# HILBERT SPACES OF MARTINGALES SUPPORTING CERTAIN SUBSTITUTION-DYNAMICAL SYSTEMS 

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Dedicated to the memory of Shizuo Kakutani


#### Abstract

Let $X$ be a compact Hausdorff space. We study finite-to-one mappings $r: X \rightarrow X$, onto $X$, and measures on the corresponding projective limit space $X_{\infty}(r)$. We show that certain quasi-invariant measures on $X_{\infty}(r)$ correspond in a one-to-one fashion to measures on $X$ which satisfy two identities. Moreover, we identify those special measures on $X_{\infty}(r)$ which are associated via our correspondence with a function $V$ on $X$, a Ruelle transfer operator $R_{V}$, and an equilibrium measure $\mu_{V}$ on $X$.


## Contents

1. Introduction
2. Projective limits
2.1. Definitions
2.2. A scaling equation
3. Induction of measures on $X_{\infty}(r)$ from measures on $X$
3.1. Martingales
3.2. From $\mu$ on $X$ to $\hat{\mu}$ on $X_{\infty}(r)$
3.3. Quasi-invariance
4. A fixed-point problem
5. Transformations of measures
6. Extreme points

Acknowledgments
References

## 1. Introduction

This paper is motivated by our desire to apply wavelet methods to some nonlinear problems in symbolic and complex dynamics. Recent research by many authors (see, e.g., [4, 2, 6]) on iterated function systems with affine scaling has suggested that the scope of the multiresolution method is wider than the more traditional wavelet context.

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In this paper we exploit the connections between multiscale-resolution analysis (wavelet theory) and dynamics, in particular iteration theory. For this purpose, we study certain quasi-invariant measures, extending the familiar scaling equation from wavelet theory. Our proofs use martingale theory, a transfer operator, called the Ruelle transfer operator, and a fixed point theorem. A by-product of our analysis is Theorem 5.9 generalizing a theorem of Perron and Frobenius.

As suggested first by S. Mallat [29, a particularly productive approach to the construction of wavelet bases in $L^{2}(\mathbb{R})$ may be based on the notion of resolution from optics. In this model, a fixed image-resolution corresponds to a certain closed subspace $V_{0}$ in $L^{2}(\mathbb{R})$. A coarser resolution will be a proper subspace of $V_{0}$, and a finer resolution is a subspace containing $V_{0}$. What results is a $\mathbb{Z}$-indexed scale of nested Hilbert subspaces $V_{n}, n \in \mathbb{Z}$, in $L^{2}(\mathbb{R})$ such that the intersection is $\{0\}$ and the union is dense. Moreover the operation of dyadic scaling transforms each $V_{n}$ to the next $V_{n+1}$. This is called a multiresolution approach to wavelets, see [12], and it is based on the interplay between the two abelian groups $\mathbb{T}$ (the circle group $=$ one-torus), and $\mathbb{R}$ (the real line), with $\mathbb{T}$ representing a period interval placed on the line $\mathbb{R}$.

Our present paper is based on the observation that multiresolutions really are martingales; and on our exploiting this fact. Thus we are able to adapt the geometric notion of scaling subspaces $\left(V_{n}\right)$, or resolutions, to nonlinear dynamics in a variety of applications where such a pair of groups is not available. In the context of discrete dynamical systems, we have instead an endomorphism $r: X \rightarrow X$ of a compact space. And then successive iterations of $r$ may serve to define a certain selfsimilarity (for example statistical similarity up to scale in time or in space), mirroring the more familiar notion of scale-similarity which is such a powerful tool in wavelet theory. The usefulness of this geometric viewpoint is illustrated for example in multivariable operator theory (e.g., [10], [23], [25], [26], [33], and [34]), as well as in recent work on encoding of quantum stochastic processes; see, e.g., [1], [3], [13], [18, 28, and 30.

The notion of multiresolution (from wavelets) is in fact closely related to more familiar concepts from dynamics, e.g., to sequences of partitions or tower constructions in dynamics, such as laminations [39, or successive subdivisions 40]; to Pesin boxes [9], to attractors [38].

Let $r$ be an endomorphism in a compact metric space $X$ (for example the Julia set of a given rational map), and suppose $r$ is onto $X$ and finite-to-one. Form a projective space $X_{\infty}$ such that $r$ induces an automorphism $\hat{r}$ of $X_{\infty}$. Let $V$ be a Borel function on $X$ (naturally extended to a function on $X_{\infty}$ ). Generalizing the more traditional approach to scaling functions, we give in Theorem4.1 a complete classification of measures on $X_{\infty}$ which are quasi-invariant under $\hat{r}$ and have RadonNikodym derivative equal to $V$; see also (1.1) and Section 2.2 for further details.

The interest in quasi-invariance derives in part from the observation that in the more familiar wavelet case, this invariance notion plays a fundamental role, that of a scaling equation (see [12, Chapter 5]). But in the wavelet case, the measure is prescribed a priori; it is the Lebesgue measure. Not so for other scaling problems in dynamics. Starting with a possibly nonlinear system $(X, r)$, we must examine measures $\mu$ on $X$ which are suitably invariant for the inverse branches of the given endomorphism $r: X \rightarrow X$. This invariance amounts to equation (2.4) below, and we typically assume that $\mu$ (if it exists) is normalized. If a probability measure $\mu$
exists, solving equation (2.4), we say that $\mu$ is strongly invariant. But existence of strongly invariant measures is now a question. (If for example $r: X \rightarrow X$ is known to have strictly contractive inverse branches [20], or if $X$ is the Julia set of some rational map, $w=r(z)$, then we know that strongly invariant measures $\mu$ exist [11], but not in general.) However, assuming that a strongly invariant measure $\mu$ does exist, then the statement of the scaling equation involves $\hat{\mu}$, and it takes the form of

$$
\begin{equation*}
\frac{d(\hat{\mu} \circ \hat{r})}{d \hat{\mu}}=V \circ \theta_{0} \tag{1.1}
\end{equation*}
$$

where $\theta_{0}$ is the $X$-coordinate function on $X_{\infty}$. (Both sides in equation (1.1) are functions on $X_{\infty}$. The function $V$ on the right-hand side is defined on $X$, so the composition is indeed a function on $X_{\infty}$.)

Hence there are three separate problems: (1) Given such a measure $\mu$, find the functions $V$ such that (1.1) holds. (2) For what functions $V$ are there measures $\mu$ on $X$ such that the corresponding $\hat{\mu}$ on $X_{\infty}$ satisfies (1.1). (3) Starting with a measure $\nu$ on $X_{\infty}$ for which the scaling equation (1.1) holds, then when is there a measure $\mu$ on $X$ such that $\nu=\hat{\mu}$.

Our earlier paper [17] addressed (1); see also Theorem 2.1 below. Here we turn instead to questions (2) and (3).

We show in Section 2 that all three problems relate to a certain eigenvalue problem for a transfer operator $R_{V}$ defined by equation (2.5) in terms of the function $V$. The operator $R_{V}$ is of general interest in operator theory; see e.g., 8], [19], [22], [32], and [37].

Our analysis of the quasi-invariant measures is based on certain Hilbert spaces of martingales (Theorem 3.3) and on a transfer operator (Section [2.2) studied first by David Ruelle. In Theorem 6.3 we give a characterization of the extreme points in the set of $V$-quasi-invariant probability measures.

Example 1.1. We now give an example of an endomorphism $r_{A}$ in a compact space $X(A)$ which is onto, and has finite pre-images. The construction is as follows:

Let $A$ be a $k \times k$ matrix with entries in $\{0,1\}$. Suppose every column in $A$ contains an entry 1. (This guarantees that the restricted shift $r_{A}$ below maps onto $X(A))$.

Set

$$
X(A):=\left\{\left(\xi_{i}\right)_{i \in \mathbb{N}} \in \prod_{\mathbb{N}}\{1, \ldots, k\} \mid A\left(\xi_{i}, \xi_{i+1}\right)=1\right\}
$$

and

$$
r_{A}\left(\xi_{1}, \xi_{2}, \ldots\right)=\left(\xi_{2}, \xi_{3}, \ldots\right) \text { for } \xi \in X(A)
$$

Then $r_{A}$ is a subshift, and the pair $\left(X(A), r_{A}\right)$ satisfies our conditions.
It is known [37] that, for each $A$, as described above, the corresponding system $r_{A}: X(A) \rightarrow X(A)$ has a unique strongly $r_{A}$-invariant probability measure, $\rho_{A}$, i.e., a probability measure on $X(A)$ such that

$$
\int_{X(A)} f d \rho_{A}=\int_{X(A)} \frac{1}{\# r_{A}^{-1}(x)} \sum_{r_{A}(y)=x} f(y) d \rho_{A}(x)
$$

for all bounded measurable functions $f$ on $X(A)$.

In this paper we analyze the connection between measures on $X$ and the induced measures on $X_{\infty}$, and we characterize those measures $X_{\infty}$ which are quasi-invariant with respect to the invertible mapping $\hat{r}$,

where

$$
\begin{gather*}
X_{\infty}:=\left\{\hat{x}=\left(x_{0}, x_{1}, \ldots\right) \in \prod_{\mathbb{N}_{0}} X \mid r\left(x_{n+1}\right)=x_{n}, n \in \mathbb{N}_{0}\right\},  \tag{1.2}\\
\hat{r}(\hat{x})=\left(r\left(x_{0}\right), x_{0}, x_{1}, \ldots\right)
\end{gather*}
$$

and

$$
\hat{r}^{-1}(\hat{x})=\left(x_{1}, x_{2}, \ldots\right)
$$

We recall the simplest instance in the wavelet context: As described in [12] Chapter 5], wavelets may be obtained from an $X_{\infty}$ construction in the special case when $X=\mathbb{T}=\{z \in \mathbb{C}| | z \mid=1\}$, and $r(z)=z^{2}$, as follows:

While the real line $\mathbb{R}$ is not $\mathbb{T}_{\infty}$, we may still build the Hilbert space $L^{2}(\mathbb{R})$ with its multiresolution wavelets from an inductive and isometric procedure based on $\mathbb{T}_{\infty}$. The intuitive idea is to get $\mathbb{R}$ in the limit by successive doubling of periods. Following [19], we are in fact using a recursive system of Hilbert spaces for this, and we note that this is really a martingale construction. To highlight the distinction between the compact group $G=\mathbb{T}_{\infty}$ (called a solenoid) and the reals $\mathbb{R}$, contrast the following two familiar short exact sequences of abelian groups, $0 \rightarrow C \rightarrow G \rightarrow$ $\mathbb{T} \rightarrow 0$, where $C$ is the Cantor group, on the one hand; and $0 \rightarrow \mathbb{Z} \rightarrow \mathbb{R} \rightarrow \mathbb{T} \rightarrow 0$, on the other. While the groups $G$ and $\mathbb{R}$ are very different, the use of $G$ helps us even in this case to build the Hilbert space $L^{2}(\mathbb{R})$; but within the category of isometries in Hilbert space. Moreover, the $X_{\infty}$ viewpoint is useful in important applications outside the group context, and this is the focus of the present paper.

We need a specific interplay between spaces $(X, r), r: X \rightarrow X$ a non-invertible endomorphism, and induced spaces $\left(X_{\infty}, \hat{r}\right)$ where $\hat{r}$ is an automorphism. This issue, and variants of it, arise in a number of areas of mathematics; first in probability theory, going back to [27]; and also more recently in a number of wavelet problems; see for example [5], 7], 8], [15], 16], 17], [19, [21], and [22]. In these applications, the problem is to carry along some isometric operator defined on a Hilbert space of functions on $X$, to the space $X_{\infty}$ (see Proposition 2.2 below). In the language of operator theory, we wish to make a covariant unitary dilation from $X$ to $X_{\infty}$. (See e.g., 24] and [30].) This means that we need to induce measures $\mu$ on $X$ to measures $\hat{\mu}$ on $X_{\infty}$ in such a way that a prescribed covariance is preserved.

To make our paper self-contained, we have recalled Doob's martingale convergence Theorem in Section 3 in the form in which we need it, but we refer the reader to the books [14] and [31] for background on martingale theory.

## 2. Projective limits

2.1. Definitions. Let $X$ be a compact Hausdorff space, and let $r: X \rightarrow X$ be a finite-to-one mapping:

$$
1 \leq \# r^{-1}(x)<\infty
$$

where $r^{-1}(x)=\{y \in X \mid r(y)=x\}$. More generally, set

$$
r^{-1}(E)=\{y \in X \mid r(y) \in E\}, \quad \text { if } E \subset X
$$

By the projective limit $X_{\infty}(r)$, we mean

$$
X \longleftarrow r_{r} X \underset{r}{\longleftarrow} X \longleftarrow X_{\infty}(r) ;
$$

see also (1.2).
We will further assume that $r$ is not invertible, i.e., that $\# r^{-1}(x)$ is not the constant function one.

It is well known that, if $r$ is continuous, then pull-backs of open sets in $X$ define a topology on $X_{\infty}(r)$ making $X_{\infty}(r)$ compact.

The restriction to $X_{\infty}(r)$ of the coordinate projections $\left(x_{0}, x_{1}, \ldots\right) \mapsto x_{n}$ will be denoted by $\theta_{n}$, and we have

$$
\begin{equation*}
r \circ \theta_{n+1}=\theta_{n} \quad\left(n \in \mathbb{N}_{0}\right) \tag{2.1}
\end{equation*}
$$

One advantage of passing from $X$ to $X_{\infty}(r)$ is that $r$ induces an invertible mapping $\hat{r}: X_{\infty}(r) \rightarrow X_{\infty}(r)$, such that

$$
r \circ \theta_{n}=\theta_{n} \circ \hat{r}=\theta_{n-1} .
$$

Moreover, both $\hat{r}$ and $\hat{r}^{-1}$ are continuous if $r$ is.
If $\mathfrak{B}$ is a sigma-algebra of subsets in $X$ (typically we will take $\mathfrak{B}$ to be the Borel subsets in $X$ ), then there are sigma-algebras

$$
\mathfrak{B}_{n}:=\theta_{n}^{-1}(\mathfrak{B})=\left\{\theta_{n}^{-1}(E) \mid E \in \mathfrak{B}\right\}
$$

and

$$
\begin{equation*}
\mathfrak{B}_{\infty}=\bigcup_{n \in \mathbb{N}_{0}} \mathfrak{B}_{n} \tag{2.2}
\end{equation*}
$$

These $\mathfrak{B}_{n}$ 's are sigma-algebras of subsets of $X_{\infty}(r)$.
Using (2.1), we get

$$
\theta_{n}^{-1}=\theta_{n+1}^{-1} \circ r^{-1}
$$

If $r$ is measurable, then $r^{-1}(\mathfrak{B}) \subset \mathfrak{B}$, and we conclude that

$$
\begin{equation*}
\mathfrak{B}_{n} \subset \mathfrak{B}_{n+1} \tag{2.3}
\end{equation*}
$$

Then both of the mappings $\hat{r}$ and $\hat{r}^{-1}$ are measurable on $X_{\infty}(r)$ with respect to the sigma-algebra $\mathfrak{B}_{\infty}$.
2.2. A scaling equation. In [17], we studied the following restrictive setup: we assumed that $X$ carries a probability measure $\mu$ which is strongly r-invariant. By this we mean that

$$
\begin{equation*}
\int_{X} f d \mu=\int_{X} \frac{1}{\# r^{-1}(x)} \sum_{y \in X, r(y)=x} f(y) d \mu(x) \quad\left(f \in L^{\infty}(X)\right) \tag{2.4}
\end{equation*}
$$

If, for example $X=\mathbb{R} / \mathbb{Z}$, and $r(x)=2 x \bmod 1$, then the Haar measure on $\mathbb{R} / \mathbb{Z}=$ Lebesgue measure on $[0,1)$, is the unique strongly $r$-invariant measure on $X$. We prescribe probabilities via a fixed function $V: X \rightarrow[0, \infty)$.

Suppose $V$ is bounded and measurable. Then define $R=R_{V}$, the Ruelle operator, by

$$
\begin{equation*}
R_{V} f(x)=\frac{1}{r^{-1}(x)} \sum_{r(y)=x} V(y) f(y) \quad\left(f \in L^{1}(X, \mu)\right) . \tag{2.5}
\end{equation*}
$$

Theorem 2.1. ([17]) Let $r: X \rightarrow X$ and $X_{\infty}(r)$ be as described in Section 2.1, and suppose that $X$ has a strongly $r$-invariant measure $\mu$. Let $V$ be a non-negative, measurable function on $X$, and let $R_{V}$ be the corresponding Ruelle operator.
(i) There is a unique measure $\hat{\mu}$ on $X_{\infty}(r)$ such that
(a) $\hat{\mu} \circ \theta_{0}^{-1} \ll \mu \quad\left(\right.$ set $\left.h:=\frac{d\left(\hat{\mu} \circ \theta_{0}^{-1}\right)}{d \mu}\right)$,
(b) $\int_{X} f d \hat{\mu} \circ \theta_{n}^{-1}=\int_{X} R_{V}^{n}(f h) d \mu \quad\left(n \in \mathbb{N}_{0}\right)$.
(ii) The measure $\hat{\mu}$ on $X_{\infty}(r)$ satisfies

$$
\begin{equation*}
\frac{d(\hat{\mu} \circ \hat{r})}{d \hat{\mu}}=V \circ \theta_{0} \tag{2.6}
\end{equation*}
$$

and

$$
R_{V} h=h
$$

In this paper we turn around the problem and take (2.6) as the fundamental axiom. We will not assume the existence of a strongly $r$-invariant measure $\mu$.

Proposition 2.2. Let $r: X \rightarrow X$ be as described above, and let $\hat{r}: X_{\infty}(r) \rightarrow X_{\infty}(r)$ be the corresponding automorphism. Let $m: X \rightarrow \mathbb{C}$ be a bounded measurable function on $X$. Then the operator

$$
\begin{equation*}
\hat{S}_{m} f:=\left(m \circ \theta_{0}\right) f \circ \hat{r} \tag{2.7}
\end{equation*}
$$

defines an isometry in the Hilbert space $L^{2}\left(X_{\infty}(r), \hat{\mu}\right)$ if and only if

$$
\begin{equation*}
\hat{\mu} \circ \hat{r} \ll \hat{\mu} \tag{2.8}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{d(\hat{\mu} \circ \hat{r})}{d \hat{\mu}}=\left|m \circ \theta_{0}\right|^{2} \tag{2.9}
\end{equation*}
$$

Proof. The isometric property for (2.7) may be stated in the form.

$$
\begin{equation*}
\int_{X_{\infty}(r)}\left|m \circ \theta_{0}\right|^{2}|f \circ \hat{r}|^{2} d \hat{\mu}=\int_{X_{\infty}(r)}|f|^{2} d \hat{\mu}, \text { for all } f \in L^{2}\left(X_{\infty}(r), \hat{\mu}\right) \tag{2.10}
\end{equation*}
$$

Setting $V:=|m|^{2}$, and $g:=|f \circ \hat{r}|^{2}$, (2.10) reads

$$
\int_{X_{\infty}(r)} V \circ \theta_{0} g d \hat{\mu}=\int_{X_{\infty}(r)} g \circ \hat{r} d \hat{\mu}
$$

But this amounts precisely to the two assertions (2.8) and (2.9). In other words, the $V$-quasi-invariance property for $V=|m|^{2}$ is equivalent to $S_{m}$ defining an $L^{2}$ isometry.

## 3. Induction of measures on $X_{\infty}(r)$ from measures on $X$

3.1. Martingales. In this subsection we introduce the sequence conditional expectations which support our martingale construction of the global Hilbert space.

If $\hat{\mu}$ is a measure defined on the sigma-algebra $\mathfrak{B}_{\infty}$ on $X_{\infty}(r)$, then there is an associated sequence of measures $\left(\mu_{n}\right)$ on $X$ defined as follows:

$$
\begin{equation*}
\mu_{n}:=\hat{\mu} \circ \theta_{n}^{-1} \quad\left(n \in \mathbb{N}_{0}\right) \tag{3.1}
\end{equation*}
$$

or more precisely

$$
\mu_{n}(E)=\hat{\mu}\left(\theta_{n}^{-1}(E)\right) \quad(E \in \mathfrak{B}),
$$

where

$$
\theta_{n}^{-1}(E)=\left\{\hat{x} \in X_{\infty}(r) \mid \theta_{n}(\hat{x}) \in E\right\} .
$$

Our measures will be assumed positive and finite, unless specified otherwise.
We now introduce the Hilbert spaces:

$$
\begin{equation*}
\hat{H}(\hat{\mu}):=L^{2}\left(X_{\infty}(r), \hat{\mu}\right) \tag{3.2}
\end{equation*}
$$

and

$$
H_{n}(\hat{\mu}):=\left\{\xi \circ \theta_{n} \mid \xi \in L^{2}\left(X, \mu_{n}\right)\right\},
$$

the orthogonal projections

$$
P_{n}: \hat{H}(\hat{\mu}) \rightarrow H_{n}(\hat{\mu})
$$

and

$$
E_{n}: \hat{H}(\hat{\mu}) \rightarrow L^{2}\left(X, \mu_{n}\right)
$$

setting

$$
P_{n}(f)=: E_{n}(f) \circ \theta_{n} \quad\left(n \in \mathbb{N}_{0}\right)
$$

then

$$
\int_{X_{\infty}(r)} \xi \circ \theta_{n} f d \hat{\mu}=\int_{X} \xi E_{n}(f) d \mu_{n}, \text { for all } n \in \mathbb{N}_{0}, \xi \in L^{2}\left(X, \mu_{n}\right)
$$

Lemma 3.1. Let $\hat{\mu}$ be a measure on $\left(X_{\infty}(r), \mathfrak{B}_{\infty}\right)$ and let $\hat{H}(\hat{\mu})$ be the corresponding Hilbert space from (3.2).
(i) Then there is an isometric isomorphism

$$
\begin{equation*}
J: \hat{H}(\hat{\mu}) \rightarrow \lim _{n} L^{2}\left(X, \mu_{n}\right) \tag{3.3}
\end{equation*}
$$

where the inductive limit in (3.3) is defined by the isometric embeddings

$$
\xi \mapsto \xi \circ r: L^{2}\left(X, \mu_{n}\right) \rightarrow L^{2}\left(X, \mu_{n+1}\right)
$$

The isomorphism $J$ in (3.3) is

$$
J f=\left(E_{n} f\right)_{n \in \mathbb{N}_{0}} \quad(f \in \hat{H}(\hat{\mu}))
$$

(ii) The Hilbert norm is given by

$$
\|f\|_{\hat{H}(\hat{\mu})}^{2}=\lim _{n \rightarrow \infty}\left\|E_{n}(f)\right\|_{L^{2}\left(X, \mu_{n}\right)}^{2}
$$

and
(iii) For every sequence $\xi_{n} \in L^{2}\left(X, \mu_{n}\right)$ such that $P_{n}\left(\xi_{n+1} \circ \theta_{n+1}\right)=\xi_{n} \circ \theta_{n}$, and

$$
\sup _{n}\left\|\xi_{n}\right\|_{L^{2}\left(X, \mu_{n}\right)}<\infty
$$

there is a unique $f \in L^{2}\left(X_{\infty}(r), \hat{\mu}\right)$ such that

$$
P_{n} f=\xi_{n} \circ \theta_{n} \text { for all } n \in \mathbb{N}_{0} .
$$

Moreover,

$$
f=\lim _{n \rightarrow \infty} \xi_{n} \circ \theta_{n}, \text { pointwise } \hat{\mu} \text {-a.e., and in } \hat{H}(\hat{\mu}) .
$$

Proof. The proof depends on Doob's martingale convergence theorem; see [14, [17], and 31 .
3.2. From $\mu$ on $X$ to $\hat{\mu}$ on $X_{\infty}(r)$. In this subsection we introduce the sequence martingale measures $\left(\mu_{n}\right)$ needed in our construction of the global Hilbert space.

Lemma 3.2. Let $\left(\mu_{n}\right)_{n \in \mathbb{N}_{0}}$ be a sequence of measures on $(X, \mathfrak{B})$. Then there is a measure $\hat{\mu}$ on $\left(X_{\infty}(r), \mathfrak{B}_{\infty}\right)$ such that

$$
\begin{equation*}
\hat{\mu} \circ \theta_{n}^{-1}=\mu_{n} \quad\left(n \in \mathbb{N}_{0}\right) \tag{3.4}
\end{equation*}
$$

if and only if

$$
\begin{equation*}
\mu_{n+1} \circ r^{-1}=\mu_{n} \quad\left(n \in \mathbb{N}_{0}\right) \tag{3.5}
\end{equation*}
$$

Proof. We first recall the definition, the measure $\mu_{n+1} \circ r^{-1}$ in (3.5). On a Borel subset $E$, it is

$$
\begin{equation*}
\mu_{n+1} \circ r^{-1}(E):=\mu_{n+1}\left(r^{-1}(E)\right) \tag{3.6}
\end{equation*}
$$

It is well known (see [35]) that (3.6) defines a measure.
Suppose $\hat{\mu}$ exists such that (3.4) holds. Let $\xi$ be a bounded $\mathfrak{B}$-measurable function on $X$. Then
$\int_{X} \xi \circ r d \mu_{n+1}=\int_{X_{\infty}(r)} \xi \circ r \circ \theta_{n+1} d \hat{\mu}=\int_{X_{\infty}(r)} \xi \circ \theta_{n} d \hat{\mu}=\int_{X} \xi d \hat{\mu} \circ \theta_{n}^{-1}=\int_{X} \xi d \mu_{n}$, which proves the identity (3.5).

Conversely, suppose $\left(\mu_{n}\right)_{n \in \mathbb{N}_{0}}$ satisfies (3.5). Note that the $\mathfrak{B}_{n}$-measurable functions are of the form $f \circ \theta_{n}$ with $f \mathfrak{B}$-measurable. The $\mathfrak{B}_{n}$-measurable functions are also $\mathfrak{B}_{n+1}$-measurable and this inclusion is given by

$$
\begin{equation*}
f \circ \theta_{n}=(f \circ r) \circ \theta_{n+1} . \tag{3.7}
\end{equation*}
$$

The relations (2.2) and (2.3) imply that the union of the algebras of bounded $\mathfrak{B}_{n}$-measurable functions is $L^{1}(\hat{\mu})$-dense in the algebra of bounded $\mathfrak{B}_{\infty}$-measurable functions on $X_{\infty}(r)$.

Now set

$$
\hat{\mu}\left(f \circ \theta_{n}\right):=\mu_{n}(f), f \text { assumed } \mathfrak{B}_{n} \text {-measurable and bounded. }
$$

To see that this is consistent, use (3.7) and (3.5) to check that

$$
\hat{\mu}\left((f \circ r) \circ \theta_{n+1}\right)=\mu_{n+1}(f \circ r)=\mu_{n}(f)=\hat{\mu}\left(f \circ \theta_{n}\right)
$$

The existence and uniqueness of $\hat{\mu}$ follows from Kolmogorov's theorem; see [17] and [27].
3.3. Quasi-invariance. In this subsection we state our theorem regarding the scaling equation (1.1) from the Introduction.

In this section we will study the problem (2.6) via our correspondence $\hat{\mu} \leftrightarrow$ $\left(\mu_{n}\right)_{n \in \mathbb{N}_{0}}$ from Lemma 3.2

The questions are two: (1) Which measures $\mu_{0}$ on $X$ admit "extensions" $\hat{\mu}$ to measures on $X_{\infty}(r)$ satisfying (2.6)?
(2) Understand the measures $\mu_{0}$ which admit solutions $\hat{\mu}$ to (2.6).

Our first theorem answers the question (1).
Theorem 3.3. Let $\left(\mu_{n}\right)_{n \in \mathbb{N}_{0}}$ be a sequence of measures as in Lemma 3.1 i.e., satisfying $\mu_{n+1} \circ r^{-1}=\mu_{n}$, and let $V: X \rightarrow[0, \infty)$ be $\mathfrak{B}$-measurable. Then the "extended" measure $\hat{\mu}$ satisfies

$$
\begin{equation*}
\frac{d(\hat{\mu} \circ \hat{r})}{d \hat{\mu}}=V \circ \theta_{0} \text { a.e. on }\left(X_{\infty}(r), \mathfrak{B}_{\infty}\right) \tag{3.8}
\end{equation*}
$$

if and only if

$$
\begin{equation*}
d \mu_{0}=\left(V d \mu_{0}\right) \circ r^{-1}, \text { and } d \mu_{n+1}=\left(V \circ r^{n}\right) d \mu_{n} \quad\left(n \in \mathbb{N}_{0}\right) \tag{3.9}
\end{equation*}
$$

Proof. Note first that (3.9) is equivalent to the following

$$
\int_{X} f d \mu_{0}=\int_{X} f \circ r V d \mu_{0}, \text { and } \int_{X} f d \mu_{n+1}=\int_{X} f V \circ r^{n} d \mu_{n} \quad\left(n \in \mathbb{N}_{0}\right)
$$

for all bounded $\mathfrak{B}$-measurable functions $f$ on $X$.
Let $\left(\mu_{n}\right)_{n \in \mathbb{N}_{0}}$ be a sequence of measures that satisfy

$$
\begin{equation*}
\mu_{n+1} \circ r^{-1}=\mu_{n} \tag{3.10}
\end{equation*}
$$

and let $\hat{\mu}$ be the corresponding measure on $X_{\infty}(r)$. Suppose first that $\hat{\mu}$ satisfies (3.8) for some $\mathfrak{B}$-measurable function $V: X \rightarrow[0, \infty)$; in other words, we have the identity

$$
\begin{equation*}
\int_{X_{\infty}(r)} f \circ \hat{r}^{-1} d \hat{\mu}=\int_{X_{\infty}(r)} f V \circ \theta_{0} d \hat{\mu} \tag{3.11}
\end{equation*}
$$

for all bounded $\mathfrak{B}_{\infty}$-measurable functions $f$ on $X_{\infty}(r)$.
Now specializing (3.11) to $f=\xi \circ \theta_{n}$ with $\xi$ some bounded $\mathfrak{B}$-measurable function on $X$, we get

$$
\begin{aligned}
\int_{X_{\infty}(r)} \xi \circ \theta_{n+1} d \hat{\mu} & =\int_{X_{\infty}(r)} \xi \circ \theta_{n} V \circ \theta_{0} d \hat{\mu} \\
& =\int_{X_{\infty}(r)} \xi \circ \theta_{n} V \circ r^{n} \circ \theta_{n} d \hat{\mu}=\int_{X} \xi V \circ r^{n} d \mu_{n}
\end{aligned}
$$

and therefore

$$
d \mu_{n+1}=\left(V \circ r^{n}\right) d \mu_{n}
$$

by (3.1).

We take a closer look at the case $n=0$. Substitute $f=\xi \circ \theta_{0}$ into (3.11). We get
$\int_{X_{\infty}(r)} \xi \circ \theta_{0} d \hat{\mu}=\int_{X_{\infty}(r)} \xi \circ \theta_{0} \circ \hat{r} V \circ \theta_{0} d \hat{\mu}=\int_{X_{\infty}(r)} \xi \circ r \circ \theta_{0} V \circ \theta_{0} d \hat{\mu}=\int_{X} \xi \circ r V d \mu_{0}$ and therefore $d \mu_{0}=\left(V d \mu_{0}\right) \circ r^{-1}$ which is (3.9).

Conversely, suppose $\left(\mu_{n}\right)_{n \in \mathbb{N}_{0}}$ satisfies (3.10) and (3.9). Lemma 3.2 yields the existence of a $\hat{\mu}$ on $\left(X_{\infty}(r), \mathfrak{B}_{\infty}\right)$ such that

$$
\mu_{n}=\hat{\mu} \circ \theta_{n}^{-1}
$$

We claim that $\hat{\mu}$ is $V$-quasi-invariant, i.e., that (3.8) holds. This amounts to the identity

$$
\begin{equation*}
\int_{X_{\infty}(r)} f \circ \hat{r}^{-1} d \hat{\mu}=\int_{X_{\infty}(r)} f V \circ \theta_{0} d \hat{\mu} \tag{3.12}
\end{equation*}
$$

By Lemma 3.1 this is equivalent to

$$
d \mu_{n+1}=V \circ r^{n} d \mu_{n}
$$

To see this, substitute $f=\xi \circ \theta_{n}$ into (3.12) for $\xi$ a bounded $\mathfrak{B}$-measurable function on $X$, and $n \in \mathbb{N}_{0}$. But this holds by (3.9), and (3.8) follows.

## 4. A FIXED-POINT PROBLEM

In Theorem [3.3], we saw that our condition (2.6), V-quasi-invariance, on measures $\hat{\mu}$ on $\left(X_{\infty}(r), \mathfrak{B}_{\infty}\right)$ entails a fixed-point property (4.1) for the corresponding measure $\mu_{0}:=\hat{\mu} \circ \theta_{0}^{-1}$. In this section, we turn the problem around. We show that this fixed-point property characterizes the measures $\hat{\mu}$ satisfying the quasiinvariance, i.e., satisfying (2.6).

If the function $V: X \rightarrow[0, \infty)$ is given, we define

$$
V^{(n)}(x):=V(x) V(r(x)) \cdots V\left(r^{n-1}(x)\right)
$$

and set $d \mu_{n}:=V^{(n)} d \mu_{0}$. Our result states that the corresponding measure $\hat{\mu}$ on $X_{\infty}(r)$ is $V$-quasi-invariant if and only if

$$
\begin{equation*}
d \mu_{0}=\left(V d \mu_{0}\right) \circ r^{-1} \tag{4.1}
\end{equation*}
$$

Theorem 4.1. Let $V: X \rightarrow[0, \infty)$ be $\mathfrak{B}$-measurable, and let $\mu_{0}$ be a measure on $X$ satisfying the following fixed-point property:

$$
\begin{equation*}
d \mu_{0}=\left(V d \mu_{0}\right) \circ r^{-1} \tag{4.2}
\end{equation*}
$$

Then there exists a unique measure $\hat{\mu}$ on $X_{\infty}(r)$ such that

$$
\begin{equation*}
\frac{d(\hat{\mu} \circ \hat{r})}{d \hat{\mu}}=V \circ \theta_{0} \tag{4.3}
\end{equation*}
$$

and

$$
\hat{\mu} \circ \theta_{0}^{-1}=\mu_{0}
$$

Proof. Let $\mu_{0}$ be a measure satisfying the fixed-point property (4.2), and set

$$
\begin{equation*}
d \mu_{n+1}:=\left(V \circ r^{n}\right) d \mu_{n} \tag{4.4}
\end{equation*}
$$

This defines a sequence $\mu_{0}, \mu_{1}, \mu_{2}, \ldots$ of measures on $X$. To establish the existence of the measure $\hat{\mu}$ on $X_{\infty}(r)$, we appeal to Lemma 3.2. We claim that (3.5) holds, or equivalently that

$$
\begin{equation*}
\int_{X} f \circ r d \mu_{n+1}=\int_{X} f d \mu_{n} \tag{4.5}
\end{equation*}
$$

for all bounded $\mathfrak{B}$-measurable functions $f$ on $X$, and for all $n \in \mathbb{N}_{0}$.
Substitution of (4.4) into (4.5) yields

$$
\int_{X} f \circ r d \mu_{n+1}=\int_{X} f \circ r V \circ r^{n} d \mu_{n}
$$

and we will prove (4.5) by induction:
First

$$
\int_{X} f \circ r d \mu_{1}=\int_{X} f \circ r V d \mu_{0}=\int_{X} f\left(V d \mu_{0}\right) \circ r^{-1}=\int_{X} f d \mu_{0}
$$

where we used (4.2) in the last step. So (4.5) holds for $n=0$.
Assume that (4.5) holds for $\mu_{m}$ and $m<n$. Then

$$
\int_{X} f \circ r d \mu_{n+1}=\int_{X} f \circ r V \circ r^{n-1} \circ r d \mu_{n}=\int_{X} f V \circ r^{n-1} d \mu_{n-1}=\int_{X} f d \mu_{n}
$$

where we used the definition (4.4) in the last step. The induction is completed.
This proves (3.5). An application of Lemma 3.2 yields the existence of a unique measure $\hat{\mu}$ on $\left(X_{\infty}(r), \mathfrak{B}_{\infty}\right)$ such that

$$
\hat{\mu} \circ \theta_{n}^{-1}=\mu_{n} \quad\left(n \in \mathbb{N}_{0}\right)
$$

Since the sequence $\left(\mu_{n}\right)_{n \in \mathbb{N}_{0}}$ satisfies the pair of conditions (3.9) by construction, we conclude from Theorem 3.3 that $\hat{\mu}$ must satisfy (4.3), i.e., $\hat{\mu}$ is $V$-quasi-invariant as claimed.

The proof of Theorem 4.1 also yields the following corollary:
Corollary 4.2. Let $V: X \rightarrow[0, \infty)$ be a bounded $\mathfrak{B}$-measurable function, and let $\mu_{0}$ be a measure on $(X, \mathfrak{B})$. Suppose

$$
d \mu_{0}=\left(V d \mu_{0}\right) \circ r^{-1}
$$

Set

$$
\begin{equation*}
d \mu_{n}:=V^{(n)} d \mu_{0}, \text { for } n \in \mathbb{N} \tag{4.6}
\end{equation*}
$$

Then

$$
\begin{equation*}
\mu_{n+1} \circ r^{-1}=\mu_{n} \tag{4.7}
\end{equation*}
$$

and

$$
\begin{equation*}
\mu_{n}(X)=\int_{X} V^{(n)} d \mu_{0}=\mu_{0}(X), \text { for all } n \in \mathbb{N} \tag{4.8}
\end{equation*}
$$

Proof. Indeed, by (4.6) we have the recursive formula $d \mu_{n+1}=\left(V \circ r^{n}\right) d \mu_{n}$, and as in the proof of Theorem 4.1) we get that (4.7) is satisfied. Applying (4.6) to the constant function 1, we get that

$$
\mu_{n}(X)=\int_{X} V^{(n)} d \mu_{0}
$$

and, with 4.7)

$$
\mu_{n+1}(X)=\int_{X} 1 d \mu_{n+1}=\int_{X} 1 \circ r d \mu_{n+1}=\int_{X} 1 d \mu_{n}=\mu_{n}(X)
$$

Now (4.8) follows by induction.

## 5. Transformations of measures

Let $X$ be a compact Hausdorff space, $\mathfrak{B}$ the sigma-algebra of all Borel subsets of $X$, and $C(X)$ the continuous functions on $X$.

For $a \in \mathbb{R}_{+}$denote by

$$
M_{a}(X):=\{\mu \mid \mu \text { is a measure on }(X, \mathfrak{B}), \mu(X)=a\}
$$

Note that $M_{a}(X)$ is equipped with the weak*-topology coming from the duality

$$
(f, \mu) \mapsto \int_{X} f d \mu=: \mu(f)
$$

where $f \in C(X)$. The neighborhoods are generated by the sets $N_{\mu_{0}}\left(f_{1}, \ldots, f_{k}, \epsilon\right)$ where $\epsilon \in \mathbb{R}_{+}, f_{1}, \ldots, f_{k} \in C(X)$ and

$$
N_{\mu_{0}}\left(f_{1}, \ldots, f_{k}, \epsilon\right)=\left\{\mu \in M(X)| | \mu\left(f_{i}\right)-\mu_{0}\left(f_{i}\right) \mid<\epsilon, i=1, \ldots, k\right\}
$$

It is known 35 that each $M_{a}(X)$ is weak*-compact, i.e., $M_{a}(X)$ is a compact convex set in the topology determined by the neighborhoods $N_{\mu_{0}}\left(f_{1}, \ldots, f_{k}, \epsilon\right)$.

Let $V: X \rightarrow[0, \infty)$ be bounded and $\mathfrak{B}$-measurable, and define

$$
T_{V}(\mu):=(V d \mu) \circ r^{-1} \quad(\mu \in M(X))
$$

Lemma 5.1. Let $V$ be as above, and assume $V$ is also continuous. Then $T_{V}: M(X)$ $\rightarrow M(X)$ is continuous.
Proof. Let $\mu \in M(X), k \in \mathbb{Z}_{+}$and $f_{i} \in C(X), i=1, \ldots, k$. Then

$$
\left\langle T_{V}(\mu) \mid f_{i}\right\rangle=\left\langle\mu \mid V f_{i} \circ r\right\rangle
$$

Setting $g_{k}:=V f_{i} \circ r$, we see that

$$
T_{V}\left(N_{\mu_{0}}\left(g_{1}, \ldots, g_{k}, \epsilon\right)\right) \subset N_{T_{V}\left(\mu_{0}\right)}\left(f_{1}, \ldots, f_{k}, \epsilon\right)
$$

and the conclusion follows if $V$ is assumed continuous.
Definition 5.2. Let $V: X \rightarrow[0, \infty)$ be bounded and $\mathfrak{B}$-measurable. We use the notation

$$
M^{V}(X):=\left\{\mu \in M(X) \mid d \mu=(V d \mu) \circ r^{-1}\right\}
$$

For measures $\hat{\mu}$ on $\left(X_{\infty}(r), \mathfrak{B}_{\infty}\right)$ we introduce

$$
M_{q i}^{V}\left(X_{\infty}(r)\right):=\left\{\hat{\mu} \in M\left(X_{\infty}(r)\right) \mid \hat{\mu} \circ \hat{r} \ll \hat{\mu} \text { and } \frac{d(\hat{\mu} \circ \hat{r})}{d \hat{\mu}}=V \circ \theta_{0}\right\}
$$

The results of the previous section may be summarized as follows:
Theorem 5.3. Let $V$ be as in Definition 5.2. For measures $\hat{\mu}$ on $X_{\infty}(r)$ and $n \in \mathbb{N}_{0}$, define

$$
C_{n}(\hat{\mu}):=\hat{\mu} \circ \theta_{n}^{-1}
$$

Then $C_{0}$ is a bijective affine isomorphism of $M_{q i}^{V}\left(X_{\infty}(r)\right)$ onto $M^{V}(X)$ that preserves the total measure, i.e., $C_{0}(\hat{\mu})(X)=\hat{\mu}\left(X_{\infty}(r)\right)$ for all $\hat{\mu} \in M_{q i}^{V}\left(X_{\infty}(r)\right)$.

Proof. We showed in Theorem 3.3 that $C_{0}$ maps $M_{q i}^{V}\left(X_{\infty}(r)\right)$ onto $M^{V}(X)$. The inverse mapping

$$
C_{0}^{-1}: M^{V}(X) \rightarrow M_{q i}^{V}\left(X_{\infty}(r)\right)
$$

may be realized using Theorem 4.1 and Lemma 4.2.
Remark 5.4. The intuitive idea behind $C_{0}^{-1}: M^{V}(X) \rightarrow M_{q i}^{V}\left(X_{\infty}(r)\right)$ is as follows:
Let $C_{0}^{-1}\left(\mu_{0}\right)=\hat{\mu}$. Then $\hat{\mu}$ is a measure on $X_{\infty}(r)$, and we view points $\hat{x}$ in $X_{\infty}(r)$ as infinite paths. Recall, if $\hat{x}=\left(x_{0}, x_{1}, \ldots\right) \in X_{\infty}(r)$, then $r\left(x_{n+1}\right)=x_{n}$. So in a random walk we choose $x_{n+1} \in r^{-1}\left(x_{n}\right)$, and the function $V$ assigns the probabilities in each step. At step $n$, there are $\# r^{-1}\left(x_{n}\right)$ choices.

The assertion in the theorem is that the measure $\hat{\mu}$ is completely determined by the prescribed measure $\mu_{0}$ at the starting point $x_{0}$ of the path, and by the function $V$.

The two measures $\hat{\mu}$ and $\mu_{0}$ are normalized so that

$$
\mu_{0}(X)=\hat{\mu}\left(X_{\infty}(r)\right)
$$

In a special case, an explicit formula for $C_{0}^{-1}\left(\mu_{0}\right)$ is known [17, Proposition 8.2].
Theorem 5.5. Let $V: X \rightarrow[0, \infty)$ be continuous. Assume that there exist some measure $\nu$ on $(X, \mathfrak{B})$ and two numbers $0<a<b$ such that

$$
\begin{equation*}
a \leq \nu(X) \leq b, \text { and } a \leq \int_{X} V^{(n)} d \nu \leq b \text { for all } n \in \mathbb{N} \tag{5.1}
\end{equation*}
$$

Then there exists a measure $\mu_{0}$ on $(X, \mathfrak{B})$ that satisfies

$$
d \mu_{0}=\left(V d \mu_{0}\right) \circ r^{-1}
$$

and there exists a $V$-quasi-invariant measure $\hat{\mu}$ on $\left(X_{\infty}(r), \mathfrak{B}_{\infty}\right)$.
Proof. Condition (5.1) guarantees that the set

$$
M_{a b}^{V}(X):=\left\{\mu \in M(X) \mid a \leq \mu(X) \leq b \text { and } a \leq \int_{X} V^{(n)} d \mu \leq b, \text { for all } n \in \mathbb{N}\right\}
$$

is non-empty. Moreover, by the Banach-Alaoglu theorem [35], this set is compact in the weak*-topology. It is clear also that the set is convex.

We claim that the operator $T_{V}$ maps $M_{a b}^{V}(X)$ into itself. Indeed, if $\mu \in M_{a b}^{V}(X)$, then

$$
T_{V}(\mu)(X)=\int_{X} V 1 \circ r d \mu \in[a, b]
$$

and, using $V V^{(n)} \circ r=V^{(n+1)}$,

$$
T_{V}(\mu)\left(V^{(n)}\right)=\int_{X} V V^{(n)} \circ r d \mu=\int_{X} V^{(n+1)} d \mu \in[a, b]
$$

Lemma 5.1 shows that $T_{V}$ is continuous so we can apply the Markov-Kakutani fixed-point theorem (see [36]) to conclude that there exists a $\mu_{0} \in M_{a b}^{V}$ such that $T_{V}\left(\mu_{0}\right)=\mu_{0}$. The $V$-quasi-invariant measure $\hat{\mu}$ can be obtained from $\mu_{0}$ using Theorem 4.1.

Definition 5.6. We now connect the above discussion with Section 2.2, Let ( $X, \mathfrak{B}$ ) and $r: X \rightarrow X$ be as described above. Suppose in addition that $(X, \mathfrak{B})$ carries a strongly $r$-invariant probability measure, $\rho$; see (2.4) for the definition.

Set

$$
M(\rho, X):=\{\mu \in M(X) \mid \mu \ll \rho\} .
$$

Lemma 5.7. Let $V: X \rightarrow[0, \infty)$ be a bounded measurable function. Then $T_{V}$ leaves $M(\rho, X)$ invariant; and $\mu_{0}$ in $M(\rho, X)$ satisfies

$$
T_{V}\left(\mu_{0}\right)=\mu_{0}
$$

if and only if there is an $h \in L^{1}(X, \rho)$ such that $d \mu_{0}=h d \rho$ and $R_{V} h=h$.
Proof. Let $\mu \in M(\rho, X)$, and write $d \mu=f d \rho, f \in L^{1}(X, \rho)$. Let $\xi$ be a bounded measurable function on $X$. Then

$$
\begin{gathered}
T_{V}(\mu)(\xi)=\int_{X} \xi \circ r V f d \rho=\int_{X} \xi(x) \frac{1}{\# r^{-1}(x)} \sum_{r(y)=x} V(y) f(y) d \rho(x) \\
=\int_{X} \xi(x)\left(R_{V} f\right)(x) d \rho(x)
\end{gathered}
$$

Stated differently,

$$
T_{V}(f d \rho)=R_{V}(f) d \rho
$$

Hence $T_{V}(f d \rho)=f d \rho$ if and only if $R_{V} f=f$ as claimed.
Before stating and proving our next result we need one more lemma.
Lemma 5.8. Let $(X, \mathfrak{B}), r: X \rightarrow X$, and $V: X \rightarrow[0, \infty)$ be as above. Suppose $(X, \mathfrak{B})$ carries a strongly $r$-invariant probability measure $\rho$. Then

$$
\begin{equation*}
\int_{X} V^{(n)} f d \rho=\int_{X} R_{V}^{n}(f) d \rho \tag{5.2}
\end{equation*}
$$

for all bounded measurable functions $f$ and $n \in \mathbb{N}$.
Proof. We prove (5.2) by induction, starting with $n=1$.
For $f \in L^{\infty}(X)$,

$$
\int_{X} V f d \rho=\int_{X} \frac{1}{\# r^{-1}(x)} \sum_{r(y)=x} V(y) f(y) d \rho(x)=\int_{X} R_{V} f d \rho
$$

where we used the definition (2.4) of strong $r$-invariance.
Suppose (5.2) holds up to $n$. Then

$$
\begin{aligned}
\int_{X} V^{(n+1)} f d \rho & =\int_{X}\left(V^{(n)} \circ r\right) V f d \rho=\int_{X} V^{(n)}(x) \frac{1}{\# r^{-1}(x)} \sum_{r(y)=x} V(y) f(y) d \rho(x) \\
& =\int_{X} V^{(n)} R_{V} f d \rho=\int_{X} R_{V}^{n} R_{V} f d \rho=\int_{X} R_{V}^{n+1} f d \rho
\end{aligned}
$$

where we used the induction hypothesis in the last step.
This completes the proof of the lemma.
Our next result is an infinite-dimensional analogue of the (finite-dimensional) Perron-Frobenius theorem for non-negative matrices.

Theorem 5.9. Let $(X, \mathfrak{B})$, and $r: X \rightarrow X$, be as described above. Suppose $V: X \rightarrow[0, \infty)$ is measurable,

$$
\frac{1}{\# r^{-1}(x)} \sum_{r(y)=x} V(y) \leq 1
$$

and that some probability measure $\nu_{V}$ on $X$ satisfies

$$
\begin{equation*}
\nu_{V} \circ R_{V}=\nu_{V} \tag{5.3}
\end{equation*}
$$

Assume also that $(X, \mathfrak{B})$ carries a strongly r-invariant probability measure $\rho$, such that

$$
\rho(\{x \in X \mid V(x)>0\})>0
$$

Then
(i) $T_{V}^{n}(d \rho)=R_{V}^{n}(\mathbf{1}) d \rho$, for $n \in \mathbb{N}$, where $\mathbf{1}$ denotes the constant function one.
(ii) $[$ Monotonicity $] \quad \cdots \leq R_{V}^{n+1}(\mathbf{1}) \leq R_{V}^{n}(\mathbf{1}) \leq \cdots \leq \mathbf{1}$, pointwise on $X$.
(iii) The limit $\lim _{n \rightarrow} R_{V}^{n}(\mathbf{1})=h_{V}$ exists, $R_{V} h_{V}=h_{V}$, and

$$
\begin{equation*}
\rho\left(\left\{x \in X \mid h_{V}(x)>0\right\}\right)>0 \tag{5.4}
\end{equation*}
$$

(iv) The measure $d \mu_{0}^{(V)}=h_{V} d \rho$ is a solution to the fixed-point problem

$$
T_{V}\left(\mu_{0}^{(V)}\right)=\mu_{0}^{(V)}
$$

(v) The sequence $d \mu_{n}^{(V)}=V^{(n)} h_{V} d \rho$ defines a unique $\hat{\mu}^{(V)}$ (as in Theorem 4.1 and Lemma 4.2); and
(vi) $\mu_{n}^{(V)}(f)=\int_{X} R_{V}^{n}\left(f h_{V}\right) d \rho$ for all bounded measurable functions $f$ on $X$, and all $n \in \mathbb{N}$.
Finally,
(vii) the measure $\hat{\mu}^{(V)}$ on $X_{\infty}(r)$ satisfying $\hat{\mu}^{(V)} \circ \theta_{n}^{-1}=\mu_{n}^{(V)}$ has total mass

$$
\hat{\mu}^{(V)}\left(X_{\infty}(r)\right)=\rho\left(h_{V}\right)=\int_{X} h_{V}(x) d \rho(x)
$$

Proof. Part (i) follows from Lemma 5.7
(ii) It is clear that

$$
R_{V}(\mathbf{1})=\frac{1}{\# r^{-1}(x)} \sum_{r(y)=x} V(y) \leq 1, \text { for } x \in X
$$

Starting the induction, suppose

$$
R_{V}^{n}(\mathbf{1}) \leq R_{V}^{n-1}(\mathbf{1})
$$

Then

$$
\begin{aligned}
R_{V}^{n+1}(\mathbf{1}) & =\frac{1}{\# r^{-1}(x)} \sum_{r(y)=x} V(y) R_{V}^{n}(\mathbf{1})(y) \\
& \leq \frac{1}{\# r^{-1}(x)} \sum_{r(y)=x} V(y) R_{V}^{n-1}(\mathbf{1})(y)=R_{V}^{n}(\mathbf{1})(x)
\end{aligned}
$$

An induction now proves (ii).
It follows that the limit $h_{V}$ in (iii) exists, and that

$$
R_{V}\left(h_{V}\right)=h_{V}
$$

Using (5.3), we get that

$$
\nu_{V}\left(R_{V}^{n}(\mathbf{1})\right)=\nu_{V}(\mathbf{1})=1, \text { for all } n \in \mathbb{N}
$$

and therefore

$$
\nu_{V}\left(h_{V}\right)=\int_{X} h_{V} d \nu_{V}=1
$$

The conclusion (5.4) follows from [22, Chapter 3].

The assertion in (iv) is that $\mu_{0}^{(V)}:=h_{V} d \rho$ is a fixed point, i.e., that

$$
T_{V}\left(h_{V} d \rho\right)=h_{V} d \rho
$$

This follows in turn from Lemma 5.7 Hence, the measures

$$
d \mu_{n}^{(V)}=V^{(n)} h_{V} d \rho
$$

extend to a unique measure $\hat{\mu}^{(V)}$ on $X_{\infty}(r)$ and, by Theorem 5.3,

$$
\hat{\mu}^{(V)} \in M_{q i}^{V}\left(X_{\infty}(r)\right)
$$

Moreover, using now Lemma 5.8, we get

$$
\mu_{n}^{(V)}(f)=\int_{X} f V^{(n)} h_{V} d \rho=\int_{X} R_{V}^{n}\left(f h_{V}\right) d \rho
$$

which yields the desired conclusion (vi) and (vii).
Indeed

$$
\hat{\mu}^{(V)}\left(X_{\infty}(r)\right)=\int_{X_{\infty}(r)} \mathbf{1} \circ \theta_{0} d \hat{\mu}^{(V)}=\int_{X} \mathbf{1} d \mu_{0}=\int_{X} h_{V} d \rho
$$

Comparing with (5.1) in Theorem [5.5, notice that

$$
\int_{X} h_{V} d \rho=\lim _{n \rightarrow \infty} \int_{X} R_{V}^{n}(\mathbf{1}) d \rho=\lim _{n \rightarrow \infty} \rho\left(V^{(n)}\right)
$$

The conclusions of Theorem 5.9 hold in the following more general form.
Corollary 5.10. Let $V, W: X \rightarrow[0, \infty)$ be measurable, and suppose the assumptions (i) and (ii) are satisfied:
(i) $\frac{1}{\# r^{-1}(x)} \sum_{r(y)=x} V(y) W(y) \leq 1, x \in X$.
(ii) There exists a non-negative function $h_{W}$ on $X$, and a measure $\rho_{W}$ such that $\rho_{W} R_{W}=\rho_{W}, R_{W} h_{W}=h_{W}$ and $\rho_{W}\left(h_{W}\right)=1$.
Then

$$
\begin{equation*}
T_{V}^{n}\left(f d \rho_{W}\right)=\left(R_{V W}^{n} f\right) d \rho_{W} \tag{5.5}
\end{equation*}
$$

and the limit

$$
\lim _{n \rightarrow \infty} R_{V W}^{n}(\mathbf{1})=h
$$

exists, and satisfies

$$
\begin{equation*}
T_{V}\left(h d \rho_{W}\right)=h d \rho_{W} \tag{5.6}
\end{equation*}
$$

Proof. We first show (5.5) for $n=1$. Let $f$ and $\xi$ be bounded, measurable functions on $X$. Then

$$
\begin{gathered}
T_{V}\left(f d \rho_{W}\right)(\xi)=\int_{X} \xi \circ r V f d \rho_{W}=\int_{X} R_{W}(\xi \circ r V f) d \rho_{W} \\
=\int_{X} \xi R_{W}(V f) d \rho_{W}=\int_{X} \xi R_{V W}(f) d \rho_{W}
\end{gathered}
$$

Hence (5.5) holds for $n=1$, and the general case follows by iteration.
The argument from the proof of the theorem shows that

$$
\cdots \leq R_{V W}^{n+1}(\mathbf{1}) \leq R_{V W}^{n}(\mathbf{1}) \leq \cdots \leq R_{V W}(\mathbf{1}) \leq \mathbf{1}
$$

so the limit

$$
h:=\lim _{n \rightarrow \infty} R_{V W}^{n}(\mathbf{1})
$$

exists, and satisfies (5.6) by the same argument.
Remark 5.11. Stated in Ruelle's thermodynamical formalism [37], the data $\rho_{W}$ (measure) and $h_{W}$ (eigenfunction) in part (ii) of Corollary 5.10 represent an equilibrium distribution where $W$ is related to a potential function. Under mild conditions on $(X, r)$ and $W$, it is known that solutions $\left(\rho_{W}, h_{W}\right)$ exist, and we say that $\lambda_{W}=1$ is the Perron-Frobenius eigenvalue of the Ruelle operator $R_{W}$. In that case

$$
1=\sup \left\{|\lambda| \mid \lambda \in \operatorname{spectrum}\left(R_{W}\right)\right\}
$$

The reader is referred to [37], [8], [32] for further details regarding the spectral theory of $R_{W}$.

Notice further that the conclusion of Ruelle's Perron-Frobenius theorem is a generalization of the classical Perron-Frobenius theorem for matrices with nonnegative entries.

## 6. Extreme points

Theorem5.3 shows that the map $\hat{\mu} \mapsto \mu_{0}:=\hat{\mu} \circ \theta_{0}^{-1}$ establishes a bijective affine correspondence between the following two sets:

$$
M_{q i, 1}^{V}\left(X_{\infty}(r)\right):=\left\{\hat{\mu} \in M\left(X_{\infty}(r)\right) \mid \hat{\mu}\left(X_{\infty}(r)\right)=1, \hat{\mu} \text { is } V \text {-quasi-invariant }\right\}
$$

and

$$
M_{1}^{V}(X):=\left\{\mu_{0} \in M(X) \mid \mu_{0}(X)=1, T_{V}\left(\mu_{0}\right)=\mu_{0}\right\}
$$

It is easy two see that both these sets are convex, and an application of the BanachAlaoglu theorem [36] shows that they are compact in the weak*-topology. Then, using the Krein-Milman theorem [36], we conclude that each of these sets is the convex weak*-closure of their extreme points. Moreover, since the correspondence is affine it preserves the extreme points.

This section is devoted to an analysis of the extreme points. Before we state and prove our main result on extreme points, we need to define the concepts of conditional expectation $E_{\mu_{0}}$, and relative ergodicity.
Proposition 6.1. Let $(X, \mathfrak{B}), r$ and $V$ be as above. Let $\mu_{0}$ be a measure in $M_{1}^{V}(X)$. Then for each bounded measurable function $g$ on $X$ there exists a bounded $r^{-1}(\mathfrak{B})$ measurable function $E_{\mu_{0}}(V g)$ such that

$$
\begin{equation*}
\int_{X} V g f \circ r d \mu_{0}=\int_{X} E_{\mu_{0}}(V g) f \circ r d \mu_{0}, \tag{6.1}
\end{equation*}
$$

for all bounded $\mathfrak{B}$-measurable functions $f$ on $X$. Moreover, this is unique up to $\mu_{0} \circ r^{-1}$-measure zero.

Proof. The positive linear functional

$$
\Lambda_{g}: f \circ r \mapsto \int_{X} V g f \circ r d \mu_{0}
$$

defines a measure on $\left(X, r^{-1}(\mathfrak{B})\right)$ which is absolutely continuous with respect to $\mu_{0}$. Indeed, if $\mu(E)=0$, then

$$
0=\int_{X} \chi_{E} d \mu_{0}=\int_{X} V \chi_{E} \circ r d \mu_{0}
$$

SO

$$
\int_{X} V g \chi_{E} \circ r d \mu_{0}
$$

Therefore, by the Radon-Nikodym theorem, there exists some $r^{-1}(\mathfrak{B})$-measurable function $E_{\mu_{0}}(V f)$ such that (6.1) holds. Since

$$
\left|\int_{X} E_{\mu_{0}}(V g) f \circ r d \mu_{0}\right| \leq\|V g\|_{\infty} \int_{X}|f \circ r| d \mu_{0}
$$

it follows that $\left|E_{\mu_{0}}(V g)\right| \leq\|V g\|_{\infty}$ holds $\mu_{0} \circ r^{-1}$ a.e.
The uniqueness is also clear from the definition (6.1).

Definition 6.2. A measure $\mu_{0} \in M_{1}^{V}(X)$ is called relatively ergodic with respect to $(r, V)$ if the only non-negative, bounded $\mathfrak{B}$-measurable functions $f$ on $X$ satisfying

$$
E_{\mu_{0}}(V f)=E_{\mu_{0}}(V) f \circ r, \text { pointwise } \mu_{0} \circ r^{-1} \text {-a.e., }
$$

are the functions which are constant $\mu_{0}$-a.e.
Theorem 6.3. Let $V: X \rightarrow[0, \infty)$ be bounded and measurable. Let $\hat{\mu} \in M_{q i, 1}^{V}\left(X_{\infty}(r)\right)$, and $\mu_{0}:=\hat{\mu} \circ \theta_{0}^{-1} \in M_{1}^{V}(X)$. The following affirmations are equivalent:
(i) $\hat{\mu}$ is an extreme point of $M_{q i, 1}^{V}\left(X_{\infty}(r)\right)$;
(ii) $V \circ \theta_{0} d \hat{\mu}$ is ergodic with respect to $\hat{r}$;
(iii) $\mu_{0}$ is an extreme point of $M_{1}^{V}(X)$;
(iv) $\mu_{0}$ is relatively ergodic with respect to $(r, V)$.

Proof. The arguments in the beginning of this section (mainly Theorem 5.3) show that (i) and (iii) are equivalent.

We now prove $(\mathrm{i}) \Rightarrow(\mathrm{ii})$. Suppose $\hat{\mu}$ is not ergodic. Then there exists two measurable subsets $A$ and $B$ of $X_{\infty}(r)$ such that

- $\int_{A} V \circ \theta_{0} d \hat{\mu}>0$ and $\int_{B} V \circ \theta_{0} d \hat{\mu}>0$
(note that then $\hat{\mu}(A)>0$ and $\hat{\mu}(B)>0$ ),
- $A \cup B=X_{\infty}(r)$, and
- $A$ and $B$ are $\hat{r}$-invariant.

Define then the measures

$$
\hat{\mu}_{A}(E)=\frac{\hat{\mu}(A \cap E)}{\hat{\mu}(A)}, \quad \hat{\mu}_{B}(E)=\frac{\hat{\mu}(B \cap E)}{\hat{\mu}(B)}
$$

for all $E \in \mathfrak{B}$.
Note that

$$
\begin{equation*}
\hat{\mu}=\hat{\mu}(A) \mu_{a}+\hat{\mu}(B) \mu_{B} . \tag{6.2}
\end{equation*}
$$

We prove next that $\mu_{A}, \mu_{B}$ are in $M_{q i, 1}^{V}\left(X_{\infty}(r)\right)$. Clearly they have total mass equal to 1 , so we only have to prove the $V$-quasi-invariance. For $f$ bounded, measurable
function on $X_{\infty}(r)$, we have:

$$
\begin{aligned}
\int_{X_{\infty}(r)} f \circ \hat{r}^{-1} d \mu_{A} & =\frac{1}{\hat{\mu}(A)} \int_{X_{\infty}(r)} f \circ \hat{r}^{-1} \chi_{A} d \hat{\mu} \\
& =\frac{1}{\hat{\mu}(A)} \int_{X_{\infty}(r)} f \circ \hat{r}^{-1} \chi_{A} \circ \hat{r}^{-1} d \hat{\mu} \quad(\hat{r}(A)=A) \\
& =\frac{1}{\hat{\mu}(A)} \int_{X_{\infty}(r)} V \circ \theta_{0} f \chi_{A} d \hat{\mu} \\
& =\int_{X_{\infty}(r)} V \circ \theta_{0} f d \mu_{A}
\end{aligned}
$$

Hence $\mu_{A}$ is $V$-quasi-invariant. The same argument works for $\mu_{B}$. Therefore $\mu_{A}$ and $\mu_{B}$ are both in $M_{q i, 1}^{V}\left(X_{\infty}(r)\right)$. Now, equation (6.2) contradicts the fact that $\hat{\mu}$ is extreme.
$($ ii $) \Rightarrow(\mathrm{i})$ : Suppose there are some measures $\mu_{1}, \mu_{2} \in M_{q i, 1}^{V}\left(X_{\infty}(r)\right)$ and some $\lambda \in[0,1]$ such that

$$
\hat{\mu}=\lambda \mu_{1}+(1-\lambda) \mu_{2}
$$

Then $\mu_{1}$ and $\mu_{2}$ are absolutely continuous with respect to $\hat{\mu}$. Therefore there exist $f_{1}$ and $f_{2}$ in $L^{1}\left(X_{\infty}(r), \hat{\mu}\right)$ such that

$$
d \mu_{1}=f_{1} d \hat{\mu}, \quad d \mu_{2}=f_{2} d \hat{\mu}
$$

Since $\mu_{1}$ and $\hat{\mu}$ are $V$-quasi-invariant, we have, for all bounded measurable functions $f$ on $X_{\infty}(r)$,

$$
\begin{aligned}
\int_{X_{\infty}(r)} f f_{1} \circ \hat{r} V \circ \theta_{0} d \hat{\mu} & =\int_{X_{\infty}(r)} f \circ \hat{r}^{-1} f_{1} d \hat{\mu} \\
& =\int_{X_{\infty}(r)} f \circ \hat{r}^{-1} d \mu_{1} \\
& =\int_{X_{\infty}(r)} V \circ \theta_{0} f d \mu_{1} \\
& =\int_{X_{\infty}(r)} V \circ \theta_{0} f f_{1} d \hat{\mu}
\end{aligned}
$$

Therefore $f_{1}=f_{1} \circ r$ pointwise $V \circ \theta_{0} d \hat{\mu}$-almost everywhere. The hypothesis implies then that $f_{1}$ is constant $V \circ \theta_{0} d \hat{\mu}$-a.e. The same argument shows that $f_{2}$ is constant $V \circ \theta_{0} d \hat{\mu}$-a.e.

Since $\hat{\mu}$ and $\mu_{1}$ are $V$-quasi-invariant we have also that

$$
1=\int_{X_{\infty}(r)} 1 \circ \hat{r}^{-1} d \mu_{1}=\int_{X_{\infty}(r)} V d \mu_{1}=\int_{X_{\infty}(r)} V f_{1} d \hat{\mu}
$$

It follows that $f_{1}=1$, pointwise $V \circ \theta_{0} d \hat{\mu}$-a.e. The same is true for $f_{2}$. Then

$$
\begin{aligned}
\int_{X_{\infty}(r)} f d \mu_{1}= & \int_{X_{\infty}(r)} V \circ \theta_{0} f \circ \hat{r} d \mu_{1}=\int_{X_{\infty}(r)} V \circ \theta_{0} f \circ \hat{r} f_{1} d \hat{\mu} \\
& =\int_{X_{\infty}(r)} V f \circ \hat{r} d \hat{\mu}=\int_{X_{\infty}(r)} f d \hat{\mu}
\end{aligned}
$$

Hence $\mu_{1}=\hat{\mu}=\mu_{2}$ and $\hat{\mu}$ is extreme.
(iii) $\Rightarrow$ (iv): Suppose $\mu_{0}$ is not relatively ergodic. Then there exists a bounded measurable function $f_{1} \geq 0$ on $X$ such that $E_{\mu_{0}}\left(V f_{1}\right)=E_{\mu_{0}}(V) f_{1} \circ r$, pointwise
$\mu_{0} \circ r^{-1}$-a.e., and $f_{1}$ is not constant $\mu_{0}$-a.e. We may assume that $\int_{X} f_{1} d \mu_{0}=1$. Define the measure $d \mu_{1}:=f_{1} d \mu_{0}$. We check that $\mu_{1}$ is in $M_{1}^{V}(X)$.

$$
\begin{gathered}
\int_{X} f d \mu_{1}=\int_{X} f f_{1} d \mu_{0}=\int_{X} V f \circ r f_{1} \circ r d \mu_{0}=\int_{X} E_{\mu_{0}}(V) f_{1} \circ r f \circ r d \mu_{0} \\
=\int_{X} E_{\mu_{0}}\left(V f_{1}\right) f \circ r d \mu_{0}=\int_{X} V f_{1} f \circ r d \mu_{0}=\int_{X} V f \circ r d \mu_{1} .
\end{gathered}
$$

Now choose some $0<\lambda<1$ such that $\lambda f_{1} \leq 1$ and define

$$
f_{2}:=\frac{1-\lambda f_{1}}{1-\lambda}, \quad d \mu_{2}=f_{2} d \mu_{0}
$$

Then $E_{\mu_{0}}\left(V f_{2}\right)=E_{\mu_{0}}(V) f_{2} \circ r$, and the same calculation as before shows that $T_{V}\left(\mu_{2}\right)=\mu_{2}$. Note also that $\mu_{2}(X)=1$. Since $\mu_{0}=\lambda \mu_{1}+(1-\lambda) \mu_{2}$ and $f_{1}$ is not constant, we have that $\mu_{1} \neq \mu_{0}$. It follows that $\mu_{0}$ is not an extreme point, thus contradicting the hypothesis and proving (iv).
(iv) $\Rightarrow$ (iii): Suppose

$$
\begin{equation*}
\mu_{0}=\lambda \mu_{1}+(1-\lambda) \mu_{2} \tag{6.3}
\end{equation*}
$$

for some $\mu_{1}, \mu_{2} \in M_{1}^{V}(X), \lambda \in[0,1]$. Then $\mu_{1}$ and $\mu_{2}$ are absolutely continuous with respect to $\mu_{0}$. Let $f_{1}, f_{2}$ be the corresponding Radon-Nikodym derivatives. The relation (6.3) implies that

$$
1=\lambda f_{1}+(1-\lambda) f_{2}, \quad \mu_{0} \text {-a.e. }
$$

In particular, $f_{1}$ and $f_{2}$ are bounded.
We know that $\mu$ and $\mu_{1}$ have the fixed-point property. Then for all bounded measurable functions $f$ on $X$,

$$
\begin{gathered}
\int_{X} E_{\mu_{0}}(V) f_{1} \circ r f \circ r d \mu_{0}=\int_{X} V f \circ r f_{1} \circ r d \mu_{0}=\int_{X} f f_{1} d \mu_{0}=\int_{X} f d \mu_{1} \\
=\int_{X} V f \circ r d \mu_{1}=\int_{X} V f \circ r f_{1} d \mu_{0}=\int_{X} E_{\mu_{0}}\left(V f_{1}\right) f \circ r d \mu_{0} .
\end{gathered}
$$

Therefore $E_{\mu_{0}}(V) f_{1} \circ r=E_{\mu_{0}}\left(V f_{1}\right), \mu_{0} \circ r^{-1}$-a.e.
The hypothesis implies that $f_{1}$ is constant $\mu_{0}$-a.e. Since $\mu_{1}(X)=\mu_{0}(X)=1$, it follows that $f_{1}=1$ and $\mu_{1}=\mu_{0}$, and therefore $\mu_{0}$ is extreme.

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## References

[1] L. Accardi, A. Frigerio, J.T. Lewis, Quantum stochastic processes, Publ. Res. Inst. Math. Sci. 18 (1982), 97-133. MR0660823 (84m:82031)
[2] A. Aldroubi, D. Larson, W.-S. Tang, E. Weber, Geometric aspects of frame representations of abelian groups, Trans. Amer. Math. Soc. 356 (2004), 4767-4786. MR2084397
[3] J. Ashley, B. Markus, S. Tuncel, The classification of one-sided Markov chains, Ergodic Theory Dynam. Systems 17 (1997), 269-295. MR1444053 (98k:28021)
[4] A. Aldroubi, Q. Sun, W.-S. Tang, Nonuniform average sampling and reconstruction in multiply generated shift-invariant spaces, Constr. Approx. 20 (2004), 173-189. MR2036639
[5] L.W. Baggett, A. Carey, W. Moran, P. Ohring, General existence theorems for orthonormal wavelets, an abstract approach, Publ. Res. Inst. Math. Sci. 31 (1995), 95-111. MR1317525 (96c:42060)
[6] L.W. Baggett, P.E.T. Jorgensen, K.D. Merrill, J.A. Packer, An analogue of Bratteli-Jorgensen loop group actions for GMRA's, in Wavelets, Frames and Operator Theory (College Park, 2003), eds. C. Heil et al., Contemp. Math., vol. 345, American Mathematical Society, 2004, pp. 11-25. MR2066819
[7] L.W. Baggett, K.D. Merrill, Abstract harmonic analysis and wavelets in $\mathbb{R}^{n}$, in The Functional and Harmonic Analysis of Wavelets and Frames (San Antonio, 1999), eds. L.W. Baggett, D.R. Larson, Contemp. Math., vol. 247, American Mathematical Society, 1999, pp. 17-27. MR $1735967 \mid(2001 \mathrm{~b}: 42043)$
[8] V. Baladi, Positive Transfer Operators and Decay of Correlations, Advanced Series in Nonlinear Dynamics, vol. 16, World Scientific, River Edge, NJ, Singapore, 2000. MR1793194 (2001k:37035)
[9] L. Barreira, Ya. B. Pesin, Lyapunov exponents and smooth ergodic theory, University Lecture Series, vol. 23, American Mathematical Society, Providence, RI, 2002. MR1862379 (2003a:37040)
[10] B. Brenken, P.E.T. Jorgensen, A family of dilation crossed product algebras, J. Operator Theory 25 (1991), 299-308. MR1203035 (94m:46103)
[11] H. Brolin, Invariant sets under iteration of rational functions, Ark. Mat. 6 (1965), 103-144. MR 0194595 (33:2805)
[12] I. Daubechies, Ten Lectures on Wavelets, CBMS-NSF Regional Conf. Ser. in Appl. Math., vol. 61, Society for Industrial and Applied Mathematics, Philadelphia, 1992. MR1162107 (93e:42045)
[13] V. Deaconu, P.S. Muhly, $C^{*}$-algebras associated with branched coverings, Proc. Amer. Math. Soc. 129 (2001), 1077-1086. MR1814145 (2002a:46075)
[14] J.L. Doob, Classical Potential Theory and Its Probabilistic Counterpart, Grundlehren der Mathematischen Wissenschaften, vol. 262, Springer-Verlag, 1984. MR0731258 (85k:31001)
[15] D.E. Dutkay, The wavelet Galerkin operator, J. Operator Theory 51 (2004), 49-70. MR 2055804 (2005c:47041)
[16] D.E. Dutkay, P.E.T. Jorgensen, Wavelets on fractals, to appear in Rev. Mat. Iberoamericana.
[17] D.E. Dutkay, P.E.T. Jorgensen, Martingales, endomorphisms, and covariant systems of operators in Hilbert space, preprint 2004, arxiv math.CA/0407330.
[18] R. Gohm, B. Kummerer, T. Lang, Non-commutative symbolic coding, Preprint 2005, Greifswald University.
[19] R.F. Gundy, Two remarks concerning wavelets: Cohen's criterion for low-pass filters and Meyer's theorem on linear independence, in The Functional and Harmonic Analysis of Wavelets and Frames (San Antonio, 1999), eds. L.W. Baggett, D.R. Larson, Contemp. Math., vol. 247, American Mathematical Society, 1999, pp. 249-258. MR1738093 (2001b:42051)
[20] J.E. Hutchinson, Fractals and self-similarity, Indiana Univ. Math. J. 30 (1981), no. 5, 713747. MR0625600 (82h:49026)
[21] P.E.T. Jorgensen, Ruelle operators: Functions which are harmonic with respect to a transfer operator, Mem. Amer. Math. Soc. 152 (2001), no. 720. MR. 1837681 (2002c:46117)
[22] P.E.T. Jorgensen, Analysis and Probability, preprint 2004.
[23] P.E.T. Jorgensen, D.W. Kribs, Wavelet representations and Fock space on positive matrices, J. Funct. Anal. 197 (2003), 526-559. MR 1960424 (2004e:42055)
[24] P.E.T. Jorgensen, P.S. Muhly, Selfadjoint extensions satisfying the Weyl operator commutation relations, J. Analyse Math. 37 (1980), 46-99. MR0583632 (82k:47058)
[25] E. Katsoulis, D.W. Kribs, Isomorphisms of algebras associated with directed graphs, Math. Ann. 330 (2004), 709-728. MR 2102309
[26] D.W. Kribs, On bilateral weighted shifts in noncommutative multivariable operator theory, Indiana Univ. Math. J. 52 (2003), 1595-1614. MR2021049 (2004k:47004)
[27] A.N. Kolmogoroff, Grundbegriffe der Wahrscheinlichkeitsrechnung (German), SpringerVerlag, 1977 (reprint of the 1933 original); Foundations of the Theory of Probability (English), Chelsea, 1950. MR0494348(58:13242)
[28] B. Kummerer, Markov dilations on $W^{*}$-algebras, J. Funct. Anal. 63 (1985), 139-177. MR 0803091 ( $87 \mathrm{~b}: 46070$ )
[29] S. G. Mallat, Multiresolution approximations and wavelet orthonormal bases of $L^{2}(R)$, Trans. Amer. Math. Soc. 315 (1989), 69-87. MR1008470|(90e:42046)
[30] P.S. Muhly, B. Solel, Quantum Markov processes (correspondences and dilations), Internat. J. Math. 13 (2002), no. 8, 863-906. MR1928802 (2003j:46100)
[31] J. Neveu, Discrete-Parameter Martingales, North-Holland, 1975. MR0402915 (53:6729)
[32] R.D. Nussbaum, S.M. Verduyn Lunel, Generalizations of the Perron-Frobenius theorem for nonlinear maps, Mem. Amer. Math. Soc. 138 (1999), no. 659. MR1470912 (99i:58125)
[33] G. Popescu, Commutant lifting, tensor algebras, and functional calculus, Proc. Edinb. Math. Soc. (2) 44 (2001), 389-406. MR.1880399|(2003f:47125)
[34] G. Popescu, Central intertwining lifting, suboptimization, and interpolation in several variables, J. Funct. Anal. 189 (2002), 132-154. MR1887631 (2002m:47013)
[35] W. Rudin, Real and Complex Analysis, third edition, McGraw-Hill, 1987. MR0924157 (88k:00002)
[36] W. Rudin, Functional Analysis, second edition, International Series in Pure and Applied Mathematics, McGraw-Hill, 1991. MR1157815 (92k:46001)
[37] D. Ruelle, The thermodynamic formalism for expanding maps, Comm. Math. Phys. 125 (1989), 239-262. MR1016871 (91a:58149)
[38] D. Ruelle, Application of hyperbolic dynamics to physics: some problems and conjectures, Bull. Amer. Math. Soc. (N.S.) 41 (2004), 275-278. MF2058287 (2005b:37038)
[39] M. Su, The information topology and true laminations for diffeomorphisms, Conform. Geom. Dyn. 8 (2004), 36-51. MR2060377
[40] D. Szasz, T. Varju, Local limit theorem for the Lorentz process and its recurrence in the plane, Ergodic Theory Dynam. Systems 24 (2004), 257-278. MR2041271

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