L}\right)^{-1}
\end{array}\right)
$$

For $z \neq 0$, the inverse of $A^{R}(\alpha, z)$ is given by

$$
A^{R}(\alpha, z)^{-1}=\left(\begin{array}{cc}
z^{-1}\left(\rho^{R}\right)^{-1} & \left(\rho^{R}\right)^{-1} \alpha \\
z^{-1}\left(\rho^{L}\right)^{-1} \alpha^{\dagger} & \left(\rho^{L}\right)^{-1}
\end{array}\right)
$$

and hence

$$
\begin{align*}
\varphi_{n}^{R} & =z^{-1} \varphi_{n+1}^{R}\left(\rho_{n}^{R}\right)^{-1}+z^{-1} \varphi_{n+1}^{L, *}\left(\rho_{n}^{L}\right)^{-1} \alpha_{n}^{\dagger},  \tag{3.17}\\
\varphi_{n}^{L, *} & =\varphi_{n+1}^{R}\left(\rho_{n}^{R}\right)^{-1} \alpha_{n}+\varphi_{n+1}^{L, *}\left(\rho_{n}^{L}\right)^{-1} . \tag{3.18}
\end{align*}
$$

### 3.3 Second Kind Polynomials

In the scalar case, second kind polynomials go back to Geronimus [89, 91, 92]. For $n \geq 1$, let us introduce the second kind polynomials $\psi_{n}^{L, R}$ by

$$
\begin{align*}
\psi_{n}^{L}(z) & =\int \frac{e^{i \theta}+z}{e^{i \theta}-z}\left(\varphi_{n}^{L}\left(e^{i \theta}\right)-\varphi_{n}^{L}(z)\right) d \mu(\theta)  \tag{3.19}\\
\psi_{n}^{R}(z) & =\int \frac{e^{i \theta}+z}{e^{i \theta}-z} d \mu(\theta)\left(\varphi_{n}^{R}\left(e^{i \theta}\right)-\varphi_{n}^{R}(z)\right) \tag{3.20}
\end{align*}
$$

For $n=0$, let us set $\psi_{0}^{L}(z)=\psi_{0}^{R}(z)=\mathbf{1}$. For future reference, let us display the first polynomials of each series:

$$
\begin{array}{ll}
\varphi_{1}^{L}(z)=\left(\rho_{0}^{L}\right)^{-1}\left(z-\alpha_{0}^{\dagger}\right), & \varphi_{1}^{R}(z)=\left(z-\alpha_{0}^{\dagger}\right)\left(\rho_{0}^{R}\right)^{-1} \\
\psi_{1}^{L}(z)=\left(\rho_{0}^{L}\right)^{-1}\left(z+\alpha_{0}^{\dagger}\right), & \psi_{1}^{R}(z)=\left(z+\alpha_{0}^{\dagger}\right)\left(\rho_{0}^{R}\right)^{-1} . \tag{3.22}
\end{array}
$$

We will also need formulas for $\psi_{n}^{L, *}$ and $\psi_{n}^{R, *}, n \geq 1$. These formulas follow directly from the above definition and from

$$
\overline{\left(\frac{e^{i \theta}+1 / \bar{z}}{e^{i \theta}-1 / \bar{z}}\right)}=-\frac{e^{i \theta}+z}{e^{i \theta}-z} .
$$

Indeed, we have

$$
\begin{align*}
\psi_{n}^{L, *}(z) & =z^{n} \int \frac{e^{i \theta}+z}{e^{i \theta}-z} d \mu(\theta)\left(\varphi_{n}^{L}(1 / \bar{z})^{\dagger}-\varphi_{n}^{L}\left(e^{i \theta}\right)^{\dagger}\right),  \tag{3.23}\\
\psi_{n}^{R, *}(z) & =z^{n} \int \frac{e^{i \theta}+z}{e^{i \theta}-z}\left(\varphi_{n}^{R}(1 / \bar{z})^{\dagger}-\varphi_{n}^{R}\left(e^{i \theta}\right)^{\dagger}\right) d \mu(\theta) . \tag{3.24}
\end{align*}
$$

Proposition 3.4. The second kind polynomials obey the recurrence relations

$$
\begin{align*}
\psi_{n+1}^{L}(z) & =\left(\rho_{n}^{L}\right)^{-1}\left(z \psi_{n}^{L}(z)+\alpha_{n}^{\dagger} \psi_{n}^{R, *}(z)\right),  \tag{3.25}\\
\psi_{n+1}^{R, *}(z) & =\left(\rho_{n}^{R}\right)^{-1}\left(z \alpha_{n} \psi_{n}^{L}(z)+\psi_{n}^{R, *}(z)\right) \tag{3.26}
\end{align*}
$$

and

$$
\begin{align*}
& \psi_{n+1}^{R}(z)=\left(z \psi_{n}^{R}(z)+\psi_{n}^{L, *}(z) \alpha_{n}^{\dagger}\right)\left(\rho_{n}^{R}\right)^{-1}  \tag{3.27}\\
& \psi_{n+1}^{L, *}(z)=\left(\psi_{n}^{L, *}(z)+z \psi_{n}^{R}(z) \alpha_{n}\right)\left(\rho_{n}^{L}\right)^{-1} \tag{3.28}
\end{align*}
$$

for $n \geq 0$.
Proof. 1. Let us check (3.25) for $n \geq 1$. Denote the right-hand side of (3.25) by $\tilde{\psi}_{n+1}^{L}(z)$. Using the recurrence relations for $\varphi_{n}^{L}, \varphi_{n}^{R, *}$ and the definition (3.19) of $\psi_{n}^{L}$, we find

$$
\psi_{n+1}^{L}(z)-\tilde{\psi}_{n+1}^{L}(z)=\int \frac{e^{i \theta}+z}{e^{i \theta}-z} A_{n}(\theta, z) d \mu(\theta)
$$

where

$$
\begin{aligned}
A_{n}(\theta, z)= & \varphi_{n+1}^{L}\left(e^{i \theta}\right)-\varphi_{n+1}^{L}(z) \\
& \quad-\left(\rho_{n}^{L}\right)^{-1}\left[z \varphi_{n}^{L}\left(e^{i \theta}\right)-z \varphi_{n}^{L}(z)+\alpha_{n}^{\dagger} z^{n} \varphi_{n}^{R}(1 / \bar{z})^{\dagger}-\alpha_{n}^{\dagger} z^{n} \varphi_{n}^{R}\left(e^{i \theta}\right)^{\dagger}\right] \\
= & \left(\rho_{n}^{L}\right)^{-1}\left[e^{i \theta} \varphi_{n}^{L}\left(e^{i \theta}\right)-\alpha_{n}^{\dagger} \varphi_{n}^{R, *}\left(e^{i \theta}\right)-z \varphi_{n}^{L}(z)+\alpha_{n}^{\dagger} \varphi_{n}^{R, *}(z)\right. \\
& \left.-z \varphi_{n}^{L}\left(e^{i \theta}\right)+z \varphi_{n}^{L}(z)-\alpha_{n}^{\dagger} \varphi_{n}^{R, *}(z)+\alpha_{n}^{\dagger} z^{n} \varphi_{n}^{R}\left(e^{i \theta}\right)^{\dagger}\right] \\
= & \left(\rho_{n}^{L}\right)^{-1}\left[\left(e^{i \theta}-z\right) \varphi_{n}^{L}\left(e^{i \theta}\right)+\alpha_{n}^{\dagger}\left(z^{n} e^{-i n \theta}-1\right) \varphi_{n}^{R, *}\left(e^{i \theta}\right)\right] .
\end{aligned}
$$

Using the orthogonality relations

$$
\begin{equation*}
\int \varphi_{n}^{L}\left(e^{i \theta}\right) d \mu(\theta) e^{-i m \theta}=\int \varphi_{n}^{R, *}\left(e^{i \theta}\right) d \mu(\theta) e^{-i(m+1) \theta}=\mathbf{0} \tag{3.29}
\end{equation*}
$$

$m=0,1, \ldots, n-1$, and the formula

$$
\frac{e^{i n \theta}-z^{n}}{e^{i \theta}-z}=e^{i(n-1) \theta}+e^{i(n-2) \theta} z+\cdots+z^{n-1}
$$

we obtain

$$
\begin{aligned}
\rho_{n}^{L} \int \frac{e^{i \theta}+z}{e^{i \theta}-z} A_{n}(\theta, z) d \mu(\theta) & =\int\left(e^{i \theta} \varphi_{n}^{L}\left(e^{i \theta}\right)-\alpha_{n}^{\dagger} \varphi_{n}^{R, *}\left(e^{i \theta}\right)\right) d \mu(\theta) \\
& =\rho_{n}^{L} \int \varphi_{n+1}^{L}\left(e^{i \theta}\right) d \mu(\theta)=\mathbf{0}
\end{aligned}
$$

2. Let us check (3.26) for $n \geq 1$. Denote the right-hand side of (3.26) by $\tilde{\psi}_{n+1}^{R, *}(z)$. Similarly to the argument above, we find

$$
\psi_{n+1}^{R, *}(z)-\tilde{\psi}_{n+1}^{R, *}(z)=\int \frac{e^{i \theta}+z}{e^{i \theta}-z} B_{n}(\theta, z) d \mu(\theta),
$$

where

$$
\begin{aligned}
B_{n}(\theta, z)= & z^{n+1} \varphi_{n+1}^{R}(1 / \bar{z})^{\dagger}-z^{n+1} \varphi_{n+1}^{R}\left(e^{i \theta}\right)^{\dagger} \\
& \quad-\left(\rho_{n}^{R}\right)^{-1}\left[z \alpha_{n} \varphi_{n}^{L}\left(e^{i \theta}\right)-z \alpha_{n} \varphi_{n}^{L}(z)+z^{n} \varphi_{n}^{R}(1 / \bar{z})^{\dagger}-z^{n} \varphi_{n}^{R}\left(e^{i \theta}\right)^{\dagger}\right] \\
= & \left(\rho_{n}^{R}\right)^{-1}\left[\varphi_{n}^{R, *}(z)-\alpha_{n} z \varphi_{n}^{L}(z)-z^{n+1} e^{-i(n+1) \theta} \varphi_{n}^{R, *}\left(e^{i \theta}\right)+z^{n+1} e^{-i n \theta} \alpha_{n} \varphi_{n}^{L}\left(e^{i \theta}\right)\right. \\
& \left.\quad-z \alpha_{n} \varphi_{n}^{L}\left(e^{i \theta}\right)+z \alpha_{n} \varphi_{n}^{L}(z)-\varphi_{n}^{R, *}(z)+z^{n} e^{-i n \theta} \varphi_{n}^{R, *}\left(e^{i \theta}\right)\right] \\
= & \left(\rho_{n}^{R}\right)^{-1}\left[z \alpha_{n}\left(z^{n} e^{-i n \theta}-1\right) \varphi_{n}^{L}\left(e^{i \theta}\right)+z^{n} e^{-i n \theta}\left(1-z e^{-i \theta}\right) \varphi_{n}^{R, *}\left(e^{i \theta}\right)\right] .
\end{aligned}
$$

Using the orthogonality relations (3.29), we get

$$
\begin{aligned}
\rho_{n}^{R} \int \frac{e^{i \theta}+z}{e^{i \theta}-z} & B_{n}(\theta, z) d \mu(\theta)= \\
& =z^{n+1} \int\left(e^{-i(n+1) \theta} \varphi_{n}^{R, *}\left(e^{i \theta}\right)-\alpha_{n} e^{-i n \theta} \varphi_{n}^{L}\left(e^{i \theta}\right)\right) d \mu(\theta) \\
& =z^{n+1} \rho_{n}^{R} \int e^{-i(n+1) \theta} \varphi_{n+1}^{R, *}\left(e^{i \theta}\right) d \mu(\theta)=\mathbf{0}
\end{aligned}
$$

3. Relations (3.25) and (3.26) can be checked for $n=0$ by a direct substitution of (3.21) and (3.22).
4. We obtain (3.27) and (3.28) from (3.25) and (3.26) by applying the $*$-operation.

Writing the above recursion in matrix form, we get

$$
\binom{\psi_{n+1}^{L}}{\psi_{n+1}^{R, *}}=A^{L}\left(-\alpha_{n}, z\right)\binom{\psi_{n}^{L}}{\psi_{n}^{R, *}}, \quad\binom{\psi_{0}^{L}}{\psi_{0}^{R, *}}=\binom{\mathbf{1}}{\mathbf{1}}
$$

for left-orthogonal polynomials and

$$
\left(\begin{array}{ll}
\psi_{n+1}^{R} & \psi_{n+1}^{L, *}
\end{array}\right)=\left(\begin{array}{ll}
\psi_{n}^{R} & \psi_{n}^{L, *}
\end{array}\right) A^{R}\left(-\alpha_{n}, z\right), \quad\left(\begin{array}{ll}
\psi_{0}^{R} & \psi_{0}^{L, *}
\end{array}\right)=\left(\begin{array}{ll}
\mathbf{1} & \mathbf{1}
\end{array}\right) .
$$

for right-orthogonal polynomials.
Equivalently,

$$
\begin{equation*}
\binom{\psi_{n+1}^{L}}{-\psi_{n+1}^{R, *}}=A^{L}\left(\alpha_{n}, z\right)\binom{\psi_{n}^{L}}{-\psi_{n}^{R, *}}, \quad\binom{\psi_{0}^{L}}{-\psi_{0}^{R, *}}=\binom{\mathbf{1}}{-\mathbf{1}} \tag{3.30}
\end{equation*}
$$

and

$$
\left(\begin{array}{ll}
\psi_{n+1}^{R} & -\psi_{n+1}^{L, *}
\end{array}\right)=\left(\begin{array}{ll}
\psi_{n}^{R} & -\psi_{n}^{L, *}
\end{array}\right) A^{R}\left(\alpha_{n}, z\right), \quad\left(\begin{array}{ll}
\psi_{0}^{R} & -\psi_{0}^{L, *}
\end{array}\right)=\left(\begin{array}{ll}
\mathbf{1} & -\mathbf{1}
\end{array}\right) .
$$

In particular, we see that the second kind polynomials $\psi_{n}^{L, R}$ correspond to Verblunsky coefficients $\left\{-\alpha_{n}\right\}$. We have the following Wronskian-type relations:
Proposition 3.5. For $n \geq 0$ and $z \in \mathbb{C}$, we have

$$
\begin{align*}
2 z^{n} \mathbf{1} & =\varphi_{n}^{L}(z) \psi_{n}^{L, *}(z)+\psi_{n}^{L}(z) \varphi_{n}^{L, *}(z),  \tag{3.31}\\
2 z^{n} \mathbf{1} & =\psi_{n}^{R, *}(z) \varphi_{n}^{R}(z)+\varphi_{n}^{R, *}(z) \psi_{n}^{R}(z)  \tag{3.32}\\
\mathbf{0} & =\varphi_{n}^{L}(z) \psi_{n}^{R}(z)-\psi_{n}^{L}(z) \varphi_{n}^{R}(z),  \tag{3.33}\\
\mathbf{0} & =\psi_{n}^{R, *}(z) \varphi_{n}^{L, *}(z)-\varphi_{n}^{R, *}(z) \psi_{n}^{L, *}(z) . \tag{3.34}
\end{align*}
$$

Proof. We prove this by induction. The four identities clearly hold for $n=0$. Suppose (3.31)-(3.34) hold for some $n \geq 0$. Then,

$$
\begin{aligned}
\varphi_{n+1}^{L} \psi_{n+1}^{L, *}+ & \psi_{n+1}^{L} \varphi_{n+1}^{L, *}= \\
= & \left(\rho_{n}^{L}\right)^{-1}\left[\left(z \varphi_{n}^{L}-\alpha_{n}^{\dagger} \varphi_{n}^{R, *}\right)\left(\psi_{n}^{L, *}+z \psi_{n}^{R} \alpha_{n}\right)\right. \\
& \left.\quad+\left(z \psi_{n}^{L}+\alpha_{n}^{\dagger} \psi_{n}^{R, *}\right)\left(\varphi_{n}^{L, *}-z \varphi_{n}^{R} \alpha_{n}\right)\right]\left(\rho_{n}^{L}\right)^{-1} \\
= & \left(\rho_{n}^{L}\right)^{-1}\left[z\left(\varphi_{n}^{L} \psi_{n}^{L, *}+\psi_{n}^{L} \varphi_{n}^{L, *}\right)-z \alpha_{n}^{\dagger}\left(\psi_{n}^{R, *} \varphi_{n}^{R}+\varphi_{n}^{R, *} \psi_{n}^{R}\right) \alpha_{n}\right. \\
& \left.\quad+\alpha_{n}^{\dagger}\left(\psi_{n}^{R, *} \varphi_{n}^{L, *}-\varphi_{n}^{R, *} \psi_{n}^{L, *}\right)+z 2\left(\varphi_{n}^{L} \psi_{n}^{R}-\psi_{n}^{L} \varphi_{n}^{R}\right) \alpha_{n}\right]\left(\rho_{n}^{L}\right)^{-1} \\
= & \left(\rho_{n}^{L}\right)^{-1}\left[2 z^{n+1}\left(\mathbf{1}-\alpha_{n}^{\dagger} \alpha_{n}\right)\right]\left(\rho_{n}^{L}\right)^{-1}=2 z^{n+1} \mathbf{1},
\end{aligned}
$$

where we used (3.31)-(3.34) for $n$ in the third step. Thus, (3.31) holds for $n+1$.
For (3.32), we note that

$$
\begin{aligned}
\psi_{n+1}^{R, *} \varphi_{n+1}^{R}+ & \varphi_{n+1}^{R, *} \psi_{n+1}^{R}= \\
= & \left(\rho_{n}^{R}\right)^{-1}\left[\left(z \alpha_{n} \psi_{n}^{L}+\psi_{n}^{R, *}\right)\left(z \varphi_{n}^{R}-\varphi_{n}^{L, *} \alpha_{n}^{\dagger}\right)\right. \\
& \left.\quad+\left(\varphi_{n}^{R, *}-z \alpha_{n} \varphi_{n}^{L}\right)\left(z \psi_{n}^{R}+\psi_{n}^{L, *} \alpha_{n}^{\dagger}\right)\right]\left(\rho_{n}^{R}\right)^{-1} \\
= & \left(\rho_{n}^{R}\right)^{-1}\left[z\left(\psi_{n}^{R, *} \varphi_{n}^{R}+\varphi_{n}^{R, *} \psi_{n}^{R}\right)-z \alpha_{n}\left(\psi_{n}^{L} \varphi_{n}^{L, *}+\varphi_{n}^{L} \psi_{n}^{L, *}\right) \alpha_{n}^{\dagger}\right. \\
& \left.\quad+z 2 \alpha_{n}\left(\psi_{n}^{L} \varphi_{n}^{R}-\varphi_{n}^{L} \psi_{n}^{R}\right)-\left(\psi_{n}^{R, *} \varphi_{n}^{L, *}-\varphi_{n}^{R, *} \psi_{n}^{L, *}\right) \alpha_{n}^{\dagger}\right]\left(\rho_{n}^{R}\right)^{-1} \\
= & \left(\rho_{n}^{R}\right)^{-1} 2 z^{n+1}\left(\mathbf{1}-\alpha_{n} \alpha_{n}^{\dagger}\right)\left(\rho_{n}^{R}\right)^{-1}=2 z^{n+1} \mathbf{1},
\end{aligned}
$$

again using (3.31)-(3.34) for $n$ in the third step.
Next,

$$
\begin{aligned}
\varphi_{n+1}^{L} \psi_{n+1}^{R}= & \psi_{n+1}^{L} \varphi_{n+1}^{R}= \\
= & \left(\rho_{n}^{L}\right)^{-1}\left[\left(z \varphi_{n}^{L}-\alpha_{n}^{\dagger} \varphi_{n}^{R, *}\right)\left(z \psi_{n}^{R}+\psi_{n}^{L, *} \alpha_{n}^{\dagger}\right)\right. \\
& \left.\quad-\left(z \psi_{n}^{L}+\alpha_{n}^{\dagger} \psi_{n}^{R, *}\right)\left(z \varphi_{n}^{R}-\varphi_{n}^{L, *} \alpha_{n}^{\dagger}\right)\right]\left(\rho_{n}^{R}\right)^{-1} \\
= & \left(\rho_{n}^{L}\right)^{-1}\left[z 2\left(\varphi_{n}^{L} \psi_{n}^{R}-\psi_{n}^{L} \varphi_{n}^{R}\right)-\alpha_{n}^{\dagger}\left(\varphi_{n}^{R, *} \psi_{n}^{L, *}-\psi_{n}^{R, *} \varphi_{n}^{L, *}\right) \alpha_{n}^{\dagger}\right. \\
& \left.\quad-z \alpha_{n}^{\dagger}\left(\varphi_{n}^{R, *} \psi_{n}^{R}+\psi_{n}^{R, *} \varphi_{n}^{R}\right)+z\left(\varphi_{n}^{L} \psi_{n}^{L, *}+\psi_{n}^{L} \varphi_{n}^{L, *}\right) \alpha_{n}^{\dagger}\right]\left(\rho_{n}^{R}\right)^{-1}
\end{aligned}
$$

$$
=0
$$

which implies first (3.33) for $n+1$ and then, by applying the $*$-operation of order $2 n+2$, also (3.34) for $n+1$. This concludes the proof of the proposition.

### 3.4 Christoffel-Darboux Formulas

Proposition 3.6. (a) (CD)-left orthogonal

$$
\begin{aligned}
(1-\bar{\xi} z) \sum_{k=0}^{n} \varphi_{k}^{L}(\xi)^{\dagger} \varphi_{k}^{L}(z) & =\varphi_{n}^{R, *}(\xi)^{\dagger} \varphi_{n}^{R, *}(z)-\bar{\xi} z \varphi_{n}^{L}(\xi)^{\dagger} \varphi_{n}^{L}(z) \\
& =\varphi_{n+1}^{R, *}(\xi)^{\dagger} \varphi_{n+1}^{R, *}(z)-\varphi_{n+1}^{L}(\xi)^{\dagger} \varphi_{n+1}^{L}(z)
\end{aligned}
$$

(b) (CD)-right orthogonal

$$
\begin{aligned}
(1-\bar{\xi} z) \sum_{k=0}^{n} \varphi_{k}^{R}(z) \varphi_{k}^{R}(\xi)^{\dagger} & =\varphi_{n}^{L, *}(z) \varphi_{n}^{L, *}(\xi)^{\dagger}-\bar{\xi} z \varphi_{n}^{R}(z) \varphi_{n}^{R}(\xi)^{\dagger} \\
& =\varphi_{n+1}^{L, *}(z) \varphi_{n+1}^{L, *}(\xi)^{\dagger}-\varphi_{n+1}^{R}(z) \varphi_{n+1}^{R}(\xi)^{\dagger}
\end{aligned}
$$

(c) (Mixed CD)-left orthogonal

$$
\begin{aligned}
(1-\bar{\xi} z) \sum_{k=0}^{n} \psi_{k}^{L}(\xi)^{\dagger} \varphi_{k}^{L}(z) & =2 \cdot \mathbf{1}-\psi_{n}^{R, *}(\xi)^{\dagger} \varphi_{n}^{R, *}(z)-\bar{\xi} z \psi_{n}^{L}(\xi)^{\dagger} \varphi_{n}^{L}(z) \\
& =2 \cdot \mathbf{1}-\psi_{n+1}^{R, *}(\xi)^{\dagger} \varphi_{n+1}^{R, *}(z)-\psi_{n+1}^{L}(\xi)^{\dagger} \varphi_{n+1}^{L}(z) .
\end{aligned}
$$

(d) (Mixed CD)-right orthogonal

$$
\begin{aligned}
(1-\bar{\xi} z) \sum_{k=0}^{n} \varphi_{k}^{R}(z) \psi_{k}^{R}(\xi)^{\dagger} & =2 \cdot \mathbf{1}-\varphi_{n}^{L, *}(z) \psi_{n}^{L, *}(\xi)^{\dagger}-\bar{\xi} z \varphi_{n}^{R}(z) \psi_{n}^{R}(\xi)^{\dagger} \\
& =2 \cdot \mathbf{1}-\varphi_{n+1}^{L, *}(z) \psi_{n+1}^{L, *}(\xi)^{\dagger}-\varphi_{n+1}^{R}(z) \psi_{n+1}^{R}(\xi)^{\dagger}
\end{aligned}
$$

Remark. Since the $\psi$ 's are themselves MOPUCs, the analogue of (a) and (b), with all $\varphi$ 's replaced by $\psi$ 's, holds.

Proof. (a) Write

$$
F_{n}^{L}(z)=\binom{\varphi_{n}^{L}(z)}{\varphi_{n}^{R, *}(z)}, \quad J=\left(\begin{array}{rr}
\mathbf{1} & \mathbf{0} \\
\mathbf{0} & -\mathbf{1}
\end{array}\right), \quad \tilde{J}=\left(\begin{array}{cc}
\bar{\xi} z \mathbf{1} & \mathbf{0} \\
\mathbf{0} & -\mathbf{1}
\end{array}\right) .
$$

Then,

$$
F_{n+1}^{L}(z)=A^{L}\left(\alpha_{n}, z\right) F_{n}^{L}(z)
$$

and

$$
\begin{aligned}
A^{L}(\alpha, \xi)^{\dagger} & J A^{L}(\alpha, z)= \\
& =\left(\begin{array}{cc}
\bar{\xi}\left(\rho^{L}\right)^{-1} & -\bar{\xi} \alpha^{\dagger}\left(\rho^{R}\right)^{-1} \\
-\alpha\left(\rho^{L}\right)^{-1} & \left(\rho^{R}\right)^{-1}
\end{array}\right)\left(\begin{array}{cc}
z\left(\rho^{L}\right)^{-1} & -\left(\rho^{L}\right)^{-1} \alpha^{\dagger} \\
z\left(\rho^{R}\right)^{-1} \alpha & -\left(\rho^{R}\right)^{-1}
\end{array}\right) \\
& =\left(\begin{array}{cc}
\bar{\xi} z\left(\rho^{L}\right)^{-2}-\bar{\xi} z \alpha^{\dagger}\left(\rho^{R}\right)^{-2} \alpha & -\bar{\xi}\left(\rho^{L}\right)^{-2} \alpha^{\dagger}+\bar{\xi} \alpha^{\dagger}\left(\rho^{R}\right)^{-2} \\
-z \alpha\left(\rho^{L}\right)^{-2}+z\left(\rho^{R}\right)^{-2} \alpha & \alpha\left(\rho^{L}\right)^{-2} \alpha^{\dagger}-\left(\rho^{R}\right)^{-2}
\end{array}\right) \\
& =\left(\begin{array}{cc}
\bar{\xi} z \mathbf{0} & \mathbf{0} \\
\mathbf{0} & -\mathbf{1}
\end{array}\right)=\tilde{J} .
\end{aligned}
$$

Thus,

$$
\begin{aligned}
F_{n+1}^{L}(\xi)^{\dagger} J F_{n+1}^{L}(z) & =F_{n}^{L}(\xi)^{\dagger} A^{L}\left(\alpha_{n}, \xi\right)^{\dagger} J A^{L}\left(\alpha_{n}, z\right) F_{n}^{L}(z) \\
& =F_{n}^{L}(\xi)^{\dagger} \tilde{J} F_{n}^{L}(z)
\end{aligned}
$$

and hence

$$
\varphi_{n+1}^{L}(\xi)^{\dagger} \varphi_{n+1}^{L}(z)-\varphi_{n+1}^{R, *}(\xi)^{\dagger} \varphi_{n+1}^{R, *}(z)=\bar{\xi} z \varphi_{n}^{L}(\xi)^{\dagger} \varphi_{n}^{L}(z)-\varphi_{n}^{R, *}(\xi)^{\dagger} \varphi_{n}^{R, *}(z)
$$

which shows that the last two expressions in (a) are equal. Denote their common value by $Q_{n}^{L}(z, \xi)$. Then,

$$
\begin{aligned}
Q_{n}^{L}(z, \xi)-Q_{n-1}^{L}(z, \xi)= & \varphi_{n}^{R, *}(\xi)^{\dagger} \varphi_{n}^{R, *}(z)-\bar{\xi} z \varphi_{n}^{L}(\xi)^{\dagger} \varphi_{n}^{L}(z) \\
& -\varphi_{n}^{R, *}(\xi)^{\dagger} \varphi_{n}^{R, *}(z)+\varphi_{n}^{L}(\xi)^{\dagger} \varphi_{n}^{L}(z) \\
= & (1-\bar{\xi} z) \varphi_{n}^{L}(\xi)^{\dagger} \varphi_{n}^{L}(z) .
\end{aligned}
$$

Summing over $n$ completes the proof since $Q_{-1}^{L}(z, \xi)=0$.
(b) The proof is analogous to (a): Write $F_{n}^{R}(z)=\left(\begin{array}{ll}\varphi_{n}^{R}(z) \quad \varphi_{n}^{L, *}(z)\end{array}\right)$. Then, $F_{n+1}^{R}(z)=$ $F_{n}^{R}(z) A^{R}\left(\alpha_{n}, z\right)$ and $A^{R}(\alpha, z) J A^{R}(\alpha, \xi)^{\dagger}=\tilde{J}$. Thus,

$$
\begin{aligned}
F_{n+1}^{R}(z) J F_{n+1}^{R}(\xi)^{\dagger} & =F_{n}^{R}(z) A^{R}\left(\alpha_{n}, z\right) J A^{R}\left(\alpha_{n}, \xi\right)^{\dagger} F_{n}^{R}(\xi)^{\dagger} \\
& =F_{n}^{R}(z) \tilde{J} F_{n}^{R}(\xi)^{\dagger}
\end{aligned}
$$

and hence

$$
\varphi_{n+1}^{R}(z) \varphi_{n+1}^{R}(\xi)^{\dagger}-\varphi_{n+1}^{L, *}(z) \varphi_{n+1}^{L, *}(\xi)^{\dagger}=\bar{\xi} z \varphi_{n}^{R}(z) \varphi_{n}^{R}(\xi)^{\dagger}-\varphi_{n}^{L, *}(z) \varphi_{n}^{L, *}(\xi)^{\dagger}
$$

which shows that the last two expressions in (b) are equal. Denote their common value by $Q_{n}^{R}(z, \xi)$. Then,

$$
Q_{n}^{R}(z, \xi)-Q_{n-1}^{R}(z, \xi)=(1-\bar{\xi} z) \varphi_{n}^{R}(z) \varphi_{n}^{R}(\xi)^{\dagger}
$$

and the assertion follows as before.
(c) Write

$$
\tilde{F}_{n}^{L}(z)=\binom{\psi_{n}^{L}(z)}{-\psi_{n}^{R, *}(z)}
$$

with the second kind polynomials $\psi_{n}^{L, R}$. As in (a), we see that

$$
\begin{aligned}
\tilde{F}_{n+1}^{L}(\xi)^{\dagger} J F_{n+1}^{L}(z) & =\tilde{F}_{n}^{L}(\xi)^{\dagger} A^{L}\left(\alpha_{n}, \xi\right)^{\dagger} J A^{L}\left(\alpha_{n}, z\right) F_{n}^{L}(z) \\
& =\tilde{F}_{n}^{L}(\xi)^{\dagger} \tilde{J} F_{n}^{L}(z)
\end{aligned}
$$

and hence

$$
\psi_{n+1}^{L}(\xi)^{\dagger} \varphi_{n+1}^{L}(z)+\psi_{n+1}^{R, *}(\xi)^{\dagger} \varphi_{n+1}^{R, *}(z)=\bar{\xi} z \psi_{n}^{L}(\xi)^{\dagger} \varphi_{n}^{L}(z)+\psi_{n}^{R, *}(\xi)^{\dagger} \varphi_{n}^{R, *}(z)
$$

Denote

$$
\tilde{Q}_{n}^{L}(z, \xi)=2 \cdot \mathbf{1}-\psi_{n+1}^{R, *}(\xi)^{\dagger} \varphi_{n+1}^{R, *}(z)-\psi_{n+1}^{L}(\xi)^{\dagger} \varphi_{n+1}^{L}(z) .
$$

Then,

$$
\begin{aligned}
\tilde{Q}_{n}^{L}(z, \xi)-\tilde{Q}_{n-1}^{L}(z, \xi)= & -\psi_{n}^{R, *}(\xi)^{\dagger} \varphi_{n}^{R, *}(z)-\bar{\xi} z \psi_{n}^{L}(\xi)^{\dagger} \varphi_{n}^{L}(z) \\
& +\psi_{n}^{R, *}(\xi)^{\dagger} \varphi_{n}^{R, *}(z)+\psi_{n}^{L}(\xi)^{\dagger} \varphi_{n}^{L}(z) \\
= & (1-\bar{\xi} z) \psi_{n}^{L}(\xi)^{\dagger} \varphi_{n}^{L}(z)
\end{aligned}
$$

and the assertion follows as before.
(d) Write $\tilde{F}_{n}^{R}(z)=\left(\begin{array}{ll}\psi_{n}^{R}(z) & \left.-\psi_{n}^{L, *}(z)\right) \text {. As in (b), we see that }\end{array}\right.$

$$
\begin{aligned}
F_{n+1}^{R}(z) J \tilde{F}_{n+1}^{R}(\xi)^{\dagger} & =F_{n}^{R}(z) A^{R}\left(\alpha_{n}, z\right) J A^{R}\left(\alpha_{n}, \xi\right)^{\dagger} \tilde{F}_{n}^{R}(\xi)^{\dagger} \\
& =F_{n}^{R}(z) \tilde{J} \tilde{F}_{n}^{R}(\xi)^{\dagger}
\end{aligned}
$$

and hence

$$
\varphi_{n+1}^{R}(z) \psi_{n+1}^{R}(\xi)^{\dagger}+\varphi_{n+1}^{L, *}(z) \psi_{n+1}^{L, *}(\xi)^{\dagger}=\bar{\xi} z \varphi_{n}^{R}(z) \psi_{n}^{R}(\xi)^{\dagger}+\varphi_{n}^{L, *}(z) \psi_{n}^{L, *}(\xi)^{\dagger} .
$$

With $\tilde{Q}_{n}^{R}(z, \xi)=2 \cdot \mathbf{1}-\varphi_{n+1}^{L, *}(z) \psi_{n+1}^{L, *}(\xi)^{\dagger}-\varphi_{n+1}^{R}(z) \psi_{n+1}^{R}(\xi)^{\dagger}$, we have

$$
\tilde{Q}_{n}^{R}(z, \xi)-\tilde{Q}_{n-1}^{R}(z, \xi)=(1-\bar{\xi} z) \varphi_{n}^{R}(z) \psi_{n}^{R}(\xi)^{\dagger}
$$

and we conclude as in (c).

### 3.5 Zeros of MOPUC

Our main result in this section is:
Theorem 3.7. All the zeros of $\operatorname{det}\left(\varphi_{n}^{R}(z)\right)$ lie in $\mathbb{D}=\{z:|z|<1\}$.
We will also prove:
Theorem 3.8. For each n,

$$
\begin{equation*}
\operatorname{det}\left(\varphi_{n}^{R}(z)\right)=\operatorname{det}\left(\varphi_{n}^{L}(z)\right) \tag{3.35}
\end{equation*}
$$

The scalar analogue of Theorem 3.7 has seven proofs in [167]! The simplest is due to Landau [135] and its MOPUC analogue is Theorem 2.13 .7 of [167]. There is also a proof in Delsarte et al. [37] who attribute the theorem to Whittle [194]. We provide two more proofs here, not only for their intrinsic interest: our first proof we need because it depends only on the recursion relation (it is related to the proof of Delsarte et al. [37]). The second proof is here since it relates zeros to eigenvalues of a cutoff CMV matrix.

Theorem 3.9. We have
(i) For $z \in \partial \mathbb{D}$, all of $\varphi_{n}^{R, *}(z), \varphi_{n}^{L, *}(z), \varphi_{n}^{R}(z), \varphi_{n}^{L}(z)$ are invertible.
(ii) For $z \in \partial \mathbb{D}, \varphi_{n}^{L}(z)\left(\varphi_{n}^{R, *}(z)\right)^{-1}$ and $\left(\varphi_{n}^{*, L}(z)\right)^{-1} \varphi_{n}^{R}(z)$ are unitary.
(iii) For $z \in \mathbb{D}, \varphi_{n}^{R, *}(z)$ and $\varphi_{n}^{L, *}(z)$ are invertible.
(iv) For $z \in \mathbb{D}, \varphi_{n}^{L}(z)\left(\varphi_{n}^{R, *}(z)\right)^{-1}$ and $\left(\varphi_{n}^{*, L}(z)\right)^{-1} \varphi_{n}^{R}(z)$ are of norm at most 1 and, for $n \geq 1$, strictly less than 1 .
(v) All zeros of $\operatorname{det}\left(\varphi_{n}^{R, *}(z)\right)$ and $\operatorname{det}\left(\varphi_{n}^{L, *}(z)\right)$ lie in $\mathbb{C} \backslash \overline{\mathbb{D}}$.
(vi) All zeros of $\operatorname{det}\left(\varphi_{n}^{R}(z)\right)$ and $\operatorname{det}\left(\varphi_{n}^{L}(z)\right)$ lie in $\mathbb{D}$.

Remark. (vi) is our first proof of Theorem 3.7.
Proof. All these results are trivial for $n=0$, so we can hope to use an inductive argument. So suppose we have the result for $n-1$.

By (3.13),

$$
\begin{equation*}
\varphi_{n}^{R, *}=\left(\rho_{n-1}^{R}\right)^{-1}\left(\mathbf{1}-z \alpha_{n-1} \varphi_{n-1}^{L}\left(\varphi_{n-1}^{R, *}\right)^{-1}\right) \varphi_{n-1}^{R, *} . \tag{3.36}
\end{equation*}
$$

Since $\left|\alpha_{n-1}\right|<1$, if $|z| \leq 1$, each factor on the right of (3.36) is invertible. This proves (i) and (iii) for $\varphi_{n}^{R, *}$ and a similar argument works for $\varphi_{n}^{L, *}$. If $z=e^{i \theta}, \varphi_{n}^{R}\left(e^{i \theta}\right)=e^{i n \theta} \varphi_{n}^{R, *}\left(e^{i \theta}\right)^{\dagger}$ is also invertible, so we have (i) and (iii) for $n$.

Next, we claim that if $z \in \partial \mathbb{D}$, then

$$
\begin{equation*}
\varphi_{n}^{R, *}(z)^{\dagger} \varphi_{n}^{R, *}(z)=\varphi_{n}^{L}(z)^{\dagger} \varphi_{n}^{L}(z) . \tag{3.37}
\end{equation*}
$$

This follows from taking $z=\xi \in \partial \mathbb{D}$ in Proposition 3.6(a). Given that $\varphi_{n}^{R, *}(z)$ is invertible, this implies

$$
\begin{equation*}
\mathbf{1}=\left(\varphi_{n}^{L}(z) \varphi_{n}^{R, *}(z)^{-1}\right)^{\dagger}\left(\varphi_{n}^{L}(z) \varphi_{n}^{R, *}(z)^{-1}\right) \tag{3.38}
\end{equation*}
$$

proving the first part of (ii) for $n$. The second part of (ii) is proven similarly by using Proposition 3.6(b).

For $z \in \overline{\mathbb{D}}$, let

$$
F(z)=\varphi_{n}^{L}(z) \varphi_{n}^{R, *}(z)^{-1} .
$$

Then $F$ is analytic in $\mathbb{D}$, continuous in $\overline{\mathbb{D}}$, and $\|F(z)\|=1$ on $\partial \mathbb{D}$, so (iv) follows from the maximum principle.

Since $\varphi_{n}^{R, *}(z)$ is invertible on $\overline{\mathbb{D}}$, its det is non-zero there, proving (v). (vi) then follows from

$$
\begin{equation*}
\operatorname{det}\left(\varphi_{n}^{R}(z)\right)=z^{n l} \overline{\operatorname{det}\left(\varphi_{n}^{R, *}(1 / \bar{z})\right)} \tag{3.39}
\end{equation*}
$$

Let $\mathcal{V}$ be the $\mathbb{C}^{l}$-valued functions on $\partial \mathbb{D}$ and $\mathcal{V}_{n}$ the span of the $\mathbb{C}^{l}$-valued polynomials of degree at most $n$, so

$$
\operatorname{dim}\left(\mathcal{V}_{n}\right)=\mathbb{C}^{l(n+1)}
$$

Let $\mathcal{V}_{\infty}$ be the set $\cup_{n} \mathcal{V}_{n}$ of all $\mathbb{C}^{l}$-valued polynomials. Let $\pi_{n}$ be the projection onto $\mathcal{V}_{n}$ in the $\mathcal{V}$ inner product (1.35).

It is easy to see that

$$
\begin{equation*}
\mathcal{V}_{n} \cap \mathcal{V}_{n-1}^{\perp}=\left\{\Phi_{n}^{R}(z) v: v \in \mathbb{C}^{l}\right\} \tag{3.40}
\end{equation*}
$$

since $\left\langle z^{l}, \Phi_{n}^{R}(z) v\right\rangle=0$ for $l=0, \ldots, n-1$ and the dimensions on the left and right of (3.40) coincide. $\mathcal{V}_{n} \cap \mathcal{V}_{n-1}^{\perp}$ can also be described as the set of $\left(v^{\dagger} \Phi_{n}^{L}(z)\right)^{\dagger}$ for $v \in \mathbb{C}^{l}$.

We define $M_{z}: \mathcal{V}_{n-1} \rightarrow \mathcal{V}_{n}$ or $\mathcal{V}_{\infty} \rightarrow \mathcal{V}_{\infty}$ as the operator of multiplication by $z$.
Theorem 3.10. For all $n$, we have

$$
\begin{equation*}
\operatorname{det}_{\mathbb{C}^{l}}\left(\Phi_{n}^{R}(z)\right)=\operatorname{det}_{\mathcal{V}_{n-1}}\left(z \mathbf{1}-\pi_{n-1} M_{z} \pi_{n-1}\right) . \tag{3.41}
\end{equation*}
$$

Remarks. 1. Since $\left\|M_{z}\right\| \leq 1$, (3.41) immediately implies zeros of $\operatorname{det}\left(\varphi_{n}^{R}(z)\right)$ lie in $\overline{\mathbb{D}}$, and a small additional argument proves they lie in $\mathbb{D}$. As we will see, this also implies Theorem 3.8.
2. Of course, $\pi_{n-1} M_{z} \pi_{n-1}$ is a cutoff CMV matrix if written in a CMV basis.

Proof. If $Q \in \mathcal{V}_{n-k}$, then by (3.40),

$$
\begin{equation*}
\pi_{n-1}\left[\left(z-z_{0}\right)^{k} Q\right]=0 \Leftrightarrow\left(z-z_{0}\right)^{k} Q=\Phi_{n}^{R}(z) v \tag{3.42}
\end{equation*}
$$

for some $v \in \mathbb{C}^{l}$. Thus writing $\operatorname{det}\left(\Phi_{n}^{R}(z)\right)=\Phi_{n}^{R}(z) v_{1} \wedge \cdots \wedge \Phi_{n}^{R}(z) v_{l}$ in a Jordan basis for $\Phi_{n}^{R}(z)$, we see that the order of the zeros of $\operatorname{det}\left(\Phi_{n}^{R}(z)\right)$ at $z_{0}$ is exactly the order of $z_{0}$ as an algebraic eigenvalue of $\pi_{n-1} M_{z} \pi_{n-1}$, that is, the order of $z_{0}$ as a zero of the right side of (3.41).

Since both sides of (3.41) are monic polynomials of degree $n l$ and their zeros including multiplicity are the same, we have proven (3.41).

Proof of Theorem 3.8. On the right side of (3.42), we can put $\left(\Phi_{n}^{L}(z) v^{\dagger}\right)^{\dagger}$ and so conclude (3.41) holds with $\Phi_{n}^{L}(z)$ on the left.

This proves (3.35) if $\varphi$ is replaced by $\Phi$. Since $\alpha_{j}^{*} \alpha_{j}$ and $\alpha_{j} \alpha_{j}^{*}$ are unitarily equivalent, $\operatorname{det}\left(\rho_{j}^{L}\right)=$ $\operatorname{det}\left(\rho_{j}^{R}\right)$. Thus, $\operatorname{det}\left(\kappa_{n}^{L}\right)=\operatorname{det}\left(\kappa_{n}^{R}\right)$, and we obtain (3.35) for $\varphi$.

It is a basic fact (Theorem 1.7.5 of [167]) that for the scalar case, any set of $n$ zeros in $\mathbb{D}$ are the zeros of a unique OPUC $\Phi_{n}$ and any monic polynomial with all its zeros in $\mathbb{D}$ is a monic OPUC. It is easy to see that any set of $n l$ zeros in $\mathbb{D}$ is the set of zeros of an OPUC $\Phi_{n}$, but certainly not unique. It is an interesting open question to clarify what matrix monic OPs are monic MOPUCs.

### 3.6 Bernstein-Szegő Approximation

Given $\left\{\alpha_{j}\right\}_{j=0}^{n-1} \in \mathbb{D}^{n}$, we use Szegő recursion to define polynomials $\varphi_{j}^{R}, \varphi_{j}^{L}$ for $j=0,1, \ldots, n$. We define a measure $d \mu_{n}$ on $\partial \mathbb{D}$ by

$$
\begin{equation*}
d \mu_{n}(\theta)=\left[\varphi_{n}^{R}\left(e^{i \theta}\right) \varphi_{n}^{R}\left(e^{i \theta}\right)^{\dagger}\right]^{-1} \frac{d \theta}{2 \pi} . \tag{3.43}
\end{equation*}
$$

Notice that (3.37) can be rewritten

$$
\begin{equation*}
\varphi_{n}^{R}\left(e^{i \theta}\right) \varphi_{n}^{R}\left(e^{i \theta}\right)^{\dagger}=\varphi_{n}^{L}\left(e^{i \theta}\right)^{\dagger} \varphi_{n}^{L}\left(e^{i \theta}\right) \tag{3.44}
\end{equation*}
$$

We use here and below the fact that the proof of Theorem 3.9 only depends on Szegő recursion and not on the a priori existence of a measure. That theorem also shows the inverse in (3.43) exists. Thus,

$$
\begin{equation*}
d \mu_{n}(\theta)=\left[\varphi_{n}^{L}\left(e^{i \theta}\right)^{\dagger} \varphi_{n}^{L}\left(e^{i \theta}\right)\right]^{-1} \frac{d \theta}{2 \pi} \tag{3.45}
\end{equation*}
$$

Theorem 3.11. The measure $d \mu_{n}$ is normalized (i.e., $\mu_{n}(\partial \mathbb{D})=\mathbf{1}$ ) and its right MOPUC for $j=0, \ldots, n$ are $\left\{\varphi_{j}^{R}\right\}_{j=0}^{n}$, and for $j>n$,

$$
\begin{equation*}
\varphi_{j}^{R}(z)=z^{j-n} \varphi_{n}^{R}(z) \tag{3.46}
\end{equation*}
$$

The Verblunsky coefficients for $d \mu_{n}$ are

$$
\alpha_{j}\left(d \mu_{n}\right)= \begin{cases}\alpha_{j}, & j \leq n,  \tag{3.47}\\ \mathbf{0}, & j \geq n+1 .\end{cases}
$$

Remarks. 1. In the scalar case, one can multiply by a constant and renormalize, and then prove the constant is 1 . Because of commutativity issues, we need a different argument here.
2. Of course, using (3.45), $\varphi_{n}^{L}$ are left MOPUC for $d \mu_{n}$.
3. Our proof owes something to the scalar proof in [80].

Proof. Let $\langle\langle\cdot \cdot \cdot\rangle\rangle_{R}$ be the inner product associated with $\mu_{n}$. By a direct computation, $\left\langle\left\langle\varphi_{n}^{R}, \varphi_{n}^{R}\right\rangle\right\rangle_{R}=$ 1, and for $j=0,1, \ldots, n-1$,

$$
\begin{aligned}
\left\langle\left\langle z^{j}, \varphi_{n}^{R}\right\rangle\right\rangle_{R} & =\frac{1}{2 \pi} \int_{0}^{2 \pi} e^{-i j \theta}\left(\varphi_{n}^{R}\left(e^{i \theta}\right)^{\dagger}\right)^{-1} d \theta \\
& =\frac{1}{2 \pi i} \oint z^{n-j-1}\left(\varphi_{n}^{R, *}(z)\right)^{-1} d z=\mathbf{0}
\end{aligned}
$$

by analyticity of $\varphi_{n}^{R, *}(z)^{-1}$ in $\mathbb{D}$ (continuity in $\overline{\mathbb{D}}$ ).
This proves $\varphi_{n}^{R}$ is a MOPUC for $d \mu_{n}$ (and a similar calculation works for the right side of (3.46) if $j \geq n$ ). By the inverse Szegő recursion and induction downwards, $\left\{\varphi_{j}^{R}\right\}_{j=0}^{n-1}$ are also OPs, and by the Szegő recursion, they are normalized. In particular, since $\varphi_{0}^{R} \equiv \mathbf{1}$ is normalized, $d \mu_{n}$ is normalized.

### 3.7 Verblunsky's Theorem

We can now prove the analogue of Favard's theorem for MOPUC; the scalar version is called Verblunsky's theorem in [167] after [192]. A history and other proofs can be found in [167]. The proof below is essentially the matrix analogue of that of Geronimus [90] (rediscovered in [37, 80]). Delsarte et al. [37] presented their proof in the MOPUC case and they seem to have been the first with a matrix Verblunsky theorem. One can extend the CMV and the Geronimus theorem proofs from the scalar case to get alternate proofs of the theorem below.

Theorem 3.12 (Verblunsky's Theorem for MOPUC). Any sequence $\left\{\alpha_{j}\right\}_{j=0}^{\infty} \in \mathbb{D}^{\infty}$ is the sequence of Verblunsky coefficients of a unique measure.

Proof. Uniqueness is easy, since the $\alpha$ 's determine the $\varphi_{j}^{R}$ 's and so the $\Phi_{j}^{R}$,s which determine the moments.

Given a sequence $\left\{\alpha_{j}\right\}_{j=0}^{\infty}$, let $d \mu_{n}$ be the measures of the last section. By compactness of $l \times l$ matrix-valued probability measures on $\partial \mathbb{D}$, they have a weak limit. By using limits, $\left\{\varphi_{j}^{R}\right\}_{j=0}^{\infty}$ are the right MOPUC for $d \mu$ and they determine the proper Verblunsky coefficients.

### 3.8 Matrix POPUC

Analogously to the scalar case (see $[15,100,119,172,197]$ ), given any unitary $\beta$ in $\mathcal{M}_{l}$, we define

$$
\begin{equation*}
\varphi_{n}^{R}(z ; \beta)=z \varphi_{n-1}^{R}(z)-\varphi_{n-1}^{L, *}(z) \beta^{\dagger} . \tag{3.48}
\end{equation*}
$$

As in the scalar case, this is related to the secular determinant of unitary extensions of the cutoff CMV matrix. Moreover,

Theorem 3.13. Fix $\beta$. All the zeros of $\left(\varphi_{n}(z ; \beta)\right)$ lie on $\partial \mathbb{D}$.
Proof. If $|z|<1, \varphi_{n-1}^{L, *}(z)$ is invertible and

$$
\varphi_{n}^{R}(z ; \beta)=-\varphi_{n-1}^{L, *}(z) \beta^{\dagger}\left(\mathbf{1}-z \beta \varphi_{n-1}^{L, *}(z)^{-1} \varphi_{n-1}^{R}(z)\right)
$$

is invertible since the last factor differs from 1 by a strict contraction. A similar argument shows invertibility if $|z|>1$. Thus, the only zeros of $\operatorname{det}(\cdot)$ lie in $\partial \mathbb{D}$.

### 3.9 Matrix-Valued Carathéodory and Schur Functions

An analytic matrix-valued function $F$ defined on $\mathbb{D}$ is called a (matrix-valued) Carathéodory function if $F(0)=\mathbf{1}$ and $\operatorname{Re} F(z) \equiv \frac{1}{2}\left(F(z)+F(z)^{\dagger}\right) \geq 0$ for every $z \in \mathbb{D}$. The following result can be found in [44, Thm. 2.2.2].

Theorem 3.14 (Riesz-Herglotz). If $F$ is a matrix-valued Carathéodory function, then there exists a unique positive semi-definite matrix measure $d \mu$ such that

$$
\begin{equation*}
F(z)=\int \frac{e^{i \theta}+z}{e^{i \theta}-z} d \mu(\theta) \tag{3.49}
\end{equation*}
$$

The measure $d \mu$ is given by the unique weak limit of the measures $d \mu_{r}(\theta)=\operatorname{Re} F\left(r e^{i \theta}\right) \frac{d \theta}{2 \pi}$ as $r \uparrow 1$. Moreover,

$$
F(z)=c_{0}+2 \sum_{n=1}^{\infty} c_{n} z^{n}
$$

where

$$
c_{n}=\int e^{-i n \theta} d \mu(\theta)
$$

Conversely, if $d \mu$ is a positive semi-definite matrix measure, then (3.49) defines a matrix-valued Carathéodory function.

An analytic matrix-valued function $f$ defined on $\mathbb{D}$ is called a (matrix-valued) Schur function if $f(z)^{\dagger} f(z) \leq \mathbf{1}$ for every $z \in \mathbb{D}$. This condition is equivalent to $f(z) f(z)^{\dagger} \leq \mathbf{1}$ for every $z \in \mathbb{D}$ and to $\|f(z)\| \leq 1$ for every $z \in \mathbb{D}$. By the maximum principle, if $f$ is not constant, the inequalities are strict. The following can be found in [167, Prop. 4.5.3]:
Proposition 3.15. The association

$$
\begin{align*}
f(z) & =z^{-1}(F(z)-\mathbf{1})(F(z)+\mathbf{1})^{-1}  \tag{3.50}\\
F(z) & =(\mathbf{1}+z f(z))(\mathbf{1}-z f(z))^{-1} \tag{3.51}
\end{align*}
$$

sets up a one-one correspondence between matrix-valued Carathéodory functions and matrix-valued Schur functions.
Proposition 3.16. For $z \in \mathbb{D}$, we have

$$
\begin{equation*}
\operatorname{Re} F(z)=\left(\mathbf{1}-\bar{z} f(z)^{\dagger}\right)^{-1}\left(\mathbf{1}-|z|^{2} f(z)^{\dagger} f(z)\right)(\mathbf{1}-z f(z))^{-1} \tag{3.52}
\end{equation*}
$$

and the non-tangential boundary values $\operatorname{Re} F\left(e^{i \theta}\right)$ and $f\left(e^{i \theta}\right)$ exist for Lebesgue almost every $\theta$.
Write $d \mu(\theta)=w(\theta) \frac{d \theta}{2 \pi}+d \mu_{\mathrm{s}}$. Then, for almost every $\theta$,

$$
\begin{equation*}
w(\theta)=\operatorname{Re} F\left(e^{i \theta}\right) \tag{3.53}
\end{equation*}
$$

and for a.e. $\theta$, $\operatorname{det}(w(\theta)) \neq 0$ if and only if $f\left(e^{i \theta}\right)^{\dagger} f\left(e^{i \theta}\right)<\mathbf{1}$.
Proof. The identity (3.52) follows from (3.51). The existence of the boundary values of $f$ follows by application of the scalar result to the individual entries of $f$. Then (3.51) gives the boundary values of $F$. We also used the following fact: Away from a set of zero Lebesgue measure, $\operatorname{det}(\mathbf{1}-z f(z))$ has non-zero boundary values by general $H^{\infty}$ theory.
(3.53) holds for $\langle\eta, F(z) \eta\rangle_{\mathbb{C}^{l}}$ and $\langle\eta, d \mu \eta\rangle_{\mathbb{C}^{l}}$ for any $\eta \in \mathbb{C}^{l}$ by the scalar result. We get (3.53) by polarization. From

$$
w(\theta)=\left(\mathbf{1}-e^{-i \theta} f\left(e^{i \theta}\right)^{\dagger}\right)^{-1}\left(\mathbf{1}-f\left(e^{i \theta}\right)^{\dagger} f\left(e^{i \theta}\right)\right)\left(\mathbf{1}-e^{i \theta} f\left(e^{i \theta}\right)\right)^{-1}
$$

it follows immediately that $f\left(e^{i \theta}\right)^{\dagger} f\left(e^{i \theta}\right)<\mathbf{1}$ implies $\operatorname{det}(w(\theta))>0$. Conversely, if $f\left(e^{i \theta}\right)^{\dagger} f\left(e^{i \theta}\right) \leq \mathbf{1}$ but not $f\left(e^{i \theta}\right)^{\dagger} f\left(e^{i \theta}\right)<1$, then $\operatorname{det}\left(\mathbf{1}-f\left(e^{i \theta}\right)^{\dagger} f\left(e^{i \theta}\right)\right)=0$ and by our earlier arguments $\operatorname{det}(\mathbf{1}-$ $\left.e^{-i \theta} f\left(e^{i \theta}\right)^{\dagger}\right)^{-1}$ and $\operatorname{det}\left(\mathbf{1}-e^{i \theta} f\left(e^{i \theta}\right)\right)^{-1}$ exist and are finite; hence $\operatorname{det}(w(\theta))=0$. All previous statements are true away from suitable sets of zero Lebesgue measure.

### 3.10 Coefficient Stripping, the Schur Algorithm, and Geronimus' Theorem

The matrix version of Geronimus' theorem goes back at least to the book of Bakonyi-Constantinescu [6]. Let $F(z)$ be the matrix-valued Carathéodory function (3.49) (with the same measure $\mu$ as the one used in the definition of $\left.\langle\langle\cdot, \cdot\rangle\rangle_{R}\right)$. Let us denote

$$
\begin{aligned}
u_{n}^{L}(z) & =\psi_{n}^{L}(z)+\varphi_{n}^{L}(z) F(z), \\
u_{n}^{R}(z) & =\psi_{n}^{R}(z)+F(z) \varphi_{n}^{R}(z) .
\end{aligned}
$$

We also define

$$
\begin{aligned}
u_{n}^{L, *}(z) & =\psi_{n}^{L, *}(z)-F(z) \varphi_{n}^{L, *}(z) \\
u_{n}^{R, *}(z) & =\psi_{n}^{R, *}(z)-\varphi_{n}^{R, *}(z) F(z)
\end{aligned}
$$

Proposition 3.17. For any $|z|<1$, the sequences $u_{n}^{L}(z), u_{n}^{R}(z), u_{n}^{L, *}(z), u_{n}^{R, *}(z)$ are square summable.

Proof. Denote

$$
f(\theta)=\frac{e^{-i \theta}+\bar{z}}{e^{-i \theta}-\bar{z}}, \quad g(\theta)=\frac{e^{i \theta}+z}{e^{i \theta}-z} .
$$

By the definitions (3.19)-(3.24), we have

$$
\begin{aligned}
u_{n}^{L}(z) & =\left\langle\left\langle f, \varphi_{n}^{L}\right\rangle\right\rangle_{L}, \\
-u_{n}^{R, *}(z) & =z^{n}\left\langle\left\langle\varphi_{n}^{R}, g\right\rangle\right\rangle_{R}, \\
-u_{n}^{L, *}(z) & =z^{n}\left\langle\left\langle\varphi_{n}^{L}, g\right\rangle_{L},\right. \\
u_{n}^{R}(z) & =\left\langle\left\langle f, \varphi_{n}^{R}\right\rangle\right\rangle_{R} .
\end{aligned}
$$

Using the Bessel inequality and the fact that $|z|<1$, we obtain the required statements.
Next we will consider sequences defined by

$$
\begin{equation*}
\binom{s_{n}}{t_{n}}=A^{L}\left(\alpha_{n-1}, z\right) \cdots A^{L}\left(\alpha_{0}, z\right)\binom{s_{0}}{t_{0}} \tag{3.54}
\end{equation*}
$$

where $s_{n}, t_{n} \in \mathcal{M}_{l}$. Similarly, we will consider the sequences

$$
\begin{equation*}
\left(s_{n}, t_{n}\right)=\left(s_{0}, t_{0}\right) A^{R}\left(\alpha_{0}, z\right) \cdots A^{R}\left(\alpha_{n-1}, z\right) \tag{3.55}
\end{equation*}
$$

Theorem 3.18. Let $z \in \mathbb{D}$ and let $f$ be the Schur function associated with $d \mu$ via (3.49) and (3.50). Then:
(i) A solution of (3.54) is square summable if and only if the initial condition is of the form

$$
\binom{s_{0}}{t_{0}}=\binom{c}{z f(z) c}
$$

for some matrix c in $\mathcal{M}_{l}$.
(ii) A solution of (3.55) is square summable if and only if the initial condition is of the form

$$
\left(s_{0}, t_{0}\right)=(c, c z f(z))
$$

for some matrix c.
Proof. We shall prove (i); the proof of (ii) is similar.

1. By Proposition 3.17 and (3.14), (3.30), we have the square summable solution

$$
\binom{s_{n}}{t_{n}}=\left(\begin{array}{ll}
u_{n}^{L}(z) & -u_{n}^{R, *}(z) \tag{3.56}
\end{array}\right), \quad\binom{s_{0}}{t_{0}}=\binom{d}{z f(z) d}, \quad d=(F(z)+\mathbf{1}) .
$$

The matrix $d=F(z)+\mathbf{1}$ is invertible. Thus, multiplying the above solution on the right by $d^{-1} c$ for any given matrix $c$, we get the "if" part of the theorem.
2. Let us check that $\varphi_{n}^{R, *} c$ is not square summable for any matrix $c \neq 0$. By the CD formula with $\xi=z$, we have

$$
\left(1-|z|^{2}\right) \sum_{k=1}^{n-1} \varphi_{k}^{L}(z)^{\dagger} \varphi_{k}^{L}(z)+\varphi_{n}^{L}(z)^{\dagger} \varphi_{n}^{L}(z)=\varphi_{n}^{R, *}(z)^{\dagger} \varphi_{n}^{R, *}(z)
$$

and so

$$
\varphi_{n}^{R, *}(z)^{\dagger} \varphi_{n}^{R, *}(z) \geq\left(1-|z|^{2}\right) \varphi_{0}^{L}(z)^{\dagger} \varphi_{0}^{L}(z)=\left(1-|z|^{2}\right) \mathbf{1} .
$$

Thus, we get

$$
\left\|\varphi_{n}^{R, *} c\right\|_{R} 2 \geq\left(1-|z|^{2}\right) \operatorname{Tr} c^{\dagger} c>0
$$

3. Let $\binom{s_{n}}{t_{n}}$ be any square summable solution to (3.54). Let us write this solution as

$$
\begin{equation*}
\binom{s_{n}}{t_{n}}=\binom{\varphi_{n}^{L} a}{\varphi_{n}^{R, *} a}+\binom{\psi_{n}^{L} b}{-\psi_{n}^{R, *} b}, \quad a=\frac{s_{0}+t_{0}}{2}, \quad b=\frac{s_{0}-t_{0}}{2} . \tag{3.57}
\end{equation*}
$$

Multiplying the solution (3.56) by $b$ and subtracting from (3.57), we get a square summable solution

$$
\binom{\varphi_{n}^{L}(z)(a-F(z) b)}{\varphi_{n}^{R, *}(a-F(z) b)} .
$$

It follows that $a=F(z) b$, which proves the "only if" part.
The Schur function $f$ is naturally associated with $d \mu$ and hence with the Verblunsky coefficients $\alpha_{0}, \alpha_{1}, \alpha_{2}, \ldots$. The Schur functions obtained by coefficient stripping will be denoted by $f_{1}, f_{2}, f_{3}, \ldots$, that is, $f_{n}$ corresponds to Verblunsky coefficients $\alpha_{n}, \alpha_{n+1}, \alpha_{n+2}, \ldots$. We also write $f_{0} \equiv f$.

Theorem 3.19 (Schur Algorithm and Geronimus' Theorem). For the Schur functions $f_{0}, f_{1}, f_{2}, \ldots$ associated with Verblunsky coefficients $\alpha_{0}, \alpha_{1}, \alpha_{2}, \ldots$, the following relations hold:

$$
\begin{align*}
f_{n+1}(z) & =z^{-1}\left(\rho_{n}^{R}\right)^{-1}\left[f_{n}(z)-\alpha_{n}\right]\left[\mathbf{1}-\alpha_{n}^{\dagger} f_{n}(z)\right]^{-1} \rho_{n}^{L}  \tag{3.58}\\
f_{n}(z) & =\left(\rho_{n}^{R}\right)^{-1}\left[z f_{n+1}(z)+\alpha_{n}\right]\left[\mathbf{1}+z \alpha_{n}^{\dagger} f_{n+1}(z)\right]^{-1} \rho_{n}^{L} . \tag{3.59}
\end{align*}
$$

Remarks. 1. See (1.81) and Theorem 1.4 to understand this result.
2. (1.83) provides an alternate way to write (3.58) and (3.59).

Proof. It clearly suffices to consider the case $n=0$. Consider the solution of (3.54) with initial condition

$$
\binom{\mathbf{1}}{z f_{0}(z)} .
$$

By Theorem 3.18, there exists a matrix $c$ such that

$$
\begin{aligned}
\binom{c}{z f_{1}(z) c} & =A^{L}\left(\alpha_{0}, z\right)\binom{\mathbf{1}}{z f_{0}(z)} \\
& =\binom{z\left(\rho_{0}^{L}\right)^{-1}-z\left(\rho_{0}^{L}\right)^{-1} \alpha_{0}^{\dagger} f_{0}(z)}{-z\left(\rho_{0}^{R}\right)^{-1} \alpha_{0}+z\left(\rho_{0}^{R}\right)^{-1} f_{0}(z)}
\end{aligned}
$$

From this we can compute $z f_{1}(z)$ :

$$
\begin{aligned}
z f_{1}(z) & =\left[-z\left(\rho_{0}^{R}\right)^{-1} \alpha_{0}+z\left(\rho_{0}^{R}\right)^{-1} f_{0}(z)\right]\left[z\left(\rho_{0}^{L}\right)^{-1}-z\left(\rho_{0}^{L}\right)^{-1} \alpha_{0}^{\dagger} f_{0}(z)\right]^{-1} \\
& =\left(\rho_{0}^{R}\right)^{-1}\left[f_{0}(z)-\alpha_{0}\right]\left[\mathbf{1}-\alpha_{0}^{\dagger} f_{0}(z)\right]^{-1} \rho_{0}^{L}
\end{aligned}
$$

which is (3.58).
Similarly, we can express $f_{0}$ in terms of $f_{1}$. From

$$
\begin{aligned}
\binom{\mathbf{1}}{z f_{0}(z)} & =A^{L}\left(\alpha_{0}, z\right)^{-1}\binom{c}{z f_{1}(z) c} \\
& =\left(\begin{array}{cc}
z^{-1}\left(\rho_{0}^{L}\right)^{-1} & z^{-1}\left(\rho_{0}^{L}\right)^{-1} \alpha_{0}^{\dagger} \\
\left(\rho_{0}^{R}\right)^{-1} \alpha_{0} & \left(\rho_{0}^{R}\right)^{-1}
\end{array}\right)\binom{c}{z f_{1}(z) c} \\
& =\binom{z^{-1}\left(\rho_{0}^{L}\right)^{-1} c+\left(\rho_{0}^{L}\right)^{-1} \alpha_{0}^{\dagger} f_{1}(z) c}{\left(\rho_{0}^{R}\right)^{-1} \alpha_{0} c+\left(\rho_{0}^{R}\right)^{-1} z f_{1}(z) c}
\end{aligned}
$$

we find that

$$
\begin{aligned}
z f_{0}(z) & =\left[\left(\rho_{0}^{R}\right)^{-1} \alpha_{0} c+\left(\rho_{0}^{R}\right)^{-1} z f_{1}(z) c\right]\left[z^{-1}\left(\rho_{0}^{L}\right)^{-1} c+\left(\rho_{0}^{L}\right)^{-1} \alpha_{0}^{\dagger} f_{1}(z) c\right]^{-1} \\
& =\left[\left(\rho_{0}^{R}\right)^{-1} \alpha_{0}+\left(\rho_{0}^{R}\right)^{-1} z f_{1}(z)\right]\left[z^{-1}\left(\rho_{0}^{L}\right)^{-1}+\left(\rho_{0}^{L}\right)^{-1} \alpha_{0}^{\dagger} f_{1}(z)\right]^{-1} \\
& =\left(\rho_{0}^{R}\right)^{-1}\left[\alpha_{0}+z f_{1}(z)\right]\left[z^{-1} \mathbf{1}+\alpha_{0}^{\dagger} f_{1}(z)\right]^{-1} \rho_{0}^{L}
\end{aligned}
$$

which gives (3.59).

### 3.11 The CMV Matrix

In this section and the next, we discuss CMV matrices for MOPUC. This was discussed first by Simon in [170], which also has the involved history in the scalar case. Most of the results in this section appear already in [170]; the results of the next section are new here - they parallel the discussion in [167, Sect. 4.4] where these results first appeared in the scalar case.

### 3.11.1 The CMV basis

Consider the two sequences $\chi_{n}, x_{n} \in \mathcal{H}$, defined by

$$
\begin{array}{ll}
\chi_{2 k}(z)=z^{-k} \varphi_{2 k}^{L, *}(z), & \chi_{2 k-1}(z)=z^{-k+1} \varphi_{2 k-1}^{R}(z) \\
x_{2 k}(z)=z^{-k} \varphi_{2 k}^{R}(z), & x_{2 k-1}(z)=z^{-k} \varphi_{2 k-1}^{L, *}(z)
\end{array}
$$

For an integer $k \geq 0$, let us introduce the following notation: $i_{k}$ is the $(k+1)$ th term of the sequence $0,1,-1,2,-2,3,-3, \ldots$, and $j_{k}$ is the $(k+1)$ th term of the sequence $0,-1,1,-2,2,-3,3, \ldots$ Thus, for example, $i_{1}=1, j_{1}=-1$.

We use the right module structure of $\mathcal{H}$. For a set of functions $\left\{f_{k}(z)\right\}_{k=0}^{n} \subset \mathcal{H}$, its module span is the set of all sums $\sum f_{k}(z) a_{k}$ with $a_{k} \in \mathcal{M}_{l}$.

Proposition 3.20. (i) For any $n \geq 1$, the module span of $\left\{\chi_{k}\right\}_{k=0}^{n}$ coincides with the module span of $\left\{z^{i_{k}}\right\}_{k=0}^{n}$ and the module span of $\left\{x_{k}\right\}_{k=0}^{n}$ coincides with the module span of $\left\{z^{j_{k}}\right\}_{k=0}^{n}$. (ii) The sequences $\left\{\chi_{k}\right\}_{k=0}^{\infty}$ and $\left\{x_{k}\right\}_{k=0}^{\infty}$ are orthonormal:

$$
\begin{equation*}
\left\langle\left\langle\chi_{k}, \chi_{m}\right\rangle\right\rangle_{R}=\left\langle\left\langle x_{k}, x_{m}\right\rangle\right\rangle_{R}=\delta_{k m} . \tag{3.60}
\end{equation*}
$$

Proof. (i) Recall that

$$
\begin{aligned}
\varphi_{n}^{R}(z) & =\kappa_{n}^{R} z^{n}+\text { linear combination of }\left\{\mathbf{1}, \ldots, z^{n-1}\right\}, \\
\varphi_{n}^{L, *}(z) & =\left(\kappa_{n}^{L}\right)^{\dagger}+\text { linear combination of }\left\{z, \ldots, z^{n}\right\}
\end{aligned}
$$

where both $\kappa_{n}^{R}$ and $\left(\kappa_{n}^{L}\right)^{\dagger}$ are invertible matrices. It follows that

$$
\begin{aligned}
& \chi_{n}(z)=\gamma_{n} z^{i_{n}}+\text { linear combination of }\left\{z^{i_{0}}, \ldots, z^{i_{n-1}}\right\}, \\
& x_{n}(z)=\delta_{n} z^{j_{n}}+\text { linear combination of }\left\{z^{j_{0}}, \ldots, z^{j_{n-1}}\right\},
\end{aligned}
$$

where $\gamma_{n}, \delta_{n}$ are invertible matrices. This proves (i).
(ii) By the definition of $\varphi_{n}^{L}$ and $\varphi_{n}^{R}$, we have

$$
\begin{gather*}
\left\langle\left\langle\varphi_{n}^{R}, \varphi_{m}^{R}\right\rangle\right\rangle_{R}=\left\langle\left\langle\varphi_{n}^{L, *}, \varphi_{m}^{L, *}\right\rangle\right\rangle_{R}=\delta_{n m},  \tag{3.61}\\
\left\langle\left\langle\varphi_{n}^{R}, z^{m}\right\rangle\right\rangle_{R}=\mathbf{0}, m=0, \ldots, n-1 ; \quad\left\langle\left\langle\varphi_{n}^{L, *}, z^{m}\right\rangle\right\rangle_{R}=\mathbf{0}, m=1, \ldots, n . \tag{3.62}
\end{gather*}
$$

From (3.61) with $n=m$, we get

$$
\left\langle\left\langle\chi_{n}, \chi_{n}\right\rangle\right\rangle_{R}=\left\langle\left\langle x_{n}, x_{n}\right\rangle\right\rangle_{R}=\mathbf{1} .
$$

Considering separately the cases of even and odd $n$, it is easy to prove that

$$
\begin{array}{rr}
\left\langle\left\langle\chi_{n}, z^{m}\right\rangle\right\rangle_{R}=\mathbf{0}, & m=i_{0}, i_{1}, \ldots, i_{n-1} \\
\left\langle\left\langle x_{n}, z^{m}\right\rangle\right\rangle_{R}=\mathbf{0}, & m=j_{0}, j_{1}, \ldots, j_{n-1} . \tag{3.64}
\end{array}
$$

For example, for $n=2 k, m \in\left\{i_{0}, \ldots, i_{2 k-1}\right\}$ we have $m+k \in\{1,2, \ldots, 2 k\}$ and so, by (3.62),

$$
\left\langle\left\langle\chi_{n}, z^{m}\right\rangle\right\rangle_{R}=\left\langle\left\langle z^{-k} \varphi_{2 k}^{L, *}, z^{m}\right\rangle\right\rangle_{R}=\left\langle\left\langle\varphi_{2 k}^{L, *}, z^{m+k}\right\rangle\right\rangle_{R}=\mathbf{0} .
$$

The other three cases are considered similarly. From (3.63), (3.64), and (i), we get (3.60) for $k \neq m$.

From the above proposition and the $\|\cdot\|_{\infty}$-density of Laurent polynomials in $C(\partial \mathbb{D})$, it follows that $\left\{\chi_{k}\right\}_{k=0}^{\infty}$ and $\left\{x_{k}\right\}_{k=0}^{\infty}$ are right orthonormal modula bases in $\mathcal{H}$, that is, any element $f \in \mathcal{H}$ can be represented as a

$$
\begin{equation*}
f=\sum_{k=0}^{\infty} \chi_{k}\left\langle\left\langle\chi_{k}, f\right\rangle\right\rangle_{R}=\sum_{k=0}^{\infty} x_{k}\left\langle\left\langle x_{k}, f\right\rangle\right\rangle_{R} . \tag{3.65}
\end{equation*}
$$

### 3.11.2 The CMV matrix

Consider the matrix of the right homomorphism $f(z) \mapsto z f(z)$ with respect to the basis $\left\{\chi_{k}\right\}$. Denote $\mathcal{C}_{n m}=\left\langle\left\langle\chi_{n}, z \chi_{m}\right\rangle\right\rangle_{R}$. The matrix $\mathcal{C}$ is unitary in the following sense:

$$
\sum_{k=0}^{\infty} \mathcal{C}_{k n}^{\dagger} \mathcal{C}_{k m}=\sum_{k=0}^{\infty} \mathcal{C}_{n k} \mathcal{C}_{m k}^{\dagger}=\delta_{n m} \mathbf{1}
$$

The proof follows from (3.65):

$$
\begin{aligned}
& \delta_{n m} \mathbf{1}=\left\langle\left\langle z \chi_{n}, z \chi_{m}\right\rangle\right\rangle_{R}=\left\langle\left\langle\sum_{k=0}^{\infty} \chi_{k}\left\langle\left\langle\chi_{k}, z \chi_{n}\right\rangle\right\rangle_{R}, z \chi_{m}\right\rangle\right\rangle_{R}=\sum_{k=0}^{\infty} \mathcal{C}_{k n}^{\dagger} \mathcal{C}_{k m}, \\
& \delta_{n m} \mathbf{1}=\left\langle\left\langle\bar{z} \chi_{n}, \bar{z} \chi_{m}\right\rangle\right\rangle_{R}=\left\langle\left\langle\sum_{k=0}^{\infty} \chi_{k}\left\langle\left\langle\chi_{k}, \bar{z} \chi_{n}\right\rangle\right\rangle_{R}, \bar{z} \chi_{m}\right\rangle\right\rangle_{R}=\sum_{k=0}^{\infty} \mathcal{C}_{n k} \mathcal{C}_{m k}^{\dagger} .
\end{aligned}
$$

We note an immediate consequence:
Lemma 3.21. Let $|z| \leq 1$. Then, for every $m \geq 0$,

$$
\sum_{n=0}^{\infty} \chi_{n}(z) \mathcal{C}_{n m}=z \chi_{m}(z), \quad \sum_{n=0}^{\infty} \mathcal{C}_{m n} \chi_{n}(1 / \bar{z})^{\dagger}=z \chi_{m}(1 / \bar{z})^{\dagger} .
$$

Proof. First note that the above series contains only finitely many non-zero terms. Expanding $f(z)=z \chi_{n}$ according to (3.65), we see that

$$
z \chi_{n}(z)=\sum_{k=0}^{\infty} \chi_{k}(z)\left\langle\left\langle\chi_{k}, z \chi_{n}\right\rangle\right\rangle_{R}=\sum_{k=0}^{\infty} \chi_{k}(z) \mathcal{C}_{k n}
$$

which is the first identity. Next, taking adjoints, we get

$$
\bar{z} \chi_{n}(z)^{\dagger}=\sum_{k=0}^{\infty} \mathcal{C}_{k n}^{\dagger} \chi_{k}(z)^{\dagger}
$$

which yields

$$
\begin{aligned}
\bar{z} \sum_{n=0}^{\infty} \mathcal{C}_{m n} \chi_{n}(z)^{\dagger} & =\sum_{n=0}^{\infty} \mathcal{C}_{m n} \sum_{k=0}^{\infty} \mathcal{C}_{k n}^{\dagger} \chi_{k}(z)^{\dagger} \\
& =\sum_{k=0}^{\infty}\left(\sum_{n=0}^{\infty} \mathcal{C}_{m n} \mathcal{C}_{k n}^{\dagger}\right) \chi_{k}(z)^{\dagger} \\
& =\sum_{k=0}^{\infty} \delta_{m k} \chi_{k}(z)^{\dagger}=\chi_{m}(z)^{\dagger}
\end{aligned}
$$

Replacing $z$ by $1 / \bar{z}$, we get the required statement.

### 3.11.3 The $\mathcal{L} \mathcal{M}$-representation

Using (3.65) for $f=\chi_{m}$, we obtain:

$$
\begin{equation*}
\mathcal{C}_{n m}=\left\langle\left\langle\chi_{n}, z \chi_{m}\right\rangle\right\rangle_{R}=\sum_{k=0}^{\infty}\left\langle\left\langle\chi_{n}, z x_{k}\right\rangle\right\rangle_{R}\left\langle\left\langle x_{k}, \chi_{m}\right\rangle\right\rangle_{R}=\sum_{k=0}^{\infty} \mathcal{L}_{n k} \mathcal{M}_{k m} . \tag{3.66}
\end{equation*}
$$

Denote by $\Theta(\alpha)$ the $2 l \times 2 l$ unitary matrix

$$
\Theta(\alpha)=\left(\begin{array}{cc}
\alpha^{\dagger} & \rho^{L} \\
\rho^{R} & -\alpha
\end{array}\right)
$$

Using the Szegő recursion formulas (3.15) and (3.18), we get

$$
\begin{align*}
z \varphi_{n}^{R} & =\varphi_{n+1}^{R} \rho_{n}^{R}+\varphi_{n}^{L, *} \alpha_{n}^{\dagger},  \tag{3.67}\\
\varphi_{n+1}^{L, *} & =\varphi_{n}^{L, *} \rho_{n}^{L}-\varphi_{n+1}^{R} \alpha_{n} . \tag{3.68}
\end{align*}
$$

Taking $n=2 k$ and multiplying by $z^{-k}$, we get

$$
\begin{aligned}
z x_{2 k} & =\chi_{2 k} \alpha_{2 k}^{\dagger}+\chi_{2 k+1} \rho_{2 k}^{R}, \\
z x_{2 k+1} & =\chi_{2 k} \rho_{2 k}^{L}-\chi_{2 k+1} \alpha_{2 k} .
\end{aligned}
$$

It follows that the matrix $\mathcal{L}$ has the structure

$$
\mathcal{L}=\Theta\left(\alpha_{0}\right) \oplus \Theta\left(\alpha_{2}\right) \oplus \Theta\left(\alpha_{4}\right) \oplus \cdots .
$$

Taking $n=2 k-1$ in (3.67), (3.68) and multiplying by $z^{-k}$, we get

$$
\begin{aligned}
\chi_{2 k-1} & =x_{2 k-1} \alpha_{2 k-1}^{\dagger}+x_{2 k} \rho_{2 k-1}^{R} \\
\chi_{2 k} & =x_{2 k-1} \rho_{2 k-1}^{L}-x_{2 k} \alpha_{2 k-1} .
\end{aligned}
$$

It follows that the matrix $\mathcal{M}$ has the structure

$$
\begin{equation*}
\mathcal{M}=\mathbf{1} \oplus \Theta\left(\alpha_{1}\right) \oplus \Theta\left(\alpha_{3}\right) \oplus \cdots \tag{3.69}
\end{equation*}
$$

Substituting this into (3.66), we obtain:

$$
\mathcal{C}=\left(\begin{array}{cccccc}
\alpha_{0}^{\dagger} & \rho_{0}^{L} \alpha_{1}^{\dagger} & \rho_{0}^{L} \rho_{1}^{L} & \mathbf{0} & \mathbf{0} & \cdots  \tag{3.70}\\
\rho_{0}^{R} & -\alpha_{0} \alpha_{1}^{\dagger} & -\alpha_{0} \rho_{1}^{L} & \mathbf{0} & \mathbf{0} & \cdots \\
\mathbf{0} & \alpha_{2}^{\dagger} \rho_{1}^{R} & -\alpha_{2}^{\dagger} \alpha_{1} & \rho_{2}^{L} \alpha_{3}^{\dagger} & \rho_{2}^{L} \rho_{3}^{L} & \cdots \\
\mathbf{0} & \rho_{2}^{R} \rho_{1}^{R} & -\rho_{2}^{R} \alpha_{1} & -\alpha_{2} \alpha_{3}^{\dagger} & -\alpha_{2} \rho_{3}^{L} & \cdots \\
\mathbf{0} & \mathbf{0} & \mathbf{0} & \alpha_{4}^{\dagger} \rho_{3}^{R} & -\alpha_{4}^{\dagger} \alpha_{3} & \cdots \\
\vdots & \vdots & \vdots & \vdots & \vdots & \ddots
\end{array}\right) .
$$

We note that the analogous formula to this in [170], namely, (4.30), is incorrect! The order of the factors below the diagonal is wrong there.

### 3.12 The Resolvent of the CMV Matrix

We begin by studying solutions to the equations

$$
\begin{align*}
& \sum_{k=0}^{\infty} \mathcal{C}_{m k} w_{k}=z w_{m}, \quad m \geq 2,  \tag{3.71}\\
& \sum_{k=0}^{\infty} \tilde{w}_{k} \mathcal{C}_{k m}=z \tilde{w}_{m}, \quad m \geq 1 . \tag{3.72}
\end{align*}
$$

Let us introduce the following functions:

$$
\begin{aligned}
\tilde{x}_{n}(z) & =\chi_{n}(1 / \bar{z})^{\dagger}, \\
\Upsilon_{2 n}(z) & =-z^{-n} \psi_{2 n}^{L, *}(z), \\
\Upsilon_{2 n-1}(z) & =z^{-n+1} \psi_{2 n-1}^{R}(z), \\
y_{2 n}(z) & =-\Upsilon_{2 n}(1 / \bar{z})^{\dagger}=z^{-n} \psi_{2 n}^{L}(z), \\
y_{2 n-1}(z) & =-\Upsilon_{2 n-1}(1 / \bar{z})^{\dagger}=-z^{-n} \psi_{2 n-1}^{R, *}(z), \\
p_{n}(z) & =y_{n}(z)+\tilde{x}_{n}(z) F(z), \\
\pi_{n}(z) & =\Upsilon_{n}(z)+F(z) \chi_{n}(z) .
\end{aligned}
$$

Proposition 3.22. Let $z \in \mathbb{D} \backslash\{0\}$.
(i) For each $n \geq 0$, a pair of values ( $\tilde{w}_{2 n}, \tilde{w}_{2 n+1}$ ) uniquely determines a solution $\tilde{w}_{n}$ to (3.72). Also, for any pair of values $\left(\tilde{w}_{2 n}, \tilde{w}_{2 n+1}\right)$ in $\mathcal{M}_{l}$, there exists a solution $\tilde{w}_{n}$ to (3.72) with these values at $(2 n, 2 n+1)$.
(ii) The set of solutions $\tilde{w}_{n}$ to (3.72) coincides with the set of sequences

$$
\begin{equation*}
\tilde{w}_{n}(z)=a \chi_{n}(z)+b \pi_{n}(z) \tag{3.73}
\end{equation*}
$$

where $a, b$ range over $\mathcal{M}_{l}$.
(iii) A solution (3.73) is in $\ell^{2}$ if and only if $a=0$.
(iv) A solution (3.73) obeys (3.72) for all $m \geq 0$ if and only if $b=0$.

Proposition 3.23. Let $z \in \mathbb{D} \backslash\{0\}$.
(i) For each $n \geq 1$, a pair of values $\left(w_{2 n-1}, w_{2 n}\right)$ uniquely determines a solution $w_{n}$ to (3.71). Also, for any pair of values $\left(w_{2 n-1}, w_{2 n}\right)$ in $\mathcal{M}_{l}$, there exists a solution $w_{n}$ to (3.71) with these values at $(2 n-1,2 n)$.
(ii) The set of solutions $w_{n}$ to (3.71) coincides with the set of sequences

$$
\begin{equation*}
w_{n}(z)=\tilde{x}_{n}(z) a+p_{n}(z) b \tag{3.74}
\end{equation*}
$$

where $a, b$ range over $\mathcal{M}_{l}$.
(iii) A solution (3.74) is in $\ell^{2}$ if and only if $a=0$.
(iv) A solution (3.74) obeys (3.71) for all $m \geq 0$ if and only if $b=0$.

Proof of Proposition 3.22. (i) The matrix $\mathcal{C}-z$ can be written in the form

$$
\mathcal{C}-z=\left(\begin{array}{ccccc}
A_{0} & B_{0} & \mathbf{0} & \mathbf{0} & \ldots  \tag{3.75}\\
\mathbf{0} & A_{1} & B_{1} & \mathbf{0} & \ldots \\
\mathbf{0} & \mathbf{0} & A_{2} & B_{2} & \ldots \\
\vdots & \vdots & \vdots & \vdots & \ddots
\end{array}\right)
$$

where

$$
\begin{gathered}
A_{0}=\binom{\alpha_{0}^{\dagger}-z}{\rho_{0}^{R}} \quad A_{n}=\left(\begin{array}{cc}
\alpha_{2 n}^{\dagger} \rho_{2 n-1}^{R} & -\alpha_{2 n}^{\dagger} \alpha_{2 n-1}-z \\
\rho_{2 n}^{R} \rho_{2 n-1}^{R} & -\rho_{2 n}^{R} \alpha_{2 n-1}
\end{array}\right) \\
B_{n}=\left(\begin{array}{cc}
\rho_{2 n}^{L} \alpha_{2 n+1}^{\dagger} & \rho_{2 n}^{L} \rho_{2 n+1}^{L} \\
-\alpha_{2 n} \alpha_{2 n+1}^{\dagger}-z & -\alpha_{2 n} \rho_{2 n+1}^{L}
\end{array}\right) .
\end{gathered}
$$

Define $\widetilde{W}_{n}=\left(\tilde{w}_{2 n}, \tilde{w}_{2 n+1}\right)$ for $n=0,1,2, \ldots$ Then (3.72) for $m=2 n+1,2 n+2$ is equivalent to

$$
\widetilde{W}_{n} B_{n}+\widetilde{W}_{n+1} A_{n+1}=\mathbf{0}
$$

It remains to prove that the $2 l \times 2 l$ matrices $A_{j}, B_{j}$ are invertible. Suppose that for some $x, y \in \mathbb{C}^{l}$, $A_{n}\binom{x}{y}=\binom{0}{0}$. This is equivalent to the system

$$
\begin{array}{r}
\alpha_{2 n}^{\dagger} \rho_{2 n-1}^{R} x-\alpha_{2 n}^{\dagger} \alpha_{2 n-1} y-z y=\mathbf{0}, \\
\rho_{2 n}^{R} \rho_{2 n-1}^{R} x-\rho_{2 n}^{R} \alpha_{2 n-1} y=\mathbf{0} .
\end{array}
$$

The second equation of this system yields $\rho_{2 n-1}^{R} x=\alpha_{2 n-1} y$ (since $\rho_{2 n}^{R}$ is invertible), and upon substitution into the first equation, we get $y=x=0$. Thus, $\operatorname{ker}\left(A_{n}\right)=\{0\}$. In a similar way, one proves that $\operatorname{ker}\left(B_{n}\right)=\{0\}$.
(ii) First note that $\tilde{w}_{n}=\chi_{n}$ is a solution to (3.72) by Lemma 3.21. Let us check that $\tilde{w}_{n}=\Upsilon_{n}$ is also a solution. If $U_{k m}=(-1)^{k} \delta_{k m}$, then $(U C U)_{k m}$ for $m \geq 1$ coincides with the CMV matrix corresponding to the coefficients $\left\{-\alpha_{n}\right\}$. Recall that $\psi_{n}^{L, R}$ are the orthogonal polynomials $\varphi_{n}^{L, R}$, corresponding to the coefficients $\left\{-\alpha_{n}\right\}$. Taking into account the minus signs in the definition of $\Upsilon_{n}$, we see that $\tilde{w}_{n}=\Upsilon_{n}$ solves (3.72) for $m \geq 1$. It follows that any $\tilde{w}_{n}$ of the form (3.73) is a solution to (3.72).

Let us check that any solution to (3.72) can be represented as (3.73). By (i), it suffices to show that for any $\tilde{w}_{0}, \tilde{w}_{1}$, there exist $a, b \in \mathcal{M}_{l}$ such that

$$
\begin{aligned}
& a \chi_{0}(z)+b \pi_{0}(z)=\tilde{w}_{0}, \\
& a \chi_{1}(z)+b \pi_{1}(z)=\tilde{w}_{1} .
\end{aligned}
$$

Recalling that $\chi_{0}=1, \Upsilon_{0}=-1, \Upsilon_{1}(z)=\left(z+\alpha_{0}^{\dagger}\right)\left(\rho_{0}^{R}\right)^{-1}, \chi_{1}(z)=\left(z-\alpha_{0}^{\dagger}\right)\left(\rho_{0}^{R}\right)^{-1}$, we see that the above system can be easily solved for $a, b$ if $z \neq 0$.
(iii) Let us prove that the solution $\pi_{n}$ is square integrable. We will consider separately the sequences $\pi_{2 n}$ and $\pi_{2 n-1}$ and prove that they both belong to $\ell^{2}$. By (3.20) and (3.23), we have

$$
\begin{align*}
\psi_{n}^{R}(z)+F(z) \varphi_{n}^{R}(z) & =\int \frac{e^{i \theta}+z}{e^{i \theta}-z} d \mu(\theta) \varphi_{n}^{R}\left(e^{i \theta}\right),  \tag{3.76}\\
\psi_{n}^{L, *}(z)-F(z) \varphi_{n}^{L, *}(z) & =-z^{n} \int \frac{e^{i \theta}+z}{e^{i \theta}-z} d \mu(\theta) \varphi_{n}^{L}\left(e^{i \theta}\right)^{\dagger} \tag{3.77}
\end{align*}
$$

Taking $n=2 k$ in (3.77) and $n=2 k-1$ in (3.76), we get

$$
\begin{equation*}
\pi_{2 k}(z)=z^{k} \int \frac{e^{i \theta}+z}{e^{i \theta}-z} d \mu(\theta) \varphi_{2 k}^{L}\left(e^{i \theta}\right)^{\dagger} \tag{3.78}
\end{equation*}
$$

$$
\begin{equation*}
\pi_{2 k-1}(z)=z^{-k+1} \int \frac{e^{i \theta}+z}{e^{i \theta}-z} d \mu(\theta) \varphi_{2 k-1}^{R}\left(e^{i \theta}\right) . \tag{3.79}
\end{equation*}
$$

As $\varphi_{2 k}^{L}$ is an orthonormal sequence, using the Bessel inequality, from (3.78) we immediately get that $\pi_{2 k}$ is in $\ell^{2}$.

Consider the odd terms $\pi_{2 k-1}$. We claim that

$$
\begin{equation*}
z^{-k+1} \int \frac{e^{i \theta}+z}{e^{i \theta}-z} d \mu(\theta) \varphi_{2 k-1}^{R}\left(e^{i \theta}\right)=\int \frac{e^{i \theta}+z}{e^{i \theta}-z} d \mu(\theta) e^{i(-k+1) \theta} \varphi_{2 k-1}^{R}\left(e^{i \theta}\right) \tag{3.80}
\end{equation*}
$$

Indeed, using the right orthogonality of $\varphi_{2 k-1}^{R}$ to $e^{i m \theta}, m=0,1, \ldots, 2 k-2$, we get

$$
\begin{aligned}
\int \frac{e^{i \theta}+z}{e^{i \theta}-z} d \mu(\theta) \varphi_{2 k-1}^{R}\left(e^{i \theta}\right) & =\left\langle\left\langle 1+2 \sum_{m=1}^{\infty} \bar{z}^{m} e^{i m \theta}, \varphi_{2 k-1}^{R}\right\rangle\right\rangle_{R} \\
& =\left\langle\left\langle 2 \sum_{m=2 k-1}^{\infty} \bar{z}^{m} e^{i m \theta}, \varphi_{2 k-1}^{R}\right\rangle\right\rangle_{R}
\end{aligned}
$$

and

$$
\begin{aligned}
& \int \frac{e^{i \theta}+z}{e^{i \theta}-z} z^{k-1} e^{i(-k+1) \theta} d \mu(\theta) \varphi_{2 k-1}^{R}\left(e^{i \theta}\right)= \\
&=\left\langle\left\langle\bar{z}^{k-1} e^{i(k-1) \theta}\left(1+2 \sum_{m=1}^{\infty} \bar{z}^{m} e^{i m \theta}\right), \varphi_{2 k-1}^{R}\right\rangle\right\rangle_{R} \\
&=\left\langle\left\langle 2 \sum_{m=2 k-1}^{\infty} \bar{z}^{m} e^{i m \theta}, \varphi_{2 k-1}^{R}\right\rangle\right\rangle_{R}
\end{aligned}
$$

which proves (3.80). The identities (3.80) and (3.79) yield

$$
\pi_{2 k-1}(z)=\left\langle\left\langle\frac{e^{-i \theta}+\bar{z}}{e^{-i \theta}-\bar{z}}, \chi_{2 k-1}\right\rangle\right\rangle_{R}
$$

and, since $\chi_{2 k-1}$ is a right orthogonal sequence, the Bessel inequality ensures that $\pi_{2 k-1}(z)$ is in $\ell^{2}$. Thus, $\pi_{k}(z)$ is in $\ell^{2}$.

Next, as in the proof of Theorem 3.18, using the CD formula, we check that the sequence $\left\|\varphi_{n}^{L, *}(z)\right\|_{R}$ is bounded below and therefore the sequence $\chi_{2 n}(z)$ is not in $\ell^{2}$. This proves the statement (iii).
(iv) By Lemma 3.21, the solution $\chi_{n}(z)$ obeys (3.72) for all $m \geq 0$. It is easy to check directly that the solution $\pi_{n}(z)$ does not obey (3.72) for $m=0$ if $z \neq 0$. This proves the required statement.

Proof of Proposition 3.23. (i) For $j=1,2, \ldots$, define $W_{j}=\left(w_{2 j-1}, w_{2 j}\right)$. Then, using the block structure (3.75), we can rewrite (3.71) for $m=2 j, 2 j+1$ as $A_{j} W_{j}+B_{j} W_{j+1}=\mathbf{0}$. By the proof of Proposition 3.22, the matrices $A_{j}$ and $B_{j}$ are invertible, which proves (i).
(ii) Lemma 3.21 ensures that $\tilde{x}_{n}(z)$ is a solution of (3.71). As in the proof of Proposition 3.22, by considering the matrix $(U C U)_{k m}$, one checks that $y_{n}(z)$ is also a solution to (3.71).

Let us prove that any solution to (3.71) can be represented in the form (3.74). By (i), it suffices to show that for any $w_{1}, w_{2}$, there exist $a, b \in \mathcal{M}_{l}$ such that

$$
\begin{aligned}
& \tilde{x}_{1}(z) a+p_{1}(z) b=w_{1}, \\
& \tilde{x}_{2}(z) a+p_{2}(z) b=w_{2} .
\end{aligned}
$$

We claim that this system of equations can be solved for $a, b$. Here are the main steps. Substituting the definitions of $\tilde{x}_{n}(z)$ and $p_{n}(z)$, we rewrite this system as

$$
\begin{aligned}
\varphi_{1}^{R, *}(a+F(z) b)-\psi_{1}^{R, *}(z) b & =z w_{1} \\
\varphi_{2}^{L}(a+F(z) b)+\psi_{2}^{L}(z) b & =z w_{2} .
\end{aligned}
$$

Using Szegő recurrence, we can substitute the expressions for $\varphi_{2}^{L}, \psi_{2}^{L}$, which helps rewrite our system as

$$
\begin{aligned}
\varphi_{1}^{R, *}(a+F(z) b)-\psi_{1}^{R, *}(z) b & =z w_{1}, \\
\varphi_{1}^{L}(a+F(z) b)+\psi_{1}^{L}(z) b & =\rho_{1}^{L} w_{2}+\alpha_{1}^{\dagger} w_{1} .
\end{aligned}
$$

Substituting explicit formulas for $\varphi_{1}^{R, *}, \varphi_{1}^{L}, \psi_{1}^{R, *}, \psi_{1}^{L}$, and expressing $F(z)$ in terms of $f(z)$, we can rewrite this as

$$
\begin{aligned}
& \left(\rho_{0}^{R}\right)^{-1}\left(1-\alpha_{0} z\right) a+2 z\left(\rho_{0}^{R}\right)^{-1}\left(f(z)-\alpha_{0}\right)(\mathbf{1}-z f(z))^{-1} b=z w_{1}, \\
& \left(\rho_{0}^{L}\right)^{-1}\left(z-\alpha_{0}^{\dagger}\right) a+2 z\left(\rho_{0}^{L}\right)^{-1}\left(\mathbf{1}-\alpha_{0}^{\dagger} f(z)\right)(\mathbf{1}-z f(z))^{-1} b=\rho_{1}^{L} w_{2}+\alpha_{1}^{\dagger} w_{1} .
\end{aligned}
$$

Denote

$$
\begin{aligned}
a_{1} & =\left(\rho_{0}^{R}\right)^{-1}\left(1-\alpha_{0} z\right) a, \\
b_{1} & =2 z\left(\rho_{0}^{L}\right)^{-1}\left(\mathbf{1}-\alpha_{0}^{\dagger} f(z)\right)(\mathbf{1}-z f(z))^{-1} b .
\end{aligned}
$$

Then in terms of $a_{1}, b_{1}$, our system can be rewritten as

$$
\begin{aligned}
& a_{1}+X_{1} b_{1}=z w_{1} \\
& X_{2} a_{1}+b_{1}=\rho_{1}^{L} w_{2}+\alpha_{1}^{\dagger} w_{1}
\end{aligned}
$$

where

$$
\begin{aligned}
& X_{1}=\left(\rho_{0}^{R}\right)^{-1}\left(f(z)-\alpha_{0}\right)\left(\mathbf{1}-\alpha_{0} f(z)\right)^{-1} \rho_{0}^{L}, \\
& X_{2}=\left(\rho_{0}^{L}\right)^{-1}\left(z-\alpha_{0}^{\dagger}\right)\left(\mathbf{1}-\alpha_{0} z\right)^{-1} \rho_{0}^{R} .
\end{aligned}
$$

Since $\|f(z)\|<1$ and $|z|<1$, we can apply Corollary 1.5 , which yields $\left\|X_{1}\right\|<1$ and $\left\|X_{2}\right\|<1$. It follows that our system can be solved for $a_{1}, b_{1}$.
(iii) As $p_{n}(z)=-\pi_{n}(1 / \bar{z})^{\dagger}$, by Proposition 3.22, we get that $p_{n}(z)$ is in $\ell^{2}$. In the same way, as $\tilde{x}_{n}(z)=\chi_{n}(1 / \bar{z})^{\dagger}$, we get that $\tilde{x}_{n}(z)$ is not in $\ell^{2}$.
(iv) By Lemma 3.21, the solution $\tilde{x}_{n}(z)$ obeys (3.71) for all $m \geq 0$. Using the explicit formula for $y_{n}(z)$, one easily checks that the solution $y_{n}(z)$ does not obey (3.71) for $m=0,1$.

Theorem 3.24. We have for $z \in \mathbb{D}$,

$$
\left[(\mathcal{C}-z)^{-1}\right]_{k, l}= \begin{cases}(2 z)^{-1} \tilde{x}_{k}(z) \pi_{l}(z), & l>k \text { or } k=l \text { even } \\ (2 z)^{-1} p_{k}(z) \chi_{l}(z), & k>l \text { or } k=l \text { odd } .\end{cases}
$$

Proof. Fix $z \in \mathbb{D}$. Write $G_{k, l}(z)=\left[(\mathcal{C}-z)^{-1}\right]_{k, l}$. Then $G_{\cdot, l}(z)$ is equal to $(\mathcal{C}-z)^{-1} \delta_{l}$, which means that $G_{k, l}(z)$ solves (3.71) for $m \neq l$. Since $G_{\cdot, l}(z)$ is $\ell^{2}$ at infinity and obeys the equation at $m=0$, we see that it is a right-multiple of $p$ for large $k$ and a right-multiple of $\tilde{x}$ for small $k$. Thus,

$$
G_{k, l}(z)= \begin{cases}\tilde{x}_{k}(z) a_{l}(z), & k<l \text { or } k=l \text { even }, \\ p_{k}(z) b_{l}(z), & k>l \text { or } k=l \text { odd }\end{cases}
$$

Similarly,

$$
G_{k, l}(z)= \begin{cases}\tilde{b}_{k}(z) \pi_{l}(z), & k<l \text { or } k=l \text { even }, \\ \tilde{a}_{k}(z) \chi_{l}(z), & k>l \text { or } k=l \text { odd } .\end{cases}
$$

Equating the two expressions, we find

$$
\begin{align*}
\tilde{x}_{k}(z) a_{l}(z) & =\tilde{b}_{k}(z) \pi_{l}(z)  \tag{3.81}\\
p_{k}(z) b_{l}(z) & =\tilde{a}_{k}(z) \chi_{l}(z)  \tag{3.82}\\
& k>l \text { or } k=l \text { or } k=l \text { odd. }
\end{align*}
$$

Putting $k=0$ in (3.81) and setting $\tilde{b}_{0}(z)=c_{1}(z)$, we find $a_{l}(z)=c_{1}(z) \pi_{l}(z)$. Putting $l=0$ in (3.82) and setting $b_{0}(z)=c_{2}(z)$, we find $p_{k}(z) c_{2}(z)=\tilde{a}_{k}(z)$. Thus,

$$
G_{k, l}(z)= \begin{cases}\tilde{x}_{k}(z) c_{1}(z) \pi_{l}(z), & k<l \text { or } k=l \text { even } \\ p_{k}(z) c_{2}(z) \chi_{l}(z), & k>l \text { or } k=l \text { odd }\end{cases}
$$

We claim that $c_{1}(z)=c_{2}(z)=(2 z)^{-1} \mathbf{1}$. Consider the case $k=l=0$. Then, on the one hand, by the definition,

$$
\begin{align*}
G_{0,0}(z) & =\int \frac{1}{e^{i \theta}-z} d \mu\left(e^{i \theta}\right) \\
& =\int(2 z)^{-1}\left[\frac{e^{i \theta}+z}{e^{i \theta}-z}-1\right] d \mu\left(e^{i \theta}\right) \\
& =(2 z)^{-1}(F(z)-\mathbf{1}) \tag{3.83}
\end{align*}
$$

and on the other hand,

$$
G_{0,0}(z)=\tilde{x}_{0}(z) c_{1}(z) \pi_{0}(z)=c_{1}(z)(F(z)-\mathbf{1}) .
$$

This shows $c_{1}(z)=(2 z)^{-1} \mathbf{1}$. Next, consider the case $k=1, l=0$. Then, on the one hand, by the definition,

$$
G_{1,0}(z)=\left\langle\left\langle\chi_{1},\left(e^{i \theta}-z\right)^{-1} \chi_{0}\right\rangle\right\rangle_{R}
$$

and on the other hand,

$$
G_{1,0}(z)=p_{1}(z) c_{2}(z) \chi_{0}(z)
$$

Let us calculate the expressions on the right-hand side. We have

$$
\begin{equation*}
p_{1}(z) c_{2}(z) \chi_{0}(z)=\left(\rho_{0}^{R}\right)^{-1}\left(-z^{-1}-\alpha_{0}+\left(z^{-1}-\alpha_{0}\right) F(z)\right) c_{2}(z) \tag{3.84}
\end{equation*}
$$

and

$$
\begin{aligned}
& \left\langle\left\langle\chi_{1},\left(e^{i \theta}-z\right)^{-1} \chi_{0}\right\rangle\right\rangle_{R}= \\
& \quad=\left(\rho_{0}^{R}\right)^{-1} \int\left(e^{-i \theta}-\alpha_{0}\right) d \mu(\theta)\left(e^{i \theta}-z\right)^{-1} \\
& \quad=\left(\rho_{0}^{R}\right)^{-1} \int\left[z^{-1}\left(e^{i \theta}-z\right)^{-1}-z^{-1} e^{-i \theta}-\alpha_{0}\left(e^{i \theta}-z\right)^{-1}\right] d \mu(\theta) \\
& \quad=\left(\rho_{0}^{R}\right)^{-1}\left[\frac{1}{2 z 2}(F(z)-\mathbf{1})-\frac{1}{2 z} \alpha_{0}(F(z)-\mathbf{1})-\frac{1}{z} \int e^{-i \theta} d \mu(\theta)\right]
\end{aligned}
$$

Taking into account the identity

$$
\int e^{-i \theta} d \mu(\theta)=\alpha_{0}
$$

(which can be obtained, e.g., by expanding $\left\langle\left\langle\varphi_{1}^{R}, \varphi_{0}^{R}\right\rangle\right\rangle_{R}=0$ ), we get

$$
\left\langle\left\langle\chi_{1},\left(e^{i \theta}-z\right)^{-1} \chi_{0}\right\rangle\right\rangle_{R}=\frac{1}{2 z}\left(\rho_{0}^{R}\right)^{-1}\left(-z^{-1}-\alpha_{0}+\left(z^{-1}-\alpha_{0}\right) F(z)\right)
$$

Comparing this with (3.84), we get $c_{2}(z)=(2 z)^{-1} \mathbf{1}$.
As an immediate corollary, evaluating the kernel on the diagonal for even and odd indices, we obtain the formulas

$$
\begin{align*}
\int \varphi_{2 n}^{L}\left(e^{i \theta}\right) \frac{d \mu(\theta)}{e^{i \theta}-z} \varphi_{2 n}^{L}\left(e^{i \theta}\right)^{\dagger} & =-\frac{1}{2 z^{2 n+1}} \varphi_{2 n}^{L}(z) u_{2 n}^{L, *}(z)  \tag{3.85}\\
\int \varphi_{2 n-1}^{R}\left(e^{i \theta}\right)^{\dagger} \frac{d \mu(\theta)}{e^{i \theta}-z} \varphi_{2 n-1}^{R}\left(e^{i \theta}\right) & =-\frac{1}{2 z^{2 n}} u_{2 n-1}^{R, *}(z) \varphi_{2 n-1}^{R}(z) \tag{3.86}
\end{align*}
$$

Combining this with (3.31) and (3.32), we find

$$
\begin{align*}
u_{n}^{L}(z) \varphi_{n}^{L, *}(z)+\varphi_{n}^{L}(z) u_{n}^{L, *}(z) & =2 z^{n}  \tag{3.87}\\
\varphi_{n}^{R, *}(z) u_{n}^{R}(z)+u_{n}^{R, *}(z) \varphi_{n}^{R}(z) & =2 z^{n} \tag{3.88}
\end{align*}
$$

### 3.13 Khrushchev Theory

Among the deepest and most elegant methods in OPUC are those of Khrushchev [125, 126, 101]. We have not been able to extend them to MOPUC! We regard their extension as an important open question; we present the first very partial steps here.

Let

$$
\Omega=\{\theta: \operatorname{det} w(\theta)>0\}
$$

Theorem 3.25. For every $n \geq 0$,

$$
\left\{\theta: f_{n}\left(e^{i \theta}\right)^{\dagger} f_{n}\left(e^{i \theta}\right)<\mathbf{1}\right\}=\Omega
$$

up to a set of zero Lebesgue measure.
Consequently,

$$
\begin{equation*}
\int\left\|f_{n}\left(e^{i \theta}\right)\right\| \frac{d \theta}{2 \pi} \geq 1-\frac{|\Omega|}{2 \pi} \tag{3.89}
\end{equation*}
$$

Proof. Recall that, by Proposition 3.16, up to a set of zero Lebesgue measure,

$$
\left\{\theta: f_{0}\left(e^{i \theta}\right)^{\dagger} f_{0}\left(e^{i \theta}\right)<\mathbf{1}\right\}=\{\theta: \operatorname{det} w(\theta)>0\}
$$

so, by induction, it suffices to show that, up to a set of zero Lebesgue measure,

$$
\left\{\theta: f_{0}\left(e^{i \theta}\right)^{\dagger} f_{0}\left(e^{i \theta}\right)<\mathbf{1}\right\}=\left\{\theta: f_{1}\left(e^{i \theta}\right)^{\dagger} f_{1}\left(e^{i \theta}\right)<\mathbf{1}\right\} .
$$

This in turn follows from the fact that the Schur algorithm, which relates the two functions, preserves the property $g^{\dagger} g<\mathbf{1}$.

Notice that away from $\Omega, f_{n}\left(e^{i \theta}\right)$ has norm one and therefore,

$$
\int\left\|f_{n}\left(e^{i \theta}\right)\right\| \frac{d \theta}{2 \pi} \geq \int_{\Omega^{c}}\left\|f_{n}\left(e^{i \theta}\right)\right\| \frac{d \theta}{2 \pi}=\int_{\Omega^{c}} 1 \frac{d \theta}{2 \pi}
$$

which yields (3.89).
Define

$$
b_{n}(z ; d \mu)=\varphi_{n}^{L}(z ; d \mu) \varphi_{n}^{R, *}(z ; d \mu)^{-1}
$$

Proposition 3.26. (a) $b_{n+1}=\left(\rho_{n}^{L}\right)^{-1}\left(z b_{n}-\alpha_{n}^{\dagger}\right)\left(1-z \alpha_{n} b_{n}\right)^{-1} \rho_{n}^{R}$.
(b) The Verblunsky coefficients of $b_{n}$ are $\left(-\alpha_{n-1}^{\dagger},-\alpha_{n-2}^{\dagger}, \ldots,-\alpha_{0}^{\dagger}, \mathbf{1}\right)$.

Proof. (a) By the Szegő recursion, we have that

$$
\begin{aligned}
b_{n+1}= & \varphi_{n+1}^{L}\left(\varphi_{n+1}^{R, *}\right)^{-1} \\
= & \left(\left(\rho_{n}^{L}\right)^{-1} z \varphi_{n}^{L}-\left(\rho_{n}^{L}\right)^{-1} \alpha_{n}^{\dagger} \varphi_{n}^{R, *}\right)\left(\left(\rho_{n}^{R}\right)^{-1} \varphi_{n}^{R, *}-z\left(\rho_{n}^{R}\right)^{-1} \alpha_{n} \varphi_{n}^{L}\right)^{-1} \\
= & \left(\rho_{n}^{L}\right)^{-1}\left(z \varphi_{n}^{L}-\alpha_{n}^{\dagger} \varphi_{n}^{R, *}\right)\left(\varphi_{n}^{R, *}-z \alpha_{n} \varphi_{n}^{L}\right)^{-1} \rho_{n}^{R} \\
= & \left(\rho_{n}^{L}\right)^{-1}\left(z \varphi_{n}^{L}\left(\varphi_{n}^{R, *}\right)^{-1}-\alpha_{n}^{\dagger} \varphi_{n}^{R, *}\left(\varphi_{n}^{R, *}\right)^{-1}\right) \\
& \quad\left(\varphi_{n}^{R, *}\left(\varphi_{n}^{R, *}\right)^{-1}-z \alpha_{n} \varphi_{n}^{L}\left(\varphi_{n}^{R, *}\right)^{-1}\right)^{-1} \rho_{n}^{R} \\
= & \left(\rho_{n}^{L}\right)^{-1}\left(z b_{n}-\alpha_{n}^{\dagger}\right)\left(\mathbf{1}-z \alpha_{n} b_{n}\right)^{-1} \rho_{n}^{R} .
\end{aligned}
$$

(b) It follows from part (a) that the first Verblunsky coefficient of $b_{n}$ is $-\alpha_{n-1}^{\dagger}$ and that its first Schur iterate is $b_{n-1}$; compare Theorem 3.19. This gives the claim by induction and the fact that $b_{0}=1$.

## 4 The Szegő Mapping and the Geronimus Relations

In this chapter, we present the matrix analogue of the Szegő mapping and the resulting Geronimus relations. This establishes a correspondence between certain matrix-valued measures on the unit circle and matrix-valued measures on the interval $[-2,2]$ and, consequently, a correspondence between Verblunsky coefficients and Jacobi parameters. Throughout this chapter, we will denote measures on the circle by $d \mu_{C}$ and measures on the interval by $d \mu_{I}$.

The scalar versions of these objects are due to Szegő [182] and Geronimus [90]. There are four proofs that we know of: the original argument of Geronimus [90] based on Szegő's formula in [182], a proof of Damanik-Killip [32] using Schur functions, a proof of Killip-Nenciu [127] using CMV matrices, and a proof of Faybusovich-Gekhtman [81] using canonical moments.

The matrix version of these objects was studied by Yakhlef-Marcellán [199] who proved Theorem 4.2 below using the Geronimus-Szegő approach. Our proof uses the Killip-Nenciu-CMV approach. In comparing our formula with [199], one needs the following dictionary (their objects on the left of the equal sign and ours on the right):

$$
\begin{aligned}
& H_{n}=-\alpha_{n+1}^{\dagger}, \\
& D_{n}=A_{n}, \\
& E_{n}=B_{n+1} .
\end{aligned}
$$

Dette-Studden [43] have extended the theory of canonical moments from OPRL to MOPRL. It would be illuminating to use this to extend the proof that Faybusovich-Gekhtman [81] gave of Geronimus relations for scalar OPUC to MOPUC.

Suppose $d \mu_{C}$ is a non-trivial positive semi-definite Hermitian matrix measure on the unit circle that is invariant under $\theta \mapsto-\theta$ (i.e., $z \mapsto \bar{z}=z^{-1}$ ). Then we define the measure $d \mu_{I}$ on the interval $[-2,2]$ by

$$
\int f(x) d \mu_{I}(x)=\int f(2 \cos \theta) d \mu_{C}(\theta)
$$

for $f$ measurable on $[-2,2]$. The map

$$
\mathrm{Sz}: d \mu_{C} \mapsto d \mu_{I}
$$

is called the Szegő mapping.
The Szegő mapping can be inverted as follows. Suppose $d \mu_{I}$ is a non-degenerate positive semidefinite matrix measure on $[-2,2]$. Then we define the measure $d \mu_{C}$ on the unit circle which is invariant under $\theta \mapsto-\theta$ by

$$
\int g(\theta) d \mu_{C}(\theta)=\int g(\arccos (x / 2)) d \mu_{I}(x)
$$

for $g$ measurable on $\partial \mathbb{D}$ with $g(\theta)=g(-\theta)$.
We first show that for the measures on the circle of interest in this section, the Verblunsky coefficients are always Hermitian.

Lemma 4.1. Suppose $d \mu_{C}$ is a non-trivial positive semi-definite Hermitian matrix measure on the unit circle. Denote the associated Verblunsky coefficients by $\left\{\alpha_{n}\right\}$. Then, $d \mu_{C}$ is invariant under $\theta \mapsto-\theta$ if and only if $\alpha_{n}^{\dagger}=\alpha_{n}$ for every $n$.

Proof. For a polynomial $P$, denote $\tilde{P}(z)=P(\bar{z})^{\dagger}$.

1. Suppose that $d \mu_{C}$ is invariant under $\theta \mapsto-\theta$. Then we have

$$
\langle\langle f, g\rangle\rangle_{L}=\left\langle\langle\tilde{g}, \tilde{f}\rangle_{R}\right.
$$

for all $f, g$. Inspecting the orthogonality conditions which define $\Phi_{n}^{L}$ and $\Phi_{n}^{R}$, we see that

$$
\begin{equation*}
\tilde{\Phi}_{n}^{L}=\Phi_{n}^{R} \text { and }\left\langle\left\langle\Phi_{n}^{L}, \Phi_{n}^{L}\right\rangle_{L}=\left\langle\left\langle\Phi_{n}^{R}, \Phi_{n}^{R}\right\rangle\right\rangle_{R} .\right. \tag{4.1}
\end{equation*}
$$

Next, we claim that

$$
\begin{equation*}
\kappa_{n}^{L}=\kappa_{n}^{R, \dagger} . \tag{4.2}
\end{equation*}
$$

Indeed, recall the definition of $\kappa_{n}^{L}, \kappa_{n}^{R}$ :

$$
\begin{aligned}
& \kappa_{n}^{L}=u_{n}\left\langle\left\langle\Phi_{n}^{L}, \Phi_{n}^{L}\right\rangle\right\rangle_{L}^{-1 / 2}, u_{n} \text { is unitary, } \kappa_{n+1}^{L}\left(\kappa_{n}^{L}\right)^{-1}>0, \\
& \kappa_{n}^{R}=\left\langle\left\langle\Phi_{n}^{R}, \Phi_{n}^{R}\right\rangle\right\rangle_{R}^{-1 / 2} v_{n}, v_{n} \text { is unitary, }\left(\kappa_{n}^{R}\right)^{-1} \kappa_{n+1}^{R}>0, \\
& \kappa_{0}^{R}=1 .
\end{aligned}
$$

Using this definition and (4.1), one can easily prove by induction that $v_{n}=u_{n}^{\dagger}$ and therefore (4.2) holds true.

Next, taking $z=0$ in (3.11), we get

$$
\begin{aligned}
& \alpha_{n}=-\left(\kappa_{n}^{R}\right)^{-1} \Phi_{n+1}^{L}(0)^{\dagger}\left(\kappa_{n}^{L}\right)^{\dagger}, \\
& \alpha_{n}=-\left(\kappa_{n}^{R}\right)^{\dagger} \Phi_{n+1}^{R}(0)^{\dagger}\left(\kappa_{n}^{L}\right)^{-1} .
\end{aligned}
$$

From here and (4.1), (4.2), we get $\alpha_{n}=\alpha_{n}^{\dagger}$.
2. Assume $\alpha_{n}^{\dagger}=\alpha_{n}$ for all $n$. Then, by Theorem 3.3(c), we have $\rho_{n}^{L}=\rho_{n}^{R}$. It follows that in this case the Szegő recurrence relation is invariant with respect to the change $\varphi_{n}^{L} \mapsto \tilde{\varphi}_{n}^{R}, \varphi_{n}^{R} \mapsto \tilde{\varphi}_{n}^{L}$. It follows that $\varphi_{n}^{L}=\tilde{\varphi}_{n}^{R}, \varphi_{n}^{R}=\tilde{\varphi}_{n}^{L}$. In particular, we get

$$
\begin{equation*}
\left\langle\left\langle\varphi_{n}^{L}, \varphi_{m}^{L}\right\rangle\right\rangle_{L}=\left\langle\left\langle\tilde{\varphi}_{m}^{R}, \tilde{\varphi}_{n}^{R}\right\rangle\right\rangle_{R} \tag{4.3}
\end{equation*}
$$

Now let $f$ and $g$ be any polynomials; we have

$$
f(z)=\sum_{n} f_{n} \varphi_{n}^{L}(z), \quad \tilde{f}(z)=\sum_{n} \tilde{\varphi}_{n}^{L}(z) f_{n}^{\dagger}
$$

and a similar expansion for $g$. Using these expansions and (4.3), we get

$$
\left\langle\langle f, g\rangle_{L}=\langle\langle\tilde{g}, \tilde{f}\rangle\rangle_{R} \text { for all polynomials } f, g\right.
$$

From here it follows that the measure $d \mu_{C}$ is invariant under $\theta \mapsto-\theta$.
Now consider two measures $d \mu_{C}$ and $d \mu_{I}=\mathrm{Sz}\left(d \mu_{C}\right)$ and the associated CMV and Jacobi matrices. What are the relations between the parameters of these matrices?

Theorem 4.2. Given $d \mu_{C}$ and $d \mu_{I}=\mathrm{Sz}\left(d \mu_{C}\right)$ as above, the coefficients of the associated CMV and Jacobi matrices satisfy the Geronimus relations:

$$
\begin{align*}
& B_{k+1}=\sqrt{\mathbf{1 - \alpha _ { 2 k - 1 }}} \alpha_{2 k} \sqrt{\mathbf{1 - \alpha _ { 2 k - 1 }}}-\sqrt{\mathbf{1}+\alpha_{2 k-1}} \alpha_{2 k-2} \sqrt{\mathbf{1}+\alpha_{2 k-1}}  \tag{4.4}\\
& A_{k+1}=\sqrt{\mathbf{1}-\alpha_{2 k-1}} \sqrt{\mathbf{1}-\alpha_{2 k}^{2}} \sqrt{\mathbf{1}+\alpha_{2 k+1}} \tag{4.5}
\end{align*}
$$

Remarks. 1. For these formulas to hold for $k=0$, we set $\alpha_{-1}=\mathbf{1}$.
2. There are several proofs of the Geronimus relations in the scalar case. We follow the proof given by Killip and Nenciu in [127].
3. These $A$ 's are, in general, not type 1 or 2 or 3 .

Proof. For a Hermitian $l \times l$ matrix $\alpha$ with $\|\alpha\|<1$, define the unitary $2 l \times 2 l$ matrix $S(\alpha)$ by

$$
S(\alpha)=\frac{1}{\sqrt{2}}\left(\begin{array}{cc}
\sqrt{1-\alpha} & -\sqrt{1+\alpha} \\
\sqrt{1+\alpha} & \sqrt{1-\alpha}
\end{array}\right) .
$$

Since $\alpha^{\dagger}=\alpha$, the associated $\rho^{L}$ and $\rho^{R}$ coincide and will be denoted by $\rho$. We therefore have

$$
\Theta(\alpha)=\left(\begin{array}{cc}
\alpha & \rho \\
\rho & -\alpha
\end{array}\right)
$$

and hence, by a straightforward calculation,

$$
\begin{aligned}
S(\alpha) \Theta(\alpha) S(\alpha)^{-1} & =\frac{1}{2}\left(\begin{array}{cc}
\sqrt{\mathbf{1}-\alpha} & -\sqrt{\mathbf{1 + \alpha}} \\
\sqrt{\mathbf{1}+\alpha} & \sqrt{\mathbf{1}-\alpha}
\end{array}\right)\left(\begin{array}{cc}
\alpha & \rho \\
\rho & -\alpha
\end{array}\right)\left(\begin{array}{cc}
\sqrt{\mathbf{1}-\alpha} & \sqrt{\mathbf{1 + \alpha}} \\
-\sqrt{\mathbf{1}+\alpha} & \sqrt{\mathbf{1}-\alpha}
\end{array}\right) \\
& =\left(\begin{array}{cc}
-\mathbf{1} & \mathbf{0} \\
\mathbf{0} & \mathbf{1}
\end{array}\right) .
\end{aligned}
$$

Thus, if we define

$$
\mathcal{S}=\mathbf{1} \oplus S\left(\alpha_{1}\right) \oplus S\left(\alpha_{3}\right) \oplus \ldots
$$

and

$$
\mathcal{R}=\mathbf{1} \oplus(-\mathbf{1}) \oplus \mathbf{1} \oplus(-\mathbf{1}) \oplus \ldots
$$

it follows that (see (3.69))

$$
\mathcal{S M S} \mathcal{S}^{\dagger}=\mathcal{R}
$$

The matrix $\mathcal{L} \mathcal{M}+\mathcal{M} \mathcal{L}$ is unitarily equivalent to

$$
\mathcal{A}=\mathcal{S}(\mathcal{L M}+\mathcal{M} \mathcal{L}) \mathcal{S}^{\dagger}=\mathcal{S} \mathcal{L} \mathcal{S}^{\dagger} \mathcal{R}+\mathcal{R} \mathcal{S} \mathcal{L} \mathcal{S}^{\dagger}
$$

Observe that $\mathcal{A}$ is the direct sum of two block Jacobi matrices. Indeed, it follows quickly from the explicit form of $\mathcal{R}$ that the even-odd and odd-even block entries of $\mathcal{A}$ vanish. Consequently, $\mathcal{A}$ is the direct sum of its odd-odd and even-even block entries. We will call these two block Jacobi matrices $J$ and $\tilde{J}$, respectively.

Consider $\mathcal{C}$ as an operator on $\mathcal{H}_{v}$. Then $d \mu_{C}$ is the spectral measure of $\mathcal{C}$ in the following sense:

$$
\left[\mathcal{C}^{m}\right]_{0,0}=\int e^{i m \theta} d \mu_{C}(\theta)
$$

see (3.83). Then, by the spectral theorem, $d \mu_{C}(-\theta)$ is the spectral measure of $\mathcal{C}^{-1}$ and so $d \mu_{I}$ is the spectral measure of $\mathcal{C}+\mathcal{C}^{-1}=\mathcal{L} \mathcal{M}+(\mathcal{L M})^{-1}=\mathcal{L} \mathcal{M}+\mathcal{M} \mathcal{L}$. Since $\mathcal{S}$ leaves $\left(\begin{array}{lllll}1 & \mathbf{0} & \mathbf{0} & \mathbf{0} & \cdots\end{array}\right)^{t}$ invariant, we see that $d \mu_{I}$ is the spectral measure of $J$.

To determine the block entries of $J$, we only need to compute the odd-odd block entries of $\mathcal{A}$. For $k \geq 0$, we have

$$
\begin{aligned}
\mathcal{A}_{2 k+1,2 k+1} & =\left(\begin{array}{ll}
\sqrt{\mathbf{1}+\alpha_{2 k-1}} & \sqrt{\mathbf{1}-\alpha_{2 k-1}}
\end{array}\right)\left(\begin{array}{cc}
-\alpha_{2 k-2} & \mathbf{0} \\
\mathbf{0} & \alpha_{2 k}
\end{array}\right)\binom{\sqrt{\mathbf{1}+\alpha_{2 k-1}}}{\sqrt{\mathbf{1}-\alpha_{2 k-1}}} \\
& =\sqrt{\mathbf{1}-\alpha_{2 k-1}} \alpha_{2 k} \sqrt{\mathbf{1}-\alpha_{2 k-1}}-\sqrt{\mathbf{1}+\alpha_{2 k-1}} \alpha_{2 k-2} \sqrt{\mathbf{1}+\alpha_{2 k-1}}
\end{aligned}
$$

and

$$
\begin{aligned}
\mathcal{A}_{2 k+1,2 k+3} & =\left(\begin{array}{ll}
\sqrt{\mathbf{1}+\alpha_{2 k-1}} & \sqrt{\mathbf{1}-\alpha_{2 k-1}}
\end{array}\right)\left(\begin{array}{cc}
\mathbf{0} & \mathbf{0} \\
\rho_{2 k} & \mathbf{0}
\end{array}\right)\binom{\sqrt{\mathbf{1}+\alpha_{2 k+1}}}{\sqrt{\mathbf{1}-\alpha_{2 k+1}}} \\
& =\sqrt{\mathbf{1}-\alpha_{2 k-1}} \sqrt{\mathbf{1}-\alpha_{2 k}^{2}} \sqrt{\mathbf{1}+\alpha_{2 k+1}} .
\end{aligned}
$$

The result follows.
As in $[127,168]$, one can also use this to get Geronimus relations for the second Szegő map.

## 5 Regular MOPRL

The theory of regular (scalar) OPs was developed by Stahl-Totik [180] generalizing a definition of Ullman [190] for $[-1,1]$. (See Simon [171] for a review and more about the history.) Here we develop the basics for MOPRL; it is not hard to do the same for MOPUC.

### 5.1 Upper Bound and Definition

Theorem 5.1. Let $d \mu$ be a non-trivial $l \times l$ matrix-valued measure on $\mathbb{R}$ with $E=\operatorname{supp}(d \mu)$ compact. Then (with $C(E)=$ logarithmic capacity of $E$ )

$$
\begin{equation*}
\limsup _{n \rightarrow \infty}\left|\operatorname{det}\left(A_{1} \cdots A_{n}\right)\right|^{1 / n} \leq C(E)^{l} . \tag{5.1}
\end{equation*}
$$

Remarks. 1. $\left|\operatorname{det}\left(A_{1} \cdots A_{n}\right)\right|$ is constant over equivalent Jacobi parameters.
2. For the scalar case, this is a result of Widom [195] (it might be older) whose proof extends to the matrix case.

Proof. Let $T_{n}$ be the Chebyshev polynomials for $E$ (see [171, Appendix B] for a definition) and let $T_{n}^{(l)}$ be $T_{n} \otimes \mathbf{1}$, that is, the $l \times l$ matrix polynomial obtained by multiplying $T_{n}(x)$ by $\mathbf{1}$. $T_{n}^{(l)}$ is monic so, by (2.12) and (2.34),

$$
\begin{aligned}
\left|\operatorname{det}\left(A_{1} \cdots A_{n}\right)\right|^{1 / n} & \leq\left|\operatorname{det}\left(\int\left|T_{n}^{(l)}(x)\right|^{2} d \mu(x)\right)\right|^{1 / 2 n} \\
& \leq \sup _{n}\left|T_{n}(x)\right|^{l / n}
\end{aligned}
$$

By a theorem of Szegő $[183], \sup _{n}\left|T_{n}(x)\right|^{1 / n} \rightarrow C(E)$, so (5.1) follows.
Definition. Let $d \mu$ be a non-trivial $l \times l$ matrix-valued measure with $E=\operatorname{supp}(d \mu)$ compact. We say $\mu$ is regular if equality holds in (5.1).

### 5.2 Density of Zeros

The following is a simple extension of the scalar results (see [74] or [171, Sect. 2]):

Theorem 5.2. Let $d \mu$ be a regular measure with $E=\operatorname{supp}(d \mu)$. Let $d \nu_{n}$ be the zero counting measure for $\operatorname{det}\left(p_{n}^{L}(x)\right)$, that is, if $\left\{x_{j}^{(n)}\right\}_{j=1}^{n l}$ are the zeros of this determinant (counting degenerate zeros multiply), then

$$
\begin{equation*}
d \nu_{n}=\frac{1}{n l} \sum_{j=1}^{n l} \delta_{x_{j}^{(n)}} . \tag{5.2}
\end{equation*}
$$

Then $d \nu_{n}$ converges weakly to $d \rho_{E}$, the equilibrium measure for $E$.
Remark. For a discussion of capacity, equilibrium measure, quasi-every (q.e.), etc., see [171, Appendix A] and [115, 136, 158].

Proof. By (2.80) and (2.93),

$$
\begin{equation*}
\left|\operatorname{det}\left(p_{n}^{R}(x)\right)\right| \geq\left(\frac{d}{D}\right)^{l}\left(1+\left(\frac{d}{D}\right)^{2}\right)^{(n-1) l / 2} \tag{5.3}
\end{equation*}
$$

so, in particular,

$$
\begin{equation*}
\underset{n}{\lim \inf }\left|\operatorname{det}\left(p_{n}^{R}(x)\right)\right|^{1 / n l} \geq\left(1+\left(\frac{d}{D}\right)^{2}\right)^{1 / 2} \geq 1 \tag{5.4}
\end{equation*}
$$

But

$$
\begin{equation*}
\left|\operatorname{det}\left(p_{n}^{R}(x)\right)\right|=\left|\operatorname{det}\left(A_{1} \cdots A_{n}\right)\right|^{-1} \exp \left(-n l \Phi_{\nu_{n}}(x)\right) \tag{5.5}
\end{equation*}
$$

where $\Phi_{\nu}$ is the potential of the measure $\nu$. Let $\nu_{\infty}$ be a limit point of $\nu_{n}$ and use equality in (5.1) and (5.4) to conclude, for $x \notin \operatorname{cvh}(E)$,

$$
\exp \left(-\Phi_{\nu_{\infty}}(x)\right) \geq C(E)
$$

which, as in the proof of Theorem 2.4 of [171], implies that $\nu_{\infty}=\rho_{e}$.
The analogue of the almost converse of this last theorem has an extra subtlety relative to the scalar case:

Theorem 5.3. Let $d \mu$ be a non-trivial $l \times l$ matrix-valued measure on $\mathbb{R}$ with $E=\operatorname{supp}(d \mu)$ compact. If $d \nu_{n} \rightarrow d \rho_{E}$, then either $\mu$ is regular or else, with $d \mu=M(x) d \mu_{\operatorname{tr}}(x)$, there is a set $S$ of capacity zero, so $\operatorname{det}(M(x))=0$ for $d \mu_{\text {tr }}$-a.e. $x \notin S$.

Remark. By taking direct sums of regular point measures, it is easy to find regular measures where $\operatorname{det}(M(x))=0$ for $d \mu_{\mathrm{tr}^{-}}$-a.e. $x$.

Proof. For a.e. $x$ with $\operatorname{det}(M(x)) \neq 0$, we have (see Lemma 5.7 below)

$$
p_{n}^{R}(x) \leq C(n+1) \mathbf{1}
$$

The theorem then follows from the proof of Theorem 2.5 of [171].

### 5.3 General Asymptotics

The following generalizes Theorem 1.10 of [171] from OPRL to MOPRL-it is the matrix analogue of a basic result of Stahl-Totik [180]. Its proof is essentially the same as in [171]. By $\sigma_{\text {ess }}(\mu)$, we mean the essential spectrum of the block Jacobi matrix associated to $\mu$.
Theorem 5.4. Let $E \subset \mathbb{R}$ be compact and let $\mu$ be an $l \times l$ matrix-valued measure of compact support with $\sigma_{\mathrm{ess}}(\mu)=E$. Then the following are equivalent:
(i) $\mu$ is regular, that is, $\lim _{n \rightarrow \infty}\left|\operatorname{det}\left(A_{1} \cdots A_{n}\right)\right|^{1 / n}=C(E)^{l}$.
(ii) For all $z$ in $\mathbb{C}$, uniformly on compacts,

$$
\begin{equation*}
\lim \sup \left|\operatorname{det}\left(p_{n}^{R}(z)\right)\right|^{1 / n} \leq e^{G_{E}(z)} \tag{5.6}
\end{equation*}
$$

(iii) For q.e. $z$ in $E$, we have

$$
\begin{equation*}
\lim \sup \left|\operatorname{det}\left(p_{n}^{R}(z)\right)\right|^{1 / n} \leq 1 \tag{5.7}
\end{equation*}
$$

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

$$
\begin{equation*}
\lim _{n \rightarrow \infty}\left|\operatorname{det}\left(p_{n}^{R}(z)\right)\right|^{1 / n}=e^{G_{E}(z)} . \tag{5.8}
\end{equation*}
$$

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

$$
\begin{equation*}
\limsup _{n \rightarrow \infty}\left|\operatorname{det}\left(p_{n}^{R}(z)\right)\right|^{1 / n}=1 \tag{5.9}
\end{equation*}
$$

Remarks. 1. $G_{E}$, the potential theorists' Green's function for $E$, is defined by $G_{E}(z)=$ $-\log (C(E))-\Phi_{\rho_{E}}(z)$.
2. There is missing here one condition from Theorem 1.10 of [171] involving general polynomials. Since the determinant of a sum can be much larger than the sum of the determinants, it is not obvious how to extend this result.

### 5.4 Weak Convergence of the CD Kernel and Consequences

The results of Simon in [174] extend to the matrix case. The basic result is:
Theorem 5.5. The measures $d \nu_{n}$ and $\frac{1}{(n+1) l} \operatorname{Tr}\left(K_{n}(x, x)\right) d \mu(x)$ have the same weak limits. In particular, if $d \mu$ is regular,

$$
\begin{equation*}
\frac{1}{(n+1) l} \operatorname{Tr}\left(K_{n}(x, x) d \mu(x)\right) \xrightarrow{w} d \rho_{E} . \tag{5.10}
\end{equation*}
$$

As in [174], $\left(\pi_{n} M_{x} \pi_{n}\right)^{j}$ and $\left(\pi_{n} M_{x}^{j} \pi_{n}\right)$ have a difference of traces which is bounded as $n \rightarrow \infty$, and this implies the result. Once one has this, combining it with Theorem 2.20 leads to:
Theorem 5.6. Let $I=[a, b] \subset E \subset \mathbb{R}$ with $E$ compact. Let $\sigma_{\text {ess }}(d \mu)=E$ for an $l \times l$ matrix-valued measure, and suppose $d \mu$ is regular for $E$ and

$$
\begin{equation*}
d \mu=W(x) d x+d \mu_{\mathrm{s}} \tag{5.11}
\end{equation*}
$$

where $d \mu_{\mathrm{s}}$ is singular and $\operatorname{det}(W(x))>0$ for a.e. $x \in I$. Then,

$$
\begin{align*}
& \lim _{n \rightarrow \infty} \int_{I} p_{n}^{R}(x)^{\dagger} d \mu_{\mathrm{s}}(x) p_{n}^{R}(x) \rightarrow \mathbf{0}  \tag{1}\\
& \int_{I}\left\|\frac{1}{n+1} \sum_{j=0}^{n} p_{j}^{R}(x)^{\dagger} W(x) p_{j}(x)-\rho_{E}(x) \mathbf{1}\right\| d x \rightarrow 0 \tag{2}
\end{align*}
$$

### 5.5 Widom's Theorem

Lemma 5.7. Let $d \mu$ be an $l \times l$ matrix-valued measure supported on a compact $E \subset \mathbb{R}$ and let $d \eta$ be a scalar measure on $\mathbb{R}$ so

$$
\begin{equation*}
d \mu(x)=W(x) d \eta(x)+d \mu_{\mathrm{s}}(x) \tag{5.14}
\end{equation*}
$$

where $d \mu_{\mathrm{s}}$ is $d \eta$ singular. Suppose for $d \eta$-a.e. $x$,

$$
\begin{equation*}
\operatorname{det}(W(x))>0 . \tag{5.15}
\end{equation*}
$$

Then for $d \eta$-a.e. $x$, there is a positive real function $C(x)$ so that

$$
\begin{equation*}
\left\|p_{n}^{R}(x)\right\| \leq C(x)(n+1) \mathbf{1} \tag{5.16}
\end{equation*}
$$

In particular,

$$
\begin{equation*}
\left|\operatorname{det}\left(p_{n}^{R}(x)\right)\right| \leq C(x)^{l}(n+1)^{l} \tag{5.17}
\end{equation*}
$$

Proof. Since $\left\|p_{n}^{R}\right\|_{R}^{2}=1$, we have

$$
\begin{equation*}
\sum_{n=0}^{\infty}(n+1)^{-2}\left\|p_{n}^{R}\right\|_{R}^{2}<\infty \tag{5.18}
\end{equation*}
$$

so

$$
\sum_{n=0}^{\infty}(n+1)^{-2} \operatorname{Tr}\left(p_{n}^{R}(x)^{\dagger} W(x) p_{n}^{R}(x)\right)<\infty
$$

for $d \eta$-a.e. $x$. Since (5.15) holds, for a.e. $x$,

$$
W(x) \geq b(x) \mathbf{1}
$$

for some scalar function $b(x)$. Thus, for a.e. $x$,

$$
\sum_{n=0}^{\infty}(n+1)^{-2} \operatorname{Tr}\left(p_{n}^{R}(x)^{\dagger} p_{n}^{R}(x)\right) \leq C(x)^{2} .
$$

Since $\|A\|^{2} \leq \operatorname{Tr}\left(A^{\dagger} A\right)$, we find (5.16), which in turn implies (5.17).
This lemma replaces Lemma 4.1 of [171] and then the proof there of Theorem 1.12 extends to give (a matrix version of the theorem of Widom [195]):

Theorem 5.8. Let $d \mu$ be an $l \times l$ matrix-valued measure with $\sigma_{\mathrm{ess}}(d \mu)=E \subset \mathbb{R}$ compact. Suppose

$$
\begin{equation*}
d \mu(x)=W(x) d \rho_{E}(x)+d \mu_{\mathrm{s}}(x) \tag{5.19}
\end{equation*}
$$

with $d \mu_{\mathrm{s}}$ singular with respect to $d \rho_{E}$. Suppose for $d \rho_{E}$-a.e. $x, \operatorname{det}(W(x))>0$. Then $\mu$ is regular.

### 5.6 A Conjecture

We end our discussion of regular MOPRL with a conjecture - an analog of a theorem of Stahl-Totik [180]; see also Theorem 1.13 of [171] for a proof and references. We expect the key will be some kind of matrix Remez inequality. For direct sums, this conjecture follows from Theorem 1.13 of [171].

Conjecture 5.9. Let $E$ be a finite union of disjoint closed intervals in $\mathbb{R}$. Suppose $\mu$ is an $l \times l$ matrix-valued measure on $\mathbb{R}$ with $\sigma_{\text {ess }}(d \mu)=E$. For each $\eta>0$ and $m=1,2, \ldots$, define

$$
\begin{equation*}
S_{m, \eta}=\left\{x: \mu\left(\left[x-\frac{1}{m}, x+\frac{1}{m}\right]\right) \geq e^{-\eta m} \mathbf{1}\right\} . \tag{5.20}
\end{equation*}
$$

Suppose that for each $\eta$ (with $|\cdot|=$ Lebesgue measure)

$$
\lim _{m \rightarrow \infty}\left|E \backslash S_{m, \eta}\right|=0
$$

Then $\mu$ is regular.
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