# THE CHARACTERIZATION OF CONTINUOUS, FOUR-COEFFICIENT SCALING FUNCTIONS AND WAVELETS* 

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#### Abstract

We examine four-coefficient dilation equations and give results converse to a theorem of Daubechies-Lagarias. These results complete the characterization of those four-coefficient dilation equations having a continuous solution.


A wavelet basis for $L^{2}(\mathbf{R})$ is an orthonormal basis $\left\{2^{n / 2} \psi\left(2^{n} x-m\right)\right\}_{n, m \in \mathbf{Z}}$ generated from a single function $\psi$, the wavelet. The classical example is the Haar system, where $\psi=\chi_{[0,1 / 2)}-\chi_{[1 / 2,1)}$. Wavelet bases have many applications, e.g., image and speech processing. The variety of applications demands that wavelet bases having specific properties be available. It is therefore important to have means by which wavelets with desired properties can be constructed. One method is to solve a dilation equation or two-scale difference equation $f(t)=\sum_{k=-\infty}^{\infty} c_{k} f(2 t-k)$. A functional solution $f$ to a dilation equation is called a scaling function. The wavelet $\psi$ can be realized as $\psi(t)=\sum(-1)^{k} c_{1-k} f(2 t-k)$ whenever the scaling function $f$ defines a multiresolution analysis, cf. Example 1. Recently, numerous papers have addressed the issue of constructing scaling functions, e.g., the various subdivision schemes in [1], [2], [3], [4], [5], [6], [7]. The purpose of this paper is to characterize four-coefficient dilation equations having continuous solutions. One part of this characterization (Theorem 1) was provided by Daubechies and Lagarias [4], [5]. We obtained converse results (Theorem 2) which completes this characterization and which shows that the bounds for smoothness obtained in [5] are exact [8].

In the case of four-coefficient systems, it is sufficient to consider only those $\left\{c_{k}\right\}$ for which at most $c_{0}, c_{1}, c_{2}, c_{3}$ are nonzero, i.e., dilation equations of the form

$$
\begin{equation*}
f(t)=c_{0} f(2 t)+c_{1} f(2 t-1)+c_{2} f(2 t-2)+c_{3} f(2 t-3) \tag{1}
\end{equation*}
$$

We are interested in only those solutions $f$ to (1) which are integrable, hence supp $f \subset[0,3]$. We further require that the coefficients be real and satisfy $c_{0}+c_{2}=c_{1}+c_{3}=1$. This condition is necessary for the existence of a multiresolution analysis. Thus, we have a two-parameter family of equations, and we select the independent parameters to be $c_{0}$ and $c_{3}$. Each scaling function $f$ is then associated with a point in the $\left(c_{0}, c_{3}\right)$-plane. Since the time-reversed function $f(3-t)$ is a solution to the system determined by $\left(c_{3}, c_{0}\right)$, any result obtained regarding this plane will be symmetric about the line $c_{0}=c_{3}$. Finally, recall that a function $f$ is Hölder continuous with Hölder exponent $\alpha$ if there exists a constant $C$ for which $|f(x)-f(y)| \leq C|x-y|^{\alpha}$ for all $x, y$. An extension of this definition to vector-valued functions is immediate.

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Fig. 1. Circle of orthogonality, ellipse, dashed/solid line, and triangle.
Example 1. To illustrate the range of phenomena that can occur even for the fourcoefficient setting, we exhibit in Figure 1 several geometric objects in the ( $c_{0}, c_{3}$ )-plane: the ellipse $c_{0}^{2}+c_{3}^{2}-c_{0}-c_{3}+c_{0} c_{3}=1$; the circle $\left(c_{0}-1 / 2\right)^{2}+\left(c_{3}-1 / 2\right)^{2}=1 / 2$; a triangle; and a solid/dashed line. In [9] we show that there are no integrable solutions to (1) for ( $c_{0}, c_{3}$ ) on or outside the ellipse, with the exception of the point $\left(c_{0}, c_{3}\right)=(1,1)$. There is an integrable solution to (1) for every point on and inside the circle. Neither of these results is sharp, and they are improved upon in [9]. Furthermore, Cohen [10] and Lawton [11] have characterized the $N$-coefficient systems $\left\{c_{k}\right\}_{k=0}^{N}$ that give a multiresolution analysis. In the $\left(c_{0}, c_{3}\right)$-plane, these are identified as all points on the circle except $(1,1)$. We refer to the circle as the circle of orthogonality and associated scaling functions as orthogonal scaling functions. We see below that the triangle gives a first approximation to the region of points resulting in continuous scaling functions. The line identifies a necessary condition in order that a scaling function be differentiable. The intersection of this line and the circle is the well-known Daubechies scaling function $D_{4}$. Differentiable solutions to (1) do exist, and are indicated by the solid line segment in Figure 1. No four-coefficient orthogonal scaling function is differentiable ( $D_{4}$ is differentiable "almost everywhere"). Finally, $\left(c_{0}, c_{3}\right)=(1,0)$ gives rise to the scaling function $\chi_{[0,1)}$ whose wavelet basis is the Haar system. The points $(0,1)$ and $(0,0)$ are translates of this system, and $(1,1)$ is a "stretched" version resulting in the non-orthogonal scaling function $\chi_{[0,3)}$.

We begin by outlining related work of [4], [5], where a sufficient condition is given
for the continuity of scaling functions. We assume familiarity with the notation and main results in those papers. Given a dilation equation (1), continuity of the scaling function $f$ can be examined in the following way. First define the matrix

$$
M=\left(\begin{array}{ll}
c_{1} & c_{0} \\
c_{3} & c_{2}
\end{array}\right)
$$

The eigenvalues of $M$ are 1 and $1-c_{0}-c_{3}$. In order for $f$ to be continuous we must have $f(0)=f(3)=0$. Using (1) with $t=1,2$, we obtain the values for $f(1)$ and $f(2)$ as components of the eigenvector associated with the eigenvalue 1 for $M$, cf. [12], therefore $f(1)=c_{0}$ and $f(2)=c_{3}$. (When $c_{0}+c_{3} \neq 0$, it is customary to normalize these values so that $f(1)+f(2)=1$, cf. [5].) For $f$ to be differentiable, $M$ must also have the eigenvalue $1 / 2$, so that $c_{0}+c_{3}=1 / 2$. This is the equation of the line plotted in Figure 1. No four-coefficient scaling function can have more than one continuous derivative.

Next define the matrices

$$
T_{0}=\left(\begin{array}{ccc}
c_{0} & 0 & 0 \\
c_{2} & c_{1} & c_{0} \\
0 & c_{3} & c_{2}
\end{array}\right) \quad \text { and } \quad T_{1}=\left(\begin{array}{ccc}
c_{1} & c_{0} & 0 \\
c_{3} & c_{2} & c_{1} \\
0 & 0 & c_{3}
\end{array}\right)
$$

and note that $M$ is a submatrix of both $T_{0}$ and $T_{1}$. Assuming $f$ exists, we let $v:[0,1] \rightarrow \mathbf{R}^{3}$ be the vector-valued function

$$
v(x)=\left(\begin{array}{c}
f(x)  \tag{2}\\
f(x+1) \\
f(x+2)
\end{array}\right)
$$

Then $v(0), v(1)$ are right eigenvectors associated with the eigenvalue 1 for $T_{0}, T_{1}$, respectively. Let $\tau$ denote the left shift operator on $[0,1]$, so that $\tau x=2 x$ for $0 \leq x<1 / 2$ and $\tau x=2 x-1$ for $1 / 2 \leq x \leq 1$. Expanding $x$ into its binary representation $x=. d_{1} d_{2} \ldots$, we have $v(x)=T_{d_{1}} v(\tau x)$ for all $x \in[0,1]$. The definition $\tau(1 / 2)=0$ can be replaced by $\tau(1 / 2)=1$ without affecting this property of $v$ since $T_{0} v(1)=T_{1} v(0)$. By induction, $v(x)=T_{d_{1}} \cdots T_{d_{m}} v\left(\tau^{m} x\right)$ for every $m>0$. Using the matrices $T_{0}, T_{1}$ and the vector representation $v$, we can then determine the value of $f$ at any dyadic point $\left\{k / 2^{n}\right\}_{k, n \in \mathbf{z}}$. For example, $\left(f\left(2^{-m}\right), f\left(1+2^{-m}\right), f\left(2+2^{-m}\right)\right)^{\mathrm{t}}=v\left(2^{-m}\right)=T_{0}^{m-1} T_{1} v(0)=T_{0}^{m} v(1)$. The scaling function $f$ is continuous if and only if $v$ is continuous and satisfies

$$
\begin{align*}
& v_{1}(0)=v_{3}(1)=0  \tag{3}\\
& v_{i+1}(0)=v_{i}(1) \quad \text { for } i=1,2,3  \tag{4}\\
& v(x)=T_{d_{1}} v(\tau x) \quad \text { for } x \in[0,1] \tag{5}
\end{align*}
$$

In the case that $f$ is continuous, $v\left(2^{-m}\right) \rightarrow v(0)=\left(0, c_{0}, c_{3}\right)$. We compute

$$
v\left(2^{-m}\right)-v(0)=\frac{1}{\left(1-2 c_{0}-c_{3}\right)}\left(\begin{array}{c}
\left(1-2 c_{0}-c_{3}\right) c_{0}^{m+1}  \tag{6}\\
\left(2 c_{0}-1\right) c_{0}^{m+1}+c_{3}\left(1-c_{0}-c_{3}\right)^{m+1} \\
c_{3} c_{0}^{m+1}+c_{3}\left(1-c_{0}-c_{3}\right)^{m+1}
\end{array}\right)
$$

and therefore $v$ continuous at zero forces $\left|c_{0}\right|,\left|1-c_{0}-c_{3}\right|<1$. By symmetry, the continuity of $v$ at 1 implies $\left|c_{3}\right|,\left|1-c_{0}-c_{3}\right|<1$. Since the eigenvalues of $T_{0}$ and $T_{1}$ combined are $1, c_{0}, c_{3}$, and $1-c_{0}-c_{3}$, continuity of $f$ implies that all eigenvalues of $T_{0}, T_{1}$ other than 1 are less than one in absolute value and that the Hölder exponent of $v$ is at most $-\log _{2}\left(\max \left\{\left|c_{0}\right|,\left|c_{3}\right|,\left|1-c_{0}-c_{3}\right|\right\}\right)$. The region of points $\left\{\left(c_{0}, c_{3}\right):\left|c_{0}\right|,\left|c_{3}\right|,\left|1-c_{0}-c_{3}\right|<1\right\}$ is the interior of the triangle shown in Figure 1. Some orthogonal scaling functions are therefore discontinuous, and, in fact, unbounded.

Conversely, a scaling function $f$ can sometimes be constructed from a vector function $v$. Given $\left(c_{0}, c_{3}\right)$, set $v(0)=\left(0, c_{0}, c_{3}\right)^{\mathrm{t}}$ and $v(1)=\left(c_{0}, c_{3}, 0\right)^{\mathrm{t}}$, the right eigenvectors for the eigenvalue 1 for $T_{0}, T_{1}$. Then define $v$ on the set of dyadic points on $[0,1]$ using (5). Conditions on the coefficients $\left\{c_{k}\right\}$ that give uniform continuity for $v$ on the dyadics allows us to extend $v$ to a continuous function on $[0,1]$. If $v$ also satisfies (3)-(4), the continuous scaling function $f$ is then obtained from $v$ by (2). More specifically, we seek bounds on the Hölder exponent of $v$, which is identical to the Hölder exponent of $f$. The function $f$ constructed in this way is the unique integrable solution to (1), up to multiplication by a constant, and we say that $f$, respectively $v$, is the scaling function, respectively, vector scaling function, associated with the point $\left(c_{0}, c_{3}\right)$.

Obtaining the appropriate bounds for the Hölder exponent requires an examination of the joint spectral radius (defined below) for the matrix operators $T_{0}, T_{1}$ restricted to the two-dimensional subspace $V=\left\{u \in \mathbf{R}^{3}: u_{1}+u_{2}+u_{3}=0\right\}$, cf. [5]. In particular, since $(1,1,1)$ is a common left eigenvalue for $T_{0}$ and $T_{1}$, the action of $T_{i}$ on $V$ is equivalent to the $2 \times 2$ matrices $S_{0}, S_{1}$ acting on $\mathbf{R}^{2}$, where

$$
S_{0}=\left(\begin{array}{cc}
c_{0} & 0 \\
-c_{3} & 1-c_{0}-c_{3}
\end{array}\right) \quad \text { and } \quad S_{1}=\left(\begin{array}{cc}
1-c_{0}-c_{3} & -c_{0} \\
0 & c_{3}
\end{array}\right)
$$

Except for the eigenvalue 1, the eigenvalues of $S_{0}, S_{1}$ are the same as those of $T_{0}, T_{1}$, respectively. Before we can state the result of [5], we need some definitions.

Let $\|\cdot\|$ be a norm on $\mathbf{R}^{3}$, with corresponding operator norm $\|A\|=\sup \{\|A u\| /\|u\|$ : $u \neq 0\}$ for a matrix $A$. Typical norms are $\|u\|_{p}=\left(\left|u_{1}\right|^{p}+\left|u_{2}\right|^{p}+\left|u_{3}\right|^{p}\right)^{1 / p}$ for fixed $p \geq 1$, and $\|u\|_{\infty}=\max \left\{\left|u_{1}\right|,\left|u_{2}\right|,\left|u_{3}\right|\right\}$ for $p=\infty$.

Definition 1. The joint spectral radius $\hat{\rho}\left(A_{0}, A_{1}\right)$ of two matrices $A_{0}, A_{1}$ is defined by $\hat{\rho}\left(A_{0}, A_{1}\right)=\limsup \lambda_{m}$, where $\lambda_{m}=\max _{d_{j}=0,1}\left\|A_{d_{1}} \cdots A_{d_{m}}\right\|^{1 / m}$.

This definition generalizes the usual spectral radius of a single matrix $A$, which is given by $\rho(A)=\limsup \left\|A^{m}\right\|^{1 / m}=\max \{|\lambda|: \lambda$ is an eigenvalue of $A\}$. An extension to larger collections of matrices is made in the obvious way. To our knowledge, the joint spectral radius was first discussed by Rota and Strang [13]. Some recent papers on the subject include [14], [15]. The joint spectral radius is independent of the choice of norm $\|\cdot\|$, and of the choice of basis, that is, $\hat{\rho}\left(A_{0}, A_{1}\right)=\hat{\rho}\left(B_{0}, B_{1}\right)$ whenever $B_{i}=B A_{i} B^{-1}$ for any fixed matrix $B$, cf. [5]. We are interested in $\hat{\rho}\left(S_{0}, S_{1}\right)=\hat{\rho}\left(\left.T_{0}\right|_{V},\left.T_{1}\right|_{V}\right)$. The importance of the joint spectral radius in this work is not accidental, for example, cf. [6], where a probabilistic approach is considered.

The first step in characterizing the continuous scaling functions is provided by the following theorem.

Theorem 1 [5]. Given $\left(c_{0}, c_{3}\right)$, if $\hat{\rho}\left(S_{0}, S_{1}\right)=\hat{\rho}\left(\left.T_{0}\right|_{V},\left.T_{1}\right|_{V}\right)<1$ then the associated scaling function is Hölder continuous with Hölder exponent $\alpha \geq-\log _{2} \hat{\rho}\left(S_{0}, S_{1}\right)-\varepsilon$ for every $\varepsilon>0$.

Let $\sigma_{m}=\max _{d_{j}=0,1} \rho\left(S_{d_{1}} \cdots S_{d_{m}}\right)^{1 / m}$, so that $\sigma_{m}$ is the largest absolute eigenvalue that occurs among all products of length $m$ of the matrices $S_{0}, S_{1}$. Then $\sigma_{m} \leq \hat{\rho}\left(S_{0}, S_{1}\right) \leq$ $\lambda_{m}$ for every $m$. Berger and Wang [14] prove that $\hat{\rho}\left(S_{0}, S_{1}\right)=\lim \sup \sigma_{m}$. It therefore follows that $\sup \sigma_{m}=\hat{\rho}\left(S_{0}, S_{1}\right)=\lim \lambda_{m}$. In particular, $\hat{\rho}\left(S_{0}, S_{1}\right) \geq \sigma_{1}$ and in general one expects $\hat{\rho}\left(S_{0}, S_{1}\right)$ to be strictly larger than $\sigma_{1}$. Since $\sigma_{1}=\max \left\{\left|c_{0}\right|,\left|c_{3}\right|,\left|1-c_{0}-c_{3}\right|\right\}$, the set $\left\{\left(c_{0}, c_{3}\right): \sigma_{1}<1\right\}$ is the triangular region in Figure 1.

Now consider the set $C S=\left\{\left(c_{0}, c_{3}\right): \hat{\rho}\left(S_{0}, S_{1}\right)<1\right\}$. Theorem 1 implies that the dilation equation (1) has a continuous solution for every point $\left(c_{0}, c_{3}\right) \in C S$. A major goal of this paper is to (numerically) determine the region where continuous solutions exist. We do this by bounding $\hat{\rho}\left(S_{0}, S_{1}\right)$ by $\lambda_{m}$ and $\sigma_{m}$ for various $m$. Since $\lambda_{m}$ is norm dependent, let $\lambda_{m}(p)$ denote that value of $\lambda_{m}$ obtained using the norm $\|\cdot\|_{p}$, and define the set $C_{m}(p)=\left\{\left(c_{0}, c_{3}\right): \lambda_{m}(p)<1\right\}$. Every $C_{m}(p)$ is a subset of $C S$, so each point in $C_{m}(p)$ is associated with a continuous scaling function. Since for each fixed $p, C_{m}(p)$ increases to the set $C S$, we can approximate $C S$ with $C_{m}(p)$ for various $m$ and $p$. The regions $C_{1}(1), C_{1}(2), C_{1}(\infty)$, determined by considering products of length one and the $\|\cdot\|_{1}$-norm, $\|\cdot\|_{2}$-norm, and $\|\cdot\|_{\infty}$-norm, respectively, are plotted in [16]. Figure 2 shows a numerical estimation of the set $C_{16}(1)$.


Fig. 2. The set $C_{16}(1)=\left\{\left(c_{0}, c_{3}\right): \lambda_{16}<1\right\}$ (shaded area) .

The calculation of $\lambda_{m}$ or $\sigma_{m}$ can become computationally intensive even for small values of $m$. This makes the exact determination of the joint spectral radius nearly impossible, except in special cases. One case when it is possible is the following.

Lemma 1. If $\left\{S_{0}, S_{1}\right\}$ can be simultaneously symmetrized, i.e., there exists an invertible matrix $B$ such that $B S_{i} B^{-1}$ is real and symmetric for $i=0,1$, then $\hat{\rho}\left(S_{0}, S_{1}\right)=\sigma_{1}$.

Tedious but elementary calculations allow us to determine those points of the $\left(c_{0}, c_{3}\right)$ plane for which simultaneous symmetrization is possible.

Proposition 1. If $\sigma_{1}<1$ and if $R=c_{0} c_{3}\left(1-c_{0}-c_{3}\right)\left(1-2 c_{0}-2 c_{3}\right) \geq 0$ then the scaling function associated with $\left(c_{0}, c_{3}\right)$ is continuous with Hölder exponent exactly $-\log _{2} \sigma_{1}$.

Lemma 1 holds also when $S_{0}, S_{1}$ can be simultaneously Hermitianized. One might then hope that the hypothesis of Proposition 1 can be weakened to include that case as well. However, for the matrices $S_{0}, S_{1}$ simultaneous symmetrization is possible exactly when simultaneous Hermitianization is possible.

Example 2. The region $S S$, defined by $\sigma_{1}<1$ and $R \geq 0$, is depicted in Figure 3.


Fig. 3. The region $S S$ where simultaneous symmetrization is possible and leads to continuous scaling functions (shaded area).

If simultaneous symmetrization is not possible and it is not practical to compute $\lambda_{m}$ directly, it may still be possible to bound $\hat{\rho}\left(S_{0}, S_{1}\right)$ from above by computing a small number of matrix product norms. The following result can be found in [5].

Lemma 2. Let $\left\{P_{j}\right\}$ be a set of building blocks of products of the matrices $S_{i}$, i.e.,
(a) each $P_{j}$ is some product of $m_{j}$ of the matrices $S_{i}$, and
(b) there is some $r \geq 0$ such that if $P$ is any product of the matrices $S_{i}$ then $P=$ $P_{j_{1}} \cdots P_{j_{k}} Q$, where $Q$ is some product of at most $r$ of the matrices $S_{i}$. Then $\hat{\rho}\left(S_{0}, S_{1}\right) \leq \sup \left\|P_{j}\right\|^{1 / m_{j}}$.

The following recursive algorithm can be used to implement Lemma 2.
Algorithm 1. Given $\hat{\rho}>\hat{\rho}\left(S_{0}, S_{1}\right)$. For each of the matrices $S_{0}, S_{1}$ in turn, implement the following recursion.

Given a product $P=S_{d_{1}} \cdots S_{d_{m}}$. If $\|P\|^{1 / m}<\hat{\rho}$ then keep $P$ as a building block. Otherwise, repeat this step with each of the products $P S_{0}, P S_{1}$.

This recursion will clearly terminate since there must exist some $\lambda_{m}$ with $\hat{\rho}\left(S_{0}, S_{1}\right) \leq$ $\lambda_{m}<\hat{\rho}$. The resulting set of building blocks will satisfy Lemma 2 with $r=l-1$, where $l$ is the length of the longest product in the set of building blocks. An analogous recursive algorithm for bounding $\hat{\rho}\left(S_{0}, S_{1}\right)$ from below based on the eigenvalues of products of matrices is not possible, since the eigenvalues of a product are not, in general, related to the eigenvalues of the matrices in the product.

Example 3. We consider points on the circle of orthogonality. If we could always simultaneously symmetrize $S_{0}$ and $S_{1}$ then the resulting Hölder exponent would be exactly $-\log _{2} \sigma_{1}$ for every point. On the circle, this expression is maximized when $\sigma_{1}$ is minimized, at $\left(c_{0}, c_{3}\right)=(0.6,-0.2)$. Then, $\sigma_{1}=0.6$, so the Hölder exponent would be $\approx 0.737$ (note the Hölder exponent of $D_{4}$ is $\approx 0.550$ ). Unfortunately, at $(0.6,-0.2)$ we have $R=$ $-0.0144<0$, so symmetrization is not possible. We estimate the Hölder exponent by direct bounding of $\hat{\rho}\left(S_{0}, S_{1}\right)$. Considering all $2^{13}$ products of length 13, we find that $\hat{\rho}\left(S_{0}, S_{1}\right) \leq \lambda_{13}(1) \approx 0.682455<1$. The associated scaling function is therefore continuous, with Hölder exponent $\alpha \geq-\log _{2} \lambda_{13}(1) \approx 0.551$, i.e., $f$ is smoother than $D_{4}$. (Note $-\log _{2} \lambda_{m}(1)<0.550$ for $m=1, \ldots, 12$.) Since the only orthogonal scaling function which lies on the line $c_{0}+c_{3}=1 / 2$, and so the only one which "could be" differentiable, is $D_{4}$, it is surprising that there exists a four-coefficient, orthogonal scaling function smoother than $D_{4}$.

Algorithm 1 can be used to improve the above estimate. Allowing a maximum product length of 73 in that algorithm results in the bound $\hat{\rho}\left(S_{0}, S_{1}\right) \leq 0.660500$, which implies $\alpha \geq 0.598370$. This search required 14156 matrix norm computations, of which 7079 matrix products were selected as building blocks. A direct computation of $\lambda_{73}(1)$ would not improve this estimate (e.g., $\left\|S_{0}^{2} S_{1} S_{0}^{14} S_{1} S_{0}^{14} S_{1} S_{0}^{13} S_{1} S_{0}^{12} S_{1} S_{0}^{12} S_{1}\right\|^{1 / 73}=0.663200$ ), and would, in fact, be impossible to carry out.

Since the above computation involved products of significant length, it is difficult to believe that the Hölder exponent of $f$ could be much greater than 0.598370 . In order to show this, we bound $\alpha$ from above by bounding $\hat{\rho}\left(S_{0}, S_{1}\right)$ from below. Direct computation of $2^{31}-1$ matrix products reveals that $\max \left\{\sigma_{1}, \ldots, \sigma_{30}\right\}=\sigma_{13}=\rho\left(S_{1} S_{0}^{12}\right)^{1 / 13} \approx 0.659679$. Therefore, by Theorem 2 below, $\alpha \leq-\log _{2} \sigma_{13} \approx 0.600164$. Furthermore, since the Hölder exponent of continuity for any scaling function is at least $-\log _{2} \sigma_{1}$, we can easily identify points on the circle of orthogonality whose scaling functions cannot have Hölder exponent
greater than 0.598370 , the lower bound for the Hölder exponent for $f$. In particular, the smoothest orthogonal scaling function must come from a point somewhere near ( $0.6,-0.2$ ). We do not know if there exists an orthogonal scaling function smoother than $f$.

Theorem 1 implies continuity for the scaling function whenever $\hat{\rho}\left(S_{0}, S_{1}\right)$ is strictly less than one. We now consider the converse problem, i.e., whether this joint spectral radius condition exactly characterizes the continuous scaling functions.

Fix the point $\left(c_{0}, c_{3}\right)$ and assume $f$ is an integrable solution to (1). Note first that if $f(0) \neq 0$ (so $f$ is discontinuous), then $c_{0}$, one of the eigenvalues for $T_{0}$, is one. Thus $\hat{\rho}\left(\left.T_{0}\right|_{V},\left.T_{1}\right|_{V}\right) \geq 1$. A similar argument holds for $f(3)$ and $c_{3}$. We therefore assume $f(0)=f(3)=0$, and that $c_{0}, c_{3} \neq 1$.

Since our results depend on the eigenvalues of matrices, and since eigenvalues are in general complex, we need to consider the complex space $\mathbf{C}^{3}$. Let $A$ be a matrix with real or complex entries and having eigenvalue $\lambda$. Then $U_{\lambda}=\left\{u \in \mathbf{C}^{3}:(A-\lambda)^{k} u=0\right.$ for some $k>$ $0\}$ is an $A$-invariant subspace of $\mathbf{C}^{3}$, for if $u \in U_{\lambda}$ then $A u=(A-\lambda) u+\lambda u \in U_{\lambda}$. The standard Jordan decomposition for $A$ ensures that every vector $v \in \mathbf{C}^{3}$ can be written uniquely as a sum of vectors $v_{\lambda} \in U_{\lambda}$ for the distinct eigenvalues $\lambda$ of $A$. We say that a vector $v \in \mathbf{C}^{3}$ has a component in $U_{\lambda}$ if $v_{\lambda} \neq 0$.

THEOREM 2. Let $v$ be the vector scaling function associated with the point $\left(c_{0}, c_{3}\right)$. Suppose $T=T_{d_{1}} \cdots T_{d_{m}}$ is any fixed product of the matrices $T_{0}, T_{1}$, and $\lambda$ is an eigenvalue for $\left.T\right|_{V}$. Let $x$ be that point whose binary decimal expansion is $x=. d_{1} \ldots d_{m} d_{1} \ldots d_{m} \ldots$.
(a) If $|\lambda|>1$, and if there exists some $z \in[0,1]$ such that $v(z)$ has a component in $U_{\lambda}$ then $v$ is unbounded.
(b) If $v$ is continuous and there exists some $z \in[0,1]$ such that $v(x)-v(z)$ has a component in $U_{\lambda}$ then $|\lambda|<1$ and the Hölder exponent of $v$ is at most $-\log _{2}|\lambda|^{1 / m}$.

For four-coefficient systems, the requirement in Theorem 2 that $v(z)$, respectively $v(x)-v(z)$, has a component in the eigenspace $U_{\lambda}$ for some $z$ is always satisfied whenever $\lambda \neq 0$. As we shall see, Proposition 2 below implies this fact for all but one special case. When $\lambda=0$, part (a) above does not apply and part (b) remains valid with essentially no bound on the Hölder exponent.

Proposition 2. Assume that $\{v(z): z \in[0,1]\}$ is not contained in any line in $\mathbf{C}^{3}$. Let $T$ and $x$ be defined as in Theorem 2, and let $\lambda$ be any eigenvalue of $\left.T\right|_{V}$. Then there exists a $z \in[0,1]$ such that $v(z)$, respectively $v(x)-v(z)$, has a component in $U_{\lambda}$.

We now proceed to show $\{v(z): z \in[0,1]\}$ is not contained in a line in $\mathbf{C}^{3}$ unless $1-c_{0}-c_{3}=0$. Recalling that $v(0)=\left(0, c_{0}, c_{3}\right)^{\mathrm{t}}$ and $v(1)=\left(c_{0}, c_{3}, 0\right)^{\mathrm{t}}$, we compute $v(1 / 2)=T_{0} v(1)=\left(c_{0}^{2}, c_{0}\left(1-c_{0}\right)+c_{3}\left(1-c_{3}\right), c_{3}^{2}\right)^{\mathrm{t}}$. If it were the case that $v(0), v(1 / 2)$, and $v(1)$ did lie on a line in $\mathbf{C}^{3}$ then $v(0)-v(1 / 2)$ would be a multiple of $v(0)-v(1)$. This implies immediately that $1-c_{0}-c_{3}=0$.

When $1-c_{0}-c_{3}=0$ then $\{v(z): z \in[0,1]\}$ does indeed lie on a line in $\mathbf{C}^{3}$. For example, if $c_{0}=c_{3}=1 / 2$ then $v(z)=(z, 1,1-z)^{\mathrm{t}}$. We consider this case separately. Let $1-c_{0}-c_{3}=0$, and suppose $T=T_{d_{1}} \cdots T_{d_{m}}$. Then 1 is an eigenvalue for $T$ since $(1,1,1)$ is a common left eigenvector for $T_{0}$ and $T_{1}$ for this eigenvalue. Let $w_{1}$ be a right
eigenvector for $T$ for the eigenvalue 1. Since $\left(c_{3}^{2},-c_{0} c_{3}, c_{0}^{2}\right)$ is a common left eigenvector for $T_{0}$ and $T_{1}$ for the eigenvalue 0 , that value is also an eigenvalue for $T$. Let $w_{3}$ be a right eigenvector for $T$ for the eigenvalue 0 . To find the remaining eigenvalue of $T$, note that $w_{2}=\left(-c_{0}, c_{0}-c_{3}, c_{3}\right)^{\mathrm{t}}$ is a right eigenvector for $T_{0}$ for the eigenvalue $c_{0}$ and is a right eigenvector for $T_{1}$ for the eigenvalue $c_{3}$. Therefore if $k$ of the $d_{i}$ are 0 and the remaining $m-k$ are 1 then $w_{2}$ is a right eigenvector for $T$ for the eigenvalue $c_{0}^{k} c_{3}^{m-k}$. The three eigenvalues of $T$ are distinct since $c_{0}, c_{3} \neq 0,1$, and so $w_{1}, w_{2}, w_{3}$ is a basis for $\mathbf{C}^{3}$. Therefore, there exist scalars $\alpha_{i}, \beta_{i}, \gamma_{i}$ such that $v(0)=\alpha_{0} w_{1}+\beta_{0} w_{2}+\gamma_{0} w_{3}$ and $v(1)=\alpha_{1} w_{1}+\beta_{1} w_{2}+\gamma_{1} w_{3}$. However, $v(0)-v(1)=\left(0, c_{0}, c_{3}\right)^{\mathrm{t}}-\left(c_{0}, c_{3}, 0\right)^{\mathrm{t}}=w_{2}$, and so $\alpha_{0}=\alpha_{1}, \beta_{0}=\beta_{1}+1$, and $\gamma_{0}=\gamma_{1}$. Hence, either $v(0)$ or $v(1)$ must have a component in the eigenspace of $T$ corresponding to the eigenvalue $c_{0}^{k} c_{3}^{m-k}$. This is one eigenvalue of $\left.T\right|_{V}$. The other eigenvalue of $\left.T\right|_{V}$ is zero, and although neither $v(0)$ nor $v(1)$ need have a component in the zero eigenspace, this eigenspace is irrelevant for the conclusions of Theorem 2.

A summary of the results for four-coefficient systems is given in the next theorem.
ThEOREM 3. Let $f$ be the scaling function associated with the point $\left(c_{0}, c_{3}\right)$.
(a) If $\hat{\rho}\left(S_{0}, S_{1}\right)<1$ then $f$ is continuous.
(b) If $\hat{\rho}\left(S_{0}, S_{1}\right)>1$ then $f$ is unbounded.
(c) If $f$ is continuous then $\hat{\rho}\left(S_{0}, S_{1}\right) \leq 1$ with $\sigma_{m}<1$ for every $m$ and $f$ is Hölder continuous with Hölder exponent $-\log _{2} \hat{\rho}\left(S_{0}, S_{1}\right)-\varepsilon \leq \alpha \leq-\log _{2} \hat{\rho}\left(S_{0}, S_{1}\right)$ for every $\varepsilon>0$.

Theorem 3(c) is almost a complete converse of Theorem 3(a). If it was the case that a solution existed to a dilation equation with $\sup \sigma_{m}=\hat{\rho}\left(\left.T_{0}\right|_{V},\left.T_{1}\right|_{V}\right)=1$ and $\sigma_{m}<1$ for every $m$, then Theorem 3(a) would not imply that this solution was continuous, and Theorem 3(c) would not imply that it was discontinuous. However, it is shown in [8] that a solution for this special case must be discontinuous, and therefore continuity occurs if and only if $\hat{\rho}\left(S_{0}, S_{1}\right)<1$. We do not know if this special case is possible.

Example 4. For each $m$, define $E_{m}=\left\{\left(c_{0}, c_{3}\right): \sigma_{m} \geq 1\right\}$. The set $E_{1}$ is the boundary and exterior of the triangle in Figure 1. By Theorem 2, no scaling function determined by a point in $E_{m}$ can be continuous. In Figure 4 we display the curve which represents the boundary of the set $E_{16}$. This curve serves as an outer approximation for $C S$, i.e., every point on or outside the curve has $\hat{\rho}\left(S_{0}, S_{1}\right) \geq 1$. The solid region in Figure 4 is the largest subset of $C S$ which we have explicitly determined, namely, the union of the sets $S S, C_{16}(1)$, and $C_{16}(2)$. The boundary of this region serves as an inner approximation for $C S$. The actual boundary of $C S$ therefore lies somewhere between the boundary of $E_{16}$ and the boundary of the region $S S \cup C_{16}(1) \cup C_{16}(2)$.

Generalizations of Theorem 2 and Proposition 2 remain valid for $N$-coefficient systems and are proved in [8]. The work of Berger and Wang [6] indicate that the hypotheses regarding the components of $v(z)$ in Theorem 2 cannot be removed when $N>3$, e.g., consider the remarks on "stretched dilation equations" found in that paper. It is conjectured in [8] that for each fixed $N$ the collection of coefficients $\left\{c_{0}, c_{1}, \ldots, c_{N}\right\}$ where this condition fails is a set of measure zero.


Fig. 4. The set $S S \cup C_{16}(1) \cup C_{16}(2)$ (shaded area); boundary of the set $E_{16}=\left\{\left(c_{0}, c_{3}\right)\right.$ : $\left.\sigma_{16} \geq 1\right\}$ (solid line).

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