# Operator monotone functions and Löwner functions of several variables 

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#### Abstract

We prove generalizations of Löwner's results on matrix monotone functions to several variables. We give a characterization of when a function of $d$ variables is locally monotone on $d$-tuples of commuting self-adjoint $n$-by- $n$ matrices. We prove a generalization to several variables of Nevanlinna's theorem describing analytic functions that map the upper half-plane to itself and satisfy a growth condition. We use this to characterize all rational functions of two variables that are operator monotone.


## 1. Introduction

In 1934, K. Löwner published a very influential paper [25] studying functions on an open interval $E \subseteq \mathbb{R}$ that are matrix monotone; i.e., functions $f$ with the property that whenever $S$ and $T$ are self-adjoint matrices whose spectra are in $E$, then

$$
\begin{equation*}
S \leq T \quad \Rightarrow \quad f(S) \leq f(T) \tag{1.1}
\end{equation*}
$$

This property is equivalent (see Section 1.3) to being locally matrix monotone; i.e., if $S(t)$ is a $C^{1}$-arc of self-adjoint matrices with $\sigma(S(t)) \subset E$, then

$$
\begin{equation*}
S^{\prime}(t) \geq 0 \quad \Rightarrow \quad \frac{d}{d t} f(S(t)) \geq 0 \tag{1.2}
\end{equation*}
$$

Roughly speaking, Löwner showed that if one fixes a dimension $n$ and wants (1.1) or (1.2) to hold for $n$-by- $n$ self-adjoint matrices, then certain matrices derived from the values of $f$ must all be positive semi-definite. As $n$ increases, the conditions become more restrictive. In the limit as $n \rightarrow \infty$ (equivalently, if one passes to self-adjoint operators on an infinite-dimensional Hilbert space), then a necessary and sufficient condition is that the function $f$ must have an analytic continuation to a function $F$ that maps the upper half-plane to itself.

[^0]The goal of this paper is to extend the above notions to several variables. In particular, we want to study functions of $d$ variables applied to $d$-tuples of commuting self-adjoint operators. Given two $d$-tuples $S=\left(S^{1}, \ldots, S^{d}\right)$ and $T=\left(T^{1}, \ldots, T^{d}\right)$, we shall say that $S \leq T$ if and only if $S^{r} \leq T^{r}$ for every $1 \leq r \leq d$. We want to study functions that satisfy (1.1) or (1.2) for $d$-tuples.

Before we can describe our results, we must first give a more detailed description of the one-dimensional case. We recommend the book [14] by W. Donoghue for a well-written account from a modern perspective. See also the paper [31].

Note that there is another approach to extending Löwner's results to several variables where the operators $S^{1}, \ldots, S^{d}$ act on different spaces $\mathcal{H}^{1}, \ldots, \mathcal{H}^{d}$, and $f(S)$ is interpreted to act on $\mathcal{H}^{1} \otimes \cdots \otimes \mathcal{H}^{d}$. We refer the reader to the papers [16], [36], [23] and references therein.
1.1. Dimension one. Let $E$ be an open set in $\mathbb{R}$, and let $n \geq 2$ be a natural number. The Löwner class $\mathcal{L}_{n}^{1}(E)$ is the set of $C^{1}$-functions $f: E \rightarrow \mathbb{R}$ with the property that, whenever $\left\{x_{1}, \ldots, x_{n}\right\}$ is a set of $n$ distinct points in $E$, then the matrix $A$, defined by

$$
A_{i j}= \begin{cases}\frac{f\left(x_{j}\right)-f\left(x_{i}\right)}{x_{j}-x_{i}} & \text { if } i \neq j, \\ f^{\prime}\left(x_{i}\right) & \text { if } i=j,\end{cases}
$$

is positive semi-definite.
We shall let $M_{n}$ denote the $n$-by- $n$ complex matrices, $\mathrm{SAM}_{n}$ the selfadjoint $n$-by- $n$ matrices, and SA the bounded self-adjoint operators on an infinite-dimensional separable Hilbert space.

Definition 1.1. A function $f$ is locally $n$-matrix monotone on the open set $E \subset \mathbb{R}$ if, whenever $S$ is in $\mathrm{SAM}_{n}$ with $\sigma(S)$ consisting of $n$ distinct points in $E$, and $S(t)$ is a $C^{1}$-curve in $\mathrm{SAM}_{n}$ with $S(0)=S$ and $\left.\frac{d}{d t} S(t)\right|_{t=0} \geq 0$, then $\left.\frac{d}{d t} f(S(t))\right|_{t=0} \geq 0$.

Remark 1.2. This definition is slightly different from the one in the first paragraph, where the eigenvalues were not required to be distinct. We use this definition to be consistent with the multivariable Definition 1.7 below. However, using formula (6.6.31) in [18] for $\frac{d}{d t} f(S(t))$, it is easy to show that in the one variable case the two different definitions are equivalent. ${ }^{1}$

[^1]where $\Delta$ means the matrix of divided differences and o denotes the Schur product.

We shall say that $f$ is $n$-matrix monotone on $E$, or $M_{n}$-monotone, if, whenever $S$ and $T$ are in $\mathrm{SAM}_{n}$ and all their eigenvalues lie in $E$, then (1.1) holds. To emphasize the difference from locally monotone, we shall also call $n$-matrix monotone functions globally $M_{n}$-monotone. Replacing $\mathrm{SAM}_{n}$ by SA, we get the definitions of locally operator monotone and operator monotone.

Theorem 1.3 (Löwner). Let $E \subseteq \mathbb{R}$ be open, and let $f \in C^{1}(E)$. Then $f$ is locally n-matrix monotone on $E$ if and only if $f$ is in $\mathcal{L}_{n}^{1}(E)$.

We shall use $\Pi$ to denote the upper half-plane, $\{z \in \mathbb{C}: \operatorname{Im} z>0\}$.
Definition 1.4. Let $E \subseteq \mathbb{R}$ be open. The Pick class on $E$, denoted $\mathcal{P}(E)$, is the set of real-valued functions $f$ on $E$ for which there exists an analytic function $F$ that maps $\Pi$ to $\bar{\Pi}$, the closure of $\Pi$, and such that $F$ extends analytically across $E$ and $^{2}$

$$
\lim _{y \searrow 0} F(x+i y)=f(x) \quad \forall x \in E .
$$

Theorem 1.5 (Löwner). Let $E \subseteq \mathbb{R}$ be open, and let $f \in C^{1}(E)$. The following are equivalent:
(i) The function $f$ is locally operator monotone on $E$.
(ii) The function $f$ is in $\mathcal{L}_{n}^{1}(E)$ for all $n$.
(iii) The function $f$ is in $\mathcal{P}(E)$.
1.2. Dimension $d \geq 2$ : Local results. We shall let $\operatorname{CSAM}_{n}^{d}$ denote the set of $d$-tuples of commuting self-adjoint $n$-by- $n$ matrices and CSA ${ }^{d}$ be the set of $d$-tuples of commuting self-adjoint bounded operators. If $S$ is a commuting $d$-tuple of self-adjoint operators acting on the Hilbert space $\mathcal{H}$ and $f$ is a realvalued continuous (indeed, measurable) function on the joint spectrum of $S$ in $\mathbb{R}^{d}$, then $f(S)$ is a well-defined self-adjoint operator on $\mathcal{H}$.

Definition 1.6. Let $E$ be an open set in $\mathbb{R}^{d}$ and let $f$ be a real-valued $C^{1}$-function on $E$. Say $f$ is locally operator monotone on $E$ if, whenever $S$ is in $\mathrm{CSA}^{d}$ with $\sigma(S) \subset E$, and $S(t)$ is a $C^{1}$-curve in $\mathrm{CSA}^{d}$ with $S(0)=S$ and $\left.\frac{d}{d t} S(t)\right|_{t=0} \geq 0$, then $\left.\frac{d}{d t} f(S(t))\right|_{t=0}$ exists and is $\geq 0$.

We shall not concern ourselves in this paper with what conditions on $f$ guarantee that $f(S(t))$ is differentiable; for these see, e.g., [28].

Definition 1.7. Let $E$ be an open set in $\mathbb{R}^{d}$ and $f$ be a real-valued $C^{1}$ function on $E$. We say $f$ is locally $M_{n}$-monotone on $E$ if, whenever $S$ is in $\operatorname{CSAM}_{n}^{d}$ with $\sigma(S)=\left\{x_{1}, \ldots, x_{n}\right\}$ consisting of $n$ distinct points in $E$, and $S(t)$ is a $C^{1}$-curve in $\operatorname{CSAM}_{n}^{d}$ with $S(0)=S$ and $\left.\frac{d}{d t} S(t)\right|_{t=0} \geq 0$, then $\left.\frac{d}{d t} f(S(t))\right|_{t=0}$ exists and is $\geq 0$.

[^2]We define the Löwner classes in $d$ variables, $\mathcal{L}_{n}^{d}(E)$, by
Definition 1.8. Let $E$ be an open subset of $\mathbb{R}^{d}$. The set $\mathcal{L}_{n}^{d}(E)$ consists of all real-valued $C^{1}$-functions on $E$ that have the following property: whenever $\left\{x_{1}, \ldots, x_{n}\right\}$ are $n$ distinct points in $E$, there exist positive semi-definite $n$-by- $n$ matrices $A^{1}, \ldots, A^{d}$ so that

$$
A^{r}(i, i)=\left.\frac{\partial f}{\partial x^{r}}\right|_{x_{i}}
$$

and

$$
f\left(x_{j}\right)-f\left(x_{i}\right)=\sum_{r=1}^{d}\left(x_{j}^{r}-x_{i}^{r}\right) A^{r}(i, j) \quad \forall 1 \leq i, j \leq n .
$$

Theorem 1.9. Let $E$ be an open set in $\mathbb{R}^{d}$ and $f$ be a real-valued $C^{1}$-function on $E$. Then $f$ is locally $M_{n}$-monotone if and only if $f$ is in $\mathcal{L}_{n}^{d}(E)$.

In generalizations of Theorem 1.5, there turns out to be a difference between the case $d=2$ and $d>2$.

Definition 1.10. The Löwner class, $\mathcal{L}^{d}$, is the set of functions $F: \Pi^{d} \rightarrow \bar{\Pi}$ with the property that there exist $d$ positive semi-definite kernel functions $A^{r}, 1 \leq r \leq d$, on $\Pi^{d}$ such that

$$
F(z)-\overline{F(w)}=\left(z^{1}-\bar{w}^{1}\right) A^{1}(z, w)+\cdots+\left(z^{d}-\bar{w}^{d}\right) A^{d}(z, w)
$$

When $d=1$ or 2 , the Löwner class coincides with the set of all analytic functions from $\Pi^{d}$ to $\bar{\Pi}$, but for $d \geq 3$ it is a proper subset (see Remark 2.18).

Definition 1.11. Let $E \subseteq \mathbb{R}^{d}$ be open. The class $\mathcal{L}(E)$ is the set of realvalued functions $f$ on $E$ for which there exists an analytic function $F$ in $\mathcal{L}^{d}$ such that $F$ extends analytically across $E$ and

$$
\lim _{y \searrow 0} F\left(x^{1}+i y, \ldots, x^{d}+i y\right)=f\left(x^{1}, \ldots, x^{d}\right) \quad \forall x \in E .
$$

We prove the following result as Theorem 8.1.
Theorem 1.12. Let $E$ be an open set in $\mathbb{R}^{d}$ and $f$ be a real-valued $C^{1}$-function on $E$. The following are equivalent:
(i) The function $f$ is locally operator monotone on $E$.
(ii) The function $f$ is in $\mathcal{L}_{n}^{d}(E)$ for all $n$.
(iii) The function $f$ is in $\mathcal{L}(E)$.
1.3. Local to global. In one variable, provided $E$ is an interval, local monotonicity implies global monotonicity immediately. Indeed, suppose $S \leq T$, and let $S(t)=(1-t) S+t T$. Then $S^{\prime}(t)=T-S \geq 0$, so

$$
\begin{equation*}
f(T)-f(S)=\int_{0}^{1} \frac{d}{d t} f(S(t)) d t \geq 0 \tag{1.13}
\end{equation*}
$$

If $E$ is not convex, this argument fails. Indeed, the function $-1 / x$ is locally $n$-matrix monotone on $\mathbb{R} \backslash\{0\}$ for all $n$; but it is only globally monotone on sets that lie entirely on one side of 0 . A result of Chandler [11] says that functions that are globally operator monotone on a set $E$ always extend to be globally monotone on the convex hull of $E$.

For intervals, (1.13) shows that the word "locally" can be dropped in both Theorems 1.3 and 1.5. One problem in going to several variables is that this simple argument no longer works, because one may not be able to connect $S$ and $T$ by a path of commuting $d$-tuples. Indeed, the following example shows that there need not be any commuting tuples between two given ones.

Example 1.14. Let $S$ and $T$ be pairs in $\mathrm{CSAM}_{2}^{2}$ given by

$$
\begin{aligned}
& S=\left(\left(\begin{array}{ll}
0 & 0 \\
0 & 5
\end{array}\right),\left(\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right)\right) \\
& T=\left(\left(\begin{array}{ll}
4 & 2 \\
2 & 6
\end{array}\right),\left(\begin{array}{ll}
2 & 2 \\
2 & 4
\end{array}\right)\right) .
\end{aligned}
$$

If $R$ is in $\mathrm{CSAM}_{2}^{2}$ and $S \leq R \leq T$, it can be shown that either $R=S$ or $R=T$.

We have been unable to resolve the question of whether the $n$-matrix monotone functions on a connected open set $E$ are a proper subset of the locally $n$-matrix monotone functions on $E$. However, as $n$ tends to infinity and we pass to locally operator monotone functions, analyticity enters the picture, and makes the problem more tractable; see Section 1.5.
1.4. The Nevanlinna representation. To prove (iii) $\Rightarrow$ (i) in Theorem 1.5, one must understand analytic functions that map the upper half-plane to itself. A key fact is a characterization due to R. Nevanlinna [26] which says that, provided they have some regularity at infinity, they are all Cauchy transforms of measures on the line.

Theorem 1.15 (Nevanlinna). If $F: \Pi \rightarrow \Pi$ is analytic and satisfies

$$
\limsup _{y \rightarrow \infty} y|F(i y)-C|<\infty
$$

for some $C \in \mathbb{R}$, then there exists a unique finite positive Borel measure $\nu$ on $\mathbb{R}$ so that

$$
\begin{equation*}
F(z)=C+\int \frac{d \nu(t)}{t-z} \tag{1.16}
\end{equation*}
$$

Nevanlinna's theorem was used by M. Stone to prove the spectral theorem [37], but one can adopt the reverse viewpoint and write (1.16) in terms of
the resolvent of a self-adjoint. Indeed, let $X$ be the self-adjoint operator of multiplication by the independent variable on $L^{2}(\nu)$, and let $v$ be the vector in $L^{2}(\nu)$ that is 1 almost everywhere. Then (1.16) can be rewritten as

$$
\begin{equation*}
F(z)=C+\left\langle(X-z)^{-1} v, v\right\rangle . \tag{1.17}
\end{equation*}
$$

This representation turns out to be useful in studying operator monotonicity, because then

$$
\begin{equation*}
F(S)=C I+R_{v}^{*}(I \otimes X-S \otimes I)^{-1} R_{v} \tag{1.18}
\end{equation*}
$$

where $R_{v}: \mathcal{H} \rightarrow \mathcal{H} \otimes \mathcal{M}$ is given by $R_{v}: \xi \mapsto \xi \otimes v$.
There is a several variable analogue of Theorem 1.15. It may require first perturbing $F$.

Definition 1.19. For each real number $t$, define

$$
\rho_{t}(z)=\frac{z+t}{1-t z} .
$$

For $F \in \mathcal{L}^{d}$, define $F_{t}$ by

$$
F_{t}\left(z^{1}, \ldots, z^{d}\right):=\rho_{t} \circ F\left(\rho_{t}\left(z^{1}\right), \ldots, \rho_{t}\left(z^{d}\right)\right)
$$

The following theorem follows from Theorem 6.33. We shall say that a function $F$ on $\Pi^{d}$ is analytic on a neighborhood of infinity if the function $F\left(1 / z^{1}, \ldots, 1 / z^{d}\right)$ extends to be analytic on a neighborhood of the origin. In Theorem 6.33, a weaker assumption is placed on $F$ than being analytic in a neighborhood of infinity.

Theorem 1.20. Let $F$ be in $\mathcal{L}^{d}$, and assume that $F$ is analytic in a neighborhood of infinity. Then for all sufficiently small $t$, with at most countably many exceptions, the function $F_{t}$ has the following representation. There is a Hilbert space $\mathcal{M}$, a densely defined self-adjoint operator $X$ on $\mathcal{M}$, a vector $v$ in $\mathcal{M}$, a real constant $C$, and d orthogonal projections $P^{1}, \ldots, P^{d}$ with $\sum_{r=1}^{d} P^{r}=I_{\mathcal{M}}$ so that

$$
\begin{equation*}
F_{t}(z)=C+\left\langle\left(X-\sum_{r=1}^{d} z^{r} P^{r}\right)^{-1} v, v\right\rangle . \tag{1.21}
\end{equation*}
$$

1.5. Dimension $d \geq 2$ : Global operator monotonicity. If $E$ is an open set in $\mathbb{R}^{d}$, we shall say that a real-valued function $f$ defined on $E$ is globally operator monotone, or just operator monotone for short, if, for every $n$, whenever $S$ and $T$ are in $\operatorname{CSAM}_{n}^{d}$, with $S \leq T$, and the joint eigenvalues of both $S$ and $T$ lie in $E$, then

$$
f(S) \leq f(T)
$$

Using the representation (1.21), we can prove results on (global) operator monotonicity. With notation as in Theorem 1.20 , let us say that the $\mu$-resolvent
of $X$ is the set of points

$$
\left\{\left(z^{1}, \ldots, z^{d}\right) \in \mathbb{C}^{d}: X-\sum_{r=1}^{d} z^{r} P^{r} \text { has a bounded inverse }\right\}
$$

We prove the following result as Theorem 9.2.
Theorem 1.22. Let $X$ be a densely-defined self-adjoint operator on a Hilbert space $\mathcal{M}$, let $v$ be a vector in $\mathcal{M}$, let $C$ be a real constant, and let $P^{1}, \ldots, P^{d}$ be projections with orthogonal ranges that sum to the identity. Let $F$ be given by

$$
F(z)=C+\left\langle\left(X-\sum_{r=1}^{d} z^{r} P^{r}\right)^{-1} v, v\right\rangle .
$$

Let $E$ be an open box in $\mathbb{R}^{d}$ that is in the $\mu$-resolvent of $X$. Then $F$ is globally operator monotone on $E$.

As an application, we can give a complete characterization of the rational functions of two variables that are operator monotone on rectangles. This is Theorem 9.6.

Theorem 1.23. Let $F$ be a rational function of two variables. Let $\Gamma$ be the zero-set of the denominator of $F$. Assume $F$ is real-valued on $\mathbb{R}^{2} \backslash \Gamma$. Let $E$ be an open rectangle in $\mathbb{R}^{2} \backslash \Gamma$. Then $F$ is globally operator monotone on $E$ if and only if $F$ is in $\mathcal{L}(E)$; that is, if and only if $F$ is the restriction to $E$ of an analytic function from $\Pi^{2}$ to $\bar{\Pi}$ that extends analytically across $E$.

## 2. Some notation

Let $\mathbb{D}$ be the unit disk in the complex plane, $\Pi$ the upper half-plane $\{z: \operatorname{Im}(z)>0\}$, and $\mathbb{H}$ the right half-plane $\{z: \operatorname{Re}(z)>0\}$. Let $\alpha$ denote the linear fractional map

$$
\begin{equation*}
\alpha(\lambda)=i \frac{1+\lambda}{1-\lambda} \tag{2.1}
\end{equation*}
$$

that maps $\mathbb{D}$ to $\Pi$, and let

$$
\begin{equation*}
\beta(z)=\frac{z-i}{z+i} \tag{2.2}
\end{equation*}
$$

denote its inverse.
We shall let $d$ denote the number of variables. If $z$ is a point in $\Pi^{d}$, we shall use $z^{1}, \ldots, z^{d}$ to denote its components; likewise, $\lambda=\left(\lambda^{1}, \ldots, \lambda^{d}\right)$ will be a point in $\mathbb{D}^{d}$. We shall write $S=\left(S^{1}, \ldots, S^{d}\right)$ for a $d$-tuple of matrices or operators and use $\|S\|$ for $\max _{1 \leq r \leq d}\left\|S^{r}\right\|$. We shall also use $\alpha$ and $\beta$ to denote the maps from $\mathbb{D}^{d}$ to $\Pi^{d}$ and back again that are defined by applying $\alpha$ and $\beta$ coordinate-wise.

A kernel on a set $E$ is a map $K: E \times E \rightarrow \mathbb{C}$ with the property that for every finite set $\left\{\lambda_{1}, \ldots, \lambda_{N}\right\}$ of distinct points in $E$, the matrix $\left[K\left(\lambda_{j}, \lambda_{i}\right)\right]$ is positive semi-definite.

Definition 2.3. The Pick class, $\mathcal{P}^{d}$, is the set of analytic functions $F$ : $\Pi^{d} \rightarrow \bar{\Pi}$.

Definition 2.4. The Schur class, $\mathcal{S}^{d}$, is the set of analytic functions $\varphi$ : $\mathbb{D}^{d} \rightarrow \overline{\mathbb{D}}$.

Definition 2.5. The Carathéodory class, $\mathcal{C}^{d}$, is the set of analytic functions $\psi: \mathbb{D}^{d} \rightarrow \overline{\mathbb{H}}$.

Definition 2.6. The Löwner class, $\mathcal{L}^{d}$, is the set of functions $F: \Pi^{d} \rightarrow \bar{\Pi}$ with the property that there exist $d$ kernel functions $A^{r}, 1 \leq r \leq d$ on $\Pi^{d}$ such that

$$
\begin{equation*}
F(z)-\overline{F(w)}=\left(z^{1}-\bar{w}^{1}\right) A^{1}(z, w)+\cdots+\left(z^{d}-\bar{w}^{d}\right) A^{d}(z, w) \tag{2.7}
\end{equation*}
$$

Definition 2.8. The Schur-Agler class, $\mathcal{A}^{d}$, is the set of functions $\varphi: \mathbb{D}^{d} \rightarrow \overline{\mathbb{D}}$ with the property that there exist $d$ kernel functions $B^{r}, 1 \leq r \leq d$ on $\mathbb{D}^{d}$ such that

$$
\begin{equation*}
1-\varphi(\lambda) \overline{\varphi(\mu)}=\left(1-\lambda^{1} \bar{\mu}^{1}\right) B^{1}(\lambda, \mu)+\cdots+\left(1-\lambda^{d} \bar{\mu}^{d}\right) B^{d}(\lambda, \mu) . \tag{2.9}
\end{equation*}
$$

When the dimension is clear, we shall drop the superscript $d$.
Remark 2.10. If we exclude the constant function 1 from $\mathcal{S}$, we have the identification

$$
\begin{equation*}
F \in \mathcal{P} \Longleftrightarrow \beta \circ F \circ \alpha \in \mathcal{S} \Longleftrightarrow-i F \circ \alpha \in \mathcal{C} . \tag{2.11}
\end{equation*}
$$

Moreover, we also have (again excluding the constant function 1)

$$
\begin{equation*}
F \in \mathcal{L} \Longleftrightarrow \beta \circ F \circ \alpha \in \mathcal{A} \tag{2.12}
\end{equation*}
$$

(see Lemma 2.13). As all our results are trivial for constant functions, we shall use (2.11) and (2.12) without explicitly mentioning the exclusion of the constant function 1.

The following change of variables formula is in [20]. A function is in $\mathcal{A}^{d}$ if and only if it is analytic and maps $d$-tuples of commuting strict contractions to contractions; a function is in $\mathcal{L}^{d}$ if and only if it is analytic and maps $d$-tuples of commuting operators with strictly positive imaginary parts ${ }^{3}$ to operators with positive imaginary parts.

[^3]Lemma 2.13. The function $F: \Pi^{d} \rightarrow \mathbb{C}$ is in the Löwner class if and only if $\varphi:=\beta \circ F \circ \alpha$ is in the Schur-Agler class $\mathcal{A}^{d}$.

Proof. Define $\varphi=\beta \circ F \circ \alpha$. Then $\varphi$ is in $\mathcal{A}^{d}$ if and only if there are kernels $B^{r}$ on $\mathbb{D}$ such that

$$
\begin{equation*}
1-\varphi(\lambda) \overline{\varphi(\mu)}=\sum_{r=1}^{d}\left(1-\lambda^{r} \bar{\mu}^{r}\right) B^{r}(\lambda, \mu) . \tag{2.14}
\end{equation*}
$$

When $z=\alpha(\lambda)$ and $w=\alpha(\mu),(2.14)$ becomes

$$
\begin{equation*}
1-\beta \circ F(z) \overline{\beta \circ F(w)}=\sum_{r=1}^{d}\left(1-\left[\frac{z^{r}-i}{z^{r}+i}\right] \overline{\left[\frac{w^{r}-i}{w^{r}+i}\right]}\right) B^{r}(\beta(z), \beta(w)) . \tag{2.15}
\end{equation*}
$$

Rearranging (2.15), we get

$$
\begin{equation*}
F(z)-\overline{F(w)}=\sum_{r=1}^{d}\left(z^{r}-\bar{w}^{r}\right) \frac{F(z)+i}{z^{r}+i} \frac{\overline{F(w)}-i}{\overline{w^{r}}-i} B^{r}(\beta(z), \beta(w)) . \tag{2.16}
\end{equation*}
$$

If $A^{r}$ is defined for $r=1, \ldots, d$ by

$$
A^{r}(z, w)=\frac{F(z)+i}{z^{r}+i} \frac{\overline{F(w)}-i}{\overline{w^{r}}-i} B^{r}(\beta(z), \beta(w)),
$$

then (2.16) becomes

$$
\begin{equation*}
F(z)-\overline{F(w)}=\sum_{r=1}^{d}\left(z^{r}-\bar{w}^{r}\right) A^{r}(z, w) \tag{2.17}
\end{equation*}
$$

which means $F$ is in $\mathcal{L}^{d}$. Reversing the argument gives the converse.
Remark 2.18. It is known that $\mathcal{A}^{d}=\mathcal{S}^{d}$ for $d=1$ or 2 and that for $d \geq 3$, $\mathcal{A}^{d} \subsetneq \mathcal{S}^{d}$ [13], [38], [2]. It follows similarly that the Löwner class equals the Pick class in dimensions 1 and 2 , and is strictly contained in it for $d \geq 3$. By Theorem 5.5.1 of [32], rational inner functions are dense in the unit ball of $\mathcal{S}^{d}$ in the topology of uniform convergence on compacta. Therefore there must be rational inner functions in $\mathcal{S}^{d} \backslash \mathcal{A}^{d}$ for $d \geq 3$. By (2.11) and (2.12), it follows that for each $d \geq 3$, there are rational functions that are real on $\mathbb{R}^{d}$ and that are in $\mathcal{P}^{d} \backslash \mathcal{L}^{d}$.

## 3. Models, $B$-points and $C$-points

For a function $\varphi$ in $\mathcal{A}^{d}$, we can take the representation (2.9) and decompose the $B^{r}$ 's as Gramians to get a Hilbert space model for $\varphi$. That means we find a separable Hilbert space $\mathcal{M}$, an orthogonal decomposition of $\mathcal{M}$,

$$
\begin{equation*}
\mathcal{M}=\mathcal{M}^{1} \oplus \cdots \oplus \mathcal{M}^{d} \tag{3.1}
\end{equation*}
$$

and an analytic map $u: \mathbb{D}^{d} \rightarrow \mathcal{M}$ such that

$$
\begin{equation*}
1-\overline{\varphi(\mu)} \varphi(\lambda)=\sum_{r=1}^{d}\left(1-\overline{\mu^{r}} \lambda^{r}\right)\left\langle u_{\lambda}^{r}, u_{\mu}^{r}\right\rangle_{\mathcal{M}^{r}} \tag{3.2}
\end{equation*}
$$

for all $\lambda, \mu \in \mathbb{D}^{d}$, where we write $u_{\lambda}$ for $u(\lambda), P^{r}$ for the projection onto $\mathcal{M}^{r}$, and $u_{\lambda}^{r}$ for $P^{r}\left[u_{\lambda}\right]$.

We shall view (3.1) interchangeably as a graded Hilbert space (i.e., one with a given orthogonal decomposition) or as a single Hilbert space with $d$ given projections $P^{1}, \ldots, P^{d}$ that are orthogonal and add up to the identity.

In general, if $\eta \in \mathcal{M}$, we set $\eta^{r}=P^{r}[\eta]$. If $\lambda \in \mathbb{C}^{d}$, we may regard $\lambda$ as an operator on $\mathcal{M}$ by letting

$$
\begin{equation*}
\lambda \eta=\lambda^{1} \eta^{1}+\cdots+\lambda^{d} \eta^{d} . \tag{3.3}
\end{equation*}
$$

Equation (3.2) can then be rewritten as

$$
\begin{equation*}
1-\overline{\varphi(\mu)} \varphi(\lambda)=\left\langle\left(1-\mu^{*} \lambda\right) u_{\lambda}, u_{\mu}\right\rangle . \tag{3.4}
\end{equation*}
$$

A lurking isometry argument yields the following result [2].
Theorem 3.5. If $(\mathcal{M}, u)$ is a model of $\varphi \in \mathcal{A}^{d}$, then there exist $a \in \mathbb{C}$, vectors $\beta, \gamma \in \mathcal{M}$ and a linear operator $D: \mathcal{M} \rightarrow \mathcal{M}$ such that the operator

$$
\left[\begin{array}{cc}
a & 1 \otimes \beta \\
\gamma \otimes 1 & D
\end{array}\right]
$$

is a contraction on $\mathbb{C} \oplus \mathcal{M}$ and, for all $\lambda \in \mathbb{D}^{d}$,

$$
\begin{align*}
(1-D \lambda) u_{\lambda} & =\gamma,  \tag{3.6}\\
\varphi(\lambda) & =a+\left\langle\lambda u_{\lambda}, \beta\right\rangle . \tag{3.7}
\end{align*}
$$

With notation as in Theorem 3.5, we shall call $(a, \beta, \gamma, D)$ a realization of ( $\mathcal{M}, u$ ).

One can rewrite (3.6) and (3.7) as

$$
\begin{equation*}
\varphi(\lambda)=a+\beta^{*} \lambda(I-D \lambda)^{-1} \gamma . \tag{3.8}
\end{equation*}
$$

How one can go from (3.8) to (3.2) is discussed in [9] and [7].
If we start instead with the representation (2.7) of a function $F$ in $\mathcal{L}^{d}$, we can decompose the kernels $A^{r}$ as the Gramians of some vectors $v^{r}$ in auxiliary separable Hilbert spaces $\mathcal{N}^{r}$. Then we get, in the analogous notation to above,

$$
\begin{align*}
F(z)-\overline{F(w)} & =\sum_{r=1}^{d}\left(z^{r}-\bar{w}^{r}\right) A^{r}(z, w)  \tag{3.9}\\
& =\sum_{r=1}^{d}\left(z^{r}-\bar{w}^{r}\right)\left\langle v_{z}^{r}, v_{w}^{r}\right\rangle_{\mathcal{N}^{r}} \\
& =\left\langle\left(z-w^{*}\right) v_{z}, v_{w}\right\rangle_{\mathcal{N}} .
\end{align*}
$$

This decomposition leads to a lurking self-adjoint argument, which we shall discuss in Section 6.

In [4] we introduced the concept of a $B$-point for $\mathcal{S}$. Let us give a unified definition for each of the classes $\mathcal{S}, \mathcal{P}$, and $\mathcal{C}$; notice that it depends on the codomain of the function.

Definition 3.10. Let $U$ and $V$ be fixed domains and $f: U \rightarrow \bar{V}$ be an analytic function. A point $\tau$ in $\partial U$ is called a $B$-point of $f$ if there is a sequence $\lambda_{n}$ of points in $U$ that converge to $\tau$ and such that

$$
\begin{equation*}
\frac{\operatorname{dist}\left(f\left(\lambda_{n}\right), \partial V\right)}{\operatorname{dist}\left(\lambda_{n}, \partial U\right)} \tag{3.11}
\end{equation*}
$$

is bounded.
So, for example, a point $\tau$ in $\partial \Pi^{d}$ is a $B$-point for a function $F$ in $\mathcal{P}^{d}$ if there exists some sequence $z_{n}$ in $\Pi^{d}$ that tends to $\tau$ and such that the quantity

$$
\frac{\operatorname{Im} F\left(z_{n}\right)}{\min _{r \in\{1, \ldots, d\}}\left(\operatorname{Im} z_{n}^{r}\right)}
$$

is bounded.
For a function in $\mathcal{L}^{d}\left(\right.$ resp. $\left.\mathcal{A}^{d}\right)$ we shall call a point $\tau$ a $B$-point if it is a $B$-point for the function thought of as an element of $\mathcal{P}^{d}$ (resp. $\mathcal{S}^{d}$ ).

For each of the three classes $\mathcal{S}, \mathcal{P}$, and $\mathcal{C}$, it follows from results of F. Jafari [19] and M. Abate [1] that if $\tau$ is a $B$-point, then the ratio (3.11) remains bounded for every sequence $\lambda_{n}$ that tends to $\tau$ nontangentially. Moreover, the function $f$ will then have a nontangential limit at $\tau$. (A sequence $\lambda_{n}$ in $U$ tends to the point $\tau$ nontangentially if $\lambda_{n}$ tends to $\tau$ and $\frac{\operatorname{dist}\left(\lambda_{n}, \tau\right)}{\operatorname{dist}\left(\lambda_{n}, \partial U\right)}$ is bounded.)

The following result was proved in [4] for $d=2$, but the proof generalizes to any $d$. We shall need it in the proof of Theorem 6.26.

Lemma 3.12. Let $\varphi \in \mathcal{A}^{d}$ and $\tau \in \mathbb{T}^{d}$. Let $(\mathcal{M}, u)$ be a model for $\varphi$ and $(a, \beta, \gamma, D)$ be a realization. The following are equivalent:
(i) $\tau$ is a $B$-point for $\varphi$.
(ii) For some sequence $\lambda_{n}$ converging to $\tau$ nontangentially, the sequence $\left\|u_{\lambda_{n}}\right\|$ is bounded.
(iii) For any sequence $\lambda_{n}$ converging to $\tau$ nontangentially, the sequence $\left\|u_{\lambda_{n}}\right\|$ is bounded.
(iv) The vector $\gamma$ is in the range of $(I-D \tau)$.

Moreover, if $u_{\lambda_{n}}$ converges to a vector weakly as $\lambda_{n}$ tends to $\tau$ nontangentially, then $u_{\lambda_{n}}$ converges in norm. The vector $u_{\tau}:=\lim _{r}{ }_{11} u_{r \tau}$ exists for every $B$-point $\tau$.

A stronger condition than being a $B$-point is being a $C$-point.

Definition 3.13. A point $x \in \mathbb{R}^{d}$ is a $C$-point for $F \in \mathcal{L}$ if there are complex numbers $\eta^{1}, \ldots, \eta^{d}$ and a real number $c$ so that

$$
\begin{equation*}
F(z)-c-\sum_{r=1}^{d} \eta^{r}\left(z^{r}-x^{r}\right)=o(\|z-x\|) \tag{3.14}
\end{equation*}
$$

as $z$ tends to $x$ nontangentially.
In particular, if $F$ is differentiable at $x$ and $F(x)$ is real, then $x$ is a $C$-point for $F$, since (3.14) then holds as $z$ tends to $x$ from any direction.

The following result was proved in [4].
Proposition 3.15. Suppose $F \in \mathcal{L}$ has a model as in equation (3.9). If $x$ is a $C$-point for $F$, then as $z$ converges to $x$ nontangentially from $\Pi^{d}$, the vectors $v_{z}$ converge in norm to some vector $v_{x}$ in $\mathcal{N}$.

## 4. Analytically continuing Pick functions

Suppose $F$ is analytic on $\Pi^{d}$ and $E$ is an open set in $\mathbb{R}^{d}$. What conditions on $F$ guarantee that it can be analytically continued across $E$ ? The edge-of-the-wedge theorem (see Theorem 4.10 below) is a common tool to give such extensions. Checking the hypotheses, however, requires knowledge of the values of $F$ as one approaches points of $E$ not just nontangentially but also tangentially. If $F$ is in the Pick class $\mathcal{P}^{d}$, Theorem 4.8 says that it suffices to know that every point of $E$ is a $B$-point (which can be checked by looking at the values of $F$ on the inward-pointing normal).

As we are using bars to denote closure, we shall use stars for the complex conjugate of a set and write $\Pi^{*}$ for the lower half-plane.
4.1. One dimension. To understand the situtation, let us first consider the one-dimensional case. Let $\psi: \mathbb{D} \rightarrow \mathbb{H}$ be nonconstant. Then $\psi$ has a Herglotz representation; if we assume $\psi(0)$ is positive, then

$$
\begin{equation*}
\psi(z)=\int_{0}^{2 \pi} \frac{e^{i \theta}+z}{e^{i \theta}-z} d \mu(\theta) \tag{4.1}
\end{equation*}
$$

for some positive measure $\mu$. There is an elegant analysis of when $\psi$ has $B$-points in the paper [35] by D. Sarason, where the following two propositions are proved. Proposition 4.2 is originally due to M. Riesz [30], and Proposition 4.4 is due to R. Nevanlinna [27].

Proposition 4.2. Let $\psi$ be given by equation (4.1), and let $\tau$ be a point in $\mathbb{T}$. Then $\psi$ has a B-point at $\tau$ if and only if

$$
\begin{equation*}
\int \frac{1}{\left|e^{i \theta}-\tau\right|^{2}} d \mu(\theta)<\infty \tag{4.3}
\end{equation*}
$$

If $\varphi=\beta \circ(i \psi)$ is the Cayley transform of $\psi$, there is a distinction between $B$-points where $\varphi(\tau)$ equals 1 , corresponding to $\psi(\tau)=\infty$, and all other cases.

Proposition 4.4. Let $\varphi=\frac{\psi-1}{\psi+1}$, where $\psi$ is given by equation (4.1), and let $\tau$ be a point in $\mathbb{T}$. Then $\varphi$ has a B-point at $\tau$ with $\varphi(\tau) \neq 1$ if and only if inequality (4.3) holds. The function $\varphi$ has a B-point at $\tau$ with $\varphi(\tau)=1$ if and only if $\tau$ is a mass point of $\mu$.

Suppose now that $\varphi: \mathbb{D} \rightarrow \mathbb{D}$ has an open $\operatorname{arc} I$ of $B$-points. Can $\varphi$ be extended analytically across $I$ ? If we know that $\varphi$ omits a value on $I$, then the answer is yes. Indeed, after a Möbius map, we can assume that $\varphi$ is the Cayley transform of some $\psi$ as (4.1). If condition (4.3) holds on an open arc $I$, then $\mu$ must vanish on $I$ by Lemma 4.5. But then the formula (4.1) gives an analytic function on the extended plane less $\mathbb{T} \backslash I$.

However, without the assumption that $\varphi$ omits a value, the answer may be no, as Example 4.6 shows.

Lemma 4.5. Suppose $\mu$ is a measure on $[-\pi, \pi$ ) and inequality (4.3) holds for $\tau=e^{i x}$ for every $x$ in an open arc $I \subset[-\pi, \pi)$. Then $\mu(I)=0$.

Proof. For $\mu$ almost every point $x$ in $I$, there is a constant $c>0$ such that

$$
\mu\left[x-\frac{1}{k}, x+\frac{1}{k}\right] \geq c \frac{2}{k}
$$

by [34, Thms. 8.6 and 8.10$]$. For such an $x$, we have

$$
\begin{aligned}
\int_{-\pi}^{\pi} \frac{1}{\left|e^{i \theta}-\tau\right|^{2}} d \mu(\theta) & \geq \int_{x-1 / k}^{x+1 / k} \frac{1}{\left|e^{i \theta}-\tau\right|^{2}} d \mu(\theta) \\
& \geq \frac{1}{\left|1-e^{i / k}\right|^{2}} \mu\left[x-\frac{1}{k}, x+\frac{1}{k}\right] \\
& \geq c k
\end{aligned}
$$

Letting $k$ tend to infinity, the integral would be infinite; so $\mu$ must put no mass on $I$.

Example 4.6. Here is an example of a function in the Schur class of one variable that has $B$-points at every point of $\mathbb{T}$ but that cannot be analytically continued across every arc.

Let $\tau_{n}=e^{i x_{n}}$ be a sequence in $\mathbb{T}$ that converges to 1 . Let $c_{n}=2^{-n}\left|1-\tau_{n}\right|^{2}$. Then for every $\lambda \in \mathbb{T}$, the quantity $\left|\frac{1-\tau_{n}}{\lambda-\tau_{n}}\right|$ is less than or equal to 1 for all but finitely many values of $n$. Therefore,

$$
\begin{equation*}
\sum 2^{-n}\left|\frac{1-\tau_{n}}{\lambda-\tau_{n}}\right|^{2}<\infty \tag{4.7}
\end{equation*}
$$

for every $\lambda$.

Let $\mu=\sum 2^{-n} \delta_{x_{n}}$, let $F$ be the Herglotz transform of $\mu$, and let $\varphi=\frac{F-1}{F+1}$. By Proposition 4.4, we have that $\varphi$ has every point of $\mathbb{T}$ as a $B$-point, but $\varphi$ cannot be analytically continued across any arc containing 1 , as it takes the value 1 infinitely often on any such arc.
4.2. $d$ dimensions. Our goal is to prove the following analytic continuation theorem.

Theorem 4.8. Let $E$ be an open subset of $\mathbb{R}^{d}$. Then there is an open set $U$ in $\mathbb{C}^{d}$ that contains $\Pi^{d} \cup E \cup \Pi^{* d}$ with the following property: whenever $F$ is in the Pick class, and every point of $E$ is a B-point for $F$, then there is an analytic function $G$ on $U$ that agrees with $F$ on $\Pi^{d}$.

This theorem immediately implies the omit-a-value theorem. Let us say that a subset $E^{\prime}$ of $\mathbb{T}^{d}$ is a $B$-set for $\varphi$ in $\mathcal{S}^{d}$ if every point of $E^{\prime}$ is a $B$-point for $\varphi$.

Theorem 4.9. Let $E^{\prime}$ be an open subset of $\mathbb{T}^{d}$. Then there is an open set $U$ in $\mathbb{C}^{d}$ containing $\mathbb{D}^{d} \cup E^{\prime} \cup\{\mathbb{C} \backslash \overline{\mathbb{D}}\}^{d}$ such that the following two statements are equivalent for any $\varphi$ in the Schur class:
(1) there is an analytic function $\psi$ on $U$ that agrees with $\varphi$ on $\mathbb{D}^{d}$;
(2) the set $E^{\prime}$ is a $B$-set for $\varphi$, and for every point $\tau$ in $E^{\prime}$, there exists a neighborhood $V$ of $\tau$ in $\mathbb{T}^{d}$ and a point $\omega$ in $\mathbb{T}$ such that no nontangential limit of $\varphi$ at any point of $V$ is equal to $\omega$.

Condition (2) says that every point of $E^{\prime}$ has a neighborhood where the nontangential limits of $\varphi$ omit some value in $\mathbb{T}$.

We need a version of the edge-of-the-wedge theorem (Theorem C from [33]). We write $\mathbb{R}_{+}$for the interval $(0, \infty)$.

Theorem 4.10 (Edge-of-the-wedge). Let $E$ be an open subset of $\mathbb{R}^{d}$. Then there is an open set $U$ in $\mathbb{C}^{d}$ that contains $\Pi^{d} \cup E \cup \Pi^{* d}$ and is such that whenever $H$ is an analytic function on $\Pi^{d}$ with the property that for every $g$ in $C_{c}^{\infty}(E)$,

$$
\begin{equation*}
\lim _{\mathbb{R}_{+}^{d} \ni y \rightarrow 0} \int_{E} g(x) \operatorname{Im} H(x+i y) d x=0 \tag{4.11}
\end{equation*}
$$

then there is an analytic function $G$ on $U$ that agrees with $H$ on $\Pi^{d}$.
Proposition 4.12. Let $E$ be an open subset of $\mathbb{R}^{d}$. Then there is an open set $U$ in $\mathbb{C}^{d}$ that contains $\Pi^{d} \cup E \cup \Pi^{* d}$ with the following property: if $J$ is a nonempty interval in $\mathbb{R}, F$ is in the Pick class, $F$ has nontangential limits at almost every point of $E$, and these limits are all in $\mathbb{R} \backslash J$, then there is an analytic function $G$ on $U$ that agrees with $F$ on $\Pi^{d}$.

Proof. Precomposing $F$ with a Möbius transformation of $\Pi$ if necessary, we can assume that $J$ is an interval about infinity, so the nontangential limits are in some compact set $[-M, M]$ almost everywhere.

Let $H(z)=\log (1+M+F(z))$. Then $H$ maps $\Pi^{d}$ into $\{z \in \mathbb{C}: 0<\operatorname{Im} z<\pi\}$, and

$$
\lim _{y \rightarrow 0} H(x+i y) \in[0, \log (2 M+1)] \text { a.e. } x \in E .
$$

As $H$ has bounded imaginary part, we can pass the limit inside the integral on the left-hand side of (4.11), and as $H$ has real boundary values, we get that the limit is 0 . Therefore, by Theorem 4.10 we get an analytic extension of $H$, and hence $F$, to the desired open set $U$.

Proof of Theorem 4.8. We can extend $F$ to $\Pi^{* d}$ by letting $F(z)=\overline{F(\bar{z})}$ on $\Pi^{* d}$. The difficulty is in showing that the definitions of $F$ on the two disjoint domains $\Pi^{d}$ and $\Pi^{* d}$ are analytic continuations of each other across $E$. This is a local property. If we can show that every point of $E$ has a neighborhood on which the boundary values of $F$ take values in a bounded set, we can apply Proposition 4.12 to conclude that the reflection of $F$ is an analytic continuation of $F$ across this neighborhood in $E$. Since this is true at every point, the conclusion of the theorem will follow.

For convenience, we will change variables and consider the function $\varphi(\lambda)=$ $(-i) F \circ \alpha$, which is in $\mathcal{C}^{d}$.

We can normalize to assume that $\varphi(0, \ldots, 0)=1$ and that the point of interest for $\varphi$ is $\beta(0, \ldots, 0)=(-1, \ldots,-1)$. So for some $0<c<\frac{\pi}{5}$, the set

$$
\begin{equation*}
\left\{\left(e^{i \theta_{1}}, \ldots, e^{i \theta_{d}}\right): \forall 1 \leq r \leq d,\left|\theta_{r}\right| \geq \pi-5 c\right\} \tag{4.13}
\end{equation*}
$$

consists of $B$-points for $\varphi$. In what follows, we shall choose arg to take values in $[-\pi, \pi)$.

For each $\tau$ in the set

$$
\left\{\tau \in \mathbb{T}^{d-1}:\left|\arg \left(\tau^{j}\right)\right|<2 c, \forall 1 \leq j \leq d-1\right\},
$$

define $g_{\tau}$ in $\mathcal{C}^{1}$ by

$$
g_{\tau}(z)=\varphi\left(z, \tau^{1} z, \tau^{2} z, \ldots, \tau^{d-1} z\right) .
$$

Then for each $\tau$, the set

$$
I_{3 c}=\{\sigma \in \mathbb{T}:|\arg (\sigma)|>\pi-3 c\}
$$

is a set of $B$-points for $g_{\tau}$, and $g_{\tau}(0)=1$. Each $g_{\tau}$ has a Herglotz representation, and by the results of Section 4.1, the corresponding measure is supported off the set $I_{3 c}$. So

$$
g_{\tau}(z)=\int_{-\pi+3 c}^{\pi-3 c} \frac{e^{i \theta}+z}{e^{i \theta}-z} d \mu_{\tau}(\theta)
$$

for some probability measure $\mu_{\tau}$. Therefore, if $\sigma$ is in the $\operatorname{arc} I_{c}=\{\sigma \in \mathbb{T}$ : $|\arg (\sigma)|>\pi-c\}$, then

$$
\begin{aligned}
\left|g_{\tau}(\sigma)\right| & \leq \int_{-\pi+3 c}^{\pi-3 c}\left|\frac{e^{i \theta}+\sigma}{e^{i \theta}-\sigma}\right| d \mu_{\tau}(\theta) \\
& \leq \sec c
\end{aligned}
$$

Therefore, on the set $\left(I_{c}\right)^{d}$ we conclude that the nontangential limits of $\varphi$ take values in the bounded $\operatorname{set}[-\sec c, \sec c]$.

Notice that $c$ does not depend on $F$ : we have shown that for any $F$, normalized to have $F(i, \ldots, i)=i$, if $F$ has $B$-points on the set $\alpha\left(\left(I_{5 c}\right)^{d}\right)$, then $F$ is bounded on $\alpha\left(\left(I_{c}\right)^{d}\right)$. By Proposition 4.12, this latter set now has a neighborhood to which $F$ can be analytically extended, and this neighborhood can be chosen independently of $F$. So every point $x$ in $E$ has a neighborhood $U_{x}$ to which all functions $F$ in the Pick class with $B$-points on $E$ can be extended; let $U$ be the union of all the $U_{x}$ as $x$ ranges over $E$.

## 5. The Löwner classes

We shall single out functions that have a representation on subsets of $\mathbb{R}^{d}$ as in (2.7).

Definition 5.1. Let $E \subseteq \mathbb{R}^{d}$ be a nonempty open set, and let $n$ be a positive integer. We define $\mathcal{L}_{n}(E)$ to be the set of real valued differentiable functions that have the following property: whenever $x_{1}, \ldots, x_{n}$ are $n$ distinct points in $E$, there exist positive semi-definite $n$-by- $n$ matrices $A^{1}, \ldots, A^{d}$ so that

$$
\begin{align*}
A^{r}(i, i) & =\left.\frac{\partial f}{\partial x^{r}}\right|_{x_{i}},  \tag{5.2}\\
f\left(x_{j}\right)-f\left(x_{i}\right) & =\sum_{r=1}^{d}\left(x_{j}^{r}-x_{i}^{r}\right) A^{r}(i, j) \forall 1 \leq i, j \leq n . \tag{5.3}
\end{align*}
$$

We shall give an alternative description of $\mathcal{L}(E)$. We shall temporarily call it $\mathcal{L}_{\partial}(E)$, but we shall show in Proposition 5.11 that it coincides with the set $\mathcal{L}(E)$ from Definition 1.11.

Definition 5.4. Let $E \subseteq \mathbb{R}^{d}$ be a nonempty open set. We shall let $\mathcal{L}_{\partial}(E)$ denote the set of differentiable real valued functions $f$ on $E$ for which there exist positive semi-definite functions $A^{1}, \ldots, A^{d}: E \times E \rightarrow \mathbb{C}$ so that

$$
\begin{equation*}
A^{r}(z, z)=\left.\frac{\partial f}{\partial x^{r}}\right|_{z} \tag{5.5}
\end{equation*}
$$

and

$$
\begin{equation*}
f(z)-f(w)=\sum_{r=1}^{d}\left(z^{r}-w^{r}\right) A^{r}(z, w) \tag{5.6}
\end{equation*}
$$

If $E \subseteq \mathbb{R}^{d}$, a function $f$ in $\mathcal{L}_{\partial}(E)$ can be extended to a function $F$ in $\mathcal{L}$ that has $f$ as its nontangential boundary values on $E$. We shall use the abbreviation nt for nontangential.

Proposition 5.7. Let $E \subseteq \mathbb{R}^{d}$ be open, and let $f \in \mathcal{L}_{\partial}(E)$. Then there exists $F \in \mathcal{L}^{d}$ such that every point of $E$ is a $B$-point for $F$ and such that

$$
\begin{equation*}
\lim _{z^{\mathrm{nt}} \rightarrow t} F(z)=f(t) \quad \forall t \in E \tag{5.8}
\end{equation*}
$$

Proposition 5.7 follows immediately from the corresponding result on the polydisk, Theorem 5.9, which was proved by J. A. Ball and V. Bolotnikov [6]. (We are changing their language slightly; they did not explicitly use the notion of $B$-point.)

Theorem 5.9 (Ball-Bolotnikov). Let $E^{\prime} \subseteq \mathbb{T}^{d}$, and let $\psi: E^{\prime} \rightarrow \mathbb{C}$. Suppose there are positive semi-definite functions $B^{1}, \ldots, B^{d}: E^{\prime} \times E^{\prime} \rightarrow \mathbb{C}$ such that, for all $\lambda, \mu$ in $E^{\prime}$,

$$
\begin{equation*}
1-\psi(\lambda) \overline{\psi(\mu)}=\left(1-\lambda^{1} \bar{\mu}^{1}\right) B^{1}(\lambda, \mu)+\cdots+\left(1-\lambda^{d} \bar{\mu}^{d}\right) B^{d}(\lambda, \mu) \tag{5.10}
\end{equation*}
$$

Then there is a function $\varphi$ in $\mathcal{A}$ such that every point of $E^{\prime}$ is a B-point for $\varphi$ and

$$
\lim _{\substack { \mathrm{nt} \\
\begin{subarray}{c}{ \mathrm { nt } \\
\begin{subarray} { c } { \\
} }\end{subarray}} \varphi(\lambda)=\psi(\tau) \quad \forall \tau \in E^{\prime}
$$

Moreover, if $\varphi$ is defined to equal $\psi$ on $E^{\prime}$, the kernels $B^{r}$ can be extended to $E^{\prime} \cup \mathbb{D}^{d}$ so that, for all $\lambda, \mu$ in $E^{\prime} \cup \mathbb{D}^{d}$,

$$
1-\varphi(\lambda) \overline{\varphi(\mu)}=\left(1-\lambda^{1} \bar{\mu}^{1}\right) B^{1}(\lambda, \mu)+\cdots+\left(1-\lambda^{d} \bar{\mu}^{d}\right) B^{d}(\lambda, \mu)
$$

We can pass back and forth between regarding functions in $\mathcal{L}(E)$ as functions in the Löwner class $\mathcal{L}^{d}$ that have $B$-points on $E$ (and so can be analytically extended across $E$ ), and as functions that are characterized by their values on $E$ by (5.4) and can then be analytically extended into $\Pi^{d}$.

Proposition 5.11. Let $E \subseteq \mathbb{R}^{d}$ be a nonempty open set. The following four sets coincide:
(i) $\cap_{n=1}^{\infty} \mathcal{L}_{n}(E)$.
(ii) The set $\mathcal{L}_{\partial}(E)$ defined by Definition (5.4).
(iii) The set $\mathcal{L}(E)$ defined by Definition (1.11).
(iv) The functions $f$ on $E$ for which there exists a function $F$ in $\mathcal{L}^{d}$ such that every point $x$ of $E$ is a $B$-point of $F$ and the nontangential limit of $F$ at $x$ is $f(x)$.

Proof. It is immediate that (ii) $\subseteq(\mathrm{i})$. Theorem 4.8 asserts that (iii) $=(\mathrm{iv})$. Proposition 5.7 says that (ii) $\subseteq$ (iv).

To show (iii) $\subseteq$ (ii), choose a model for $F$ so that equation (3.9) holds on $\Pi^{d} \times \Pi^{d}$ :

$$
F(Z)-\overline{F(W)}=\left\langle\left(Z-W^{*}\right) v_{Z}, v_{W}\right\rangle \quad \forall Z, W \in \Pi^{d} .
$$

As every point in $E$ is a $C$-point for $F$, we can let $Z$ and $W$ tend to points in $E$ nontangentially, $z$ and $w$ respectively. By Proposition 3.15, the vectors $v_{Z}$ and $v_{W}$ converge to $v_{z}$ and $v_{w}$. Let

$$
A^{r}(z, w)=\left\langle v_{z}^{r}, v_{w}^{r}\right\rangle,
$$

and one gets equation (5.6). To get equation (5.5), let $z$ be in $E$ and let $W$ in $\Pi^{d}$ tend to $z$ nontangentially. As $F$ is analytic at $z$, we have

$$
\begin{equation*}
F(W)-F(z)=\left.\sum\left(W^{r}-z^{r}\right) \frac{\partial f}{\partial x^{r}}\right|_{z}+o(\|z-W\|) . \tag{5.12}
\end{equation*}
$$

From the model,

$$
\begin{align*}
F(W)-F(z)= & \left\langle(W-z) v_{W}, v_{z}\right\rangle  \tag{5.13}\\
== & \left\langle(W-z) v_{z}, v_{z}\right\rangle+ \\
& \left\langle(W-z)\left(v_{W}-v_{z}\right), v_{z}\right\rangle .
\end{align*}
$$

The second term on the right of equation (5.13) is $o(\|z-W\|)$, so comparing with equation (5.12), we conclude that

$$
\left.\frac{\partial f}{\partial x^{r}}\right|_{z}=\left\langle v_{z}^{r}, v_{z}^{r}\right\rangle,
$$

and hence equation (5.5) holds.
To prove (i) $\subseteq$ (ii), we need to show that if equations (5.5) and (5.6) hold on every finite set, with perhaps a different choice of $A^{r}$ 's each time, then we can make one choice for the $A^{r}$ 's that works everywhere.

Let $f \in \cap_{n=1}^{\infty} \mathcal{L}_{n}(E)$. Consider any finite set $\left\{z_{1}, \ldots, z_{n}\right\}$ of distinct points in $E$. By Definition 5.1 there exist kernels $A^{1}, \ldots, A^{d}$ on $E$ such that equations (5.2), (5.3) hold, and we have

$$
A^{r}(i, i) \leq \frac{\partial f}{\partial x^{r}}\left(z_{i}\right), \quad i=1, \ldots, n, r=1, \ldots, d
$$

Since the matrices $A^{r}$ are positive semi-definite, we also obtain bounds on the off-diagonal entries of all the $A^{r}$. Hence the set $K$ of all $d$-tuples $\left(A^{1}, \ldots, A^{d}\right)$ for which equations (5.2), (5.3) hold is a compact nonempty subset of $M_{n}^{d}$.

Moreover, if $\left(B^{1}, \ldots, B^{d}\right)$ is a $d$-tuple of kernels on any finite superset $Z$ of $\left\{z_{1}, \ldots, z_{n}\right\}$ for which the analogs of equations (5.2), (5.3) hold, then the choice of $A^{r}$ to be the principal submatrix of $B^{r}$ corresponding to $\left\{z_{1}, \ldots, z_{n}\right\}$ gives a $d$-tuple that belongs to $K$. Therefore, by Kurosh's theorem [5, p. 74] or [3, p. 30], there is a $d$-tuple $\left(A^{1}, \ldots, A^{d}\right)$ of kernels on $E$ such that equations (5.2) and (5.3) hold for all points $z_{i}, z_{j} \in E$.

## 6. The $\mu$-spectral theorem

A function in the Pick class of one variable, i.e., an analytic function from $\Pi$ to $\Pi$, has an integral representation that can be obtained from the Herglotz representation (4.1) of functions from $\mathbb{D}$ to $\mathbb{H}$ by a change of variables [17].

Theorem 6.1 (Herglotz). An analytic function $F: \Pi \rightarrow \Pi$ has a unique representation of the form

$$
\begin{equation*}
F(z)=c+d z+\int \frac{1+z t}{t-z} d \mu(t) \tag{6.2}
\end{equation*}
$$

where $\operatorname{Im} c \geq 0$ and $d \geq 0$, and $\mu$ is a finite positive Borel measure on $\mathbb{R}$. Conversely any function of this form is in the Pick class of one variable.

If in addition $F$ decays up the imaginary axis, one gets that $F$ is the Cauchy transform of a finite measure on $\mathbb{R}$. This is called the Nevanlinna representation and was proved by R. Nevanlinna [26].

Theorem 6.3 (Nevanlinna). If $F: \Pi \rightarrow \Pi$ is analytic and satisfies

$$
\begin{equation*}
\limsup _{y \rightarrow \infty}|y F(i y)|<\infty, \tag{6.4}
\end{equation*}
$$

then there exists a unique finite positive Borel measure $\nu$ on $\mathbb{R}$ so that

$$
F(z)=\int \frac{d \nu(t)}{t-z}
$$

Remark 6.5. If one considers $\psi=-i F \circ \alpha: \mathbb{D} \rightarrow \mathbb{H}$, then $d$ in Theorem 6.1 is the mass assigned to the point 1 in the Herglotz representation of $\psi$. Nevanlinna's condition (6.4) in Theorem 6.3 is equivalent to saying that $\varphi=\beta \circ F \circ \alpha: \mathbb{D} \rightarrow \mathbb{D}$ has a $B$-point at 1 with $\varphi(1)=-1$.

One can prove the spectral theorem for a (possibly unbounded) self-adjoint operator by showing that if $R_{z}$ is the resolvent, then for any vector $u$, the function $\left\langle R_{z} u, u\right\rangle$ is in the one variable Pick class and satisfies Nevanlinna's growth condition. Then Theorem 6.3 gives the scalar spectral measure. See [14, Chap. V] or [24, Chap. 32]. Conversely, if $X$ is the operator of multiplication by the independent variable on $L^{2}(\mu)$ and $v$ is the constant function 1 , then (6.2) becomes

$$
\begin{equation*}
F(z)=c+d z+\left\langle(1+z X)(X-z)^{-1} v, v\right\rangle \tag{6.6}
\end{equation*}
$$

In several variables, there is also a connection between Pick functions and self-adjoint operators, which could be called a $\mu$-spectral theorem (Theorem 6.9 below).

Definition 6.7. Let $\mathcal{M}$ be a Hilbert space, with a fixed decomposition as $\mathcal{M}=\mathcal{M}^{1} \oplus \cdots \oplus \mathcal{M}^{d}$. Let $T$ be a densely defined linear operator on $\mathcal{M}$. For
$z=\left(z^{1}, \ldots, z^{d}\right)$ in $\mathbb{C}^{d}$, define the $\mu$-resolvent of $T$ at $z$ to be

$$
(T-z)^{-1}=\left(T-\left[z^{1} I_{\mathcal{M}^{1}} \oplus \cdots \oplus z^{d} I_{\mathcal{M}^{d}}\right]\right)^{-1} .
$$

The $\mu$-spectrum of $T$ is the complement of the set of points in $\mathbb{C}^{d}$ for which the $\mu$-resolvent is bounded.

The expressions " $\mu$-resolvent" and " $\mu$-spectrum" are not standard, but they are suggested by usage in control engineering. The notion of $\mu$-analysis provides an approach to robust stabilization in the presence of "structured uncertainty" [15]. Corresponding to the projections $P^{1}, \ldots, P^{d}$, one defines the "cost function" $\mu(X)$ by
$\frac{1}{\mu(X)}=\inf \left\{\|T\|: T \in B(\mathcal{M})\right.$, each $P^{r} \mathcal{M}$ reduces $T$, and $1-T X$ is singular $\}$.
In what follows, we shall write $z$ for $z^{1} I_{\mathcal{M}^{1}} \oplus \cdots \oplus z^{2} I_{\mathcal{M}^{d}}$ and $z^{*}$ for $\bar{z}^{1} I_{\mathcal{M}^{1}} \oplus \cdots \oplus \bar{z}^{d} I_{\mathcal{M}^{d}}$. Let us recall Definition 1.19.

Definition 1.19. For each real number $t$, define

$$
\begin{equation*}
\rho_{t}(z)=\frac{z+t}{1-t z} . \tag{6.8}
\end{equation*}
$$

For $F \in \mathcal{L}^{d}$, define

$$
F_{t}:=\rho_{t} \circ F \circ \rho_{t} .
$$

Note that, similarly to the maps $\alpha$ and $\beta$, we use $\rho_{t}$ on $\mathbb{C}^{d}$ to mean the component-wise action.

Theorem 6.9. Let $F$ be in $\mathcal{L}^{d}$ and $z_{0}$ be a point in $\Pi^{d}$. For all except at most a countable number of real numbers $t$, there is a Hilbert space $\mathcal{M}=$ $\mathcal{M}^{1} \oplus \cdots \oplus \mathcal{M}^{d}$, a self-adjoint operator $X$ on $\mathcal{M}$, a vector $v$ in $\mathcal{M}$, and a real number $c$ so that

$$
\begin{equation*}
F_{t}(z)=c+\langle z v, v\rangle+\left\langle\left(z-z_{0}^{*}\right)(X-z)^{-1}\left(z-z_{0}\right) v, v\right\rangle . \tag{6.10}
\end{equation*}
$$

Conversely, if $z_{0}$ is a point in $\Pi^{d}$, $c$ is a real number, $X$ is a densely defined self-adjoint operator on a Hilbert space $\mathcal{M}=\mathcal{M}^{1} \oplus \cdots \oplus \mathcal{M}^{d}$, and $v$ is a vector in $\mathcal{M}$, then the function of $z$ given by the right-hand side of equation (6.10) is in $\mathcal{L}^{d}$.

Proof. $(\Rightarrow)$ Let $\varphi=\beta \circ F \circ \alpha$ in $\mathcal{A}^{d}$, where $\alpha$ and $\beta$ are defined in (2.1) and (2.2). Choose a model for $\varphi$ so that (3.2) holds:

$$
\begin{align*}
1-\overline{\varphi(\mu)} \varphi(\lambda) & =\sum_{r=1}^{d}\left(1-\overline{\mu^{r}} \lambda^{r}\right)\left\langle u_{\lambda}^{r}, u_{\mu}^{r}\right\rangle_{\mathcal{M}^{r}}  \tag{6.11}\\
& =\left\langle\left(1-\mu^{*} \lambda\right) u_{\lambda}, u_{\mu}\right\rangle_{\mathcal{M}} .
\end{align*}
$$

Define a linear operator $V$ by

$$
V:\binom{1}{\lambda u_{\lambda}} \mapsto\binom{\varphi(\lambda)}{u_{\lambda}}
$$

and extend it by linearity to finite linear combinations of vectors of the form

$$
\binom{1}{\lambda_{i} u_{\lambda_{i}}}
$$

where the points $\lambda_{i}$ range over $\mathbb{D}^{d}$.
$V$ is defined on a subspace of $\mathbb{C} \oplus \mathcal{M}$, and by (6.11) it is isometric on its domain. If the codimensions of the closures of the domain and range of $V$ are the same, $V$ can be extended to a unitary $U$. If they are different, after the addition of a separable infinite-dimensional summand to one of the spaces $\mathcal{M}^{r}$, the codimensions become equal, and one can then extend $V$ to a unitary $U$. So we can assume that we have a unitary $U: \mathbb{C} \oplus \mathcal{M} \rightarrow \mathbb{C} \oplus \mathcal{M}$ such that

$$
\begin{equation*}
U:\binom{1}{\lambda u_{\lambda}} \mapsto\binom{\varphi(\lambda)}{u_{\lambda}} \tag{6.12}
\end{equation*}
$$

Now, let $\tau$ be a point in the unit circle that is not in the point spectrum of $U$, and let

$$
t=-i \frac{1-\tau}{1+\tau}
$$

As $\mathbb{C} \oplus \mathcal{M}$ is separable, the point spectrum of $U$ is countable, so all but countably many real numbers $t$ will arise in this way.

Let

$$
Y=-i(U-\tau)^{-1}(U+\tau):(U-\tau) \eta \mapsto-i(U+\tau) \eta
$$

Then $Y$ is densely defined and self-adjoint. Its domain $\mathcal{D}$ is $\operatorname{ran}(U-\tau)$. Moreover, by definition,

$$
Y:\binom{\varphi(\lambda)-\tau}{(1-\tau \lambda) u_{\lambda}} \mapsto\binom{-i(\varphi(\lambda)+\tau)}{-i(1+\tau \lambda) u_{\lambda}}
$$

Therefore,

$$
\begin{equation*}
Y:\binom{1}{\frac{1-\tau \lambda}{\varphi(\lambda)-\tau} u_{\lambda}} \mapsto\binom{-i \frac{\varphi(\lambda)+\tau}{\varphi(\lambda)-\tau}}{-i \frac{1+\tau \lambda}{\varphi(\lambda)-\tau} u_{\lambda}} \tag{6.13}
\end{equation*}
$$

Let

$$
\tilde{v}_{z}=\frac{1-\tau \beta(z)}{\beta \circ F(z)-\tau} u_{\beta(z)}
$$

Then one can rewrite (6.13) as

$$
\begin{equation*}
Y:\binom{1}{\tilde{v}_{z}} \mapsto\binom{\rho_{t} \circ F(z)}{-\rho_{-t}(z) \tilde{v}_{z}} \tag{6.14}
\end{equation*}
$$

Now let

$$
v_{z}=\tilde{v}_{\rho_{t}(z)}
$$

Then (6.14) becomes

$$
\begin{equation*}
Y:\binom{1}{v_{z}} \mapsto\binom{F_{t}(z)}{-z v_{z}} \tag{6.15}
\end{equation*}
$$

As $Y$ is self-adjoint, (6.15) implies that

$$
\begin{equation*}
F_{t}(z)-F_{t}(w)^{*}=\left\langle\left(z-w^{*}\right) v_{z}, v_{w}\right\rangle \tag{6.16}
\end{equation*}
$$

Let $X$ be the compression of $-Y$ to $\mathcal{M}$. By Lemma $6.24, X$ is self-adjoint with dense domain equal to $\mathcal{D} \cap \mathcal{M}$.

If $\gamma$ is in $\mathcal{D} \cap \mathcal{M}$, then

$$
Y\binom{0}{\gamma}=\binom{L(\gamma)}{-X \gamma}
$$

for some linear functional $L$.
Define $v=v_{z_{0}}$, and let $a=F_{t}\left(z_{0}\right)$. Then

$$
Y\binom{1}{v}=\binom{a}{-z_{0} v}
$$

For $z \in \Pi^{d}$, let

$$
\begin{aligned}
Y\binom{1}{v_{z}} & =Y\binom{1}{v}+Y\binom{0}{v_{z}-v} \\
& =\binom{a}{-z_{0} v}+\binom{L\left(v_{z}-v\right)}{-X\left(v_{z}-v\right)}
\end{aligned}
$$

By (6.15), we get the equations

$$
\begin{align*}
F_{t}(z) & =a+L\left(v_{z}-v\right)  \tag{6.17}\\
z v_{z} & =z_{0} v+X\left(v_{z}-v\right) \tag{6.18}
\end{align*}
$$

For $\gamma \in \mathcal{D} \cap \mathcal{M}$,

$$
\begin{aligned}
\left\langle\binom{ 0}{\gamma},\binom{a}{-z_{0} v}\right\rangle & =\left\langle\binom{ 0}{\gamma}, Y\binom{1}{v}\right\rangle \\
& =\left\langle Y\binom{0}{\gamma},\binom{1}{v}\right\rangle \\
& =\left\langle\binom{ L(\gamma)}{-X \gamma},\binom{1}{v}\right\rangle
\end{aligned}
$$

Therefore,

$$
\begin{equation*}
L(\gamma)=-\left\langle\gamma, z_{0} v\right\rangle+\langle X \gamma, v\rangle \tag{6.19}
\end{equation*}
$$

If $z$ is in $\Pi^{d}$, then by Lemma 6.25, $X-z$ is invertible, so (6.18) yields

$$
\begin{equation*}
v_{z}-v=(X-z)^{-1}\left(z-z_{0}\right) v . \tag{6.20}
\end{equation*}
$$

Combining equations (6.17) to (6.20), we get

$$
\begin{align*}
F_{t}(z)= & a-\left\langle v_{z}-v, z_{0} v\right\rangle+\left\langle z v_{z}-z_{0} v, v\right\rangle  \tag{6.21}\\
= & a-\left\langle v_{z}-v, z_{0} v-z^{*} v\right\rangle+\left\langle\left(z-z_{0}\right) v, v\right\rangle \\
= & a+\left\langle(X-z)^{-1}\left(z-z_{0}\right) v,\left(z^{*}-z_{0}\right) v\right\rangle \\
& +\left\langle\left(z-z_{0}\right) v, v\right\rangle .
\end{align*}
$$

By (6.16),

$$
a-\bar{a}=F_{t}\left(z_{0}\right)-\overline{F_{t}\left(z_{0}\right)}=\left\langle\left(z_{0}-z_{0}^{*}\right) v, v\right\rangle,
$$

so $c:=a-\left\langle z_{0} v, v\right\rangle$ is real. Then (6.21) becomes (6.10), as desired.
$(\Leftarrow)$ To prove the converse, suppose $X$ is a self-adjoint operator on $\mathcal{M}$ with dense domain $\mathcal{D}^{\prime}$. Let $F(z)$ be given by the right-hand side of (6.10). Define $v_{z}$ by (6.20); i.e.,

$$
\begin{equation*}
v_{z}=v+(X-z)^{-1}\left(z-z_{0}\right) v . \tag{6.22}
\end{equation*}
$$

Define a linear functional $L$ on $\mathcal{D}^{\prime}$ by

$$
L(\gamma)=-\left\langle\gamma, z_{0} v\right\rangle+\langle X \gamma, v\rangle
$$

Let $\mathcal{D}$ be the linear span in $\mathbb{C} \oplus \mathcal{M}$ of the vector $\binom{1}{v}$ and the vector space $0 \oplus \mathcal{D}^{\prime}$. Let $a=c+\left\langle z_{0} v, v\right\rangle$. Finally, define $Y$ on $\mathcal{D}$ by

$$
Y\binom{t}{\gamma+t v}=t\binom{a}{-z_{0} v}+\binom{L(\gamma)}{-X(\gamma)}
$$

for $t$ in $\mathbb{C}$. It is routine to verify that $Y$ is symmetric. Moreover, by (6.22), $\left(v_{z}-v\right)$ is in the domain of $X-z$ and therefore in $\mathcal{D}^{\prime}$. So $\binom{1}{v_{z}}$ is in $\mathcal{D}$ for every $z$ in $\Pi$, and therefore

$$
\begin{equation*}
\left\langle Y\binom{1}{v_{z}},\binom{1}{v_{w}}\right\rangle=\left\langle\binom{ 1}{v_{z}}, Y\binom{1}{v_{w}}\right\rangle . \tag{6.23}
\end{equation*}
$$

Expanding (6.23) and rearranging, one gets

$$
F(z)-\overline{F(w)}=\left\langle\left(z-w^{*}\right) u_{z}, u_{w}\right\rangle,
$$

and hence $F$ is in $\mathcal{L}^{d}$.
Lemma 6.24. With notation as in the proof of the forward direction of Theorem 6.9, the domain of $X$ is $\mathcal{D} \cap \mathcal{M}$. This domain is dense, and the operator $X$ is self-adjoint.

Proof. Since $\mathcal{D}$ is dense in $\mathbb{C} \oplus \mathcal{M}$, there are vectors $\xi_{n}$ in $\mathcal{M}$ that converge to zero and such that $\binom{1}{\xi_{n}}$ are in $\mathcal{D}$. If $\gamma$ is any vector in $\mathcal{M}$, there are vectors $\binom{a_{n}}{\eta_{n}}$ in $\mathcal{D}$ that converge to $\binom{0}{\gamma}$, hence so do the vectors

$$
\binom{a_{n}}{\eta_{n}}-a_{n}\binom{1}{\xi_{n}}=\binom{0}{\eta_{n}-a_{n} \xi_{n}} .
$$

Therefore $\mathcal{D} \cap \mathcal{M}$ is dense in $\mathcal{M}$.
Let $P$ be the projection from $\mathbb{C} \oplus \mathcal{M}$ onto $\mathcal{M}$. Let $X=-\left.P Y\right|_{\mathcal{M}}$ with domain $\mathcal{D}^{\prime}=\mathcal{D} \cap \mathcal{M}$. Then for $\gamma, \eta$ in $\mathcal{D}^{\prime}$, we have

$$
\begin{aligned}
\langle X \gamma, \eta\rangle & =-\langle P Y \gamma, \eta\rangle \\
& =-\langle Y \gamma, \eta\rangle \\
& =-\langle\gamma, Y \eta\rangle \\
& =\langle\gamma, X \eta\rangle .
\end{aligned}
$$

So $X$ is symmetric.
To prove $X$ is self-adjoint, it remains to show that $X$ and $X^{*}$ have the same domain. Assume that there is some vector $\eta$ in $\mathcal{M}$ such that

$$
|\langle X \gamma, \eta\rangle| \leq C\|\gamma\|
$$

for all $\gamma \in \mathcal{D}^{\prime}$. Then for every vector $\binom{c}{\delta}$ in $\mathcal{D}$ of norm at most one, we have

$$
\begin{aligned}
\left\lvert\,\left\langle Y\binom{c}{\delta},\binom{0}{\eta}\right\rangle\right. & =\left|\left\langle Y\left[\binom{0}{\delta-c \xi_{1}}+c\binom{1}{\xi_{1}}\right],\binom{0}{\eta}\right\rangle\right| \\
& \leq C\left\|\delta-\xi_{1}\right\|+|c|\left\|Y\binom{1}{\xi_{1}}\right\|\|\eta\| \leq C^{\prime}
\end{aligned}
$$

So $\binom{0}{\eta}$ is in $\mathcal{D}$, and therefore $\eta$ is in $\mathcal{D}^{\prime}$.
Lemma 6.25. Let $X$ be a densely defined self-adjoint operator on $\mathcal{M}$. The $\mu$-spectrum of $X$ is disjoint from $\Pi^{d} \cup \Pi^{d^{*}}$. Moreover,

$$
\left\|(X-z)^{-1}\right\| \leq 1 / \min _{1 \leq r \leq d}\left(\left|\operatorname{Im} z^{r}\right|\right) \quad \forall z \in \Pi^{d} \cup \Pi^{d^{*}}
$$

Proof. Let $X$ be self-adjoint on $\mathcal{M}=\mathcal{M}^{1} \oplus \cdots \oplus \mathcal{M}^{d}$, and let $z=\left(x^{1}+\right.$ $\left.i y^{1}, \ldots, x^{d}+i y^{d}\right)$ be a point in $\Pi^{d} \cup \Pi^{d^{*}}$. Then for any $v=v^{1} \oplus \cdots \oplus v^{d}$ in $\mathcal{M}$,

$$
\langle(X-z) v, v\rangle=\langle(X-x) v, v\rangle-i\left(y^{1}\left\|v^{1}\right\|^{2}+\cdots+y^{d}\left\|v^{d}\right\|^{2}\right) .
$$

The first summand on the right is real, so $X-z$ is bounded below by $\min \left(\left|y^{r}\right|\right)$, and therefore has a left inverse. Applying the same argument to $z^{*}$, we get that $X-z^{*}$ has a left inverse, and taking adjoints we get that $X-z$ has a right inverse also.

When $F$ decays at infinity, we can sharpen Theorem 6.9 to get a theorem like Nevanlinna's Theorem 6.3. (It has long been known that the limsup in Nevanlinna's theorem can be replaced by a liminf.)

Let us write $\mathbb{1}$ for $(1,1, \ldots, 1)$, and $s \mathbb{1}$ for $(s, s, \ldots, s)$, etc.
Theorem 6.26. Let $F$ be in $\mathcal{L}^{d}$, and assume $F$ has a representation as in (6.10) with $t=0$. Then the following are equivalent:
(i) $\liminf _{y \rightarrow \infty} y|F(i y \mathbb{1})|<\infty$.
(ii) There exists a vector $v_{1}$ in $\mathcal{M}$ so that

$$
\begin{equation*}
F(z)=\left\langle(X-z)^{-1} v_{1}, v_{1}\right\rangle \quad z \in \Pi^{d} . \tag{6.27}
\end{equation*}
$$

(iii) The function $\varphi=\beta \circ F \circ \alpha$ in $\mathcal{S}$ has a $B$-point at $\mathbb{1}$ and $\varphi(\mathbb{1})=-1$.
(iv) $\lim _{y \rightarrow \infty} F(i y \mathbb{1})=0$ and the vectors $v_{(i y \mathbb{1})}$ defined by (6.22) satisfy

$$
\liminf _{y \rightarrow \infty} y\left\|v_{(i y \mathbb{1})}\right\|<\infty
$$

(v) The vector $v$ is in the domain of $X$, and $\lim _{y \rightarrow \infty} F(i y \mathbb{1})=0$.

Proof. (i) $\Rightarrow$ (iii) Let $\varphi=\beta \circ F \circ \alpha$. Condition (i) becomes

$$
\begin{equation*}
\liminf _{s \rightarrow 1} \frac{1+s}{1-s}\left|\frac{1+\varphi(s \mathbb{1})}{1-\varphi(s \mathbb{1})}\right|<\infty . \tag{6.28}
\end{equation*}
$$

The left-hand side of (6.28) dominates

$$
\frac{1-|\varphi(s \mathbb{1})|}{1-s}
$$

so $\mathbb{1}$ is a $B$-point. In order for (6.28) to hold, we must have $\varphi(\mathbb{1})=-1$.
(iii) $\Leftrightarrow$ (iv) As the proof of Proposition 2.13 shows, one can pass between a model $(\mathcal{M}, u)$ for $\varphi$ and a $\operatorname{model}(\mathcal{M}, v)$ for $F$ by letting

$$
\begin{align*}
v_{z}^{r} & =\left(\frac{F(z)+i}{z^{r}+i}\right) u_{\beta(z)}^{r},  \tag{6.29}\\
u_{\lambda}^{r} & =\left(\frac{1-\varphi(\lambda)}{1-\lambda^{r}}\right) v_{\alpha(\lambda)}^{r} \quad r=1, \ldots, d .
\end{align*}
$$

By Lemma $3.12, \varphi$ having a $B$-point at $\mathbb{1}$ is equivalent to $u_{(r \mathbb{1})}$ being bounded as $r \rightarrow 1^{-}$. Moreover, $\varphi(r \mathbb{1})$ tending to -1 is the same as $F(i y \mathbb{1})$ tending to 0 as $y \rightarrow \infty$. And as long as $F(i y \mathbb{1})$ has a finite limit, (6.29) says that $u_{(r \mathbb{1})}$ is bounded if and only if $\left[y v_{(i y \mathbb{1})}\right]$ is.
(iv) $\Rightarrow(\mathrm{v})$ As $X$ is densely defined and self-adjoint, it is closed. By (6.18), the vectors $v_{z}-v$ all lie in $\mathcal{D}^{\prime}$, the domain of $X$. Let $z=(i y \mathbb{1})$, and let $y \rightarrow \infty$. Then $v-v_{(i y \mathbb{1})}$ tends to $v$. Moreover, $X\left(v-v_{(i y \mathbb{1})}\right)=z_{0} v-i y v_{(i y \mathbb{1})}$ contains a bounded sequence as $y \rightarrow \infty$ and, therefore, a subsequence that converges
weakly to some vector, $w$ say. So $(v, w)$ is in the weak closure of the graph of $X$, therefore in the graph of $X$, and hence $v$ is in $\mathcal{D}^{\prime}$.
(v) $\Rightarrow$ (ii) If $v \in \mathcal{D}^{\prime}$, then (6.19) becomes

$$
\begin{equation*}
L(\gamma)=\left\langle\gamma, X v-z_{0} v\right\rangle . \tag{6.30}
\end{equation*}
$$

Let $v_{1}=\left(X-z_{0}\right) v$. Then (6.18) says

$$
\begin{equation*}
(X-z) v_{z}=v_{1} . \tag{6.31}
\end{equation*}
$$

Combining (6.17), (6.30), and (6.31), we get

$$
\begin{equation*}
F_{t}(z)=a-\left\langle v, v_{1}\right\rangle+\left\langle(X-z)^{-1} v_{1}, v_{1}\right\rangle . \tag{6.32}
\end{equation*}
$$

Now let $z=(i y \mathbb{1})$ in (6.32), and let $y \rightarrow \infty$. By Lemma 6.25, the last term on the right tends to zero, so we must have $a-\left\langle v, v_{1}\right\rangle=0$.
(ii) $\Rightarrow$ (i) Lemma 6.25 implies that

$$
\|F(z)\| \leq\left\|v_{1}\right\|^{2} / \min \left(\left|\operatorname{Im} z^{r}\right|\right)
$$

and so (i) follows.
For later use, let us record a slight variant of Theorem 6.26; it is proved in the same way.

Theorem 6.33. Let $F$ be in $\mathcal{L}^{d}$, and assume $F$ has a representation as in (6.10) with $t=0$. Then the following are equivalent:
(i) There exists a constant $C \in \mathbb{R}$ so that

$$
\liminf _{y \rightarrow \infty} y|F(i y \mathbb{1})-C|<\infty .
$$

(ii) There exists a vector $v_{1}$ in $\mathcal{M}$ and a constant $C$ in $\mathbb{R}$ so that

$$
F(z)=C+\left\langle(X-z)^{-1} v_{1}, v_{1}\right\rangle \quad z \in \Pi^{d} .
$$

(iii) The function $\varphi=\beta \circ F \circ \alpha$ in $\mathcal{S}$ has a $B$-point at $\mathbb{1}$, and $\varphi(\mathbb{1}) \neq 1$.
(iv) $\lim _{y \rightarrow \infty} F(i y \mathbb{1})=C \in \mathbb{R}$, and

$$
\liminf _{y \rightarrow \infty} y\left\|v_{(i y \mathbb{1})}\right\|<\infty .
$$

(v) The vector $v$ is in the domain of $X$.

## 7. Locally matrix monotone functions

Recall the definition of locally $n$-matrix monotone.

Definition 1.7. Let $E$ be an open set in $\mathbb{R}^{d}$ and $f$ be a real-valued $C^{1}$-function on $E$. Say $f$ is locally $M_{n}$-monotone on $E$ if, whenever $S$ is in $\operatorname{CSAM}_{n}^{d}$ with $\sigma(S)=\left\{x_{1}, \ldots, x_{n}\right\}$ consisting of $n$ distinct points in $E$, and $S(t)$ is a $C^{1}$-curve in $\operatorname{CSAM}_{n}^{d}$ with $S(0)=S$ and $\left.\frac{d}{d t} S(t)\right|_{t=0} \geq 0$, then $\left.\frac{d}{d t} f(S(t))\right|_{t=0}$ exists and is $\geq 0$.

If $S$ is in $\mathrm{CSAM}_{n}^{d}$, we can choose an orthonormal basis of eigenvectors that diagonalize all the $S^{r}$ 's simultaneously, so

$$
S^{r}=\left(\begin{array}{ccc}
x_{1}^{r} & &  \tag{7.1}\\
& \ddots & \\
& & x_{n}^{r}
\end{array}\right) \quad \forall 1 \leq r \leq d
$$

If $S(t)$ is a $C^{1}$-curve of commuting self-adjoints, then $S(0)+\left.t \frac{d}{d t} S(t)\right|_{t=0}$ commutes to first order; i.e., the commutators are all $o(t)$.

For any $X \in M_{n}$, we define diag $X$ to be the diagonal matrix in $M_{n}$ with diagonal entries $X_{i i}$, and for any $\Delta \in \mathrm{SAM}_{n}^{d}$, we define diag $\Delta$ to be $\left(\operatorname{diag} \Delta^{1}, \ldots, \operatorname{diag} \Delta^{d}\right)$.

Definition 7.2. We shall say that $S$ in $\mathrm{CSAM}_{n}^{d}$ is generic if its spectrum consists of $n$ distinct points.

Lemma 7.3. Let $S$ be in $\operatorname{CSAM}_{n}^{d}$ and $\Delta$ be in $\operatorname{SAM}_{n}^{d}$, with $S$ generic. Then there exists a $C^{1}$-curve $S(t)$ of commuting self-adjoints with $S(0)=S$ and $S^{\prime}(0)=\Delta$ if and only if

$$
\begin{equation*}
\left[S^{r}, \Delta^{s}\right]=\left[S^{s}, \Delta^{r}\right] \quad \forall 1 \leq r \neq s \leq d \tag{7.4}
\end{equation*}
$$

Proof. $(\Rightarrow)$ : If $S(t)=S+t \Delta+o(t)$ is commutative, calculate

$$
\left[S^{r}(t), S^{s}(t)\right]=t\left(\left[S^{r}, \Delta^{s}\right]-\left[S^{s}, \Delta^{r}\right]\right)+o(t)
$$

The coefficient of $t$ must vanish, giving (7.4).
$(\Leftarrow)$ : Suppose $S$ is as in (7.1), suppose and (7.4) holds. This means

$$
\begin{equation*}
\Delta_{i j}^{s}\left(x_{j}^{r}-x_{i}^{r}\right)=\Delta_{i j}^{r}\left(x_{j}^{s}-x_{i}^{s}\right) \quad \forall r \neq s, \tag{7.5}
\end{equation*}
$$

so

$$
\begin{equation*}
\Delta_{i j}^{r} \frac{1}{x_{j}^{r}-x_{i}^{r}}=\Delta_{i j}^{s} \frac{1}{x_{j}^{s}-x_{i}^{s}} \quad \text { if } x_{j}^{r}-x_{i}^{r} \neq 0 \neq x_{j}^{s}-x_{i}^{s} . \tag{7.6}
\end{equation*}
$$

Define a skew-selfadjoint matrix $Y$ by

$$
\begin{equation*}
Y_{i j}=\Delta_{i j}^{r} \frac{1}{x_{j}^{r}-x_{i}^{r}} \quad \text { for any } r \text { such that } x_{j}^{r}-x_{i}^{r} \neq 0 . \tag{7.7}
\end{equation*}
$$

For any $i \neq j$, there is some $r$ with $x_{j}^{r}-x_{i}^{r} \neq 0$, so (7.7) defines $Y_{i j}$; (7.6) says it does not matter which $r$ we choose. Let all the diagonal terms of $Y$ be 0 .

Define

$$
\begin{equation*}
S^{r}(t)=e^{t Y}\left(S^{r}+t \operatorname{diag} \Delta^{r}\right) e^{-t Y} \tag{7.8}
\end{equation*}
$$

Since $e^{t Y}$ is a unitary matrix and $S^{r}+t$ diag $\Delta^{r}$ is diagonal, $S(t) \in \operatorname{CSAM}_{n}^{d}$ and

$$
\left.\frac{d}{d t} S^{r}(t)\right|_{t=0}=\left[Y, S^{r}\right]+\operatorname{diag} \Delta^{r}=\Delta^{r}
$$

If $S$ and $\Delta$ satisfy (7.4) and $S$ is generic, then for any function $f$ that is $C^{1}$ on a neighborhood of $\sigma(S)$, we define the directional derivative of $f$ at $S$ in direction $\Delta$ by

$$
\begin{equation*}
D_{\Delta} f(S)=\left.\frac{d}{d t} f(S(t))\right|_{t=0} \tag{7.9}
\end{equation*}
$$

where $S(t)$ is the curve given by equations (7.8) and (7.7). We shall show in Proposition 7.18 that (7.9) is actually unchanged if $S(t)$ is replaced by any other curve that agrees with it to first order. First, let us show that the right-hand side of (7.9) exists. Indeed,

$$
\begin{equation*}
f(S(t))=e^{t Y} f(S+t \operatorname{diag} \Delta) e^{-t Y} \tag{7.10}
\end{equation*}
$$

Since $S+t$ diag $\Delta$ is diagonal, $f(S+t \operatorname{diag} \Delta)$ is diagonal, with the $i$ th entry

$$
f\left(x_{i}+t \Delta_{i i}\right)=f\left(x_{i}\right)+t \sum_{r=1}^{d} \Delta_{i i}^{r} \frac{\partial f}{\partial x^{r}}\left(x_{i}\right)+o(t) .
$$

In other words,

$$
f(S+t \operatorname{diag} \Delta)=f(S)+t \sum_{r=1}^{d}\left(\operatorname{diag} \Delta^{r}\right) \frac{\partial f}{\partial x^{r}}(S)+o(t)
$$

Hence, on differentiating equation (7.10) at 0 , we obtain

$$
\left.\frac{d}{d t} f(S(t))\right|_{t=0}=[Y, f(S)]+\sum_{r=1}^{d}\left(\operatorname{diag} \Delta^{r}\right) \frac{\partial f}{\partial x^{r}}(S) .
$$

We have shown the following.
Proposition 7.11. Let $S$ be a generic d-tuple of commuting self-adjoint matrices in $M_{n}$. Fix an orthonormal basis of eigenvectors so that every $S^{r}$ is diagonal:

$$
S^{r}=\left(\begin{array}{ccc}
x_{1}^{r} & & \\
& \ddots & \\
& & x_{n}^{r}
\end{array}\right)
$$

Let $\Delta$ be a d-tuple of self-adjoints satisfying (7.4). Let $f$ be $C^{1}$ on a neighborhood of $\left\{x_{1}, \ldots, x_{n}\right\}$ in $\mathbb{R}^{d}$, where each $x_{j}$ is the d-tuple $\left(x_{j}^{1}, \ldots, x_{j}^{d}\right)$. Then

$$
\left[D_{\Delta} f(S)\right]_{i j}= \begin{cases}\Delta_{i j}^{r} \frac{f\left(x_{j}\right)-f\left(x_{i}\right)}{x_{j}^{r}-x_{i}^{r}} & \text { if } i \neq j, \text { where } x_{j}^{r} \neq x_{i}^{r},  \tag{7.12}\\ \left.\sum_{r=1}^{d} \Delta_{i i}^{r} \frac{\partial f}{\partial x^{r}}\right|_{x_{i}} & \text { if } i=j .\end{cases}
$$

Corollary 7.13. For $S, \Delta$ as in Proposition 7.11, if $f, g$ are $C^{1}$-functions that agree to first order on $\sigma(S)$, then $D_{\Delta} f(S)=D_{\Delta} g(S)$.

Lemma 7.14. Let $R$ and $S$ be in $\mathrm{CSAM}_{n}^{d}$. For every point $\mu$ in the joint spectrum of $R$ there is an $x_{p}$ in the joint spectrum of $S$ with

$$
\begin{equation*}
\left\|\mu-x_{p}\right\| \leq \sqrt{d n}\|R-S\| . \tag{7.15}
\end{equation*}
$$

Proof. Choose an orthonormal basis that diagonalizes $S$, so that $S$ is as in (7.1). Let $\Delta=R-S$. Let $\mu$ be a joint eigenvalue of $R$ with corresponding eigenvector $\xi=\left(\xi_{1}, \ldots, \xi_{n}\right)^{t}$. Choose $p$ so that $\left|\xi_{p}\right| \geq\left|\xi_{j}\right|$ for all $1 \leq j \leq n$.

Then for each $1 \leq r \leq d$, we have

$$
R^{r} \xi=\mu^{r} \xi=\left(S^{r}+\Delta^{r}\right) \xi .
$$

So, in particular,

$$
\sum_{j=1}^{n} R_{p j}^{r} \xi_{j}=\mu^{r} \xi_{p}
$$

Therefore,

$$
\left(\mu^{r}-x_{p}^{r}\right) \xi_{p}=\sum_{j=1}^{n} \Delta_{p j}^{r} \xi_{j}
$$

So

$$
\begin{aligned}
\left|\mu^{r}-x_{p}^{r}\right| & \leq \sum_{j=1}^{n}\left|\Delta_{p j}^{r}\right| \\
& \leq \sqrt{n} \sqrt{\sum_{j}\left|\Delta_{p j}^{r}\right|^{2}} \\
& \leq \sqrt{n}\left\|\Delta^{r}\right\|
\end{aligned}
$$

and hence

$$
\sum_{r=1}^{d}\left|\mu^{r}-x_{p}^{r}\right|^{2} \leq d n\|\Delta\|^{2}
$$

Lemma 7.16. If $R(t)$ is a Lipschitz path in $\operatorname{CSAM}_{n}^{d}, 0 \leq t<1$, with $R(0)=S$ generic, then there exists $\varepsilon>0$ and Lipschitz maps $X_{1}, \ldots, X_{n}$ : $[0, \varepsilon) \rightarrow \mathbb{R}^{d}$ such that $\sigma(R(t))=\left\{X_{j}(t): j=1, \ldots, n\right\}$.

Proof. Choose an orthonormal basis that diagonalizes $S$ so that $S$ is as in (7.1). The joint eigenvalues of $S$ are the points $x_{i}=\left(x_{i}^{1}, \ldots, x_{i}^{d}\right)$, and genericity means $\left\|x_{i}-x_{j}\right\|>0$ if $i \neq j$. Choose $\varepsilon$ so that for all $0 \leq t \leq \varepsilon$,

$$
\begin{equation*}
\sqrt{d n}\|R(t)-S\| \leq \frac{1}{3} \min _{i \neq j}\left\|x_{i}-x_{j}\right\| \tag{7.17}
\end{equation*}
$$

By Lemma 7.14, for every joint eigenvalue $x$ of $S$, there is a joint eigenvalue $\mu$ of $R(t)$ within distance $\sqrt{d n}\|R(t)-S\|$ of it. By (7.17), this means that $R(t)$ is also generic, and it makes sense to talk of the joint eigenvalue of $R(t)$ that is closest to $x_{j}$. Let us call these joint eigenvalues $X_{j}(t)$. We have proved that

$$
\left\|X_{j}(t)-x_{j}\right\| \leq \sqrt{d n}\|R(t)-S\| \quad \forall 0 \leq t \leq \varepsilon
$$

Repeating the argument with $R\left(t_{1}\right)$ in place of $S$, we get

$$
\left\|X_{j}\left(t_{2}\right)-X_{j}\left(t_{1}\right)\right\| \leq \sqrt{d n}\left\|R\left(t_{2}\right)-R\left(t_{1}\right)\right\| \quad \forall 0 \leq t_{1}, t_{2} \leq \varepsilon .
$$

As $R$ is assumed to be Lipschitz, we get that each $X_{j}$ is Lipschitz also.
Proposition 7.18. If $S$ is generic in $\mathrm{CSAM}_{n}^{d}, \Delta$ is in $\mathrm{SAM}_{n}^{d}$, and they satisfy the commutation relations (7.4), then for any $C^{1}$-path $R(t) \in \operatorname{CSAM}_{n}^{d}$ such that $R(0)=S, R^{\prime}(0)=\Delta$, and any $f \in C^{1}$,

$$
\begin{equation*}
\left.\frac{d}{d t} f(R(t))\right|_{t=0}=D_{\Delta} f(S) \tag{7.19}
\end{equation*}
$$

Proof. If $g$ is a monomial, then a simple calculation shows that

$$
\left.\frac{d}{d t} g(R(t))\right|_{t=0}
$$

exists and depends only on $g, S$, and $\Delta$. It follows that for any polynomial $g$,

$$
\begin{equation*}
\left.\frac{d}{d t} g(R(t))\right|_{t=0}=\left.\frac{d}{d t} g(S(t))\right|_{t=0}=D_{\Delta} g(S) \tag{7.20}
\end{equation*}
$$

Consider any $f \in C^{1}$, and pick a polynomial $g$ that agrees with $f$ to first order on $\sigma(S)$. By Corollary 7.13,

$$
\begin{equation*}
D_{\Delta} f(S)=D_{\Delta} g(S) \tag{7.21}
\end{equation*}
$$

We claim that

$$
\begin{equation*}
\left.\frac{d}{d t} g(R(t))\right|_{t=0}=\left.\frac{d}{d t} f(R(t))\right|_{t=0} \tag{7.22}
\end{equation*}
$$

By Lemma 7.16, there exist Lipschitz functions $X_{1}, \ldots, X_{n}:[0, \varepsilon) \rightarrow \mathbb{R}^{d}$ such that $\sigma(R(t))=\left\{X_{1}(t), \ldots, X_{n}(t)\right\}$ for all $t$. Then $f(S)=g(S)$ and

$$
\begin{aligned}
\|(f-g)(R(t))\| & =\max _{i}\left|(f-g)\left(X_{i}(t)\right)\right| \\
& =o\left(\max _{i}\left\|X_{i}(t)-X_{i}(0)\right\|\right) \\
& =o(t)
\end{aligned}
$$

Hence

$$
\left\|\frac{f(R(t))-f(S)}{t}-\frac{g(R(t))-g(S)}{t}\right\| \rightarrow 0 \text { as } t \rightarrow 0
$$

In view of equation (7.20),

$$
\left.\frac{f(R(t))-f(S)}{t} \rightarrow \frac{d}{d t} g(R(t))\right|_{t=0}=D_{\Delta} g(S) \text { as } t \rightarrow 0
$$

On combining this relation with equation (7.21), we obtain equation (7.19).
Corollary 7.23. A real-valued $C^{1}$-function $f$ on an open set $E \subseteq \mathbb{R}^{d}$ is locally $M_{n}$-monotone if and only if

$$
D_{\Delta} f(S) \geq 0
$$

for every generic $S$ in $\operatorname{CSAM}_{n}^{d}$ with spectrum in $E$ and every $\Delta$ in $\mathrm{SAM}_{n}^{d}$ such that $\Delta \geq 0$ and

$$
\left[S^{r}, \Delta^{s}\right]=\left[S^{s}, \Delta^{r}\right] \quad \forall 1 \leq r \neq s \leq d
$$

The statement follows immediately from Definition 1.7 and Proposition 7.18.
If one drops the genericity assumption on $S$, the situation becomes significantly more complicated. Nonetheless, K. Bickel has recently analyzed this situation and has shown that both Lemma 7.3 and Proposition 7.18 remain true without the generic assumption on the eigenvalues of $S$ [10]. Passing from differentiating functions of matrices to differentiating functions of operators brings up another set of issues. See [29] for a discussion of these.

We can now characterize locally matrix monotone functions.
Theorem 7.24. Let $E$ be an open set in $\mathbb{R}^{d}$ and $f$ be a real-valued $C^{1}$-function on $E$. Then $f$ is locally $M_{n}$-monotone if and only if $f$ is in $\mathcal{L}_{n}^{d}(E)$.

Proof. $(\Leftarrow)$ We must show: if $S$ is generic with $\sigma(S) \subset E$, if $\Delta$ is a positive $d$-tuple, and if $\left[S^{r}, \Delta^{s}\right]=\left[S^{s}, \Delta^{r}\right]$ for all $r, s$, then $D_{\Delta} f(S) \geq 0$.

Let $\sigma(S)=\left\{x_{1}, \ldots, x_{n}\right\}$. Choose $A^{r}$ as in Definition 5.1. For $i \neq j$, assume without loss of generality that $x_{j}^{1} \neq x_{i}^{1}$. Then

$$
\begin{aligned}
{\left[D_{\Delta} f(S)\right]_{i j} } & =\Delta_{i j}^{1} \frac{f\left(x_{j}\right)-f\left(x_{i}\right)}{x_{j}^{1}-x_{i}^{1}} \\
& =\frac{\Delta_{i j}^{1}}{x_{j}^{1}-x_{i}^{1}}\left(\sum_{r=1}^{d}\left(x_{j}^{r}-x_{i}^{r}\right) A^{r}(i, j)\right) \\
& =\sum_{r=1}^{d} \Delta_{i j}^{r} A^{r}(i, j) .
\end{aligned}
$$

(We get the last line by using (7.6).) By (7.12) the same formula holds for $\left[D_{\Delta} f(S)\right]_{i j}$ when $i=j$, so $D_{\Delta} f(S)$ is the sum of the Schur products of $\Delta^{r}$ with $A^{r}$, so is positive.
$(\Rightarrow)$ Let $f$ be locally $M_{n}$-monotone, and fix $\left\{x_{1}, \ldots, x_{n}\right\}$ distinct points in $E$. Let $S$ be given by (7.1). We wish to find positive semi-definite matrices $A^{r}$ such that (5.3) and (5.2) hold.

Let $\mathcal{G}$ be the set of all skew-symmetric real $n$-by- $n$ matrices $\Gamma$ with the property that there exists a $d$-tuple $A$ of real positive semi-definite matrices satisfying

$$
\begin{align*}
A^{r}(i, i) & =\left.\frac{\partial f}{\partial x^{r}}\right|_{x_{i}} & & 1 \leq i \leq n, 1 \leq r \leq d,  \tag{7.25}\\
\sum_{r=1}^{d}\left(x_{j}^{r}-x_{i}^{r}\right) A^{r}(i, j) & =\Gamma_{i j} & & 1 \leq i \neq j \leq n . \tag{7.26}
\end{align*}
$$

Let $\Lambda$ be the matrix $\Lambda_{i j}=f\left(x_{i}\right)-f\left(x_{j}\right)$. We wish to show $\Lambda$ is in $\mathcal{G}$.
Notice that $\mathcal{G}$ is a closed convex set. Moreover, it is nonempty, because $\left.\frac{\partial f}{\partial x^{r}}\right|_{x_{i}}$ is always greater than or equal to 0 . (This last assertion can be seen by letting $\Delta$ be 0 except in the $r^{\text {th }}$ slot, where it is $I$, and calculating $D_{\Delta} f(S)$.)

So if $\Lambda$ is not in $\mathcal{G}$, by the Hahn-Banach theorem there is a real linear functional $L: M_{n} \rightarrow \mathbb{R}$ that is nonnegative on $\mathcal{G}$ and negative on $\Lambda$. Any such linear functional is of the form $L(T)=\operatorname{tr}(T K)$ for some matrix $K$. Replacing $K$ by $\left(K-K^{t}\right) / 2$ will not change the value of $L$ on skew-symmetric real matrices, so we can assume that there is a real skew-symmetric matrix $K$ such that $\operatorname{tr}(\Gamma K) \geq 0$ for all $\Gamma$ in $\mathcal{G}$, and $\operatorname{tr}(\Lambda K)<0$.

Define $\Delta$ by

$$
\Delta_{i j}^{r}=\left(x_{j}^{r}-x_{i}^{r}\right) K_{j i} .
$$

Then $\Delta$ is in $\operatorname{SAM}_{n}^{d}$, and

$$
\left[\Delta^{s}, S^{r}\right]_{i j}=\left(x_{j}^{s}-x_{i}^{s}\right) K_{j i}\left(x_{j}^{r}-x_{i}^{r}\right)=\left[\Delta^{r}, S^{s}\right],
$$

so $\Delta$ satisfies (7.4).
Moreover $\Delta \geq 0$. Indeed, fix $s$ between 1 and $d$, and let $c_{1}, \ldots, c_{n}$ be complex numbers. We want to show that

$$
\begin{equation*}
\sum_{i, j=1}^{n} c_{i} \bar{c}_{j} \Delta_{i j}^{s} \geq 0 \tag{7.27}
\end{equation*}
$$

For $r \neq s$, let $A^{r}$ be the diagonal matrix with entries given by (7.25). Define $A^{s}$ to be the sum of the diagonal matrix from (7.25) with the rank one matrix $\left[c_{i} \bar{c}_{j}\right]$. Define $\Gamma$ by (7.26). Then $\Gamma$ is in $\mathcal{G}$, and since $K$ and $\Delta^{s}$ both vanish on the diagonal,

$$
\begin{aligned}
\operatorname{tr}(\Gamma K) & =\sum_{i, j=1}^{n}\left(x_{j}^{s}-x_{i}^{s}\right) A^{s}(i, j) K_{j i} \\
& =\sum_{i, j} \Delta_{i j}^{s} c_{i} \bar{c}_{j} \\
& \geq 0,
\end{aligned}
$$

yielding (7.27).

As $f$ is locally $M_{n}$ monotone, we must then have that $D_{\Delta} f(S) \geq 0$. But

$$
\begin{aligned}
0>\operatorname{tr}(\Lambda K) & =\sum_{1 \leq i \neq j \leq n}\left[f\left(x_{j}\right)-f\left(x_{i}\right)\right] K_{j i} \\
& =\sum_{1 \leq i \neq j \leq n} \frac{f\left(x_{j}\right)-f\left(x_{i}\right)}{x_{j}^{r}-x_{i}^{r}} \Delta_{i j}^{r} \\
& =\sum_{1 \leq i, j \leq n}\left[D_{\Delta} f(S)\right]_{i j}
\end{aligned}
$$

which is a contradiction.
As the dimension of the matrices increases, the condition that a function $f$ be locally monotone becomes more stringent. On an infinite-dimensional Hilbert space, the requirement becomes that $f$ be in the Löwner class, as we shall see in the next section.

## 8. Locally operator monotone functions

We defined locally operator monotone functions in Definition 1.6. We shall show that being locally operator monotone is the same as being locally $M_{n}$-monotone for all $n$, which in turn is the same as being in the Löwner class $\mathcal{L}(E)$.

TheOrem 8.1. Let $E$ be an open set in $\mathbb{R}^{d}$ and $f$ be a real-valued $C^{1}$-function on $E$. The following are equivalent:
(i) The function $f$ is locally $M_{n}$-monotone on $E$ for all $n \geq 1$.
(ii) The function $f$ is in $\mathcal{L}(E)$.
(iii) The function $f$ is locally operator monotone on $E$.

The equivalence of (i) and (ii) follows from Theorem 7.24 and Proposition 5.11. The implication (iii) $\Rightarrow$ (i) is obvious. We need to prove that (ii) $\Rightarrow$ (iii). First we need some preliminary results.

Proposition 8.2. Let $E$ be an open set in $\mathbb{R}^{d}$, and let $f \in \mathcal{L}(E)$. Then there is a model $(\mathcal{M}, v)$ for $f$ such that $v_{z}$ is locally Lipschitz on $E$.

Proof. By Proposition 5.11 we can extend $f$ to a function $F$ in $\mathcal{L}$ that extends analytically across $E$ and agrees with $f$ on $E$. For this $F$ we have a model $(\mathcal{M}, v)$ so that

$$
\begin{equation*}
F(z)-\overline{F(w)}=\left\langle\left(z-w^{*}\right) v_{z}, v_{w}\right\rangle_{\mathcal{M}} \quad \forall z, w \in E \cup \Pi^{d} \tag{8.3}
\end{equation*}
$$

and by Proposition 3.15, if $w$ is in $E$, then $v_{w}$ is the limit of $v_{z}$ as $z$ tends to $w$ nontangentially from inside $\Pi^{d}$.

Fix $w$ in $E$ (so $F(w)$ is real). Then, by analyticity, for $z$ close to $w$, we have

$$
\begin{equation*}
F(z)-F(w)=\left.\sum_{r=1}^{d} \frac{\partial f}{\partial x^{r}}\right|_{w}\left(z^{r}-w^{r}\right)+\left.\sum_{|\alpha| \geq 2} \frac{\partial^{\alpha} f}{\partial x^{\alpha}}\right|_{w} \frac{(z-w)^{\alpha}}{\alpha!} . \tag{8.4}
\end{equation*}
$$

From (8.3), we get

$$
\begin{aligned}
F(z)-F(w) & =\left\langle(z-w) v_{z}, v_{w}\right\rangle \\
& =\sum_{r=1}^{d}\left(z^{r}-w^{r}\right)\left\langle v_{w}^{r}, v_{w}^{r}\right\rangle+\sum_{r=1}^{d}\left(z^{r}-w^{r}\right)\left\langle v_{z}^{r}-v_{w}^{r}, v_{w}^{r}\right\rangle .
\end{aligned}
$$

As $z$ tends to $w$ nontangentially, the second term is $o(\|z-w\|)$, so comparing with (8.4), we see that

$$
\begin{equation*}
\left\|v_{w}^{r}\right\|^{2}=\left.\frac{\partial f}{\partial x^{r}}\right|_{w} \quad \forall 1 \leq r \leq d \tag{8.5}
\end{equation*}
$$

Now let $z$ and $w$ both be in $E$. Comparing (8.3) and (8.4), we get

$$
\begin{equation*}
\left\langle(z-w) v_{z}, v_{w}\right\rangle-\left\langle(z-w) v_{w}, v_{w}\right\rangle=\left.\sum_{|\alpha| \geq 2} \frac{\partial^{\alpha} f}{\partial x^{\alpha}}\right|_{w} \frac{(z-w)^{\alpha}}{\alpha!} . \tag{8.6}
\end{equation*}
$$

Swapping $z$ and $w$, we get

$$
\begin{equation*}
\left\langle(w-z)\left(v_{w}-v_{z}\right), v_{z}\right\rangle=\left.\sum_{|\alpha| \geq 2} \frac{\partial^{\alpha} f}{\partial x^{\alpha}}\right|_{z} \frac{(w-z)^{\alpha}}{\alpha!} \tag{8.7}
\end{equation*}
$$

Subtracting (8.6) from (8.7), we get

$$
\begin{equation*}
\left\langle(z-w)\left(v_{z}-v_{w}\right),\left(v_{z}-v_{w}\right)\right\rangle=\sum_{|\alpha| \geq 2}\left(\left.(-1)^{|\alpha|} \frac{\partial^{\alpha} f}{\partial x^{\alpha}}\right|_{z}-\left.\frac{\partial^{\alpha} f}{\partial x^{\alpha}}\right|_{w}\right) \frac{(z-w)^{\alpha}}{\alpha!} . \tag{8.8}
\end{equation*}
$$

But since $f$ is analytic,

$$
\left(\left.\frac{\partial^{\alpha} f}{\partial x^{\alpha}}\right|_{z}-\left.\frac{\partial^{\alpha} f}{\partial x^{\alpha}}\right|_{w}\right)=O(\|z-w\|)
$$

and so the right-hand side of (8.8) is $O\left(\|z-w\|^{3}\right)$. Therefore,

$$
\begin{equation*}
\left\langle(z-w)\left(v_{z}-v_{w}\right),\left(v_{z}-v_{w}\right)\right\rangle=O\left(\|z-w\|^{3}\right) . \tag{8.9}
\end{equation*}
$$

If all the differences $\left|z^{r}-w^{r}\right|$ are comparable, we can conclude immediately that $\left\|v_{z}-v_{w}\right\|=O(\|z-w\|)$. If they are not, we can get round this difficulty by connecting $z$ to $w$ by two line segments.

Indeed, suppose $\max _{1 \leq r \leq d}\left|z^{r}-w^{r}\right|=\varepsilon$. Choose numbers $a^{r}$ and $b^{r}$ with modulus between $1 / 2$ and 2 so that

$$
z^{r}-w^{r}=\left(a^{r}-b^{r}\right) \varepsilon \quad \forall 1 \leq r \leq d
$$

Let

$$
x^{r}=w^{r}+a^{r} \varepsilon=z^{r}+b^{r} \varepsilon .
$$

Then applying (8.9) to the pairs $(z, x)$ and $(x, w)$, we get

$$
\begin{aligned}
\left\|v_{z}-v_{w}\right\| & \leq\left\|v_{z}-v_{x}\right\|+\left\|v_{x}-v_{w}\right\| \\
& =O(\|z-x\|+\|x-w\|) \\
& =O(\|z-w\|)
\end{aligned}
$$

as desired.
Suppose now $E, f$, and $(\mathcal{M}, v)$ are as in Proposition 8.2. So $v: z \mapsto v_{z}$ is a map from $E$ to $\mathcal{M}$. Let $S$ be a $d$-tuple of bounded commuting self-adjoint operators on a Hilbert space $\mathcal{H}$, with $\sigma(S) \subset E$. We want to define an operator $\tilde{v}(S) \in B(\mathcal{H}, \mathcal{H} \otimes \mathcal{M})$.

We do this by choosing an orthonormal basis for $\mathcal{M}$, and writing

$$
v(z):=v_{z}=\left(\begin{array}{c}
v_{1}(z)  \tag{8.10}\\
v_{2}(z) \\
\vdots
\end{array}\right)
$$

We must caution the reader that the subscripts on the right-hand side of (8.10) run over the dimension of $\mathcal{M}$, and they are not interchangeable with the superscripts on $v$ that identify which piece of $\mathcal{M}^{1}, \ldots, \mathcal{M}^{d}$ one is in. After using the orthonormal basis to identify $\mathcal{M}$ with $\ell^{2}$, we have

$$
v_{z}=\left[\begin{array}{c}
v_{z}^{1} \\
\vdots \\
v_{z}^{d}
\end{array}\right]=\left(\begin{array}{c}
v_{1}(z) \\
v_{2}(z) \\
\vdots
\end{array}\right)
$$

Then

$$
\tilde{v}(S):=\left(\begin{array}{c}
v_{1}(S)  \tag{8.11}\\
v_{2}(S) \\
\vdots
\end{array}\right): \mathcal{H} \rightarrow \mathcal{H} \otimes \mathcal{M}
$$

The operator $\tilde{v}(S)$ is bounded, because if $S$ has spectral measure $\Lambda$ and $\xi$ is a unit vector in $\mathcal{H}$, then

$$
\begin{align*}
\|\tilde{v}(S) \xi\|^{2} & =\sum_{j} \int_{\sigma(S)}\left|v_{j}\right|^{2} d\langle\Lambda \xi, \xi\rangle  \tag{8.12}\\
& =\int_{\sigma(S)} \sum_{r=1}^{d} \frac{\partial f}{\partial x^{r}} d\langle\Lambda \xi, \xi\rangle \\
& \leq\left.\sup _{z \in \sigma(S)} \sum_{r=1}^{d} \frac{\partial f}{\partial x^{r}}\right|_{z} \tag{8.13}
\end{align*}
$$

and the last sum is finite because $\sigma(S)$ is compact and $f$ is $C^{1}$.

The operator $\tilde{v}(S)$ does not depend on the choice of orthonormal basis in $\mathcal{M}$. A simple calculation shows that for any $h \in \mathcal{H}$ and $m \in \mathcal{M}$,

$$
\tilde{v}(S)^{*}(h \otimes m)=\langle m, v(\cdot)\rangle(S) h .
$$

This gives a coordinate-free expression for $\tilde{v}(S)^{*}$, hence also for $\tilde{v}(S)$.
Lemma 8.14. Let $E, f$ and $(\mathcal{M}, v)$ be as in Proposition 8.2, and let $S$ be a d-tuple of bounded commuting self-adjoint operators on a Hilbert space $\mathcal{H}$, with $\sigma(S) \subset E$. Then

$$
\begin{equation*}
\|\tilde{v}(S)\| \leq\left(\sup _{\sigma(S)} \sum_{r=1}^{d} \frac{\partial f}{\partial x^{r}}\right)^{\frac{1}{2}} \tag{8.15}
\end{equation*}
$$

Moreover, $\tilde{v}$ is continuous.
Proof. Inequality (8.15) has been proved in (8.13). To prove continuity of $\tilde{v}$, let $K$ be a compact subset of $E$ with $\sigma(S) \subset \operatorname{int}(K) \subset E$. Let $\varepsilon>0$.

As $v$ is continuous on $K$ and $K$ is compact, there exists $N$ such that $\sum_{j=N+1}^{\infty}\left|v_{j}(z)\right|^{2} \leq \varepsilon^{2} / 9$ on $K$. For $1 \leq j \leq N$, there is a polynomial $p_{j}$ such that $\left\|p_{j}-v_{j}\right\|_{\infty} \leq \varepsilon / 9 N$ on $K$. There exists $\delta>0$ so that if $\left\|T^{r}-S^{r}\right\| \leq \delta$, then $\sigma(T) \subseteq K$ and $\left\|p_{j}(T)-p_{j}(S)\right\| \leq \varepsilon / 9 N$ for each $1 \leq j \leq N$.

Let $\tilde{v}_{N}(S)$ be the operator

$$
\left(\begin{array}{c}
v_{1}(S) \\
\vdots \\
v_{N}(S) \\
0 \\
\vdots
\end{array}\right)
$$

Then $\left\|\tilde{v}_{N}(S)-\tilde{v}(S)\right\| \leq \varepsilon / 3$ by (8.12), and similarly $\left\|\tilde{v}_{N}(T)-\tilde{v}(T)\right\| \leq \varepsilon / 3$. As

$$
\left\|v_{j}(T)-v_{j}(S)\right\| \leq\left\|v_{j}(T)-p_{j}(T)\right\|+\left\|p_{j}(T)-p_{j}(S)\right\|+\left\|p_{j}(S)-v_{j}(S)\right\|
$$

and each entry is at most $\varepsilon / 9 N$,

$$
\left\|\tilde{v}_{N}(S)-\tilde{v}_{N}(T)\right\| \leq N\left(\frac{\varepsilon}{9 N}+\frac{\varepsilon}{9 N}+\frac{\varepsilon}{9 N}\right)=\frac{\varepsilon}{3},
$$

and hence

$$
\|\tilde{v}(T)-\tilde{v}(S)\| \leq \varepsilon
$$

We assume that $\mathcal{M}$ is decomposed as $\mathcal{M}=\mathcal{M}^{1} \oplus \cdots \oplus \mathcal{M}^{d}$ and that $P^{r}$ is the orthogonal projection from $\mathcal{M}$ onto $\mathcal{M}^{r}$. If $S=\left(S^{1}, \ldots, S^{d}\right)$ is a $d$-tuple of operators on $\mathcal{H}$, we shall write

$$
\begin{equation*}
S \odot I:=S^{1} \otimes P^{1} \oplus \cdots \oplus S^{d} \otimes P^{d} \tag{8.16}
\end{equation*}
$$

which is an operator on $\mathcal{H} \otimes \mathcal{M}$.

Proposition 8.17. Let $E$ be open in $\mathbb{R}^{d}$, let $f \in \mathcal{L}(E)$, and assume $(\mathcal{M}, v)$ is a model of $f$ for which $v$ is continuous. Let $S$ and $T$ be d-tuples of commuting self-adjoint operators on a Hilbert space $\mathcal{H}$ with spectrum in $E$. Then

$$
\begin{equation*}
f(T)-f(S)=\tilde{v}(S)^{*}[T \odot I-S \odot I] \tilde{v}(T) \tag{8.18}
\end{equation*}
$$

Proof. First assume that $S$ and $T$ are (separately) diagonalizable. Let $\xi$ be an eigenvector of $S$ and $\eta$ be an eigenvector of $T$, so for some numbers $z^{r}, w^{r}$, we have

$$
\begin{aligned}
& S^{r} \xi=w^{r} \xi, \\
& T^{r} \eta=z^{r} \eta \quad \forall 1 \leq r \leq d .
\end{aligned}
$$

Then

$$
\langle[f(T)-f(S)] \eta, \xi\rangle_{\mathcal{H}}=\left\langle\left[f(z)-f(w)^{*}\right] \eta, \xi\right\rangle_{\mathcal{H}} .
$$

Also,

$$
\begin{aligned}
& \left\langle\tilde{v}(S)^{*}[T \odot I-S \odot I] \tilde{v}(T) \eta, \xi\right\rangle_{\mathcal{H}} \\
& \quad=\langle[T \odot I-S \odot I] \eta \otimes v(z), \xi \otimes v(w)\rangle_{\mathcal{H} \otimes \mathcal{M}} \\
& \quad=\sum_{r=1}^{d}\langle\eta, \xi\rangle_{\mathcal{H}}\left\langle\left(z^{r}-\bar{w}^{r}\right) v^{r}(z), v^{r}(w)\right\rangle_{\mathcal{M}^{r}} \\
& \quad=(f(z)-f(w))\langle\eta, \xi\rangle_{\mathcal{H}} .
\end{aligned}
$$

So both sides of (8.18) agree if you apply them to an eigenvector of $T$ and then take the inner product with an eigenvector of $S$. By linearity, this is true also for linear combinations of eigenvectors, and as these are assumed dense, we get that (8.18) holds.

If $S$ and $T$ are not diagonalizable, by the spectral theorem we can approximate them in norm by operators that are, and as $\tilde{v}$ and $f$ are both continuous, one gets (8.18) in the limit.

Proof of Theorem 8.1. Assume $f$ is in $\mathcal{L}(E)$, and $S(t)$ is a curve of commuting self-adjoint $d$-tuples with $S(0)=S$ and $S^{\prime}(0)=\Delta \geq 0$. Choose a model $(\mathcal{M}, v)$ with $v$ locally Lipschitz. Then by Proposition 8.17,

$$
f(S(t))-f(S)=\tilde{v}(S)^{*}[(S(t)-S) \odot I] \tilde{v}(S(t))
$$

As

$$
S(t)=S+t \Delta+o(t)
$$

we get

$$
\begin{aligned}
\left.\frac{d}{d t} f(S(t))\right|_{0} & =\lim _{t \rightarrow 0} \tilde{v}(S)^{*}[\Delta \odot I] \tilde{v}(S(t))+\lim _{t \rightarrow 0} \tilde{v}(S)^{*}[o(1)] \tilde{v}(S(t)) \\
& =\tilde{v}(S)^{*}[\Delta \odot I] \tilde{v}(S)
\end{aligned}
$$

Hence $f(S(t))$ is differentiable at 0 , and its derivative is a positive semi-definite operator.

## 9. Globally operator monotone functions

Definition 9.1. Let $E$ be an open set in $\mathbb{R}^{d}$ and $f$ be a real-valued $C^{1}$-function on $E$. Say $f$ is globally operator monotone on $E$ if, whenever $S$ and $T$ are $d$-tuples of commuting bounded self-adjoint operators on a Hilbert space with $\sigma(S) \cup \sigma(T) \subset E$, and $S \leq T$, then $f(S) \leq f(T)$.

If $F$ has the form in Theorem 6.26, then it is globally monotone on boxes in the $\mu$-resolvent of $X$.

Theorem 9.2. Let $X$ be a densely-defined self-adjoint operator on a graded Hilbert space $\mathcal{M}=\mathcal{M}^{1} \oplus \cdots \oplus \mathcal{M}^{d}$, let $v \in \mathcal{M}$, and let $F$ be given by

$$
F(z)=\left\langle(X-z)^{-1} v, v\right\rangle .
$$

Let $E$ be an open box in $\mathbb{R}^{d}$ that is in the $\mu$-resolvent of $X$. Then $F$ is globally operator monotone on $E$.

Proof. First observe that if $S$ is a commuting $d$-tuple of self-adjoint operators on $\mathcal{H}$ and $\sigma(S) \subset E$, then

$$
\begin{equation*}
F(S)=R_{v}^{*}\left(I_{\mathcal{H}} \otimes X-S \odot I\right)^{-1} R_{v}, \tag{9.3}
\end{equation*}
$$

where $\odot$ is as in equation (8.16) and

$$
\begin{aligned}
R_{v}: & \mathcal{H} \\
& \rightarrow \mathcal{H} \otimes \mathcal{M} \\
& h \mapsto h \otimes v .
\end{aligned}
$$

Thus equation (9.3) means that for any vectors $\xi$ and $\eta$ in $\mathcal{H}$,

$$
\begin{equation*}
\langle F(S) \eta, \xi\rangle_{\mathcal{H}}=\left\langle\left(I_{\mathcal{H}} \otimes X-\sum_{r=1}^{d} S^{r} \otimes P^{r}\right)^{-1} \eta \otimes v, \xi \otimes v\right\rangle_{\mathcal{H} \otimes \mathcal{M}} . \tag{9.4}
\end{equation*}
$$

Indeed, if $\eta$ is an eigenvector of $S$ with eigenvalues $a^{r}$, then $F(S) \eta=F(a) \eta$, so the left-hand side of (9.4) is $F(a)\langle\eta, \xi\rangle$. But we have

$$
\left(I_{\mathcal{H}} \otimes X-\sum_{r=1}^{d} S^{r} \otimes P^{r}\right)^{-1} \eta \otimes v=\eta \otimes(X-a)^{-1} v
$$

as one can verify by applying $\left(I_{\mathcal{H}} \otimes X-\sum_{r=1}^{d} S^{r} \otimes P^{r}\right)$ to both sides. So the right-hand side of (9.4) is

$$
\langle\eta, \xi\rangle\left\langle(X-a)^{-1} v, v\right\rangle,
$$

which is the same as the left-hand side of (9.4). If $S$ has a spanning set of eigenvectors, our claim is proved. If it does not, one can approximate it in norm by a $d$-tuple that does, and the claim follows by continuity.

Now let $S$ and $T$ be $d$-tuples of commuting self-adjoint operators with $\sigma(S) \cup \sigma(T) \subset E$ and $\Delta:=T-S \geq 0$. Let

$$
R^{r}(t)=(1-t) S^{r}+t T^{r}, \quad 1 \leq r \leq d .
$$

Then for $t$ in the range $(0,1)$, the $d$-tuple $R(t)$ will consist of self-adjoint operators that need not commute with each other. Nonetheless, letting

$$
Y(t)=\left(I_{\mathcal{H}} \otimes X-R(t) \odot I\right),
$$

then $R_{v}^{*} Y(t)^{-1} R_{v}$ still makes sense by Lemma 9.5. Moreover,

$$
\frac{d}{d t} Y(t)^{-1}=Y(t)^{-1}(\Delta \odot I) Y(t)^{-1}
$$

and so is positive. Therefore,

$$
\begin{aligned}
F(T)-F(S) & =R_{v}^{*} Y(1)^{-1} R_{v}-R_{v}^{*} Y(0)^{-1} R_{v} \\
& =R_{v}^{*} \int_{0}^{1} \frac{d}{d t} Y(t)^{-1} d t R_{v} \\
& =\int_{0}^{1} R_{v}^{*} Y(t)^{-1}(\Delta \odot I) Y(t)^{-1} R_{v} d t \\
& \geq 0 .
\end{aligned}
$$

Lemma 9.5. Let $a^{r}<b^{r}, 1 \leq r \leq d$, and let $X$ be a densely defined self-adjoint operator on a graded Hilbert space $\mathcal{M}=\oplus_{r=1}^{d} \mathcal{M}^{r}$. Suppose that for every $t$ in $(0,1)$, the point $\lambda_{t}=(1-t) a+t b$ is not in the $\mu$-spectrum of $X$. Let $S=\left(S^{1}, \ldots, S^{d}\right)$ be a d-tuple of bounded self-adjoint operators on a Hilbert space $\mathcal{H}$, with $\sigma\left(S^{r}\right) \subset\left(a^{r}, b^{r}\right)$ for each $r$. Then $I \otimes X-\sum_{r=1}^{d} S^{r} \otimes P^{r}$ is invertible, with a bounded inverse.

Proof. First, suppose $a^{r}=-1$ and $b^{r}=1$ for each $r$. Then $(-1,1) \cap \sigma(X)$ is empty, and so is $(-1,1) \cap \sigma(I \otimes X)$. So $\left\|I_{\mathcal{H}} \otimes X \xi\right\| \geq\|\xi\|$ for every $\xi$ in $\mathcal{H} \otimes \mathcal{M}$. But if $\sigma\left(S^{r}\right) \subset(-1,1)$ for each $r$, the operator $\sum S^{r} \otimes P^{r}$ has norm less than one. Therefore $I \otimes X-\sum_{r=1}^{d} S^{r} \otimes P^{r}$ is invertible.

In the general case, let $m^{r}$ be the midpoint and $c^{r}$ be half the length of the interval $\left(a^{r}, b^{r}\right)$, so $a^{r}=m^{r}-c^{r}, b^{r}=m^{r}+c^{r}$. Let

$$
Y=\left(\sum_{r=1}^{d} \frac{1}{\sqrt{c^{r}}} I \otimes P^{r}\right)\left(I \otimes X-\sum_{r=1}^{d} m^{r} I_{\mathcal{H}} \otimes P^{r}\right)\left(\sum_{r=1}^{d} \frac{1}{\sqrt{c^{r}}} I \otimes P^{r}\right) .
$$

If $\sigma_{\mu}$ denotes the $\mu$-spectrum,

$$
\sigma_{\mu}(Y)=c^{-1}\left(\sigma_{\mu}(X)-m\right),
$$

and hence the point $(1-t)(-\mathbb{1})+t \mathbb{1}$ lies in the $\mu$-resolvent set of $Y$ for $0<t<1$. Let $T^{r}=\left(1 / c^{r}\right)\left(S^{r}-m^{r} I_{\mathcal{H}}\right)$. Then $T^{r}$ is a strict contraction, and so, by the
previous case, $Y-\sum T^{r} \otimes P^{r}$ is invertible. As

$$
\begin{aligned}
& Y-\sum_{r=1}^{d} T^{r} \otimes P^{r} \\
& =\left(\sum \frac{1}{\sqrt{c^{r}}} I \otimes P^{r}\right)\left(I \otimes X-\sum S^{r} \otimes P^{r}\right)\left(\sum \frac{1}{\sqrt{c^{r}}} I \otimes P^{r}\right)
\end{aligned}
$$

we get the desired result.
We can now prove a global result for rational functions of two variables.
TheOrem 9.6. Let $F$ be a rational function of two variables. Let $\Gamma$ be the zero-set of the denominator of $F$. Assume $F$ is real-valued on $\mathbb{R}^{2} \backslash \Gamma$. Let $E$ be an open rectangle in $\mathbb{R}^{2} \backslash \Gamma$. Then $F$ is globally operator monotone on $E$ if and only if $F$ is in $\mathcal{L}(E)$.

Proof. Necessity follows from Theorem 8.1. For sufficiency, by Lemma 9.7 it is sufficient to prove the theorem for $F_{t_{n}}$ with $t_{n} \searrow 0$, where $F_{t}=\rho_{t} \circ F \circ \rho_{t}$. Suppose the degree of $F$ is $n^{1}$ in $z^{1}$ and $n^{2}$ in $z^{2}$. Let $\varphi=\beta \circ F \circ \alpha$. By a result of G. Knese [21], there is a model for $\varphi$ in a Hilbert space $\mathcal{M}=\mathcal{M}^{1} \oplus \mathcal{M}^{2}$ with $\operatorname{dim}\left(\mathcal{M}^{r}\right)=n^{r}$ for $r=1$ and 2 ; see also the paper [8] by J. A. Ball, C. Sadosky, and V. Vinnikov.

Accordingly, in Theorem 6.9, we obtain a realization of $F_{t}$ on $\mathcal{M}^{1} \oplus \mathcal{M}^{2}$ of the form (6.10); since $\mathcal{M}$ is finite-dimensional, the vector $v$ is in the domain of $X$. By Theorem $6.33(\mathrm{v}) \Rightarrow(\mathrm{ii})$, for some $v_{1} \in \mathcal{M}$,

$$
F_{t}(z)=C+\left\langle(X-z)^{-1} v_{1}, v_{1}\right\rangle_{\mathcal{M}}
$$

where $\operatorname{dim}\left(\mathcal{M}^{r}\right)=n^{r}$ for $r=1$ and 2 , and $F_{t}(\infty, \infty)=C<\infty$. Then the pole-set $\Gamma_{t}$ of $F_{t}$ is contained in the zero-set of $\operatorname{det}(X-z)$, which is a rational function of degree $\left(\operatorname{dim}\left(\mathcal{M}^{1}\right), \operatorname{dim}\left(\mathcal{M}^{2}\right)\right)$. As these two algebraic sets have the same degree, they must be equal. So the $\mu$-resolvent of $X$ is $\mathbb{R}^{2} \backslash \Gamma_{t}$, and now the result follows from Theorem 9.2.

Let $\rho_{t}$ be as in (6.8). The following lemma is elementary.
Lemma 9.7. Let $t>0$. Let $U$ be an open set in $\mathbb{R}^{d}$. Then
(i) The function $F$ is globally operator monotone on $U \cap(-1 / t, \infty)^{d}$ if and only if $F \circ \rho_{t}$ is globally operator monotone on $\rho_{t}^{-1}(U) \cap(-\infty, 1 / t)^{d}$.
(ii) The function $F$ is globally operator monotone on $U \cap F^{-1}(-1 / t, 1 / t)$ if and only if $\rho_{t} \circ F$ is globally operator monotone on the same set.

What happens to Theorem 9.6 in $d \geq 3$ variables? It is still true that rational Löwner functions have finite-dimensional models [12], [8]. However, a recent example of Knese [22] shows that the minimal dimension $n^{r}$ needed may be strictly greater than the degree of $F$ in $z^{r}$. So we cannot rule out the possibility that the $\mu$-spectrum of $X$ contains some other algebraic sets in $\mathbb{R}^{d}$ than just the zero set of the denominator of $F$.

We now give an example of a nonrational globally operator monotone function.

Example 9.8. For $0 \leq s \leq 1 / 2$, the function $\left(z^{1} z^{2}\right)^{s}$ is operator monotone on $(0, \infty) \times(0, \infty)$. Indeed, if $(0,0)<\left(A^{1}, A^{2}\right) \leq\left(B^{1}, B^{2}\right)$ and $s$ is between 0 and $1 / 2$, then

$$
\left\|\left(A^{r}\right)^{s}\left(B^{r}\right)^{-s}\right\| \leq 1 \quad \text { for } r=1,2
$$

Therefore, the norm of

$$
\left(B^{1}\right)^{-s}\left(A^{1}\right)^{s}\left(A^{2}\right)^{s}\left(B^{2}\right)^{-s}
$$

is less than or equal to 1 , so the largest eigenvalue is less than or equal to 1 , and therefore the largest eigenvalue of

$$
\begin{equation*}
\left(B^{2}\right)^{-s / 2}\left(B^{1}\right)^{-s / 2}\left(A^{1}\right)^{s}\left(A^{2}\right)^{s}\left(B^{1}\right)^{-s / 2}\left(B^{2}\right)^{-s / 2} \tag{9.9}
\end{equation*}
$$

is also less than or equal to 1 . But (9.9) is self-adjoint, so less than or equal to the identity. Therefore,

$$
\left(A^{1} A^{2}\right)^{s} \leq\left(B^{1} B^{2}\right)^{s} .
$$

We do not know if $\left(z^{1} z^{2}\right)^{s}$ can be approximated by rational functions in the Löwner class.

Let us close with some questions.

- Is Theorem 9.6 true for rational functions of more than two variables?
- Can $E$ be an arbitrary open set in Theorem 9.2?
- Is every function in $\mathcal{L}(E)$ globally operator monotone on $E$ ?
- Is every function in $\mathcal{L}_{n}(E) M_{n}$-monotone on $E$ ?


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[^1]:    ${ }^{1}$ This formula says that for a $C^{1}-\operatorname{arc} S(t)=U(t) \Lambda(t) U^{*}(t)$, with $U(t)$ unitary and $\Lambda(t)$ diagonal with diagonal entries $\lambda_{1}(t), \ldots, \lambda_{n}(t)$, and a $C^{1}$ function $f$, then

    $$
    \frac{d}{d t} f(S(t))=U(t)\left(\left[\Delta f\left(\lambda_{i}(t), \lambda_{j}(t)\right)\right] \circ\left[U(t)^{*} S^{\prime}(t) U(t)\right]\right) U(t)^{*}
    $$

[^2]:    ${ }^{2}$ The notation $y \searrow 0$ means $y$ decreases to 0 . The notation $r \nearrow 1$ means $r$ increases to 1 .

[^3]:    ${ }^{3}$ We say an operator $T$ has strictly positive imaginary part if there exists $\alpha>0$ such that $\left(T-T^{*}\right) / 2 i \geq \alpha I$.

