# ON CONTINUOUS COLLECTIONS OF MEASURES 

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## 1. Introduction

In this paper we will consider a problem whose formulation in probability terms is essentially as follows: when can one construct a stochastic process $\{Z(t) ; t \in M\}$ having continuous paths and preassigned one dimensional distributions. We always take the state space $X$ for the process to be a metric space and the parameter set $M$ to be a compact topological space. An obvious necessary condition is that the preassigned distributions for the individual $\boldsymbol{Z}(t)$ vary continuously with $t$. This condition is also sufficient [1], if $X$ is complete metric and $M$ is zero dimensional, the Cantor set for example. If $M$ is anything else, for example an interval on the real line, the simple necessary condition is no longer sufficient as one easily sees, and further conditions on $X$ and on the desired distributions for the $Z(t)$ are needed. We will give here a theorem which treats the case in which $M$ is arbitrary.

First we must introduce some notation and a precise statement of the problem, which also makes it look more like the sort of thing one ordinarily considers.

If $Y$ is a topological space, let $P(Y)$ denote the set of all probability measures on the Borel sets of $Y$. Let $C(Y)$ denote the continuous bounded real valued functions on $Y$. We give $C(Y)$ the uniform topology and $P(Y)$ the topology generated by the functions $\mu \rightarrow \int f d \mu, f \in C(Y)$. Given a measure space $(\Omega, \mathscr{F}, \mu)$ and a mapping $\varphi$ from $\Omega$ to $Y$ which is measurable relative to $\mathscr{F}$ and the Borel sets of $Y$ let $\varphi \mu$ denote the measure on the Borel sets of $Y$ defined by $\varphi \mu(A)=$ $\mu\left(\varphi^{-1}(A)\right)$. From now on let $M$ and $X$ be compact metric spaces and let $C(M, X)$ denote the continuous functions from $M$ to $X$ under the uniform topology. Each $t$ in $M$ defines, by evaluation at $t$, a continuous function from $C(M, X)$ to $X$. We denote this mapping simply by $t$, so that $t f=f(t)$ for $f \in C(M, X)$. Let $\mu$ be an element of $P(C(M, X))$; then in the notation we have just introduced, $t \mu$ defines an element of $P(X)$, namely, $t \mu(A)=\mu(\{f: f(t) \in A\})$. Moreover the mapping $t \rightarrow t \mu$ from $M$ to $P(X)$ is continuous.

In this paper we shall consider the converse construction; that is, given a continuous function $T$ from $M$ to $P(X)$, when is there a measure $\mu$ in $P(C(M, X)$ ) such that $t \mu=T(t)$ for all $t \in M$. Note that any such $T$ defines, via the formula $\left(T^{*} f\right)(t)=\int f d T(t)$, a continuous linear mapping from $C(X)$ to $C(M)$ such that $T^{*}(1)=1=\left\|T^{*}\right\|$; any such mapping $T^{*}$ arises in this way. The mappings of
this sort which are also multiplicative are those of the form $T^{*} f(t)=f(g(t))$ for some $g \in C(M, X)$, and so the existence of $\mu \in P(C(M, X))$ with $T(t)=t \mu$ is equivalent to the existence of an integral representation for $T^{*}$ as an average of multiplicative operators. Such integral representations are objects of general interest. Finally one can check immediately that given a $T$, a measure $\mu$ of the desired sort exists if and only if there is a probability space $(\Omega, \mathscr{F}, P)$ and a function $Z$ from $M \times \Omega$ to $X$ such that $\omega \rightarrow Z(t, \omega)$ is $\mathscr{F}$ measurable for each $t \in M$, and $t \rightarrow Z(t, \omega)$ is continuous for each $\omega \in \Omega$ (that is, a stochastic process with continuous paths), and such that $P(Z(t) \in A)=T(t)(A)$ for all $t \in M$ and Borel sets $A \subset X$.

## 2. The main theorem

In the first few paragraphs we shall state our result and give a description of what we require for the proof. The remainder of the section shall be devoted to the rather complicated details. In Section 3 we shall discuss the necessity of our hypotheses. We always use $M$ and $X$ to denote compact metric spaces, and $I$ to denote the closed unit interval.

Theorem 2.1. Let $X$ be connected and locally connected. Let $T: M \rightarrow P(X)$ be continuous and have the property that for each $t \in M$ the support of the measure $T(t)$ is all of $X$. Then there is a $\mu \in P(C(M, X))$ such that $t \mu=T(t)$ for all $t \in M$.

Lemma 2.1. Theorem 2.1 is valid in the case $X=I$.
Proof. Let $\Omega=(0,1)$, let $\mathscr{F}=$ Borel sets of $(0,1)$, and let $P=$ Lebesgue measure on $\mathscr{F}$. Let $F_{t}$ be the distribution function of $T(t)$, that is, $F_{t}(x)=$ $T(t)([0, x])$ for $x \in[0,1]$, and for $\omega \in \Omega$ define $Z(t, \omega)=\inf \left\{x: F_{t}(x) \geqq \omega\right\}$. Then the following facts are easy to check: (a) $\omega \rightarrow Z(t, \omega)$ is $\mathscr{F}$ measurable and $P\{\omega: Z(t, \omega) \in A\}=T^{\prime}(t)(A)$ for all $t \in M$ and Borel sets $A \subset X$; (b) for each $\omega \in \Omega$ the mapping $t \rightarrow Z(t, \omega)$ is continuous at each point $r$ such that $\left\{x: F_{r}(x)=\right.$ $\omega\}$ consists of at most one point. Since the support of $T(t)$ is by hypothesis all of $I$, it follows that $t \rightarrow Z(t, \omega)$ is continuous for all $\omega$. In view of the remarks at the end of Section 1 the proof is complete.

For the next statement let $P_{0}(X)$ denote the subspace of $P(X)$ consisting of those measures whose support is all of $X$.

Theorem 2.2. Let $X$ be connected and locally connected. Then there is a continuous function $\tilde{\varphi}$ from $P_{0}(X)$ to $P_{0}(I)$ and a continuous function $\varphi$ from $I$ to $X$ such that $\varphi(\tilde{\varphi} \mu)=\mu$ for all $\mu$ in $P_{0}(X)$.

Remarks. First of all, Theorem 2.1 is an immediate consequence of Lemma 2.1 and Theorem 2.2. Indeed if $\varphi$ and $\tilde{\varphi}$ are as in Theorem 2.2 and $T$ is continuous from $M$ into $P_{0}(X)$, then $\widetilde{T}(t)=\tilde{\varphi}(T(t))$ defines a continuous function $\tilde{T}$ from $M$ into $P_{0}(I)$. By Lemma 2.1, there is a measure $\tilde{\mu}$ in $P(C(M, I))$ such that $t \tilde{\mu}=$ $\widetilde{T}(t)$. Now $\varphi: I \rightarrow X$ defines a mapping, also called $\varphi$, from $C(M, I)$ to $C(M, X)$ by $\varphi f(t)=\varphi(f(t)), f \in C(M, I)$. If we define $\mu \in P(C(M, X))$ by $\mu=\varphi \tilde{\mu}$, then it is trivial to check that $\mu$ satisfies the conclusion of Theorem 2.1. Secondly, we should comment on the hypotheses on $X$ and on the long proof of Theorem 2.2:
according to a theorem of Hahn and Mazurkiewicz, $X$ is compact metric, connected, and locally connected if and only if there is a continuous function $\varphi$ from $I$ onto $X$. Now simply from the existence of such a $\varphi$, it follows immediately that given any $\theta \in P(X)$ there is a $\tilde{\theta} \in P(I)$ such that $\varphi \tilde{\theta}=\theta$. But unless $\varphi$ has some additional properties it will not be true that $\tilde{\theta}$ can be chosen to depend on $\theta$ in a continuous manner. Much of the detail in the proof which follows stems from this fact.

Before coming to the proof of Theorem 2.2 we must develop some additional facts. A finite partition of unity on $X$ is an indexed collection $F=\left\{f_{1}, \cdots, f_{n}\right\}$ of nonnegative continuous functions (not identically 0 ) on $X$ such that $\Sigma_{i=1}^{n} f_{i}=1$. It is possible that $f_{i} \neq f_{j}$ for some $i=j$; but we wish to regard functions with different indices as different elements of the partition. It will simplify the reading if this possibility is ignored. We shall first construct some partitions of unity with additional properties.

Lemma 2.2. There is a sequence $\left\{F_{n}\right\}_{n} \geqq 0$ of finite partitions of unity and a sequence $\left\{\pi_{n}\right\}_{n \geqq 0}$ of mappings, $\pi_{n}$ mapping $F_{n+1}$ onto $F_{n}$, such that (i) for each $n$ and $f \in F_{n}$ the support of $f$ is a connected set of diameter no greater than $1 / n$; (ii) for each $n$ and $f \in F_{n}$ we have $f=\Sigma g_{i}$, the sum being over those $g \in F_{n+1}$ such that $\pi_{n}(g)=f$.

Proof. Let $\mathscr{V}_{0}=\{X\}$ and for each $n \geqq 1$ let $\mathscr{V}_{n}$ be a finite open cover of $X$ by open connected subsets each of diameter no greater than $2^{-(n+2)}$. Since $X$ is locally connected such covers exist. Now given any subset $A$ of $X$, define $s_{n}(A)$ to be the union of those sets $V \in \mathscr{V}_{n}$ such that $V \cap A \neq \phi$, and define $s_{n, \infty}(A)$ to be $\cup_{k=1}^{\infty} s_{n+k} \cdots\left(s_{n}(A)\right)$. Then $s_{n, \infty}(A)$ is open, it is connected if $A$ is connected, and if $A$ has diameter not exceeding $\varepsilon$ then $s_{n, \infty}(A)$ has diameter not exceeding $\varepsilon+2^{-n}$. In addition it is easy to check that $s_{n, \infty}(A \cup B)=s_{n, \infty}(A) \cup s_{n, \infty}(B)$. Now define a sequence of covers $\left\{\mathscr{U}_{n}\right\}$ of $X$ by taking $\mathscr{U}_{n}$ to consist of the sets $s_{n+1, \infty}(V)$ as $V$ ranges over the elements of $\mathscr{V}_{n}$. By what we have just said, each $\mathscr{U}_{n}$ is a finite cover of $X$ by open connected sets each of diameter no greater than $1 / n$ and in addition each set in $\mathscr{U}_{n}$ is the union of the sets in $\mathscr{U}_{n+1}$ which are contained in it.

Now to construct the partitions of unity and mappings take $F_{0}$ to consist of the function 1. Suppose $F_{n}$ has been constructed and for each $f \in F_{n},\{f>0\}$ is an element $U(f)$ of $\mathscr{U}_{n}$. Let $Z_{1}, \cdots, Z_{n(f)}$ denote the elements of $\mathscr{U}_{n+1}$ contained in $U(f)$ and define

$$
\begin{equation*}
g_{f, i}=f \frac{f_{Z_{i}}}{f_{Z_{1}}+\cdots+f_{Z_{n(f)}}}, \tag{2.1}
\end{equation*}
$$

where $f_{V}$ denotes any nonnegative continuous function such that $\left\{f_{V}>0\right\}$ is the open subset $V$ of $X$. We take $F_{n+1}$ to consist of the functions $g_{f, i}, f \in F_{n}$, $i \leqq n(f)$, and we define $\pi_{n}$ by $\pi_{n} g_{f, i}=f$. This defines inductively sequences with the desired properties, so the proof is complete. We note once again that the same function may appear more than once in a partition.

Let $F=\left\{f_{1}, \cdots, f_{n}\right\}$ be an indexed collection of continuous functions on $X$. By the 1-complex associated with $F$ we mean a linearly independent set of $n$ points, which we call $f_{1}, \cdots, f_{n}$ in a Euclidean space of high enough dimension and closed line segments, denoted $\left[f_{i}, f_{j}\right]$, joining those points $f_{i}$ and $f_{j}$ for which the function $f_{i} f_{j}$ is not identically zero. The 1 -complex is topologized as a subset of the Euclidean space. Clearly the choice of Euclidean space and linearly independent points is unimportant, so we will not specify these further. It will be convenient at times to denote the point $f_{i}$ also by $\left[f_{i}, f_{i}\right]$. The symbol $\left(f_{i}, f_{j}\right)$ will denote the line segment $\left[f_{i}, f_{j}\right]$ with the end points deleted. Any two segments in a 1 -complex do not meet, or else meet in a single point which is an end point of each. Frequently we refer to the points $f$ as vertices. Now we pick once and for all a sequence $\left\{F_{n}\right\}$ of finite partitions of unity, and a sequence $\left\{\pi_{n}\right\}$ of mappings satisfying the requirements of Lemma 2.2. Let $X_{n}$ denote the 1-complex associated with $F_{n}$. We shall use $\pi_{n}: F_{n+1} \rightarrow F_{n}$ to define a mapping also called $\pi_{n}$, from $X_{n+1}$ to $X_{n}$ as follows: for a vertex $f$ in $X_{n+1}$ we set $\pi_{n}(f)=\left(\pi_{n}(f)\right)$ and we map a line segment $\left[f_{i}, f_{j}\right]$ in $X_{n+1}$ linearly onto the line segment $\left[\pi_{n}\left(f_{i}\right), \pi_{n}\left(f_{j}\right)\right]$ in $X_{n}$. We call attention to the following fact whose proof is immediate from the properties of $\left\{F_{n}\right\}$ and $\left\{\pi_{n}\right\}$.

Lemma 2.3. Each $X_{n}$ is connected; $\pi_{n}: X_{n+1} \rightarrow X_{n}$ is onto. For each segment (or point) $[f, g] \subset X_{n}, \pi_{n}^{-1}([f, g])$ is a connected 1-complex contained in $X_{n+1}$.

Let $J=[a, b]$ be a closed interval on the line and let $Y$ be a 1 -complex associated with functions $f_{1}, \cdots, f_{n}$. A mapping $\varphi: J \rightarrow Y$ is called piecewise linear if there is a finite subset $V$ of $J$ with $a$ and $b$ in $V$ such that for each $v \in V, \varphi(v)$ is a vertex of $Y$ and $\varphi$ is linear on each interval between adjacent members of $V$. Now suppose that $Y$ is connected. Then one proves easily that given two vertices $f_{i}$ and $f_{j}$ of $Y$ there is a piecewise linear mapping of $J$ onto $Y$ such that $\varphi(a)=f_{i}$ and $\varphi(b)=f_{j}$. In fact one sees without difficulty that if $Y$ is connected and $f_{i}$ and $f_{j}$ are vertices, then a piecewise linear $\varphi: J \rightarrow Y$ exists with the following properties:
(a) $\varphi$ is onto $Y, \quad \varphi(a)=f_{i}, \quad \varphi(b)=f_{j}$;
(b) for each vertex or segment $\left[f_{r}, f_{s}\right]$ of $Y, \varphi^{-1}\left(\left[f_{r}, f_{s}\right]\right)$ is a union of nondegenerate closed intervals in $J$;
(c) $\varphi\left(\left[a, a^{\prime}\right]\right)=f_{i}, \quad \varphi\left(\left[b^{\prime}, b\right]\right)=f_{j}$
where $a^{\prime}=a+\frac{1}{3}(b-a), b^{\prime}=b-\frac{1}{3}(b-a)$ (property (c) is needed only for a slight, technical reason and might better be ignored).

Now to avoid trivialities assume that $Y$ consists of more than a single point. Any piecewise linear $\varphi$ with properties (a), (b), (c) breaks up $J$ in a natural way into nondegenerate nonoverlapping closed intervals $J_{1}, \cdots, J_{r}$ and $I_{1}, \cdots, I_{r-1}$ such that $\varphi$ maps each $J$ onto a vertex in $Y$ and each $I_{k}$ linearly onto a nondegenerate segment in $Y$. We assume the numbering has been done so that (with an obvious notation) $J_{1} \leqq I_{1} \leqq J_{2} \leqq \cdots \leqq I_{r-1} \leqq J_{r}$. If we break each of $J_{2}, \cdots, J_{r-1}$ into a left half $J_{i, 1}$, and a right half $J_{i, 2}$, then $J$ is broken up into units $J_{1} I_{1} J_{2,1}, J_{2,2} I_{2} J_{3,1}, J_{3,2} I_{3} J_{4,1}, \cdots, J_{r-1,2} I_{r-1} J_{r}$, and in each unit $\varphi$ is constant on the $J$ and linear and nonconstant on the $I$. Given any $\varphi$ satisfying (b) and any decomposition of $J$ into nondegenerate intervals $J_{1} \leqq I_{1} \leqq J_{2,1} \leqq$
$J_{2,2} \leqq I_{2} \leqq J_{3,1} \leqq \cdots$, such that $\varphi$ is constant on the $J$ and nonconstant linear on the $I$, we shall refer to the setup as a piecewise linear mapping with a specification of units and we shall refer to the triples $J_{i, 2} I_{i} J_{i+1,1}$ as units. The specific decomposition we just gave is not the only one possible; for example, we could have broken up the intermediate $J$ into unequal parts. However any $\varphi$ satisfying (a), (b), (c) together with the specific decomposition we gave will be denoted by $\varphi\left[J, Y, f_{i}, f_{j}\right]$.

Now we return to the partitions $F_{n}$ and associated 1-complexes $X_{n}$ and define a sequence $\left\{\varphi_{n}\right\}$ of piecewise linear mappings of $I$ onto $X_{n}$ with certain additional properties. We are going to assume that for each $n$ and $f \in F_{n}, \pi_{n}^{-1}(f)$ consists of more than one function. This involves no loss of generality for we may replace each element $g \in F_{n+1}$ by two elements, each being the function $\frac{1}{2} g$. Define $\varphi_{0}$ by $\varphi_{0}(x)=1$ for all $x \in I$ (of course 1 denotes the single point of $X_{0}$ ). Suppose $\varphi_{n}: I \rightarrow X_{n}$ has been defined so that $\varphi_{n}$ satisfies property (b) with $J=I$ and $Y=X_{n}$, and that we have specified a decomposition $J_{1} \leqq I_{1} \leqq J_{2,1} \leqq$ $J_{2,2} \leqq \cdots$, of $[0,1]$ into units. Let $J$ denote any interval of type $J$ or $I$ in the decomposition and recall that by Lemma $2.3, \pi_{n}^{-1}\left[\varphi_{n}(J)\right]$ is a connected 1-complex contained in $X_{n+1}$. To define $\varphi_{n+1}$ pick any points $f_{1}$ and $f_{1} \in$ $\pi_{n}^{-1}\left[\varphi_{n}\left(J_{1}\right)\right]$ and define $\varphi_{n+1}$ on $J_{1}$ to be any $\varphi\left[J_{1}, \pi_{n}^{-1}\left[\varphi_{n}\left(J_{1}\right)\right], f_{1}, f_{2}\right]$. Then pick any $g \in \pi_{n}^{-1}\left[\varphi_{n}\left(J_{2,1}\right)\right]$ and define $\varphi_{n+1}$ on $I_{1}$ to be any $\varphi\left[I_{1}, \pi_{n}^{-1}\left[\varphi_{n}\left(I_{1}\right)\right]\right.$, $\left.f_{2}, g\right]$. Note that $f_{2}$ and $g$ are indeed both in $\pi_{n}^{-1}\left[\varphi_{n}\left(I_{1}\right)\right]$ so this makes sense. We continue in this manner until $\varphi_{n+1}$ is defined on all of $I$. We take the units for $\varphi_{n+1}$ to be those arising from the decomposition of the segments making up the units for $\varphi_{n}$. It is clear that each unit for $\varphi_{n}$ is a finite union of units for $\varphi_{n+1}$, and $\varphi_{n+1}$ satisfies (b) with $J=I$ and $Y=X_{n+1}$ and
(d) if $[f, g]$ is a vertex or segment in $X_{n}$ and $J$ is one of the intervals making up $\varphi_{n}^{-1}([f, g])$ then $\varphi_{n+1}(J)=\pi_{n}^{-1}([f, g])$.

Thus we have defined inductively a sequence $\left\{\varphi_{n}\right\}$ of piecewise linear maps. This sequence will be held fixed from now on. Consider the units specified for $\varphi_{n}$ and let $\varepsilon_{n}$ denote the length of the longest one. From (c) it follows easily that $\varepsilon_{n} \rightarrow 0$ as $n \rightarrow \infty$.

Now we shall use $\left\{\varphi_{n}\right\}$ to define a continuous function $\varphi$ from $I$ onto $X$. Given $n$ and $z \in X_{n}$ set

$$
X_{n}(z)= \begin{cases}\overline{\{f>0\}} & \text { if } z \text { is a vertex } f  \tag{2.2}\\ \{\overline{\{f>0\}} \cup \overline{\{g>0\}} & \text { if } z \in(f, g)\end{cases}
$$

Each $X_{n}(z)$ is a closed subset of $X$ having diameter no greater than $2 / n$. For $a \in I$ define $\Phi_{n}(a)=X_{n}\left(\varphi_{n}(a)\right)$. Using property (d) the reader will verify without difficulty that $\Phi_{n+1}(a) \subset \Phi_{n}(a)$. We define $\varphi$ by $\varphi(a)=\cap_{n}\left\{\Phi_{n}(a)\right\}$. It is simple to check that $\varphi$ is a continuous function from $I$ onto $X$. Note for later use that if $K$ is a closed subset of $X$ then $\left\{a: \Phi_{n}(a) \subset K\right\}$ is a closed subset of $I$. In addition suppose $O$ is an open subset of $X$ and $\left\{K_{n}\right\}$ is an increasing sequence of compact subsets of $O$ such that every compact subset of $O$ is contained in one of the $K_{n}$. Then the sequence $\left\{a: \Phi_{n}(a) \subset K_{n}\right\}$ increases to $\varphi^{-1}(O)$.

Finally we are ready to define the mapping $\tilde{\varphi}$ required for Theorem 2.2. Given a measure $\mu \in P(X)$ define a measure $\mu_{n} \in P\left(X_{n}\right)$ as follows: $\mu_{n}$ puts mass $\int f^{2} d \mu$ at each vertex $f \in X_{n}$ and on each open segment $(f, g), \mu_{n}$ is linear Lebesgue measure of mass $\int 2 f g d \mu$. It follows immediately from the definition of $\pi_{n}: X_{n+1} \rightarrow X_{n}$ that for each $n$, we have $\pi_{n} \mu_{n+1}=\mu_{n}$. Clearly the support of $\mu_{n}$ is all of $X_{n}$ if the support of $\mu$ is all of $X$. In addition it is important to note that if $\mu$ varies continuously in $P(X)$ then $\mu_{n}(A)$ varies continuously for every Borel subset $A$ of $X_{n}$. We shall suppose from now on that $\mu \in P_{0}(X)$.

We are going to define a sequence $\left\{\mu^{n}\right\}$ of measures in $P_{0}(I)$ such that $\varphi_{n} \mu^{n}=$ $\mu_{n}$ and having additional properties. For each vertex $f \in X_{1}$, let $\mu^{1}$ put linear Lebesgue measure of mass $\mu_{1}(\{f\})$ on the set $\varphi_{1}^{-1}(\{f\})$ (recall this set is a union of nondegenerate intervals) and for each open segment ( $f, g$ ) contained in $X_{1}$, let $\mu^{1}$ put linear Lebesgue measure of mass $\mu_{1}((f, g))$ on the set $\varphi_{1}^{-1}((f, g))$. Suppose $\mu^{n}$ has been defined and has the following properties (which do indeed hold in case $n=1$ ): (i) $\varphi_{n} \mu^{n}=\mu_{n}$, (ii) for each closed interval $J$ appearing in any unit for $\varphi_{n} ; \mu^{n}$ restricted to $J$ is a nonzero multiple of Lebesgue measure on $J$; and (iii) for each $J$ as in (b) $\mu^{n}(J)$ varies continuously with $\mu$. We shall describe how to construct $\mu^{n+1}$ so that (i), (ii), and (iii) hold at the ( $n+1$ )st stage.

As described in the construction of $\varphi_{n}$ we have $I$ broken up into units $I_{1} \leqq \cdots \leqq I_{r}$ and each $I_{i}$ is made up of three nondegenerate closed intervals $I_{i, 1} \leqq I_{i, 2} \leqq I_{i, 3}$ such that $\varphi_{n}$ is constant on $I_{i, 1}$ and $I_{i, 3}$ and $\varphi_{n}$ is linear and nonconstant on $I_{i, 2}$. Let us fix an $i$, and denote $I_{i}$ by $J$ and $I_{i, 1}, I_{i, 2}, I_{i, 3}$ by $J_{1} J_{2}, J_{3}$. Suppose that $\varphi_{n}\left(J_{1}\right)=f$ and $\varphi_{n}\left(J_{3}\right)=g$. Let $\theta_{1}$ be the measure that puts, for each vertex $f^{\prime} \in \pi_{n}^{-1}(\{f\})$, linear Lebesgue measure of mass $\mu_{n+1}\left(\left\{f^{\prime}\right\}\right) / \mu_{n}(\{f\})$ on $J \cap \varphi_{n+1}^{-1}\left(\left\{f^{\prime}\right\}\right)$, and puts, for each $\left(f^{\prime}, f^{\prime \prime}\right) \subset \pi_{n}^{-1}(\{f\})$, linear Lebesgue measure of mass $\mu_{n+1}\left(\left(f^{\prime}, f^{\prime \prime}\right)\right) / \mu_{n}(\{f\})$ on $J \cap \varphi_{n+1}^{-1}\left(\left(f^{\prime}, f^{\prime \prime}\right)\right)$. Define $\theta_{3}$ similarly using $g$ in place of $f$. Define $\theta_{2}$ to be the measure that puts, for each interval $\left(f^{\prime}, g^{\prime}\right) \subset \pi_{n}^{-1}((f, g))$, linear Lebesgue measure of mass $\mu_{n+1}\left(\left(f^{\prime}, g^{\prime}\right)\right) /$ $\mu_{n}((f, g))$ on the set $J \cap \varphi_{n+1}^{-1}\left(\left(f^{\prime}, g^{\prime}\right)\right)$. Finally, recalling that $J=I_{i}$, put

$$
\begin{equation*}
\mu_{i}^{n+1}=\mu^{n}\left(J_{1}\right) \theta_{1}+\mu^{n}\left(J_{2}\right) \theta_{2}+\mu^{n}\left(J_{3}\right) \theta_{3}, \tag{2.3}
\end{equation*}
$$

and set

$$
\begin{equation*}
\mu^{n+1}=\mu_{1}^{n+1}+\cdots+\mu_{r}^{n+1} \tag{2.4}
\end{equation*}
$$

It is a slightly tedious but straightforward matter to verify that $\mu^{n+1}$ satisfies properties (i) through (iii). In addition we have the important fact that if $J$ is one of the units for $\varphi_{n}$, then $\mu^{n+1}(J)=\mu^{n}(J)$. Now the construction of $\tilde{\varphi}$ is practically complete-if $J$ is a unit for $\varphi_{n}$ then $\mu^{m}(J)=\mu^{n}(J)$ for $m \geqq n$. (Recall a unit for $\varphi_{n}$ is a union of units for $\varphi_{n+1}$ ). Since the units for $\varphi_{n}$ approach 0 in length as $n \rightarrow \infty$, it follows that as $n \rightarrow \infty, \mu^{n}$ approaches a measure which we define to be $\tilde{\varphi} \mu$. From (ii) it follows that the support of $\tilde{\varphi} \mu$ is all of $I$, and from (iii) it follows that $\tilde{\varphi} \mu$ varies continuously with $\mu$. The last thing we must check is that $\varphi(\tilde{\varphi} \mu)=\mu$. Call $\tilde{\varphi} \mu=\nu$. Since $v$ and $\mu$ are both probability measures it is enough to check the $v\left(\varphi^{-1}(O)\right) \leqq \mu(O)$ for every open subset $O$ of $X$. Let $K_{n}$ consist of
those points in $X$ distant from $X-O$ by at least $1 / n$, and let $U_{n}=\left\{a: \Phi_{n}(a) \subset K_{n}\right\}$. If $A$ is a compact subset of $\varphi^{-1}(O)$ then $A \subset U_{n}$ for all large $n$. Since $\mu^{n}$ converges to $v$ it follows that

$$
\begin{equation*}
v\left(\varphi^{-1}(O)\right) \leqq \lim \inf \mu^{n}\left(U_{n}\right) \tag{2.5}
\end{equation*}
$$

But from the definition of $\Phi_{n}(a)$ it follows that

$$
\begin{equation*}
\mu^{n}\left(U_{n}\right) \leqq \sum_{i} \int f_{i}^{2} d \mu+2 \sum_{i<j} \int f_{i} f_{j} d \mu \tag{2.6}
\end{equation*}
$$

where the sum involves only functions $f$ in $F_{n}$ such that $f$ vanishes outside $O$. Consequently $\mu^{n}\left(U_{n}\right) \leqq \mu(O)$, so the proof is complete.

## 3. Comment on the hypotheses

Assuming that $X$ is compact metric, the condition that $X$ be connected and locally connected is in fact necessary for the conclusion of Theorem 2.1 as the following slightly stronger statement shows.

Theorem 3.1. If for every continuous mapping $T: I \rightarrow P_{0}(X)$ there is a $\mu \in P(C(I, X))$ such that $t \mu=T(t)$ for all $t \in I$, then $X$ is connected and locally connected.

Proof. It is immediate that $X$ must be arcwise connected. Indeed if $x_{1}$ and $x_{2}$ are points of $X$, and for $0 \leqq t \leqq 1$ we set

$$
\begin{equation*}
T(t)=\frac{3}{4}\left\{(1-t) \varepsilon_{x_{1}}+t \varepsilon_{x_{2}}\right\}+\frac{1}{4} \theta \tag{3.1}
\end{equation*}
$$

where $\varepsilon_{x}$ denotes unit mass at the point $x$ and $\theta$ is any measure in $P_{0}(X)$, then $T$ is a continuous mapping of $I$ into $P_{0}(X)$. If $\mu$ is a measure on $C(I, X)$ with $t \mu=T(t)$, then the two sets $\left\{f: f(0)=x_{1}\right\}$ and $\left\{f: f(1)=x_{2}\right\}$ each have $\mu$ measure no less than $\frac{3}{4}$. Thus their intersection is nonempty and so there is a curve in $X$ joining $x_{1}$ and $x_{2}$.

As to the local connectedness, it is an exercise in general topology to show that a space $X$ is locally connected if it has the following property: given any point $x \in X$ and neighborhood $U$ of $x$ there is a neighborhood $V$ of $x$ such that for each point $y \in V$ there is a continuous function $f: I \rightarrow U$ with $f(0)=x$ and $f(1)=y$. Consequently if $X$ is not locally connected there is a point $x \in X$, a neighborhood $U$ of $x$ and a sequence $\left\{x_{n}\right\}$ of points approaching $x$ such that any continuous function from an interval $[a, b]$ to $X$ having $x$ and some $x_{n}$ in its range must also take on values outside of $U$. In particular if

$$
\begin{equation*}
A_{n}=\left\{f \in C(I, X): f(0)=x, f\left(\frac{1}{n}\right)=x_{n}\right\} \tag{3.2}
\end{equation*}
$$

then $\cap_{k=1}^{\infty} \cup_{n=k}^{\infty} A_{n}$ is the empty set. To define $T(t)$ pick a measure $\theta \in P_{0}(X)$, define $T(0)$ to be $\frac{3}{4} \varepsilon_{x}+\frac{1}{4} \theta$ and

$$
\begin{equation*}
T(t)=\frac{3}{4} n(n-1)\left\{\left(\frac{1}{n-1}-t\right) \varepsilon_{x_{n}}+\left(t-\frac{1}{n}\right) \varepsilon_{x_{n-1}}\right\}+\frac{1}{4} \theta \tag{3.3}
\end{equation*}
$$

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if $1 / n \leqq t \leqq 1 /(n-1)$. Then $T$ maps $I$ continuously into $P_{0}(X)$. If $\mu$ is any probability measure on $C(I, X)$ then $\mu\left(\cup_{n=k}^{\infty} A_{n}\right) \rightarrow 0$ as $k \rightarrow \infty$. But if $\mu$ is such that $t \mu=T(t)$ for all $t$, then $\mu\left(A_{k}\right) \geqq \frac{1}{2}$ for all $k$. Thus there can be no such $\mu$ and the proof is complete.

## REFERENCE

[1] R. M. Blumenthal and H. H. Corson, "On continuous collections of measures," Ann. Inst. Fourier (Grenoble), Vol. 20 (1970), pp. 193-199.

