# Gradient Approximation of Vector Fields* 

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Communicated by R. C. Buck


#### Abstract

In this paper, we consider an approximation problem which arose in optimal control theory. We seek conditions on a compact subset $K$ of Euclidean $n$-space such that every continuous vector field on $K$ may be uniformly approximated (on $K$ ) by vector fields with strictly positive integrating factors. We prove that such approximation is possible for all $K$ in a particular subclass of the compact sets with topological dimension less than or equal to 1 .


## Introduction

The following problem arose in connection with a problem of optimal control in the Lagrange form, where the results of this paper were used to prove existence theorems (see [2]). Let $K$ be a compact subset of a Euclidean space $E^{n}$, and let $C(K)^{n}$ be the Banach space of continuous vector fields $b(x)=\left(b_{1}, \ldots, b_{n}\right)$ (with $n$ real components) defined on $K$ with

$$
\|b\|=\sup _{x \in K}|b(x)|=\sup _{x \in K}\left[b_{1}(x)^{2}+\cdots+b_{n}(x)^{2}\right]^{1 / 2}
$$

We let $\tilde{F}(K)$ be the set of vector fields $g$, defined in an open neighborhood of $K$, which are of the form $g(x)=c(x) \nabla G(x)$ for some continuous, strictly positive function $c(x)$ and some continuously differentiable function $G(x)$. Here $\nabla G(x)$ is the gradient of $G$ at $x$. Thus, the set $\tilde{F}(K)$ consists of vector fields defined in a neighborhood of $K$ which have strictly positive integrating factors. The set of vector functions $\tilde{F}(K)$ induces a subset of $C(K)^{n}$ by restriction, and we denote by $F(K)$ the norm closure of this set in $C(K)^{n}$. We note that any element $g$ of $\tilde{F}(K)$ is gradient-like, for, if $g(x)=c(x) \nabla G(x)$, then the inner product, $g(x) \cdot \nabla G(x)$, is strictly positive for all $x$ such that $g(x) \neq 0$. Thus, we say that the elements of $F(K)$ are weakly gradient-like. We now pose the problem.

Given a vector field $b$ in $C(K)^{n}$, under what conditions on $b($ and $K)$ is $b$

[^0]weakly gradient-like? As it stands, this problem has been too general to admit satisfactory solution. Rather, a restricted problem has proved more susceptible to attack. Namely, for what compact sets $K$ does $F(K)=C(K)^{n}$ ? This is the approach we take in this paper.

In Section 1, we define a subset $G(K)$ of $F(K)$ which is useful in the approximation problem under consideration, and we discuss the relations between $G(K), F(K)$ and $C(K)^{n}$. Section 2 is devoted to the case $n=2$, i.e., we seek conditions on compact subsets $K$ of the plane so that $F(K)=C(K)^{2}$. In Section 3, we review some definitions from topological dimension theory, and we use these concepts in our study of the case $n>2$.

## 1

We define the subset $\tilde{G}(K)$ of $\tilde{F}(K)$ to be the set of all continuous vector functions $g$, defined in an open neighborhood of $K$, which are of the form $g(x)=\nabla G(x)$ for some continuously differentiable function $G(x)$. The set $\widetilde{G}(K)$ defines a subset of $F(K)$ by restriction to $K$, and we define $G(K)$ to be the norm closure of this set. Thus, $G(K) \subseteq F(K) \subseteq C(K)^{n}$ for any compact set $K$ in $E^{n}$. The advantage in working with the set $G(K)$ is that it is a linear subspace of $C(K)^{n}$ (since $\tilde{G}(K)$ is closed under multiplication by scalars and under addition), while, in general, $F(K)$ is not closed under addition. In fact, we have the following.

Theorem. If $K$ is such that $F(K)$ is a linear subspace of $C(K)^{n}$, then $F(K)=C(K)^{n}$.

Proof. Let $\nu=\left(\nu_{1}, \ldots, \nu_{n}\right)$ be a Borel measure on $K$ with $n$ real components so that if $g=\left(g_{1}, \ldots, g_{n}\right)$ is an element of $F(K)$, then $\int_{K} g \cdot d \nu=$ $\sum_{i=1}^{n} \int_{K} g_{i} \cdot d \nu_{i}=0$. That is, $d \nu$ annihilates the subspace $F(K)$. By definition, for any function $c(x)$, continuous and strictly positive in a neighborhood of $K$, and any function $G$, continuously differentiable in a neighborhood of $K$, $c \nabla G \in F(K)$. We consider the function $G$ to be fixed (but arbitrary). Because $F(K)$ is closed, for any continuous, nonnegative function $c(x)$ defined on $K$, $c(x) \nabla G(x)$ lies in $F(K)$. Let $f(x)$ be any continuous, real-valued function on $K$. Let $f^{+}(x)=\max \{0, f(x)\}$, and let $f^{-}(x)=\min \{0, f(x)\}$. Then $f^{+}(x) \geqslant 0$ and $f^{-}(x) \leqslant 0$ for all $x$ in $K$. Furthermore, $f(x)=f^{+}(x)+f^{-}(x)$ so that

$$
f(x) \nabla G(x)=f^{+}(x) \nabla G(x)+f^{-}(x) \nabla G(x)
$$

for all $x$ in $K$. Therefore, if $H(x)=-G(x)$, then for all $x$ in $K$,

$$
f(x) \nabla G(x)=f^{+}(x) \nabla G(x)+\left[-f^{-}(x)\right] \nabla H(x)
$$

Since each element of the sum is in $F(K)$, and $F(K)$ is assumed to be closed under addition, $f(x) \nabla G(x)$ is in $F(K)$. Therefore,

$$
\int_{K} f(x) \nabla G(x) \cdot d \nu(x)=0,
$$

for all continuous, real-valued functions $f$ defined on $K$. We let $d \mu$ be the real measure $\nabla G \cdot d \nu=\sum_{i=1}^{n}\left(\partial G / \partial x_{i}\right) d \nu_{i}$. Then, we have shown that $\int_{K} f(x) d \mu(x)=0$ for all continuous, real-valued functions $f$ on $K$. It follows that $d \mu=\nabla G \cdot d \nu$ is the zero measure for all continuously differentiable real-valued functions $G$. Let $b=\left(b_{1}, \ldots, b_{n}\right)$ be an arbitrary fixed vector in $E^{n}$ and define $G(x)=b \cdot x=b_{1} x_{1}+\cdots+b_{n} x_{n}$ for all $x$ in $E^{n}$. Then, $\nabla G(x)=b$ for all $x$, and $b \cdot d \nu$ is the zero measure. That is,

$$
0=\int_{K^{\prime}} b \cdot d \nu=b \cdot v\left(K^{\prime}\right)
$$

for all Borel measurable subsets $K^{\prime}$ of $K$. Since $b$ is arbitrary, $\nu\left(K^{\prime}\right)=0$ for all Borel measurable subsets $K^{\prime}$ of $K$. Therefore, $\nu$ is the zero measure. Since the only Borel measure on $K$ which annihilates $F(K)$ is the zero measure, and $F(K)$ is a linear subspace of $C(K)^{n}, F(K)$ must equal $C(K)^{n}$, and the theorem is proved.

This result is not as useful as it is interesting, since demonstrating that $F(K)$ is linear seems to be as difficult as verifying that $F(K)=C(K)^{n}$ by more direct means.

We conclude this section with an example of a compact set $K_{1}$ in $E^{2}$ for which $G\left(K_{1}\right) \neq F\left(K_{1}\right) \neq C\left(K_{1}\right)^{2}$, and an example of a compact set $K_{2}$ in $E^{n}$ for which $G\left(K_{2}\right)=F\left(K_{2}\right)=C\left(K_{2}\right)^{n}$.

Example 1.1. We denote by $(x, y)$ the points in $E^{2}$ and we define $K$ to be the set of all $(x, y)$ in $E^{2}$ such that $x^{2}+y^{2}=1$. For $(x, y)$ in $K_{1}$, we define $g(x, y)=(-y, x)$ so that $g \in C(K)^{2}$. We will show that $g$ is not an element of $G\left(K_{1}\right)$ or $F\left(K_{1}\right)$. We define the bounded linear functional $T$ on elements $h=\left(h_{1}, h_{2}\right)$ of $C\left(K_{1}\right)^{2}$ by setting

$$
T(h)=\int_{0}^{2 \pi}\left[-h_{1}(\cos t, \sin t) \sin t+h_{2}(\cos t, \sin t) \cos t\right] d t,
$$

so that $T(h)$ is the counter-clockwise path integral of $h$ along $K_{1}$. Thus, $T(g)=2 \pi$. But $T$ annihilates the elements of $G\left(K_{1}\right)$, so $g \notin G\left(K_{1}\right)$. Now, suppose $g$ is an element of $F\left(K_{1}\right)$. Then, for each $n=1,2, \ldots$, there is a continuous, strictly positive function $c_{n}(x, y)$ and a continuously differentiable function $G_{n}(x, y)$, both defined in a neighborhood of $K_{1}$, such that

$$
\left|g(x, y)-c_{n}(x, y) \nabla G_{n}(x, y)\right|<1 / n
$$

for all $(x, y)$ in $K_{1}$. Since $|g(x, y)|=1$ for all $(x, y)$ in $K_{1}$, for $n \geqslant 2$, $\nabla G_{n}(x, y) \neq(0,0)$ for all $(x, y)$ in $K_{1}$. In particular, we look at $n=2$. Let $\gamma=\max c_{2}(x, y)$, the maximum taken for $(x, y)$ in $K_{1}$. Then, by the Schwarz inequality,

$$
\left|c_{2}(x, y) \nabla G_{2}(x, y) \cdot g(x, y)-|g(x, y)|^{2}\right|<\frac{1}{2}
$$

for all $(x, y)$ in $K_{1}$. Since $|g(x, y)|^{2}=1$, we have that

$$
c_{2}(x, y) \nabla G_{2}(x, y) \cdot g(x, y)>\frac{1}{2}
$$

for all $(x, y)$. Therefore,

$$
\nabla G_{2}(x, y) \cdot g(x, y)>\left(1 / 2 c_{2}(x, y)\right) \geqslant 1 / 2 \gamma
$$

for all $(x, y)$ in $K_{1}$. We now note that

$$
T\left(\nabla G_{2}\right)=\int_{0}^{2 \pi} \nabla G_{2}(\cos t, \sin t) \cdot g(\cos t, \sin t) d t \geqslant \pi / \gamma>0
$$

which contradicts the fact that $T$ annihilates elements of $G\left(K_{1}\right)$. Therefore, the open ball of radius $\frac{1}{2}$ in $C\left(K_{1}\right)^{2}$, centered at $g$, contains no elements of $F\left(K_{1}\right)$, so $F\left(K_{1}\right) \neq C\left(K_{1}\right)^{2}$. From the theorem proved above, it is clear that $G\left(K_{1}\right) \neq F\left(K_{1}\right)$, since $F\left(K_{1}\right)$ cannot be a linear subspace of $C\left(K_{1}\right)^{2}$. Thus, $G\left(K_{1}\right) \neq F\left(K_{1}\right) \neq C\left(K_{1}\right)^{2}$.

Example 1.2. Let $K_{2}$ be the compact subset of $E^{n}$ defined by

$$
K_{2}=\left\{\left(x_{1}, x_{2}, \ldots, x_{n}\right) \mid 0 \leqslant x_{1} \leqslant 1, x_{2}=\cdots=x_{n}=0\right\}
$$

We will show that $G\left(K_{2}\right)=C\left(K_{2}\right)^{n}$, and hence that $F\left(K_{2}\right)=C\left(K_{2}\right)^{n}$. Let $g\left(x_{1}, \ldots, x_{n}\right)=\left(g_{1}\left(x_{1}, \ldots, x_{n}\right), \ldots, g_{n}\left(x_{1}, \ldots, x_{n}\right)\right)$ be an element of $C\left(K_{2}\right)^{n}$ such that each function $g_{2}, \ldots, g_{n}$ is continuously differentiable. Such functions are dense in $C\left(K_{2}\right)^{n}$. Actually, each function $g_{i}$ depends only on the variable $x_{1}$, so we write the function $g$ as $g\left(x_{1}\right)=\left(g_{1}\left(x_{1}\right), \ldots, g_{n}\left(x_{1}\right)\right)$. Since $g_{1}$ is continuous, it is Riemann integrable and we may define $G_{1}\left(x_{1}\right)=\int_{0}^{x_{1}} g_{1}(t) d t$. We define $H\left(x_{1}, \ldots, x_{n}\right)$ by setting

$$
H\left(x_{1}, \ldots, x_{n}\right)=G_{1}\left(x_{1}\right)+\sum_{i=2}^{n} x_{i} g_{i}\left(x_{1}\right)
$$

Then, $H$ has continuous partial derivatives and

$$
H_{x_{i}}\left(x_{1}, 0, \ldots, 0\right)=g_{i}\left(x_{1}\right)
$$

for $i=1,2, \ldots, n$ and for $0 \leqslant x_{1} \leqslant 1$. Thus, each such element $g$ of $C\left(K_{2}\right)^{n}$ is an element of $G\left(K_{2}\right)$. Since the functions with continuous derivatives are dense in $C\left(K_{2}\right)^{n}, G\left(K_{2}\right)=C\left(K_{2}\right)^{n}$.

The difference between the sets $K_{1}$ and $K_{2}$ illustrates the underlying theme of this paper. We will consider (one dimensional) compact sets $K$ in $E^{n}$ which, in an appropriate sense, have the topological properties of $K_{2}$ and avoid those of $K_{1}$.

## 2

In this section, we study the case $n=2$. That is, we investigate conditions on a compact set $K$ in $E^{2}$ such that $G(K)=C(K)^{2}$. For the purposes of this section only, we regard $E^{2}$ as the complex plane: $E^{2}=\{x+i y \mid x$ and $y$ lie in $\left.E^{1}\right\}$, where $i^{2}=-1$. For $K$ a compact subset of $E^{2}$, we define $\mathrm{C}(K)$ to be the space of continuous, complex-valued functions defined on $K . \mathrm{C}(K)$ is a Banach space under the norm defined for $f=f_{1}+i f_{2}$ by

$$
\|f\|_{\infty}=\sup _{z \in K}|f(z)|=\sup _{z \in K}\left[f_{1}(z)^{2}+f_{2}(z)^{2}\right]^{1 / 2}
$$

The subspace $\mathbf{P}(K)$ of $\mathbf{C}(K)$ is defined to be the norm closure in $\mathbf{C}(K)$ of the polynomials on $K$. That is, $\mathbf{P}(K)$ is generated by functions $p(z)$ which are polynomials in the complex variable $z$. We will need the following theorem, which is proved in ([1], p. 48).

Theorem 2.1 (Mergelyan's Theorem). Let $K$ be a compact subset of $E^{2}$ whose complement is connected. If $f \in \mathbf{C}(K)$ and $f$ is analytic on the interior of $K$, then $f \in \mathbf{P}(K)$.

It is clear that $C(K)^{2}$ and $\mathbf{C}(K)$ are isometrically isomorphic. We will use this fact in conjunction with Theorem 2.1 to prove the main theorem of this section, Theorem 2.2.

Theorem 2.2. If $K$ is a compact subset of $E^{2}$ which is nowhere dense and has connected complement in $E^{2}$, then $G(K)=C(K)^{2}$.

Proof. For $f=\left(f_{1}, f_{2}\right)$ in $C(K)^{2}$, we define $U(f)=f_{1}-i f_{2}$ so that $U(f) \in \mathbb{C}(K)$. It is clear that $U$ is an isometry from $C(K)^{2}$ onto $\mathrm{C}(K)$. Let $p(z)$ be a polynomial in the complex variable $z=x+i y$, defined in the whole complex plane, and hence in a neighborhood of $K$. Then, there are real polynomials $p_{1}$ and $p_{2}$ in the two real variables $x$ and $y$ such that $p(x+i y)=$ $p_{1}(x, y)+i p_{2}(x, y)$. By the Cauchy-Riemann equations,

$$
\partial p_{1} / \partial y=-\partial p_{2} / \partial x
$$

for all real $x, y$. Thus, by Green's theorem, the differential $p_{1} d x-p_{2} d y$ is exact on the entire plane. It follows that there is a continuously differentiable function $P=P(x, y)$ such that $\partial P / \partial x=p_{1}$ and $\partial P / \partial y=-p_{2}$. We have therefore shown that $U^{-1}(p) \in G(K)$ for every complex polynomial $p$. Since
$U^{-1}$ is continuous and maps a dense subset of $\mathbf{P}(K)$ into $G(K)$, we conclude that $U^{-1}(\mathbf{P}(K)) \subseteq G(K)$. We now use our assumptions about $K$. Since $K$ is nowhere dense, $K$ has empty interior, and every function in $\mathbf{C}(K)$ is analytic on the (empty) interior of $K$. By Theorem 2.1, then, since $K$ has connected complement, if $f \in \mathbf{C}(K)$, we have that $f \in \mathbf{P}(K)$, i.e., $\mathbf{P}(K)=\mathbf{C}(K)$. Therefore, under our hypotheses on $K, U^{-1}(\mathbf{C}(K)) \subseteq G(K)$. But, $U$ is an isometry of $C(K)^{2}$ onto $\mathbf{C}(K)$, so $U^{-1}(\mathbf{C}(K))=C(K)^{2}$. Thus, $C(K)^{2} \subseteq G(K) \subseteq C(K)^{2}$, and the theorem follows.

In the next section, we attempt to imitate these results when $K$ is a compact subset of $E^{n}, n>2$. However, neither of the hypotheses stated in Theorem 2.2 on the subset $K$ of $E^{2}$ generalizes verbatim to the case $n>2$. If $K$ is a compact subset of $E^{n}$ which is nowhere dense in $E^{n}$, then it has no interior in $E^{n}$, but for $n>2, K$ could have "dimension" larger than or equal to 2 . For example, if $K=\left\{(x, y, z) \in E^{3} \mid x^{2}+y^{2} \leqslant 1, z=0\right\}$, then $K$ is nowhere dense in $E^{3}$, but, as in Example 1.1, it may be shown that $G(K) \neq F(K) \neq C(K)^{3}$. Similarly, the "connected complement" condition is not sufficient for $n>2$. For example, the set $K=\left\{(x, y, z) \in E^{3} \mid x^{2}+y^{2}=1, z=0\right\}$ has connected complement in $E^{3}$, but, as in Example 1.1, $G(K) \neq F(K) \neq C(K)^{3}$. Thus, in Section 3, we consider the appropriate topological generalizations for $n>2$.

## 3

We will need some of the concepts of topological dimension theory, and for these, we draw on [3]. Let $K$ be a compact subset of $E^{n}$ and let $\left\{U_{1}, \ldots, U_{k}\right\}$ be a covering of $K$ by open sets. The order of the covering $\left\{U_{1}, \ldots, U_{k}\right\}$ is defined to be the largest integer $N$ such that there are $N+1$ members of the covering whose intersection is nonempty. If $\left\{V_{1}, \ldots, V_{m}\right\}$ is a covering of $K$ by open sets, we say that $\left\{V_{1}, \ldots, V_{m}\right\}$ is a refinement of $\left\{U_{1}, \ldots, U_{k}\right\}$ if for each $V_{j}, j=1,2, \ldots, m$, there is an $i, 1 \leqslant i \leqslant k$, such that $V_{j} \subseteq U_{i}$. (See [3, pp. 52-53].) We may define the dimension of $K$, $\operatorname{dim} K$, as follows. We will say that $\operatorname{dim} K \leqslant N$ if every covering of $K$ by finitely many open sets has a refinement of order less than or equal to $N$. In particular, the empty set is the only set with dimension -1 . Further, if $\operatorname{dim} K=0$, every covering of $K$ by finitely many open sets has a refinement whose elements are pairwise disjoint. We thus have the following theorem.

Theorem 3.1. If $K$ is a nonempty compact subset of $E^{n}$ such that $\operatorname{dim} K=0$, then $G(K)=C(K)^{n}$.

Proof. Let $g$ be an arbitrary element of $C(K)^{n}$ and let $\epsilon>0$ be arbitrary. For each $x_{0}$ in $K$, there is an open neighborhood $U=U\left(x_{0}\right)$ of $x_{0}$ such that

$$
\left|g(x)-g\left(x_{0}\right)\right|<\epsilon / 2
$$

for all $x$ in $K \cap U\left(x_{0}\right)$. Since $\{U(x)\}_{x \in K}$ is a covering of $K$ by open sets, the family $\{U(x)\}_{x \in K}$ has a Lebesgue number $\delta$. Thus, if $V$ is an open subset of $E^{n}$ of diameter less than $\delta$ and if $V \cap K$ is nonempty, then $V \subseteq U(x)$ for some $x$ in $K$. Since $K$ is compact, $K$ is totally bounded, so there is a covering of $K$ by finitely many open sets $\left\{V_{1}, \ldots, V_{m}\right\}$ such that each $V_{i}, 1 \leqslant i \leqslant m$, has diameter less than $\delta$. Since $\operatorname{dim} K=0$, there is a refinement $\left\{U_{1}, \ldots, U_{N}\right\}$ of $\left\{V_{1}, \ldots, V_{m}\right\}$ such that $U_{i} \cap U_{j}=\phi$ if $i \neq j$. By the definition of the family $\left\{U_{1}, \ldots, U_{N}\right\}$ the diameter of $U_{i}$ is less than $\delta$ for each $i$. For $i=1,2, \ldots, N$, let $x_{i}$ be any element of $K \cap U_{i}$. Then, $U_{i} \subseteq U\left(x^{\prime}\right)$ for some $x^{\prime}$ in $K$, so that, for any $x$ in $K \cap U_{i}$,

$$
\left|g(x)-g\left(x_{i}\right)\right| \leqslant\left|g(x)-g\left(x^{\prime}\right)\right|+\left|g\left(x^{\prime}\right)-g\left(x_{i}\right)\right|<\epsilon
$$

We let $b_{i}=g\left(x_{i}\right) \in E^{n}$ for $i=1,2, \ldots, N$, and we define a real-valued function $G$ on $U=\bigcup_{i=1}^{N} U_{i}$ as follows. Let $G(x)=b_{i} \cdot x$ for all $x$ in $U_{i}$. Since the sets $\left\{U_{1}, \ldots, U_{N}\right\}$ are pairwise disjoint, $G$ is well-defined. Also, $\nabla G(x)=b_{i}$ for all $x$ in $U_{i}$. It follows that $G$ is continuously differentiable on $U$ and, for $x$ in $K \cap U_{i}$,

$$
|g(x)-\nabla G(x)|=\left|g(x)-b_{i}\right|<\epsilon
$$

Therefore, $|g(x)-\nabla G(x)|<\epsilon$ for all $x$ in $K$. Since $\epsilon>0$ was arbitrary, $g \in G(K)$. Since $g \in C(K)^{n}$ was arbitrary, $G(K)=C(K)^{n}$ and Theorem 3.1 is proved.

Our main result of the section concerns a subclass of the compact sets $K$ satisfying $\operatorname{dim} K \leqslant 1$.

We note that in Section 2, we have already proved a theorem which can be used to derive a statement of this type. For, by [3, p. 41, Theorem IV 1], $\operatorname{dim} E^{n}=n$ for all $n$, and, if $K \subseteq E^{n}$, then, by [3, p. 44, Theorem IV 3], $\operatorname{dim} K=n$ if and only if $K$ contains a nonempty subset which is open in $E^{n}$. In particular, for $n=2, K \subseteq E^{2}$ is nowhere dense in $E^{2}$ if and only if $\operatorname{dim} K \leqslant 1$. Thus, we may restate Theorem 2.2 to say: if $K$ is a compact subset of $E^{2}$ such that $\operatorname{dim} K \leqslant 1$ and $K$ has connected complement in $E^{2}$, then $G(K)=C(K)^{2}$. We will prove a similar, though weaker statement for $K \subseteq E^{n}, n>2$.

We now state some definitions. We will assume that $B_{1}, \ldots, B_{k}$ are open balls in the Euclidean space $E_{n}$ with centers $x_{1}, \ldots, x_{k}$, respectively. If $B_{i} \cap B_{j} \neq \phi$ for $i \neq j$, we denote by $L_{i j}$ the undirected line segment joining $x_{i}$ and $x_{j}$. Thus, $L_{i j}=L_{j i}$. If $B_{i} \cap B_{j}$ is empty, we let $L_{i j}=L_{j i}=\phi$, and we set $L_{i i}=\left\{x_{i}\right\}$ for each $i$. We define a point set $L\left(B_{1}, \ldots, B_{k}\right)$ by setting

$$
L\left(B_{\mathbf{1}}, \ldots, B_{k}\right)=\bigcup_{i, j=1}^{k} L_{i j}
$$

We note that $L\left(B_{1}, \ldots, B_{k}\right)$ is a finite union of piecewise linear curves and that $L\left(B_{1}, \ldots, B_{k}\right)$ has the same number of connected components as $\bigcup_{i=1}^{k} B_{i}$.

We will say that the collection $\left\{B_{1}, \ldots, B_{k}\right\}$ forms a simple chain if $L\left(B_{1}, \ldots, B_{k}\right)$ contains no simple closed curve. Let $K$ be a compact subset of $E^{n}$. We will say that $K$ has the simple chain covering property if for every $\epsilon>0$, there is a covering of $K$ by open balls $B_{1}, \ldots, B_{k}$ in $E^{n}$, each of radius less than $\epsilon$, such that the collection $\left\{B_{1}, \ldots, B_{k}\right\}$ forms a simple chain. We note that if $K$ has the simple chain covering property, then $\operatorname{dim} K \leqslant 1$, since every covering by open balls $\left\{B_{1}, \ldots, B_{k}\right\}$ which form a simple chain, has order less than or equal to 1 . Furthermore, if a set $K$ has the simple chain covering property, then clearly $K$ contains no simple closed curve. We may prove the following theorem.

Theorem 3.2. If $K$ is a compact subset of $E^{n}$ having the simple chain covering property, then $F(K)=C(K)^{n}$.

Proof. Let $g=\left(g^{1}, g^{2}, \ldots, g^{n}\right)$ be an element of $C(K)^{n}$ and let $\epsilon>0$ be given. For each $x_{0}$ in $K$, there is an open $n$-ball $B\left(x_{0}\right)$ centered at $x_{0}$ such that $\left|g\left(x_{0}\right)-g(x)\right|<\epsilon / 4$ for $x$ in $K \cap B\left(x_{0}\right)$. The open balls $B(x)$, for $x$ in $K$, cover the set $K$. Since $K$ is compact, the covering $\{B(x)\}_{x \in K}$ has a Lebesgue number $\delta$. Let $B_{1}, \ldots, B_{k}$ be a covering of $K$ by open balls in $E^{n}$ of radius less than $\delta / 2$, such that $\left\{B_{1}, \ldots, B_{k}\right\}$ forms a simple chain. Thus, for each $B_{i}$, there is an element $y_{i}$ of $K$ such that $B_{i} \subseteq B\left(y_{i}\right)$. Let $b_{i}=g\left(y_{i}\right)$ and write $b_{i}=\left(b_{i}{ }^{1}, b_{i}{ }^{2}, \ldots, b_{i}{ }^{n}\right)$. We let $r_{i}$ be the radius of $B_{i}$ and we write $c_{i}=\left(c_{i}{ }^{1}, c_{i}{ }^{2}, \ldots, c_{i}{ }^{n}\right)$ for the center of $B_{i}, i=1,2, \ldots, k$. We now note that since $K$ is compact, there is a positive distance $\delta^{\prime}$ between $K$ and $\partial\left(\bigcup_{i=1}^{k} B_{i}\right)$. In particular, since $\left\{B_{1}, \ldots, B_{k}\right\}$ forms a simple chain, $\bigcup_{i \neq j}\left(\partial B_{i} \cap \partial B_{j}\right) \subseteq$ $\partial\left(\bigcup_{i=1}^{k} B_{i}\right)$, and if $W$ is a $\delta^{\prime} / 2$-neighborhood of $\bigcup_{i \neq j}\left(\partial B_{i} \cap \partial B_{j}\right)$, then $K \cap c l W=\phi$. We define $B$ to be the set $\bigcup_{i=1}^{k} B_{i} \mid c l W$, so that $B$ is an open neighborhood of $K$ in $E^{n}$. We will define a vector function $f=\left(f^{1}, f^{2}, \ldots, f^{n}\right)$ on $B$ such that $|g(x)-f(x)|<\epsilon$ for all $x$ in $K$, and such that $f \in F(K)$. That is, we will find a continuous, strictly positive function $h(x)$ on $B$ and a real-valued, continuously differentiable function $G(x)$ on $B$ such that $f(x)=h(x) \nabla G(x)$ for all $x$ in $B$. The vector function $f$ will be defined inductively. In fact, since it is to be constant in each of the sets $\left.B_{j}\right)\left(\bigcup_{i \neq j} B_{i}\right)$, the induction step will be obvious when the construction is verified in the sets $B_{i} \cap B_{j}$, for $i \neq j$. Since the family $\left\{B_{1}, \ldots, B_{k}\right\}$ forms a simple chain and hence has order less than or equal to 1 , we need never consider the intersections of 3 or more distinct elements of the family.

For $x$ in $B_{j} \backslash\left(\bigcup_{i \neq j} B_{i} \cup c l W\right)$, we define $f(x)=b_{j}=g\left(y_{j}\right)$. Here, without loss of generality, we assume that if $B_{i} \cap B_{j} \neq \phi, i \neq j$, then $b_{j}$ is not perpendicular to the line connecting the centers $c_{i}$ and $c_{j}$. (If this not the
case, we may perturb $b_{j}$ by adding a vector of length less than $\epsilon / 4$ so that the above is satisfied. The remaining statements then proceed without difficulty if we consider that $\left|b_{j}-g\left(y_{j}\right)\right|<\epsilon / 4$ for each $j$.) Thus, for each $x$ in $B_{i} \backslash\left(\cup_{i \neq j} B_{i} \cup c l W\right)$,

$$
|f(x)-g(x)| \leqslant\left|b_{j}-g(x)\right| \leqslant\left|b_{j}-g\left(y_{j}\right)\right|+\left|g\left(y_{j}\right)-g(x)\right|<\epsilon / 2 .
$$

Now, suppose that $f(x)$ is extended to all of $B$ in such a way that for $x$ lying in $B_{i} \cap B_{j}$, for $i \neq j$, then $f(x)$ is a convex combination of $b_{i}$ and $b_{j}$. We then have the following. For $x$ in $B_{i} \cap B_{i} \cap K$, there is a $\lambda, 0 \leqslant \lambda \leqslant 1$, such that $f(x)=\lambda b_{i}+(1-\lambda) b_{j}$. Thus,

$$
\begin{aligned}
|f(x)-g(x)| \leqslant & \lambda\left|b_{i}-g\left(y_{i}\right)\right|+\lambda\left|g\left(y_{i}\right)-g(x)\right| \\
& +(1-\lambda)\left|b_{j}-g\left(y_{j}\right)\right|+(1-\lambda)\left|g\left(y_{j}\right)-g(x)\right|,
\end{aligned}
$$

which is less than $\epsilon$ since each absolute value is bounded by $\epsilon / 4$. We will therefore show how to define $f$ in the sets $B_{i} \cap B_{j}, i \neq j$, so that $f$ is a convex combination of $b_{i}$ and $b_{j}$.
We assume, therefore, that $B_{i} \cap B_{j}$ is non-empty, for some $i \neq j$. Without loss of generality, we may assume that the centers $c_{i}$ and $c_{j}$ of $B_{i}$ and $B_{j}$ both lie along the $x_{1}$-axis, and, in fact, we may assume that $c_{i}=0$ while $c_{j}=(a, 0,0, \ldots, 0)$ for some $a>0$. Thus, $\partial B_{i} \cap \partial B_{j}$ lies in a hyperplane perpendicular to the $x_{1}$-axis, say $\left\{\left(x_{1}, \ldots, x_{n}\right) \mid x_{1}=a^{\prime}\right\}$, where $0<a^{\prime}<a$. Let $a_{1}, a_{2}$ be positive numbers such that $a_{1}<a^{\prime}<a_{2}$ and such that: the intersection of the hyperplane $x_{1}=a_{1}$ and the set $\partial B_{j}$ lies in the open set $W$; the intersection of the hyperplane $x_{1}=a_{2}$ and the set $\partial B_{i}$ lies in the open set $W$.

We note that neither of the vectors $b_{i}$ and $b_{j}$ are perpendicular to the $x_{1}$ axis (by assumption). If $b_{i}=b_{j}$, we may set $h(x)=1$ and $f(x)=b_{i}=b_{j}$ for all $x$ in $B_{i} \cup B_{j}$. We therefore assume that $b_{i} \neq b_{j}$.

Let $\bar{b}_{i}$ and $\bar{b}_{i}$ be the vectors in $E^{n+1}$ formed from $b_{i}, b_{j}$ by setting $\bar{b}_{i}=\left(b_{i}, b\right), \bar{b}_{j}=\left(b_{j}, b\right)$, where $b$ is chosen to be any positive number large enough so that $\bar{b}_{i} \cdot \bar{b}_{j}=(b)^{2}+\sum_{k=1}^{n} b_{i}{ }^{k} b_{j}{ }^{k}>0$. Since $b_{i} \neq b_{j}, \bar{b}_{i}$ and $\bar{b}_{j}$ are linearly independent. We define row vectors $v_{1}, v_{2}, \ldots, v_{n+1}$ in $E^{n+1}$ as follows. We let $v_{3}, v_{4}, \ldots, v_{n+1}$ be $n-1$ row vectors in $E^{n+1}$ such that $\bar{b}_{i}, \bar{b}_{j}, v_{3}, \ldots, v_{n+1}$ are linearly independent and such that $v_{k}$ is perpendicular to $\bar{b}_{i}$ and $\bar{b}_{j}$ for all $k, 3 \leqslant k \leqslant n+1$. We let $v_{1}=(1,0,0, \ldots, 0)$ and we note that since $v_{1} \cdot \bar{b}_{i} \neq 0$ and $v_{1} \cdot \bar{b}_{j} \neq 0$, the family $v_{1}, \bar{b}_{i}, v_{3}, \ldots, v_{n+1}$ is linearly independent. We define $v_{2}=\bar{b}_{j}$. Thus, $\left\{v_{1}, v_{2}, \ldots, v_{n+1}\right\}$ forms a linearly independent set of row vectors. Let $M$ be the $(n+1) \times(n+1)$ matrix whose $k$ th row is $v_{k}$.
For $d$ in $E^{1}$, let $L=\left\{\bar{x}=\left(x_{1}, \ldots, x_{n+1}\right) \in E^{n+1} \mid x_{1}=d\right\}$. Then, we have defined $M$ so that $M\left(L_{0}\right)$ is a subset of $L_{0}$. (When we write $M(\bar{x})$ for some
vector $\bar{x}$ in $E^{n+1}$, we have in mind that $\bar{x}$ is a column vector.) Since the rows of $M$ are linearly independent, $M$ is invertible and $M\left(L_{0}\right)=L_{0}$. We write $L_{d}=(d, 0, \ldots, 0)+L_{0}$. Then, $M\left(L_{d}\right)=M(d, 0,0, \ldots, 0)+L_{0}$ so that $L_{d}$ is parallel to $M\left(L_{d}\right)$. Note, however, that $M(d, 0, \ldots, 0)=\left(d, d_{2}, \ldots, d_{n+1}\right)$ for some scalars $d_{2}, \ldots, d_{n+1}$ so that $M(d, 0, \ldots, 0)$ is an element of $(d, 0, \ldots, 0)+L_{0}$. Thus, $M\left(L_{d}\right) \cap L_{d}$ is non-empty. It follows that $M\left(L_{d}\right)=L_{d}$ for all $d$ in $E^{1}$.

We define

$$
p=M\left(\bar{b}_{i}\right)=\left(\bar{b}_{i} \cdot v_{1}, \bar{b}_{i} \cdot v_{2}, \bar{b}_{i} \cdot v_{3}, \ldots, \bar{b}_{i} \cdot v_{n+1}\right)=\left(p_{1}, p_{2}, 0,0, \ldots, 0\right)
$$

where we have constructed $M$ so that $p_{2}$ is positive. Similarly, we define $q=M\left(\bar{b}_{j}\right)=\left(q_{1}, q_{2}, 0,0, \ldots, 0\right)$, where $q_{2}$ is positive. Since $\bar{b}_{i}, \bar{b}_{j}$ are independent, so are $p$ and $q$. We will define a vector-valued function $\bar{f}$ on $E^{n+1}$ such that: $\bar{f}(\bar{x}) \in E^{n+1}$ for all $\bar{x} ; \bar{f}\left(x_{1}, \ldots, x_{n+1}\right)=p$ for $x_{1} \leqslant a_{1}$; $\bar{f}\left(x_{1}, \ldots, x_{n+1}\right)=q$ for $x_{1} \geqslant a_{2}$ and $\bar{f}\left(x_{1}, \ldots, x_{n+1}\right)$ is a convex combination of $p$ and $q$ for $a_{1}<x_{1}<a_{2}$. For $a_{1} \leqslant x_{1} \leqslant a_{2}$, we define

$$
\lambda\left(x_{1}\right)=\left(a_{2}-x_{1}\right) /\left\{\left(p_{2} / q_{2}\right)\left(x_{1}-a_{1}\right)+\left(a_{2}-x_{1}\right)\right\} .
$$

Then, $\lambda\left(a_{1}\right)=1, \lambda\left(a_{2}\right)=0$ and $0 \leqslant \lambda\left(x_{1}\right) \leqslant 1$ for all $x_{1}$ such that $a_{1} \leqslant x_{1} \leqslant a_{2}$. For $a_{1} \leqslant x_{1} \leqslant a_{2}$, we define $\bar{f}\left(x_{1}, \ldots, x_{n+1}\right)$ by setting

$$
\bar{f}\left(x_{1}, \ldots, x_{n+1}\right)=\lambda\left(x_{1}\right) p+\left(1-\lambda\left(x_{1}\right)\right) q .
$$

We define $\bar{f}$ by continuity and constancy outside the interval $a_{1} \leqslant x_{1} \leqslant a_{2}$. We now define a real-valued function $\bar{h}\left(x_{1}, \ldots, x_{n+1}\right)$ by setting

$$
\bar{h}\left(x_{1}, \ldots, x_{n+1}\right)=\left(a_{2}-a_{1}\right) /\left\{\left(p_{2} / q_{2}\right)\left(x_{1}-a_{1}\right)+\left(a_{2}-x_{1}\right)\right\}
$$

for $a_{1} \leqslant x_{1} \leqslant a_{2}$. We set $\bar{h}\left(x_{1}, \ldots, x_{n+1}\right)=1$ for $x_{1} \leqslant a_{1}$ and $\bar{h}\left(x_{1}, \ldots, x_{n+1}\right)=$ $q_{2} / p_{2}$ for $x_{1} \geqslant a_{2}$. Thus, $\bar{h}$ is continuous. Since $\bar{h}$ is monotone in the interval $a_{1} \leqslant x_{1} \leqslant a_{2}$, it is easy to see that $\bar{h}$ takes on values between 1 and $q_{2} / p_{2}$. Therefore, $\bar{h}\left(x_{1}, \ldots, x_{n+1}\right)>0$ for all $x_{1}, \ldots, x_{n+1}$.

We consider the quotient $\left(\bar{g}_{1}, \bar{g}_{2}, \ldots, \bar{g}_{n+1}\right)=\bar{f} / \bar{h}$. It is an easy computation to verify that

$$
\partial \bar{g}_{k} / \partial x_{m}=\partial \bar{g}_{m} / \partial x_{k}
$$

at all points $\left(x_{1}, \ldots, x_{n+1}\right)$ in $E^{n+1}$. Therefore, there is a real-valued, continuously differentiable function $\bar{G}$ such that $\nabla \bar{G}(\bar{x})=\bar{f}(\bar{x}) / \bar{h}(\bar{x})$, or, $\bar{f}(\bar{x})=\bar{h}(\bar{x}) \nabla \bar{G}(\bar{x})$ for all $\bar{x}$ in $E^{n+1}$.

For $\bar{x}$ in $E^{n+1}$, we define $G_{0}(\bar{x})=\bar{G}\left(M^{-1}(\bar{x})\right)$ and $h_{0}(\bar{x})=\bar{h}\left(M^{-1}(\bar{x})\right)$. Thus, $\nabla G_{0}(\bar{x})=M^{-1}\left[\nabla \bar{G}\left(M^{-1}(\bar{x})\right)\right]$. Since $M\left(L_{d}\right)=L_{d}$ implies $M^{-1}\left(L_{d}\right)=L_{d}$ for all real $d$, both $h_{0}\left(x_{1}, \ldots, x_{n+1}\right)$ and $\nabla G_{0}\left(x_{1}, \ldots, x_{n+1}\right)$ are constant for $x_{1} \leqslant a_{1}$ and for $x_{1} \geqslant a_{2}$ (as are $\bar{h}$ and $\left.\nabla \bar{G}\right)$. We will write $x=\left(x_{1}, \ldots, x_{n}\right)$ and $\bar{x}=\left(x, x_{n+1}\right)=\left(x_{1}, \ldots, x_{n+1}\right)$. For $x_{1} \leqslant a_{1}$,

$$
h_{0}(\bar{x}) \nabla G_{0}(\bar{x})=M^{-1}(p)=\bar{b}_{i}
$$

and for $x_{1} \geqslant a_{2}$,

$$
h_{0}(\bar{x}) \nabla G_{0}(\bar{x})=M^{-1}(q)=\bar{b}_{j} .
$$

Since, for $a_{1} \leqslant x_{1} \leqslant a_{2}, \bar{h}(\bar{x}) \nabla \bar{G}(\bar{x})$ is a convex combination of $p$ and $q$, in the same interval, $h_{0}(\bar{x}) \nabla G_{0}(\bar{x})$ is a convex combination of $\bar{b}_{i}$ and $\bar{b}_{j}$. Therefore, the $(n+1)$ st component of $h_{0} \nabla G_{0}$ is constant and equal to $b$. We may define a function $G$ of $n$ variables, then, by setting $G(x)=G_{0}(x, 0)$, and, it is clear that if $h(x)=h_{0}(x, 0)$, then $h(x) \nabla G(x)$ is a convex combination of $b_{i}$ and $b_{j}$ for $a_{1} \leqslant x_{1} \leqslant a_{2}$. Also, $h(x) \nabla G(x)=b_{i}$ for $x_{1} \leqslant a_{1}$ and $h(x) \nabla G(x)=b_{j}$ for $x_{1} \geqslant a_{2}$. Finally, both $h$ and $\nabla G$ are constant for $x_{1} \leqslant a_{1}$, and for $x_{1} \geqslant a_{2}$.

For $x$ in $\left(B_{i} \cap B_{j}\right) \backslash c l W$, we define $f(x)=h(x) \nabla G(x)$. From our earlier demonstration, we have that $|f(x)-g(x)|<\epsilon$ for all $x$ in $K \cap\left(B_{i} \cap B_{j}\right)$. The extensions of $f, h$ and $G$ to all of $B$ require merely an appropriate scaling of $h$. Since $\left\{B_{1}, \ldots, B_{k}\right\}$ forms a simple chain, all three functions will be welldefined on extension and we have that

$$
|f(x)-g(x)|=|h(x) \nabla G(x)-g(x)|<\epsilon \quad \text { for all } x \text { in } K
$$

Therefore, since $\epsilon>0$ was arbitrary, $g$ is an element of $F(K)$. Since $g \in C(K)^{n}$ was arbitrary, $F(K)=C(K)^{n}$, and the theorem is proved.

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[^0]:    * This research represents part of the author's Ph.D. thesis, written at the University of Michigan in 1972, under the guidance of Professor L. Cesari.

