# THE MULTIDIMENSIONAL FUNDAMENTAL THEOREM OF CALCULUS 

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#### Abstract

On compact oriented differentiable manifolds, we define a well behaved Riemann type integral which coincides with the Lebesgue integral on nonnegative functions, and such that the exterior derivative of a differentiable (not necessarily continuously) exterior form is always integrable and the Stokes formula holds.


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## 1. Introduction

The philosophical basis for the present work is a well established maxim: if we can calculate the value of an integral of a function, then such a function ought to be integrable. In the past, an application of this maxim to the bounded convergence theorem for the Riemann integral led to the Lebesgue integral, and an application to the one dimensional fundamental theorem of calculus led to the Denjoy-Perron integral. More recently, J. Mawhin and the author (see [12], [13], and [16], [17]) applied it to the divergence theorem, which is the local multidimensional fundamental theorem of calculus. While the resulting integral is well adapted to the local situation, it is coordinate bound and cannot be lifted to manifolds. In this paper, we apply the maxim to the global fundamental theorem of calculus, that is, to the Stokes theorem.

Theorem. Let $M$ be an oriented compact differentiable manifold with boundary $\partial M$, and let $\omega$ be a continuous exterior form on $M$ which is differentiable in $M-\partial M$. Then $\int_{\partial M} \omega=\int_{M} d \omega$ whenever the integral on the right exists.

If $\omega$ is continuously differentiable in $M$, then the Lebesgue integral $\int_{M} d \omega$ exists. However, in general, $d \omega$ is not Lebesgue integrable even in dimension one. On the other hand, the Lebesgue integral $\int_{\partial M} \omega$ exists and its value depends only on $d \omega$. In accordance with the maxim, under these circumstances the nonexistence of the Lebesgue integral $\int_{M} d \omega$ is not a reflection of an inherent pathology of $d \omega$, but rather a shortcoming of the Lebesgue integral. Indeed, we show that there is a well behaved integral which coincides with the Lebesgue integral on nonnegative functions, and for which the integral $\int_{M} d \omega$ exists and has the correct value.

Initially, we define the integral in an $m$-dimensional Euclidean space by substantially modifying the basic idea of Henstock and Kurzweil (see [3] and [6]). We integrate over sets from a family invariant with respect to diffeomorphisms (to achieve the coordinate independence) and employ partitions satisfying a strong Vitali condition (to prove the divergence theorem). As it is not possible to satisfy the Vitali condition along the boundary of the integration domain, we use only partitions which lie in the interior, and bridge the resulting gap by a suitable continuity requirement imposed on the indefinite integral. Using standard techniques, we lift the integral to differentiable manifolds, and obtain the Stokes theorem from the divergence theorem by means of triangulations.

This approach is quite different from that employed in [4] and [5], and it yields the Stokes theorem (Theorem 7.3) which is considerably more general.

The paper is divided into eight sections whose content is sufficiently indicated by their titles.

## 2. Preliminaries

By $\mathbf{R}$ and $\mathbf{R}_{+}$we denote the sets of all real and all positive real numbers, respectively. All functions in this paper are real-valued, and often the same letter is used to denote a function on a set $A$ as well as its restriction to a set $B \subset A$. The algebraic operations, partial order, and convergence among functions on the same set are defined pointwise.

Throughout, $m \geqslant 1$ is a fixed integer, and $\mathbf{R}^{m}$ denotes the $m$-dimensional Euclidean space equipped with the usual inner product $x \cdot y$ and the corresponding norm $|x|=(x \cdot x)^{1 / 2}$. The distance between a point $x \in \mathbf{R}^{m}$ and a set $E \subset \mathbf{R}^{m}$ is denoted by $\operatorname{dist}(x, E)$. If $E \subset \mathbf{R}^{m}$, then $E^{-}, E^{\circ}, E$; and $d(E)$ denote, respectively, the closure, interior, boundary, and diameter of $E$.

For an integer $k \geqslant 1$, we denote by $\lambda_{k}$ the outer Lebesgue measure in $\mathbf{R}^{k}$. We write $\lambda$ instead of $\lambda_{1}$, and $|E|$ instead of $\lambda_{m}(E)$ whenever $E \subset \mathbf{R}^{m}$. Unless specified otherwise, the words "outer measure", "measure", and "measurable", as well as the expression "almost all", always refer to $\lambda_{m}$.

In Section 3, we shall define a new Riemann type integral which will be denoted by the usual integral sign $\int$. Thus to avoid confusion, throughout, we denote the Lebesgue integral by the sign $(L) f$.

For $E \subset \mathbf{R}^{m}$ and $\eta>0$, we let $U(E, \eta)=\left\{x \in \mathbf{R}^{m}: \operatorname{dist}(x, E)<\eta\right\}$. We say that a set $E \subset \mathbf{R}^{m}$ is, respectively, thin or slight if it is compact and $\mid(U(E, \eta) \mid=$ $O(\eta)$ or $|U(e, \eta)|=o(\eta)$ as $\eta$ approaches zero.

Clearly, each slight set is thin, and each thin set is of measure zero. Moreover, the families of all thin and slight sets are closed with respect to finite unions, and also with respect to intersections by compact sets. For heuristic purposes it is convenient, though inaccurate, to think of thin and slight sets as compact rectifiable pieces of $(m-1)$ - and ( $m-2$ )-dimensional submanifolds of $\mathbf{R}^{m}$, respectively. More precisely, the ( $m-1$ )-dimensional upper Minkowski content of a thin set is finite, and that of a slight set is zero (see [1, Section 3.2.37]).

A bounded subset of $\mathbf{R}^{m}$ whose boundary is thin is called admissible. In view of the previous paragraph, the family $\mathscr{A}$ of all admissible sets is a ring of measurable subsets of $\mathbf{R}^{m}$, and $A \in \mathscr{A}$ whenever $A^{-}$is compact and $A$ is thin. For $E \subset \mathbf{R}^{m}$, we set $\mathscr{A}(E)=\{A \in \mathscr{A}: A \subset E\}$.

Let $A \in \mathscr{A}$. By [9, Theorems 42, 16, and 18], on $A^{\cdot}$ there is a unique finite Borel measure $\sigma_{A}$, called the surface measure, and a $\sigma_{A}$-almost everywhere unique unit Borel vector field $n_{*}$, called the unit exterior normal field, such that

$$
(L) \int_{A} \nabla \cdot v d \lambda_{m}=(L) \int_{A^{*}} v \cdot n_{A} d \sigma_{A}
$$

for each vector field $v=\left(f_{1}, \ldots, f_{m}\right)$ continuously differentiable in an open set containing $A^{-}$; here $\nabla \cdot v=\sum_{i=1}^{m} \partial f_{i} / \partial x_{i}$ is the divergence of $v$. It follows from [9, Theorems 4 and 33] that the number $\|A\|=\sigma_{A}\left(A^{\circ}\right)$, called the surface area of $A$, is positive if and only if $A^{\circ} \neq \varnothing$. Thus we can define the regularity $r(A)$ of $A$ by setting $r(A)=|A| / d(A)\|A\|$ if $A^{\circ} \neq \varnothing$, and $r(A)=0$ if $A^{\circ} \neq \varnothing$. This notion of regularity controls both the shape and surface area of $A$.

The traditional shape indicator of a set $E \subset \mathbf{R}^{m}$ is the number $r^{*}(E)$ defined by $r^{*}(E)=|E| /[d(E)]^{m}$ if $E$ contains at least two points, and $r^{*}(E)=0$ otherwise (see [19, Chapter IV, Section 2, p. 106]). For an admissible set $A$, the relationship between $r(A)$ and $r^{*}(A)$ takes a form of an isoperimetric inequality.
2.1. Proposition. If $A \in \mathscr{A}$, then $[2 r(A)]^{m} \leqslant r^{*}(A)$.

Proof. Let $A \in \mathscr{A}$ and $A^{\circ} \neq \varnothing$. If $A_{i}, i=1, \ldots, m$, is the projection of $A^{\circ}$ to the hyperplane perpendicular to the $i$ th coordinate axis, then it follows from [9,

Theorems 4 and 33] that $2 \lambda_{m-1}\left(A_{i}\right) \leqslant\|A\|$. By [8],

$$
|A|^{m-1} \leqslant \prod_{i=1}^{m} \lambda_{m-1}\left(A_{i}\right) \leqslant 2^{-m}\|A\|^{m}
$$

and hence

$$
[2 r(A)]^{m}=\frac{2^{m}|A|^{m}}{[d(A)]^{m}\|A\|^{m}} \leqslant \frac{|A|}{[d(A)]^{m}}=r^{*}(A)
$$

As $r(A)=r^{*}(A)=0$ if $A^{\circ}=\varnothing$, the proposition is proved.
2.2. Remark. If $A$ is an interval, then $\|A\| \leqslant 2 m[d(A)]^{m-1}$, and we have $r^{*}(A) \leqslant 2 m r(A)$. On the other hand, setting

$$
A_{n}=\left(\bigcup_{k=1}^{n}\left[\frac{1}{2 k}, \frac{1}{2 k-1}\right]\right) \times[0,1]^{m-1}
$$

for $n=1,2, \ldots$, we obtain $r^{*}\left(A_{n}\right) \geqslant 1 / 2 m^{1 / 2}$ and $r\left(A_{n}\right) \leqslant 1 / 2 n$.
A division of an admissible set $A$ is a finite disjoint family $\mathscr{D} \subset \mathscr{A}$ with $\cup \mathscr{D}=A$. An additive function in an admissible set $A$ is a function $F$ on $\mathscr{A}(A)$ such that $F(A)=\Sigma_{D \in \mathscr{D}} F(D)$ for each division $\mathscr{D}$ of $A$.
2.3. Definition. An additive function $F$ in an admissible set $A$ is called continuous whenever we can find a slight set $S_{F} \subset A^{-}$so that the following conditions are satisfied
(i) Given $\varepsilon>0$, there is an $\eta>0$ such that $|F(B)|<\varepsilon$ for each $B \in$ $\mathscr{A}\left[A \cap U\left(S_{F}, \eta\right)\right]$ with $\|B\|<\eta$.
(ii) Given $\varepsilon>0$ and a thin set $T \subset A^{-}-S_{F}$, there is an $\eta>0$ such that $|F(B)|<\varepsilon$ for each $B \in \mathscr{A}[A \cap U(T, \eta)]$ with $\|B\|<1 / \varepsilon$.

The family of all continuous additive functions in an admissible set $A$ is a real vector space denoted by $\mathscr{C}(A)$. If $F \in \mathscr{C}(A)$ and $B \in \mathscr{A}(A)$, then the restriction $G=F \upharpoonright \mathscr{A}(B)$ belongs to $\mathscr{C}(B)$; indeed, it suffices to let $S_{G}=S_{F} \cap B^{-}$.
2.4. Proposition. If $A \in \mathscr{A}$ and $F \in \mathscr{C}(A)$, then $F(B)=0$ for each $B \subset A$ for which $B^{-}$is thin.

Proof. Let $B \subset A$ be such that $B^{-}$is thin, and let $\varepsilon>0$. Then $B \in \mathscr{A}$ and $\|B\|=0$. Hence there is an $\eta>0$ such that $\left|F\left[B \cap U\left(S_{F}, \eta\right)\right]\right|<\varepsilon$. As [ $\left.B-U\left(S_{F}, \eta\right)\right]^{-}$is a thin subset of $A^{-}$disjoint from $S_{F}$, we see that

$$
\left|F\left[B-U\left(S_{F}, \eta\right)\right]\right|<\varepsilon
$$

Thus $|F(B)|<2 \varepsilon$ by the additivity of $F$, and the proposition follows.

Let $E \subset \mathbf{R}^{m}$. A partition in $E$ is a set $P=\left\{\left(A_{1}, x_{1}\right), \ldots,\left(A_{p}, x_{p}\right)\right\}$ where $A_{1}, \ldots, A_{p}$ are disjoint admissible sets, $x_{i} \in A_{i}^{-}$, and $A_{i}^{-} \subset E^{0}, i=1, \ldots, p$. If, in addition, $r\left(A_{i}\right) \geqslant \varepsilon>0$ for $i=1, \ldots, p$, then $P$ is called an $\varepsilon$-partition in $E$. Given $\delta: E \rightarrow \mathbf{R}_{+}$, we say that $P$ is $\delta$-fine whenever $d\left(A_{i}\right)<\delta\left(x_{i}\right), i=1, \ldots, p$. The family of all $\delta$-fine $\varepsilon$-partitions in $E$ is denoted by $\mathscr{P}(E, \varepsilon ; \delta)$. If $E^{*} \subset E$, $\delta^{*} \leqslant \delta \upharpoonright E^{*}$, and $\varepsilon^{*} \geqslant \varepsilon$, then clearly $\mathscr{P}\left(E^{*}, \varepsilon^{*} ; \delta^{*}\right) \subset \mathscr{P}(E, \varepsilon ; \delta)$.

A half-open interval is an interval of the form $\prod_{i=1}^{m}\left[a_{i}, b_{i}\right)$ where $a_{i}, b_{i} \in \mathbf{R}$ and $a_{i}<b_{i}, i=1, \ldots, m$. The following existence result is sometimes referred to as Cousin's lemma. For its proof, which is a simple compactness argument, and some historical discussion we refer to [11].
2.5. Proposition. Let $E \subset \mathbf{R}^{m}, \boldsymbol{\delta}: E \rightarrow \mathbf{R}_{+}$, and let $A$ be a half-open interval with $A^{-} \subset E^{\circ}$. Then there is a $\delta$-fine partition $\left\{\left(A_{1}, x_{1}\right), \ldots,\left(A_{p}, x_{p}\right)\right\}$ in $E$ such that $\bigcup_{i=1}^{P} A_{i}=A$ and each $A_{i}$ is similar to $A$.

## 3. Definition and basic properties of the integral

We begin with the definition of integrability.
3.1. Definition. Let $f$ be a function on an admissible set $A$. We say that $f$ is integrable in $A$ if there is a continuous additive function $F$ in $A$ satisfying the following condition: we can find a thin set $T_{F} \subset A^{-}$so that given $\varepsilon>0$, there is a $\delta: A \rightarrow \mathbf{R}_{+}$such that

$$
\sum_{i=1}^{p}\left|f\left(x_{i}\right)\right| A_{i}\left|-F\left(A_{i}\right)\right|<\varepsilon
$$

for each $\delta$-fine $\varepsilon$-partition $\left\{\left(A_{1}, x_{1}\right), \ldots,\left(A_{p}, x_{p}\right)\right\}$ in $A-T_{F}$.
The family of all functions integrable in an admissible set $A$ is denoted by $\mathscr{F}(A)$.

Let $A \in \mathscr{A}$ and $f \in \mathscr{I}(A)$. Any $F \in \mathscr{C}(A)$ satisfying the condition of Definition 3.1 is called an indefinite integral of $f$. If $F$ is an indefinite integral of $f$ and $B \in \mathscr{A}(A)$, then $G=F \upharpoonright \mathscr{A}(B)$ is an indefinite integral of $f \upharpoonright B$; indeed, it suffices to let $T_{G}=T_{F} \cap B^{-}$. In particular, we have the following proposition.
3.2. Proposition. A function integrable in an admissible set $A$ is integrable in each admissible subset of $A$.

Our first goal is to show that each integrable function has only one indefinite integral.
3.3. Lemma. Let $A \in \mathscr{A}, f \in \mathscr{I}(A)$, and let $F_{1}$ and $F_{2}$ be indefinite integrals of $f$. Then for each $\varepsilon>0$, there is a partition $\left\{\left(A_{1}, x_{1}\right), \ldots,\left(A_{p}, x_{p}\right)\right\}$ in $A$ such that

$$
\left|\sum_{i=1}^{p} f\left(x_{i}\right)\right| A_{i}\left|-F_{j}(A)\right|<\varepsilon
$$

for $j=1,2$.

Proof. (a) For $n=1,2, \ldots$, we denote by $\mathscr{K}_{n}$ the family of all cubes $\prod_{i=1}^{m}\left[k_{i} 2^{-n},\left(k_{i}+1\right) 2^{-n}\right)$ where $k_{1}, \ldots, k_{m}$ are integers. Throughout this proof $j=1,2$. Let $T_{j}=T_{F_{j}}$ and $S_{j}=S_{F_{j}}$ be thin and slight sets associated to $F_{j}$ by Definitions 3.1 and 2.3, respectively. Choose an $\varepsilon>0$ and let $\alpha=$ $\min \left(\varepsilon / 4,1 / 2 m^{3 / 2}\right)$. In three steps we shall construct a set $A_{0}$ which is a finite disjoint union of cubes, and for which $A_{0}^{-} \subset\left[A-\left(T_{1} \cup T_{2}\right)\right]^{-}$and $\left|F_{j}\left(A-A_{0}\right)\right|$ $<3 \alpha$.
(b) Let $S_{0}=S_{1} \cap S_{2}$, and choose $\eta_{0}$ so that $\left|F_{j}(B)\right|<\alpha$ for each $B \in$ $\mathscr{A}\left[A \cap U\left(S_{0}, \eta_{0}\right)\right]$ with $\|B\|<\eta_{0}$. Set $\mathscr{K}_{n, 0}=\left\{K \in \mathscr{K}_{n}: S_{0} \cap K^{-} \neq \varnothing\right\}$ and $K_{n, 0}$ $=U K_{n, 0}$. For $t>0$, let $\beta(t)=\left|U\left(S_{0}, t\right)\right| / t$. If $k_{0}$ is the number of cubes in $\mathscr{X}_{n, 0}$, then

$$
k_{0} 2^{-n m}=\left|K_{n, 0}\right| \leqslant\left|U\left(S_{0}, 2^{-n} m^{1 / 2}\right)\right|=2^{-n} m^{1 / 2} \beta\left(2^{-n} m^{1 / 2}\right)
$$

and hence

$$
\left\|K_{n, 0}\right\| \leqslant \sum\left\{\|K\|: K \in \mathscr{K}_{n, 0}\right\}=k_{0}\left(2 m 2^{-n(m-1)}\right) \leqslant 2 m^{3 / 2} \beta\left(2^{-n} m^{1 / 2}\right)
$$

Since $S_{0}$ is slight, we see that $\lim _{n \rightarrow \infty}\left\|K_{n, 0}\right\|=0$. By [14, Section 3], there is an integer $n_{0}$ such that $2^{-n_{0}} \mathrm{~m}^{1 / 2}<\eta_{0}$ and $\left\|A \cap K_{n, 0}\right\|<\eta_{0}$. Thus letting $K_{0}=K_{n, 0}$, we have $S_{0} \subset K_{0}^{\circ}$ and $\left|F_{j}\left(A \cap K_{0}\right)\right|<\alpha$.
(c) Choose $\eta_{1}>0$ so that $\left|F_{j}(B)\right|<\alpha / 2$ for each

$$
B \in \mathscr{A}\left[A \cap U\left(S_{1}-K_{0}^{\circ}, \eta_{1}\right)\right]
$$

with $\|B\|<\eta_{1}$. Such a choice is clearly possible for $j=1$; it is also possible for $j=2$ because $S_{1}-K_{0}^{\circ}$ is a thin subset of $A^{-}-S_{2}$. Analogously, choose $\eta_{2}>0$ so that $\left|F_{j}(B)\right|<\alpha / 2$ for each $B \in \mathscr{A}\left[A \cap U\left(S_{2}-K_{0}^{\circ}, \eta_{2}\right)\right]$ with $\|B\|<\eta_{2}$. Let $S=S_{1} \cup S_{2}$. Since $S_{j}-K_{0}^{\circ}$ are disjoint compact sets, there is a positive $\eta_{S} \leqslant$ $\min \left(\eta_{1}, \eta_{2}\right)$ such that $U\left(S-K_{0}^{\circ}, \eta_{S}\right)$ is a disjoint union of $U_{j}=U\left(S_{j}-K_{0}^{\circ}, \eta_{S}\right)$. Now if $B \in \mathscr{A}\left[A \cap U\left(S-K_{0}^{\circ}, \eta_{S}\right)\right]$ and $\|B\|<\eta_{S}$, then $\left\|B \cap U_{j}\right\| \leqslant\|B\|<\eta_{j}$, and we have

$$
\left|F_{j}(B)\right| \leqslant\left|F_{j}\left(B \cap U_{1}\right)\right|+\left|F_{j}\left(B \cap U_{2}\right)\right|<\alpha
$$

For $n \geqslant n_{0}$, set

$$
K_{n, s}=\bigcup\left\{K \in \mathscr{X}_{n}: S \cap K^{-} \neq \varnothing \text { and } K_{0} \cap K=\varnothing\right\}
$$

Arguing similarly as in (b), we find an integer $n_{S} \geqslant n_{0}$ such that $2^{-n_{S}} m^{1 / 2}<\eta_{S}$ and $\left\|A \cap K_{n_{s}, s}\right\|<\eta_{S}$. Thus letting $K_{S}=K_{n_{s}, s}$, we have $S \subset\left(K_{0} \cap K_{S}\right)^{\circ}$ and $\left|F_{j}\left(A \cap K_{S}\right)\right|<\alpha$.
(d) Since $T=A^{\cdot} \cup T_{1} \cup T_{2}-\left(K_{0} \cup K_{S}\right)^{\circ}$ is a thin subset of $A^{-}-S$, there is a positive $\eta_{T} \leqslant 1$ such that $\left|F_{j}(B)\right|<\alpha$ for each $B \in \mathscr{A}\left[A \cap U\left(T, \eta_{T}\right)\right]$ with $\|B\| \leqslant\|A\| m^{1 / 2}+2 m^{3 / 2} \gamma$ where

$$
\gamma=\sup \left\{\frac{|U(T, t)|}{t}: 0<t<1\right\}
$$

Find $n_{T} \geqslant n_{S}$ so that $2^{-n_{T}} m^{1 / 2}<\eta_{T}$, and set

$$
\mathscr{K}_{T}=\left\{K \in \mathscr{K}_{n_{T}}: T \cap K^{-} \neq \varnothing \text { and }\left(K_{0} \cup K_{S}\right) \cap K=\varnothing\right\} .
$$

If $K_{T}=\cup \mathscr{K}_{T}$, then $A^{\cdot} \cup T_{1} \cup T_{2} \subset\left(K_{0} \cup K_{S} \cup K_{T}\right)^{\circ}$. Denoting by $k_{T}$ the number of cubes in $\mathscr{K}_{T}$, we have

$$
k_{T} 2^{-n_{T} m}=\left|K_{T}\right| \leqslant\left|U\left(T, 2^{-n_{T}} m^{1 / 2}\right)\right| \leqslant 2^{-n_{T}} m^{1 / 2} \gamma
$$

From this and [9, Theorem 9(c) and Lemma 36] we obtain

$$
\begin{aligned}
\left\|A \cap K_{T}\right\| & \leqslant \sum\left\{\|A \cap K\|: K \in \mathscr{K}_{T}\right\} \leqslant\|A\| m^{1 / 2}+\sum\left\{\|K\|: K \in \mathscr{K}_{T}\right\} \\
& =\|A\| m^{1 / 2}+k_{T}\left(2 m 2^{-n_{T}(m-1)}\right) \leqslant\|A\| m^{1 / 2}+2 m^{3 / 2} \gamma
\end{aligned}
$$

and hence $\left|F_{j}\left(A \cap K_{T}\right)\right|<\alpha$. Now if $A_{0}=A-\left(K_{0} \cup K_{S} \cup K_{T}\right)$, then $A_{0}^{-} \subset$ $\left[A-\left(T_{1} \cup T_{2}\right)\right]^{\circ}$ and $A_{0}=\bigcup\left\{K \in \mathscr{K}_{n_{T}}: K \subset A_{0}\right\}$. Moreover,

$$
\left|F_{j}\left(A-A_{0}\right)\right| \leqslant\left|F_{j}\left(A \cap K_{0}\right)\right|+\left|F_{j}\left(A \cap K_{S}\right)\right|+\left|F_{j}\left(A \cap K_{T}\right)\right|<3 \alpha
$$

(e) Find $\delta: A \rightarrow \mathbf{R}_{+}$so that

$$
\sum_{i=1}^{q}\left|f\left(y_{i}\right)\right| B_{i}\left|-F_{j}\left(B_{i}\right)\right|<\alpha
$$

for each $\delta$-fine $\alpha$-partition $\left\{\left(B_{1}, y_{1}\right), \ldots,\left(B_{q}, y_{q}\right)\right\}$ in $A-\left(T_{1} \cup T_{2}\right)$. As $A_{0}$ is a finite disjoint union of cubes which are half-open intervals, it follows from Proposition 2.5 that there is a $\delta$-fine partition $P=\left\{\left(A_{1}, x_{1}\right), \ldots,\left(A_{p}, x_{p}\right)\right\}$ in $A-\left(T_{1} \cup T_{2}\right)$ such that $A_{1}, \ldots, A_{p}$ are cubes and $\cup_{i=1}^{p} A_{i}=A_{0}$. The regularity of a cube equals $1 / 2 m^{3 / 2}$, and so $P$ is an $\alpha$-partition by our choice of $\alpha$. Thus

$$
\begin{aligned}
\left|\sum_{i=1}^{p} f\left(x_{i}\right)\right| A_{i}\left|-F_{j}(A)\right| & \leqslant \sum_{i=1}^{p}\left|f\left(x_{i}\right)\right| A_{i}\left|-F_{j}\left(A_{i}\right)\right|+\left|F_{j}\left(A-A_{0}\right)\right| \\
& <4 \alpha<\varepsilon
\end{aligned}
$$

### 3.4. Proposition. All indefinite integrals of an integrable function are equal.

Proof. Let $A \in \mathscr{A}$ and $f \in \mathscr{I}(A)$, and assume that $F_{1}$ and $F_{2}$ are indefinite integrals of $f$ such that $F_{1}(B) \neq F_{2}(B)$ for some $B \in \mathscr{A}(A)$. Since $F_{1} \upharpoonright \mathscr{A}(B)$ and $F_{2} \uparrow \mathscr{A}(B)$ are indefinite integrals of $f \upharpoonright B$, a direct application of Lemma 3.3 yields a contradiction.

Let $A \in \mathscr{A}$ and $f \in \mathscr{I}(A)$. In view of Proposition 3.4, we can talk about the indefinite integral of $f$ denoted by $\int f$. Its value at $B \in \mathscr{A}(A)$ is called the integral of $f$ over $B$ denoted by $\int_{B} f$.
3.5. Proposition. If $A \in \mathscr{A}$, then $\mathscr{I}(A)$ is a linear space, and the map $f \mapsto \int_{A} f$ is a nonnegative linear functional on $\mathscr{I}(A)$.

Lemma 3.3 shows that $\int_{A} f \geqslant 0$ for each nonnegative function $f \in \mathscr{I}(A)$. The rest of the proposition follows directly from Definition 3.1.
3.6. Proposition. Let $f$ be a function on $A \in \mathscr{A}$, and let $\mathscr{D}$ be a division of $A$. If $f \in \mathscr{I}(D)$ for each $D \in \mathscr{I}$, then $f \in \mathscr{I}(A)$ and $\int_{A} f=\Sigma_{D \in \mathscr{D}} \int_{D} f$.

Proof. For $D \in \mathscr{D}$, let $F_{D}=\int(f \mid D)$, and let $T_{D}=T_{F_{D}}$ be a thin set associated to $F_{D}$ by Definition 3.1. Setting

$$
F(B)=\sum_{D \in \mathscr{D}} F_{D}(B \cap D)
$$

for each $B \in \mathscr{A}(A)$, it is easy to see that $F \in \mathscr{C}(A)$. Choose $\varepsilon>0$, and denote by $k$ the number of elements in $\mathscr{D}$. Given $D \in \mathscr{D}$, find a $\delta_{D}: D \rightarrow \mathbf{R}_{+}$so that

$$
\sum_{i=1}^{p}\left|f\left(x_{i}\right)\right| A_{i}\left|-F_{D}\left(A_{i}\right)\right|<\varepsilon / k
$$

for every $\left\{\left(A_{1}, x_{1}\right), \ldots,\left(A_{p}, x_{p}\right)\right\}$ in $\mathscr{P}\left(D-T_{D}, \boldsymbol{\varepsilon} ; \delta_{D}\right)$. Let $T=\cup_{D \in \mathscr{D}}\left(D^{\cdot} \cup T_{D}\right)$, and for $x \in A$ set

$$
\delta(x)=\min \left[\delta_{D}(x), \operatorname{dist}\left(x, D^{\bullet} \cup T_{D}\right)\right]
$$

if $x \in D^{\circ}-T_{D}$ and $D \in \mathscr{D}$, and $\delta(x)=1$ otherwise. Now if $P \in$ $\mathscr{P}(A-T, \varepsilon ; \delta)$, then $P_{D}=\{(B, x) \in P: x \in D\}$ belongs to $\mathscr{P}\left(D-T_{D}, \varepsilon ; \delta_{D}\right)$ for each $D \in \mathscr{D}$, and $P$ is a disjoint union of the $P_{D}$ 's. Consequently,

$$
\sum_{(B, x) \in P}|f(x)| B|-F(B)|=\sum_{D \in \mathscr{D}} \sum_{(B, x) \in P_{D}}|f(x)| B\left|-F_{D}(B)\right|<k \frac{\varepsilon}{k}=\varepsilon
$$

and we see that $F=\int f$. The proposition follows.

## 4. The relationship to the Lebesgue integral, and its consequences

If $E \subset \mathbf{R}^{\boldsymbol{m}}$ is a measurable set, we denote by $\mathscr{L}(E)$ the family of all functions $f$ on $E$ for which the finite Lebesgue integral $(L) \int_{E} f d \lambda_{m}$ exists. When no misunderstanding is possible, we write $(L) \int_{E} f$ instead of $(L) \int_{E} f d \lambda_{m}$.
4.1. Proposition. If $A \in \mathscr{A}$, then $\mathscr{L}(A) \subset \mathscr{I}(A)$ and $\int_{A} f=(L) \int_{A} f$ for each $f \in \mathscr{L}(A)$.

Proof. Let $f \in \mathscr{L}(A)$, and set $F(B)=(L) \int_{B} f$ for each $B \in \mathscr{A}(A)$. As thin sets are compact and of measure zero, it is easy to conclude from the absolute continuity of the Lebesgue integral that $F \in \mathscr{C}(A)$. We now show that $F=\int f$.

Given $\varepsilon>0$, there are functions $g$ and $h$ on $A$ which are, respectively, upper and lower semicontinuous, and such that $g \leqslant f \leqslant h$ and $(L) \int_{A}(h-g)<\varepsilon$. Find $\delta: A \rightarrow R_{+}$so that $g(y)<g(x)+\varepsilon$ and $h(y)>h(x)-\varepsilon$ for each $x, y \in A$ with $|x-y|<\delta(x)$. Now let $\left\{\left(A_{1}, x_{1}\right), \ldots,\left(A_{p}, x_{p}\right)\right\}$ be a $\delta$-fine partition in $A$. By the choice of $\delta$,
(L) $\int_{A_{i}} g-\varepsilon\left|A_{i}\right| \leqslant g\left(x_{i}\right)\left|A_{i}\right|<f\left(x_{i}\right)\left|A_{i}\right| \leqslant h\left(x_{i}\right)\left|A_{i}\right| \leqslant(L) \int_{A_{i}} h+\varepsilon\left|A_{i}\right|$,
and since $(L) \int_{A_{i}} g \leqslant F\left(A_{i}\right) \leqslant(L) \int_{A_{i}} h$, we conclude

$$
\begin{aligned}
\sum_{i=1}^{p}\left|f\left(x_{i}\right)\right| A_{i}\left|-F\left(A_{i}\right)\right| & \leqslant \sum_{i=1}^{p}\left[(L) \int_{A_{i}}(h-g)+\varepsilon\left|A_{i}\right|\right] \\
& =(L) \int_{A}(h-g)+\varepsilon|A|<\varepsilon(1+|A|)
\end{aligned}
$$

Let $A \in \mathscr{A}, x \in A^{-}$, and let $F$ be a function on $\mathscr{A}(A)$. We say that $F$ is derivable at $x$ if a finite $\lim \left[F\left(B_{n}\right) /\left|B_{n}\right|\right]$ exists for each sequence $\left\{B_{n}\right\}$ in $\mathscr{A}(A)$ such that $x \in B_{n}^{-}, \lim d\left(B_{n}\right)=0$, and $\inf r\left(B_{n}\right)>0$. If all these limits exist, then they have the same value, which is denoted by $F^{\prime}(x)$ and called the derivative of $F$ at $x$. In view of Proposition 2.1, it is clear that for each $x \in A^{\circ}$, the existence of the derivative $F^{\prime}(x)$ implies the existence of the ordinary derivative of $F$ at $x$ defined in [19, Chapter IV, Section 2, p. 106]; both derivatives have the same value.
4.2. Proposition. Let $A \in \mathscr{A}, f \in \mathscr{I}(A)$, and let $F=\int f$. Then for almost all $x \in A$ the function $F$ is derivable at $x$ and $F^{\prime}(x)=f(x)$.

Proof. Let $T$ be a thin set associated to $F$ by Definition 3.1, and let $E$ be the set of all $x \in A^{\circ}-T$ such that either $F$ is not derivable at $x$, or $F^{\prime}(x) \neq f(x)$. Then given $x \in E$, we can find an $\alpha(x)>0$ so that for each $\beta>0$ there is a closed set $B \in \mathscr{A}(A)$ with $x \in B, d(B)<\beta, r(B) \geqslant \alpha(x)$, and $\|F(B) / \mid B\|-$ $f(x) \mid \geqslant \alpha(x)$; for by Proposition 2.4, $F\left(B^{-}\right)=F(B)$ for each $B \in \mathscr{A}(A)$ with $B^{-} \subset A$. Fix an integer $n \geqslant 1$, and let $E_{n}=\{x \in E: \alpha(x) \geqslant 1 / n\}$. Choose a positive $\varepsilon \leqslant 1 / n$, and find a $\delta: A \rightarrow R_{+}$such that

$$
\sum_{i=1}^{p}\left|f\left(x_{i}\right)\right| A_{i}\left|-F\left(A_{i}\right)\right|<\frac{\varepsilon}{n}
$$

for each $\left\{\left(A_{1}, x_{1}\right), \ldots,\left(A_{p}, x_{p}\right)\right\}$ in $\mathscr{P}(A-T, \varepsilon ; \delta)$. Let $\mathscr{B}$ be the family of all closed sets $B \in \mathscr{A}\left(A^{\circ}-T\right)$ such that $r(B) \geqslant \varepsilon, d(B)<\delta\left(x_{B}\right)$ for some $x_{B} \in B$, and $\left|f\left(x_{B}\right)\right| B|-F(B)| \geqslant|B| / n$. Using Proposition 2.1, it is easy to see that $\mathscr{B}$ covers $E_{n}$ in the sense of Vitali. Thus by [19, Chapter IV, Theorem (3.1)], there are disjoint sets $B_{1}, B_{2}, \ldots$ in $\mathscr{B}$ such that $\left|E_{n}-\bigcup_{i=1}^{\infty} B_{i}\right|=0$. As $\left\{\left(B_{1}, x_{B_{1}}\right), \ldots,\left(B_{k}, x_{B_{k}}\right)\right\}$ belongs to $\mathscr{P}(A-T, \varepsilon ; \delta)$, we have

$$
\sum_{i=1}^{k}\left|B_{i}\right| \leqslant n \sum_{i=1}^{k}\left|f\left(x_{B_{i}}\right)\right| B_{i}\left|-F\left(B_{i}\right)\right|<\varepsilon
$$

for $k=1,2, \ldots$, and consequently

$$
\left|E_{n}\right| \leqslant\left|\bigcup_{i=1}^{\infty} B_{i}\right| \leqslant \sum_{i=1}^{\infty}\left|B_{i}\right| \leqslant \varepsilon .
$$

As $\varepsilon$ can be arbitrarily small, $\left|E_{n}\right|=0$, and as $E=\bigcup_{n=1}^{\infty} E_{n}$, also $|E|=0$. Since $\left|A^{\cdot} \cup T\right|=0$, the proposition is proved.

### 4.3. Corollary. Each function integrable on an admissible set is measurable.

The corollary follows immediately from Proposition 4.2 and [19, Chapter IV, Theorem (4.2)].
4.4. Proposition. Let $f$ be a function on an admissible set $A$. Then $f$ belongs to $\mathscr{L}(A)$ if and only if both $f$ and $|f|$ belong to $\mathscr{I}(A)$.

Proof. As the converse is given by Proposition 4.1, assume that $f$ and $|f|$ are integrable, and let $g_{n}=\min (|f|, n), n=1,2, \ldots$. Since $|A|<+\infty$, it follows from Corollary 4.3 that $g_{n} \in \mathscr{L}(A)$. By Propositions 4.1 and 3.5 , we have $g_{n} \in \mathscr{I}(A)$ and

$$
(L) \int_{A}|f|=\lim (L) \int_{A} g_{n}=\lim \int_{A} g_{n} \leqslant \int_{A}|f|<+\infty
$$

Thus $f \in \mathscr{L}(A)$ by Corollary 4.3.
4.5. Corollary. Let $f$ be a function on an admissible set $A$. Then $f=0$ almost everywhere if and only if $f \in \mathscr{I}(A)$ and $\int f=0$.

Indeed, if $f \in \mathscr{F}(A)$ and $\int f=0$, it follows directly from Definition 3.1 that $|f| \in \mathscr{I}(A)$ and $\int|f|=0$.
4.6. Corollary. Let $g, h$, and $f_{n}, n=1,2, \ldots$, be integrable on an admissible set $A$, and let $f=\lim f_{n}$. Suppose that either of the following conditions holds:
(i) $f_{n} \leqslant f_{n+1}, n=1,2, \ldots$, and $\lim \int_{A} f_{n}<+\infty$;
(ii) $g \leqslant f_{n} \leqslant h, n=1,2, \ldots$

Then $f \in \mathscr{F}(A)$ and $\int_{A} f=\lim \int_{A} f_{n}$.

Indeed, in view of Propositions 3.5, 4.1, and 4.4, it suffices to apply the monotone and dominated convergence theorems for the Lebesgue integral to the sequences $\left\{f_{n}-f_{1}\right\}$ and $\left\{f_{n}-g\right\}$, respectively.

Note. If $m \geqslant 2$, then a countable union of slight sets need not be a thin set. Thus it is unclear how to obtain Corollary 4.6 directly from Definition 3.1.

## 5. The divergence theorem

We say that sets $A, B \subset \mathbf{R}$ are $\lambda$-equivalent if $\lambda[(A-B) \cup(B-A)]=0$. For $z=\left(\zeta_{1}, \ldots, \zeta_{m-1}\right)$ in $\mathbf{R}^{m-1}$ and $t$ in $\mathbf{R}$, we write $(z, t)$ instead of $\left(\zeta_{1}, \ldots, \zeta_{m-1}, t\right)$. Given $E \subset \mathbf{R}^{m}$, we let $E^{z}=\{t \in \mathbf{R}:(z, t) \in E\}$ for each $z \in \mathbf{R}^{m-1}$, and we set $E^{\wedge}=\left\{z \in \mathbf{R}^{m-1}: E^{z} \neq \varnothing\right\}^{-}$. If $A \in \mathscr{A}$, we denote by $\bar{\sigma}_{A}$ the completion of the surface measure $\sigma_{A}$.

We begin by summarizing some results of J. Mařík.
5.1. Lemma. Let $A \in \mathscr{A}$ and let $n_{A}=\left(\nu_{1}, \ldots, \nu_{m}\right)$. There is a set $A^{*} \subset A^{\wedge}$ such that the following conditions are satisfied.
(i) $\lambda_{m-1}\left(A^{\wedge}-A^{*}\right)=0$.
(ii) For each $z \in A^{*}$, there is an integer $k(z) \geqslant 0$ and real numbers $a_{1}^{z}<b_{1}^{z}<$ $\cdots<a_{k(z)}^{z}<b_{k(z)}^{z}$ such that $A^{z}$ and $\bigcup_{j=1}^{k(z)}\left(a_{j}^{z}, b_{j}^{z}\right)$ are $\lambda$-equivalent. In particular, the points $\left(z, a_{j}^{z}\right)$ and $\left(z, b_{j}^{z}\right), j=1, \ldots, k(z)$, belong to $A$ :
(iii) $\nu_{m}\left(z, a_{j}^{z}\right)<0<\nu_{m}\left(z, b_{j}^{2}\right)$ for each $z \in A^{*}$ and $j=1, \ldots, k(z)$.
(iv) Iff is a bounded $\bar{\sigma}_{A}$-measurable function on $A$; and

$$
f^{*}(z)=\sum_{j=1}^{k(z)}\left[f\left(z, b_{j}^{z}\right)-f\left(z, a_{j}^{z}\right)\right]
$$

for each $z \in A^{*}$, then $f^{*}$ is Lebesgue integrable on $A^{*}$ and

$$
(L) \int_{A^{*}} f^{*} d \lambda_{m-1}=(L) \int_{A} f \nu_{m} d \bar{\sigma}_{A}
$$

In particular,

$$
(L) \int_{A^{\cdot}}\left|\nu_{m}\right| d \sigma_{A}=(L) \int_{A^{*}} 2 k d \lambda_{m-1}
$$

where $k$ denotes the function $z \mapsto k(z)$.
Proof. The existence of $A^{*} \subset A^{\wedge}$ which satisfies (i) and (ii) follows directly from [9, Theorem 33]. Moreover,

$$
(L) \int_{A^{\cdot}}\left|\nu_{m}\right| d \sigma_{A}=(L) \int_{A^{*}} 2 k d \lambda_{m-1}
$$

by [9, Remark 2 to Theorem 18]. Using [9, Theorem 20], we see that (iv) is satisfied whenever $f$ is a bounded Borel function on $A$ : In particular,

$$
(L) \int_{A^{*}}\left(\operatorname{sign} \nu_{m}\right)^{*} d \lambda_{m-1}=(L) \int_{A}\left|\nu_{m}\right| d \sigma_{A}=(L) \int_{A^{*}} 2 k d \lambda_{m-1}
$$

Since $\left(\operatorname{sign} \nu_{m}\right)^{*} \leqslant 2 k$ in $A^{*}$, replacing $A^{*}$ by $\left\{z \in A^{*}:\left(\operatorname{sign} \nu_{m}\right)^{*}(z)=2 k(z)\right\}$ proves (iii). Now let $f$ be a bounded function on $A^{*}$ which is zero $\bar{\sigma}_{A}$-almost everywhere. Then there is a nonnegative bounded Borel function $g$ on $A^{*}$ which is zero $\sigma_{A}$-almost everywhere, and for which $|f| \leqslant g$. By (iii), we have $\left|f^{*}\right| \leqslant$ $\left(g \operatorname{sign} \nu_{m}\right)^{*}$, and so
$(L) \int_{A^{*}}\left|f^{*}\right| d \lambda_{m-1} \leqslant(L) \int_{A^{*}}\left(g \operatorname{sign} \nu_{m}\right)^{*} d \lambda_{m-1}=(L) \int_{A} g\left|\nu_{m}\right| d \sigma_{A}=0$.
As the map $f \mapsto f^{*}$ is linear, it is easy to see that (iv) holds in general.
5.2. Proposition. Let $A \in \mathscr{A}$, and let $B \subset A^{\cdot}$ be $\bar{\sigma}_{A}$-measurable. Then $\bar{\sigma}_{A}(B) \leqslant$ $(m / 2) \liminf _{n \rightarrow \infty} n|U(B, 1 / n)|$.

Proof. Let $n_{A}=\left(\nu_{1}, \ldots, \nu_{m}\right)$, and let $\chi$ be the characteristic function of $B$. Then

$$
\bar{\sigma}_{A}(B)=(L) \int_{A^{*}} \chi d \bar{\sigma}_{A}=\sum_{i=1}^{m}(L) \int_{A} \chi \nu_{i}^{2} d \bar{\sigma}_{A}
$$

In the notation of Lemma 5.1, fix a $z \in A^{*}$ and observe that

$$
\left(\chi \nu_{m}\right)^{*}(z) \leqslant \sum_{j=1}^{k(z)}\left[\chi\left(z, b_{j}^{z}\right)+\chi\left(z, a_{j}^{z}\right)\right]
$$

If $U_{n}=U(B, 1 / n), n=1,2, \ldots$, then $U_{n}^{z}$ contains intervals $\left(a_{j}^{z}-1 / n, a_{j}^{z}+1 / n\right)$ and $\left(b_{j}^{2}-1 / n, b_{j}^{z}+1 / n\right)$ whenever $\chi\left(z, a_{j}^{z}\right)=1$ and $\chi\left(z, b_{j}^{z}\right)=1$, respectively. Thus

$$
\frac{2}{n} \sum_{j=1}^{k(z)}\left[\chi\left(z, b_{j}^{z}\right)+\chi\left(z, a_{j}^{z}\right)\right] \leqslant \lambda\left(U_{n}^{z}\right)
$$

for all sufficiently large $n$. It follows that $2\left(\chi \nu_{m}\right)^{*}(z) \leqslant \liminf _{n \rightarrow \infty} n \lambda\left(U_{n}^{z}\right)$, and so by Lemma 5.1, (iv), the Fatou lemma and Fubini theorem,

$$
\begin{aligned}
2(L) \int_{A} \chi \nu_{m}^{2} d \bar{\sigma}_{A} & =2(L) \int_{A^{*}}\left(\chi \nu_{m}\right)^{*} d \lambda_{m-1} \leqslant(L) \int_{A^{*}} \liminf _{n \rightarrow \infty} n \lambda\left(U_{n}^{2}\right) d \lambda_{m-1} \\
& \leqslant \liminf _{n \rightarrow \infty} n(L) \int_{A^{*}} \lambda\left(U_{n}^{2}\right) d \lambda_{m-1}=\liminf _{n \rightarrow \infty} n\left|U_{n}\right| .
\end{aligned}
$$

From this the proposition follows by symmetry.
5.3. Corollary. If $A \in \mathscr{A}$, then $\sigma_{A}(S)=0$ for each slight set $S \subset A$.
5.4. Lemma. Let $A \in \mathscr{A}$, let $S$ be a slight subset of $A^{-}$, and let $v$ be a bounded vector field on $A^{-}$which is continuous in $A^{-}-S$. Then $v \uparrow B^{*}$ is $\bar{\sigma}_{B^{-}}$-measurable for each $B \in \mathscr{A}(A)$, and the map $B \rightarrow(L) \int_{B} \cdot v \cdot n_{B} d \bar{\sigma}_{B}$ is a continuous additive function in $A$.

Proof. The $\bar{\sigma}_{B}$-measurability of $v \upharpoonright B^{\cdot}$ for each $B \in \mathscr{A}(A)$ follows from Corollary 5.3. Thus letting $F(B)=(L) \int_{B} \cdot v \cdot n_{B} d \bar{\sigma}_{B}$ for $B \in \mathscr{A}(A)$, we see from [9, Remark 2 to Theorem 14] that $F$ is an additive function in $A$. Furthermore, if $\alpha=\sup \left\{|v(x)|: x \in A^{-}\right\}$then $|F(B)| \leqslant \alpha\|B\|$ for each $B \in \mathscr{A}(A)$.

Let $T \subset A^{-}-S$ be a thin set and let

$$
\beta=\sup \left\{\frac{|U(T, \eta)|}{\eta}: 0<\eta<1\right\} .
$$

Choose a positive $\varepsilon \leqslant 1 / 2 \beta m$, and find a positive $\eta_{0}<1$ so that $C=\left[U\left(T, \eta_{0}\right)\right]^{-}$ does not meet $S$. Then $v$ is uniformly continuous in $A^{-} \cap C$, and hence there is a positive $h<\eta_{0} / m^{1 / 2}$ such that $|v(x)-v(y)|<\varepsilon^{2} / 2 m^{1 / 2}$ for each $x, y \in A^{-} \cap$ $C$ with $|x-y|<h m^{1 / 2}$. Among all cubes $\prod_{i=1}^{m}\left[k_{i} h,\left(k_{i}+1\right) h\right]$ where $k_{1}, \ldots, k_{m}$ are integers, let $K_{1}, \ldots, K_{n}$ be those which meet $T$. Then

$$
n h^{m}=\left|\bigcup_{j=1}^{n} K_{j}\right| \leqslant\left|U\left(T, h m^{1 / 2}\right)\right| \leqslant \beta h m^{1 / 2}
$$

and so

$$
\sum_{j=1}^{n}\left\|K_{j}\right\|=n\left(2 m h^{m-1}\right) \leqslant 2 \beta m^{3 / 2}
$$

Moreover, $\cup_{j=1}^{n} K_{j} \subset U\left(T, \eta_{0}\right)$ and $T \subset\left(\cup_{j=1}^{n} K_{j}\right)^{\circ}$. Find $\eta>0$ so that $U(T, \eta) \subset \bigcup_{j=1}^{n} K_{j}$, and for $j=1, \ldots, n$ select an $x_{j} \in K_{j} \cap T$. Now if $B \in$ $\mathscr{A}[A \cap U(T, \eta)]$ and $\|B\|<1 / \varepsilon$, then by [9, Theorem 36],

$$
\begin{aligned}
|F(B)| & \leqslant \sum_{j=1}^{n}\left|F\left(B \cap K_{j}\right)\right| \\
& =\sum_{j=1}^{n}\left|(L) \int_{\left(B \cap K_{j}\right)}\left[v(x)-v\left(x_{j}\right)\right] \cdot n_{B \cap K_{j}} d \bar{\sigma}_{B \cap K_{j}}(x)\right| \\
& \left.\leqslant \frac{\varepsilon^{2}}{2 m^{1 / 2}} \sum_{j=1}^{n}\left\|B \cap K_{j}\right\| \leqslant \frac{\varepsilon^{2}}{2 m^{1 / 2}}\left(\|B\| m^{1 / 2}+\sum_{j=1}^{n}\left\|K_{j}\right\|\right) \right\rvert\, \\
& <\frac{\varepsilon}{2}+\varepsilon^{2} \beta m \leqslant \varepsilon .
\end{aligned}
$$

It follows that $F$ is continuous.
Throughout this paper, the word differentiable is used in the usual sense (see, e.g. [18, Definition 9.11]). Thus differentiability of a function implies its continuity and the existence of partial derivatives which need not be continuous. The $i$ th partial derivative of a differentiable function $f$ is denoted by $\partial_{i} f, i=1, \ldots, m$. If $v=\left(f_{1}, \ldots, f_{m}\right)$ is a vector field on a set $E \subset \mathbf{R}^{m}$, we call a divergence of $v$ any function $g$ on $E$ such that $g(x)=\sum_{i=1}^{m} \partial_{i} f(x)$ for each $x \in E^{\circ}$ at which $v$ is differentiable. Each divergence of a vector field $v$ is denoted by $\nabla \cdot v$.
5.5. Lemma. Let $v$ be a bounded vector field in an open set $U \subset \mathbf{R}^{m}$, and let $x \in U$. Suppose that $v$ is differentiable at $x$, and that $v \upharpoonright B^{*}$ is $\bar{\sigma}_{B}$-measurable for each $B \in \mathscr{A}$ with $B^{-} \subset U$. Then given $\varepsilon>0$, there is $a \delta>0$ such that

$$
|\nabla \cdot v(x)| B\left|-(L) \int_{B^{\cdot}} v \cdot n_{B} d \bar{\sigma}_{B}\right|<\varepsilon|B|
$$

for each $B \in \mathscr{A}$ with $B^{-} \subset U, x \in B^{-}, d(B)<\delta$, and $r(B) \geqslant \varepsilon$.

Proof. It suffices to prove the lemma for $v=(0, \ldots, 0, f)$. Let $x=\left(\xi_{1}, \ldots, \xi_{m}\right)$, and let $y=\left(\eta_{1}, \ldots, \eta_{m}\right)$ be in $U$. By our assumption, there is a function $\alpha$ in $U$ such that

$$
f(y)-f(x)=\sum_{i=1}^{m} \partial_{i} f(x)\left(\eta_{i}-\xi_{i}\right)+|y-x| \alpha(y)
$$

and $\lim _{y \rightarrow x} \alpha(y)=0$. Given $\varepsilon>0$, choose a $\delta>0$ so that $|\alpha(y)|<\varepsilon^{2}$ whenever $|y-x|<\delta$. Select $B \in \mathscr{A}$ with $B^{-} \subset U, x \in B^{-}, d(B)<\delta$, and $r(B) \geqslant \varepsilon$. In the notation of Lemma 5.1, for $z \in B^{*}$ and $j=1, \ldots, k(z)$ we have

$$
\begin{aligned}
f\left(z, b_{j}^{z}\right)- & f\left(z, a_{j}^{z}\right)=\left[f\left(z, b_{j}^{z}\right)-f(x)\right]-\left[f\left(z, a_{j}^{z}\right)-f(x)\right] \\
= & \sum_{i=1}^{m-1} \partial_{i} f(x)\left(\zeta_{i}-\xi_{i}\right)+\partial_{m} f(x)\left(b_{j}^{z}-\xi_{m}\right)+\left|\left(z, b_{j}^{z}\right)-x\right| \alpha\left(z, a_{j}^{z}\right) \\
& -\sum_{i=1}^{m-1} \partial_{i} f(x)\left(\zeta_{i}-\xi_{i}\right)-\partial_{m} f(x)\left(a_{j}^{z}-\xi_{m}\right)-\left|\left(z, a_{j}^{z}\right)-x\right| \alpha\left(z, a_{j}^{z}\right) \\
= & \partial_{m} f(x)\left(b_{j}^{z}-a_{j}^{z}\right)+\left|\left(z, b_{j}^{z}\right)-x\right| \alpha\left(z, b_{j}^{z}\right)-\left|\left(z, a_{j}^{z}\right)-x\right| \alpha\left(z, a_{j}^{z}\right),
\end{aligned}
$$

and hence

$$
\left|f^{*}(z)-\partial_{m} f(x) \lambda\left(B^{z}\right)\right|<2 k(z) \varepsilon^{2} d(B)
$$

By Lemma 5.1(iv) and the Fubini theorem, we obtain

$$
\begin{aligned}
|\nabla \cdot v(x)| B \mid & -(L) \int_{B^{*}} v \cdot n_{B} d \bar{\sigma}_{B}\left|=\left|\partial_{m} f(x)\right| B\right|-(L) \int_{B^{*}} f \nu_{m} d \bar{\sigma}_{m} \mid \\
& =\left|\partial_{m} f(x)(L) \int_{B^{*}} \lambda\left(B^{z}\right) d \lambda_{m-1}(z)-(L) \int_{B^{*}} f^{*}(z) d \lambda_{m-1}(z)\right| \\
& \leqslant(L) \int_{B^{*}}\left|\partial_{m} f(x) \lambda\left(B^{z}\right)-f^{*}(z)\right| d \lambda_{m-1}(z) \\
& <\varepsilon^{2} d(B)(L) \int_{B^{*}} 2 k(z) d \lambda_{m-1}(z) \\
& =\varepsilon^{2} d(B)(L) \int_{B^{.}}\left|\nu_{m}\right| d \sigma_{B} \leqslant \varepsilon^{2} d(B)\|B\|=\varepsilon^{2} \frac{|B|}{r(B)} \leqslant \varepsilon|B|
\end{aligned}
$$

Next we prove the divergence theorem.
5.6. Theorem. Let $A \in \mathscr{A}$, and let $C, S$, and $T$ be, respectively, a countable, slight, and thin subset of $A^{-}$. Let $v$ be a bounded vector field on $A^{-}$which is continuous in $A^{-}-S$ and differentiable in $A^{\circ}-C \cup T$. Then $v \upharpoonright A^{\circ}$ is $\bar{\sigma}_{A}$-measurable, $\nabla \cdot v$ is integrable in $A$, and $\int_{A} \nabla \cdot v=(L) \int_{A} \cdot v \cdot n_{A} d \bar{\sigma}_{A}$.

Proof. By Lemma 5.4, $v \upharpoonright B^{\bullet}$ is $\bar{\sigma}_{B}$-measurable for each $B \in \mathscr{A}(A)$, and if $F(B)=(L) \int_{B} \cdot v \cdot n_{B} d \bar{\sigma}_{B}$ for $B \in \mathscr{A}(A)$, then $F$ is a continuous additive function in $A$. Choose an $\varepsilon>0$, and enumerate $C-S$ as $\left\{z_{1}, z_{2}, \ldots\right\}$. For each $z_{n}$
find a positive

$$
\delta\left(z_{n}\right) \leqslant \min \left(1, \frac{2^{-n-2} \varepsilon}{\left|\nabla \cdot v\left(z_{n}\right)\right|+1}\right)
$$

so that $\left|v(x)-v\left(z_{n}\right)\right|<2^{-n-2} \varepsilon^{2}$ whenever $x \in A^{-}$and $\left|x-z_{n}\right|<\delta\left(z_{n}\right)$. If $B \in \mathscr{A}(A), z_{n} \in B^{-}, d=d(B)<\delta\left(z_{n}\right)$, and $r(B) \geqslant \varepsilon$, then

$$
\begin{aligned}
& \left|\nabla \cdot v\left(z_{n}\right)\right| B|-F(B)| \\
& \quad \leqslant\left|\nabla \cdot v\left(z_{n}\right)\right| \cdot|B|+\left|(L) \int_{B^{\cdot}}\left[v(x)-v\left(z_{n}\right)\right] \cdot n_{B} d \bar{\sigma}_{B}(x)\right| \\
& \quad \leqslant\left|\nabla \cdot v\left(z_{n}\right)\right| d^{m}+2^{-n-2} \varepsilon^{2}\|B\|<\left|\nabla \cdot v\left(z_{n}\right)\right| \delta\left(z_{n}\right)+2^{-n-2} \frac{\varepsilon^{2}}{r(B)} \cdot \frac{|B|}{d} \\
& \quad \leqslant 2^{-n-2} \varepsilon+2^{-n-2} \varepsilon d^{m-1} \leqslant 2^{-n-1} \varepsilon .
\end{aligned}
$$

Given $x \in A^{\circ}-C \cup T$, use Lemma 5.5 to find a $\delta(x)>0$ such that $|\nabla \cdot v(x)| B|-F(B)|<\varepsilon|B| / 2(|A|+1)$ for each $B \in \mathscr{A}(A)$ with $x \in B^{-}, d(B)$ $<\delta(x)$, and $r(B) \geqslant \varepsilon$. Finally, set $\delta(x)=1$ whenever $x \in A-\left[\left(A^{\circ}-C \cup T\right)\right.$ $\cup(C-S)]$. Now let $\left\{\left(A_{1}, x_{1}\right), \ldots,\left(A_{p}, x_{p}\right)\right\}$ be in $\mathscr{P}(A-S \cup T, \varepsilon ; \delta)$. Then each $x_{i}, i=1, \ldots, p$, belongs to $(A-S \cup T)^{\circ}=A^{\circ}-S \cup T$, and hence either to $C-S$ or to $A^{\circ}-C \cup T$. It follows that

$$
\sum_{i=1}^{p}\left|\nabla \cdot v\left(x_{i}\right)\right| A_{i}\left|-F\left(A_{i}\right)\right|<\varepsilon \sum_{n=1}^{\infty} 2^{-n-1}+\frac{\varepsilon}{2(|A|+1)} \sum_{i=1}^{p}\left|A_{i}\right| \leqslant \varepsilon
$$

Since $S \cup T$ is a thin set, $F=\int \nabla \cdot v$ and the theorem is proved.
Letting $C=S=T=\varnothing$ in Theorem 5.6 yields the following corollary.
5.7. Corollary. Let $A \in \mathscr{A}$, and let $v$ be a continuous vector field in $A^{-}$which is differentiable in $A^{\circ}$. Then $\nabla \cdot v$ is integrable in $A$ and $\int_{A} \nabla \cdot v=(L) \int_{A} \cdot v \cdot$ $n_{A} d \sigma_{A}$.

The estimates in the proof of Theorem 5.5 suggest that the boundedness of the vector field $v$ may not be essential. Specifically, it appears that at countably many points $z \in Z^{-}$a mere growth condition such as $v(x)=o\left(|x-z|^{1-m}\right)$ when $x$ approaches $z$ should be sufficient (see [17, Theorem 5.4]). However, the next example shows that such a generalization is not possible for the integral defined in this paper (see Remark 7.5).
5.8. EXAMPLE. Let $x_{0}=(0,0)$, and set $v\left(x_{0}\right)=x_{0}$ and $v(x)=$ $x|x|^{-3 / 2} \cos \pi|x|^{-1 / 2}$ for each $x \in \mathbf{R}^{2}-\left\{x_{0}\right\}$. Then $v$ is a vector field in $\mathbf{R}^{2}$
which is continuously differentiable in $\mathbf{R}^{2}-\left\{x_{0}\right\}$, and $\left|v(x)-v\left(x_{0}\right)\right|=$ $O\left(\left|x-x_{0}\right|^{-1 / 2}\right)$ as $x$ approaches $x_{0}$. If $A_{k}=\left\{x \in \mathbf{R}^{2}:(2 k+1)^{-2} \leqslant|x| \leqslant\right.$ $\left.(2 k)^{-2}\right\}, k=1,2, \ldots$, then it is easy to see that $A=\bigcup_{k=1}^{\infty} A_{k}$ belongs to $\mathscr{A}$. For $x \in \mathbf{R}^{2}-\left\{x_{0}\right\}$, we have

$$
\nabla \cdot v(x)=\frac{1}{2}|x|^{-3 / 2} \cos \pi|x|^{-1 / 2}+\frac{\pi}{2}|x|^{-2} \sin \pi|x|^{-1 / 2} .
$$

Using the polar coordinates, we see that the finite integral

$$
\text { (L) } \int_{A}|x|^{-3 / 2} \cos \pi|x|^{-1 / 2} d \lambda_{2}(x)
$$

exists, and that

$$
(L) \int_{A}|x|^{-2} \sin \pi|x|^{-1 / 2} d \lambda_{2}(x)=+\infty
$$

Since $\sin \pi|x|^{-1 / 2} \geqslant 0$ for each $x \in A$, it follows from Propositions 3.5 and 4.4 that $\int_{A} \nabla \cdot v$ does not exist. An easy calculation also reveals that $(L) \int_{A} \cdot v$. $n_{A} d \sigma_{A}=+\infty$.

## 6. Change of variable

Let $E \subset \mathbf{R}^{m}$ and $\Phi: E \rightarrow \mathbf{R}^{m}$. We say that $\Phi$ is a regular map of $E$ if it can be extended to a $C^{1}$-diffeomorphism (also denoted by $\Phi$ ) of an open set $U \subset \mathbf{R}^{m}$ containing $E^{-}$. For a regular map $\Phi$, we denote by det $\Phi$ the determinant of its Jacobi matrix. If $\Phi$ is regular, then $\Phi$ and $\operatorname{det} \Phi$ are defined uniquely in $E^{-}$and $\left(E^{\circ}\right)^{-}$, respectively, and they both extend continuously to a neighborhood of $E^{-}$.
-6.1. Lemma. Let $\Phi$ be a regular map of $E \subset \mathbf{R}^{m}$. If $E$ is, respectively, thin or slight, then so is $\Phi(E)$.

Proof. Let $E$ be compact. There is an open set $U \subset \mathbf{R}^{m}$ containing $E$, and positive numbers $\alpha$ and $\beta$ such that $|x-y| \leqslant \alpha|\Phi(x)-\Phi(y)|$ for each $x, y \in U$ and $|\Phi(B)| \leqslant \beta|B|$ for each $B \subset U$. Since $\Phi(E)$ is a compact subset of the open set $\Phi(U)$, there is an $\eta_{0}>0$ with $U\left(\Phi(E), \eta_{0}\right) \subset \Phi(U)$. Now if $0<\eta \leqslant \eta_{0}$, then $U(\Phi(E), \eta) \subset \Phi[U(E, \alpha \eta)]$ and consequently $|U(\Phi(E), \eta)| \leqslant \beta|U(E, \alpha \eta)|$. Thus as $\eta$ approaches zero, we have $|U(\Phi(E), \eta)| \leqslant \beta O(\alpha \eta)$ or $|U(\Phi(E), \eta)| \leqslant \beta o(\alpha \eta)$ according to whether $E$ is thin or slight, respectively. The lemma follows.
6.2. Corollary. If $\Phi$ is a regular map of an admissible set $A$, then $\Phi(A)$ is admissible.

Indeed, since $[\Phi(A)]^{\circ}=\Phi\left(A^{*}\right)$.
6.3. Theorem. Let $\Phi$ be a regular map of an admissible set $A$, and let $f$ be an integrable function in $\Phi(A)$. Then $f \circ \Phi \cdot|\operatorname{det} \Phi|$ is integrable in $A$ and $\int_{A} f \circ \Phi$. $|\operatorname{det} \Phi|=\int_{\Phi(A)} f$.

Proof. There is an open set $U \subset \mathbf{R}^{m}$ containing $A^{-}$and positive numbers $\alpha, \beta$, $\beta^{*}$, and $\gamma$ such that the following conditions are satisfied:
(i) $\gamma \geqslant 1$ and $\beta^{*} / \alpha \gamma \leqslant 1 / 2$;
(ii) $|\Phi(x)-\Phi(y)| \leqslant \alpha|x-y|$ for each $x, y \in U$;
(iii) $\beta^{*}|B| \leqslant|\Phi(B)| \leqslant \beta|B|$ for each $B \subset U$;
(iv) $\|\Phi(B)\| \leqslant \gamma\|B\|$ for each $B \in \mathscr{A}(A)$.

Here condition (iv) follows from [9, Theorem 52]. Let $F=\int f$, and let $T_{F}$ and $S_{F}$ be the thin and slight sets associated to $F$ by Definitions 3.1 and 2.3 , respectively. Applying Lemma 6.1 to $\Phi^{-1}$, we see that the sets $T=\Phi^{-1}\left(T_{F}\right)$ and $S=\Phi^{-1}\left(S_{F}\right)$ are thin and slight subsets of $A^{-}$, respectively. We define an additive function $G$ in $A$ by setting $G(B)=F[\Phi(B)]$ for each $B \in \mathscr{A}(A)$.

Choose $\varepsilon_{S}>0$, and find an $\eta_{S}>0$ so that $|F(C)|<\varepsilon_{S}$ whenever $C \in$ $\mathscr{A}\left[\Phi(A) \cap U\left(S_{F}, \eta_{S}\right)\right]$ and $\|C\|<\eta_{S}$. Letting $\eta_{S}^{*}=\eta_{S} / \max (\alpha, \gamma)$ and using (ii) and (iv), it is easy to verify that $|G(B)|<\varepsilon_{S}$ for each $B \in \mathscr{A}\left[A \cap U\left(S, \eta_{S}^{*}\right)\right]$ with $\|B\|<\eta_{S}^{*}$. Now let $E \subset A^{-}-S$ be a thin set, and let $\varepsilon_{E}>0$. By Lemma 6.1, $\Phi(E)$ is a thin subset of $[\Phi(A)]^{-}-S_{F}$, and hence there is an $\eta_{E}>0$ such that $|F(C)|<\varepsilon_{E} / \gamma$ for each $C \in \mathscr{A}\left[\Phi(A) \cap U\left(\Phi(E), \alpha \eta_{E}\right)\right]$ with $\|C\|<\gamma / \varepsilon_{E}$. From this, (i), (ii), and (iv) it is easy to see that $|G(B)|<\varepsilon_{E}$ for each $B \in$ $\mathscr{A}\left[A \cap U\left(E, \eta_{E}\right)\right]$ with $\|B\|<1 / \varepsilon_{E}$. It follows that $G$ is continuous, and we show next that $G=\int f \circ \Phi \cdot|\operatorname{det} \Phi|$.

Choose $\varepsilon>0$, and for each $x \in A$ find $\eta(x)>0$ so that

$$
||\Phi(B)|-|\operatorname{det} \Phi(x)| \cdot| B \|<\frac{\varepsilon|B|}{2(f[\Phi(x)]+1)(|A|+1)}
$$

for each measurable $B \subset A$ with $x \in B^{-}$and $d(B)<\eta(x)$. This is possible, for $\operatorname{det} \Phi$ is continuous and $|\Phi(B)|=(L) \int_{B}|\operatorname{det} \Phi|$. Letting $\varepsilon^{*}=\varepsilon \beta^{*} / \alpha \gamma$, there is a $\delta^{*}: \Phi(A) \rightarrow R_{+}$such that $\sum_{i=1}^{p}\left|f\left(y_{i}\right)\right| C_{i}\left|-F\left(C_{i}\right)\right|<\varepsilon^{*}$ for each $\left\{\left(C_{1}, y_{1}\right), \ldots,\left(C_{p}, y_{p}\right)\right\}$ in $\mathscr{P}\left(\Phi(A)-T_{F}, \varepsilon^{*} ; \delta^{*}\right)$. Now set $\delta(x)=$ $\min \left(\eta(x), \delta^{*}[\Phi(x)] / \alpha\right)$ for each $x \in A$, and choose a partition
$\left\{\left(B_{1}, x_{1}\right), \ldots,\left(B_{p}, x_{p}\right)\right\}$ in $\mathscr{P}(A-T, \varepsilon ; \delta)$. If $C_{i}=\Phi\left(B_{i}\right)$ and $y_{i}=\Phi\left(x_{i}\right), i=$ $1, \ldots, p$, then $\left\{\left(C_{1}, y_{1}\right), \ldots,\left(C_{p}, y_{p}\right)\right\}$ belongs to $\mathscr{P}\left(\Phi(A)-T_{F}, \varepsilon^{*} ; \delta^{*}\right)$; for $T_{F}=$ $\Phi(T), d\left(C_{i}\right) \leqslant \alpha d\left(B_{i}\right)<\alpha \delta\left(x_{i}\right) \leqslant \delta^{*}\left(y_{i}\right)$, and $r\left(C_{i}\right) \geqslant\left(\beta^{*} / \alpha \gamma\right) r\left(B_{i}\right) \geqslant \varepsilon^{*}, i=$ $1, \ldots, p$. As $d\left(B_{i}\right)<\eta\left(x_{i}\right)$ for $i=1, \ldots, p$, we have

$$
\begin{aligned}
& \sum_{i=1}^{p}\left|f\left[\Phi\left(x_{i}\right)\right] \cdot\right| \operatorname{det} \Phi\left(x_{i}\right)|\cdot| B_{i}\left|-G\left(B_{i}\right)\right| \\
& \leqslant \sum_{i=1}^{p}\left|f\left[\Phi\left(x_{i}\right)\right]\right| \cdot| | \operatorname{det} \Phi\left(x_{i}\right)|\cdot| B_{i}\left|-\left|\Phi\left(B_{i}\right)\right|\right| \\
&+\sum_{i=1}^{p}\left|f\left(y_{i}\right)\right| C_{i}\left|-F\left(C_{i}\right)\right| \\
&<\frac{\varepsilon}{2(|A|+1)} \sum_{i=1}^{p}\left|B_{i}\right|+\varepsilon^{*}<\frac{\varepsilon}{2}+\varepsilon^{*}
\end{aligned}
$$

Since $\varepsilon^{*} \leqslant \varepsilon / 2$ by (i), the theorem is proved.

## 7. Integration on manifolds and the Stokes theorem

A certain familiarity with manifolds and exterior forms is assumed in this section. For the notation and terminology, we refer to [20] and [2].

Throughout, $\mathbf{R}^{m}$ and $\mathbf{R}^{m-1}$ are oriented by the standard volume elements $\mu^{m}=d \xi_{1} \wedge \cdots \wedge d \xi_{m}$ and $\mu^{m-1}=d \xi_{1} \wedge \cdots \wedge d \xi_{m-1}$, respectively. By a manifold we mean a paracompact $C^{1}$ manifold $M$ with boundary $\partial M$, which may be empty. If a manifold $M$ is oriented, then $\partial M$ is always given the induced orientation. For a subset $E$ of a manifold $M$, we give the symbols $E^{-}, E^{\circ}$, and $E^{\cdot}$ the obvious meaning, and we let $\partial E=E^{\cdot} \cup\left(E^{-} \cap \partial M\right)$.

Let $M$ be an $m$-dimensional manifold. We say that $E \subset M$ is an elementary admissible set if there is a chart $(U, \Phi)$ such that $E^{-} \subset U$ and $\Phi(E) \in \mathscr{A}$. Using Corollary 6.2, it is easy to show that the family $\mathscr{E}^{M}$ of all elementary admissible subsets of $M$ is closed with respect to differences. A finite union of elements of $\mathscr{E}^{M}$ is called an admissible subset of $M$. Slight and thin subsets of $M$ are defined similarly. The family $\mathscr{A}^{M}$ of all admissible subsets of $M$ is a ring which contains all compact submanifolds of $M$. It is not difficult to see that $A \in \mathscr{A}^{M}$ if and only if $A^{-}$is compact and $\partial A$ is a thin subset of $M$.

Let $A$ be an admissible subset of an oriented manifold $M$, and let $\theta$ be an $m$-form on $A$. If $A \in \mathscr{E}^{M}$, then there is a positively oriented chart $(U, \Phi)$ with $A^{-} \subset U$, and a unique function $g$ on $\Phi(A)$ such that $\left(\Phi^{-1}\right)^{*} \theta=g \mu^{m}$. We let $\int_{A} \theta=\int_{\Phi(A)} g$ whenever the integral on the right exists. In general, $A$ is a disjoint
union of elementary admissible sets $A_{1}, \ldots, A_{k}$, and we let $\int_{A} \theta=\sum_{i=1}^{k} \int_{A i} \theta$ whenever the integrals on the right exist. It follows from Theorem 6.3 and Proposition 3.6 that the value of $\int_{A} \theta$ does not depend on the choice of $(U, \Phi)$ and $A_{1}, \ldots, A_{k}$. We say that $\theta$ is integrable in $A$ whenever $\int_{A} \theta$ is defined.

Let $N$ be an oriented ( $m-1$ )-dimensional manifold. In complete analogy to the previous paragraphs, we define families $\mathscr{M}^{N}$ and $\mathscr{N}^{N}$ of measurable and null subsets of $N$, respectively, and $(L) \int_{A} \omega$ for a bounded $\mathscr{M}^{N}$-measurable ( $m-1$ )-form $\omega$ on $A \in \mathscr{M}^{N}$. We note that according to our definition the family $\mathscr{M}^{N}$ is a ring, which becomes a $\sigma$-algebra whenever $N$ is compact. Similarly, the family $\mathscr{N}^{N}$ is an ideal in $\mathscr{M}^{N}$, which becomes a $\sigma$-ideal whenever $N$ is compact. We omit further details as this is a well-known process of Lebesgue integration on manifolds.

Let $M$ be an $m$-dimensional manifold, and let $\omega$ be an ( $m-1$ )-form on a set $E \subset M$. We denote by $d \omega$ any $m$-form on $E$ such that $d \omega(x)$ is the exterior derivative of $\omega$ at each $x \in E^{\circ}$ at which $\omega$ is differentiable.

For $k=m-1, m$, let

$$
\Delta^{k}=\left\{\left(t_{1}, \ldots, t_{k}\right) \in \mathbf{R}^{k}: \sum_{i=1}^{k} t_{i} \leqslant 1 \text { and } t_{i} \geqslant 0, i=1, \ldots, k\right\} .
$$

To simplify the notation, we write $\mu$ and $\Delta$ instead of $\mu^{m}$ and $\Delta^{m}$, respectively. We define maps $\varepsilon_{j}: \mathbf{R}^{m-1} \rightarrow \mathbf{R}^{m}$ by setting

$$
\varepsilon_{0}\left(\zeta_{1}, \ldots, \zeta_{m-1}\right)=\left(1-\sum_{i=1}^{m-1} \zeta_{i}, \zeta_{1}, \ldots, \zeta_{m-1}\right)
$$

and

$$
\varepsilon_{j}\left(\zeta_{1}, \ldots, \zeta_{m-1}\right)=\left(\zeta_{1}, \ldots, \zeta_{j-1}, 0, \zeta_{j}, \ldots, \zeta_{m-1}\right)
$$

for $j=1, \ldots, m$. If $\Delta_{j}=\varepsilon_{j}\left(\Delta^{m-1}\right)$, then $\Delta^{\cdot}=\bigcup_{j=0}^{m} \Delta_{j}$. Let $n_{0}=$ $\left(1 / m^{1 / 2}, \ldots, 1 / m^{1 / 2}\right)$, and for $j=1, \ldots, m$, let $n_{j}=(0, \ldots,-1, \ldots, 0)$ where -1 is the $j$ th coordinate. Clearly $\Delta \in \mathscr{A}$, and using Lemma 5.1 and the symmetry of $\Delta$, it is not difficult to see that $n_{\Delta}(x)=n_{j}$ for $\sigma_{\Delta}$-almost all $x \in \Delta_{j}, j=0, \ldots, m$.

Given a vector field $v$ on a set $E \subset \mathbf{R}^{m}$, we define an ( $m-1$ )-form $\mu_{v}$ on $E$ by setting $\left\langle\mu_{v}(x) \mid v_{1}, \ldots, v_{m-1}\right\rangle=\left\langle\mu \mid v(x), v_{1}, \ldots, v_{m-1}\right\rangle$ for each $x \in E$ and each $v_{1}, \ldots, v_{m-1}$ in $\mathbf{R}^{m}$. If $v=\left(a_{1}, \ldots, a_{m}\right)$, then

$$
\mu_{v}=\sum_{i=1}^{m}(-1)^{i+1} a_{i} d \xi_{1} \wedge \ldots \widehat{d \xi_{i}} \ldots \wedge d \xi_{m}
$$

where $d \xi_{1} \wedge \ldots \widehat{d \xi_{i}} \ldots \wedge d \xi_{m}$ stands for $d \xi_{1} \wedge \ldots \wedge d \xi_{i-1} \wedge d \xi_{i+1} \wedge \ldots \wedge d \xi_{m}$. Thus for each ( $m-1$ )-form $\omega$ on $E \subset \mathbf{R}^{m}$ there is a unique vector field $v$ on $E$ such that $\omega=\mu_{v}$. Clearly, $\omega=\mu_{v}$ is, respectively, bounded, continuous, or differentiable if and only if $v$ is.
7.1. Lemma. Let $n$ be a unit vector in $\mathbf{R}^{m}$, and let $N=\left\{x \in \mathbf{R}^{m}: x \cdot n=0\right\}$. If $v \in \mathbf{R}^{m}$, then $\mu_{v}=(v \cdot n) \mu_{n}$ on $N$.

Proof. Let $n=\left(\nu_{1}, \ldots, \nu_{m}\right)$ and $v=\left(a_{1}, \ldots, a_{m}\right)$. On $N$, we have $\sum_{i=1}^{m} \nu_{i} \xi_{i}=0$, and by symmetry, we may assume that $\nu_{1} \neq 0$. Then

$$
\begin{aligned}
\mu_{v}= & \sum_{i=1}^{m}(-1)^{i+1} a_{i} d \xi_{1} \wedge \ldots{\widehat{d \xi_{i}}}_{i} \ldots \wedge d \xi_{m} \\
= & a_{1} d \xi_{2} \wedge \ldots \wedge d \xi_{m}+\sum_{i=2}^{m}(-1)^{i+1} a_{i}\left(-\sum_{j=2}^{m} \frac{\nu_{j}}{\nu_{1}} d \xi_{j}\right) \\
& \wedge d \xi_{2} \wedge \ldots \widehat{d \xi}_{i} \ldots \wedge d \xi_{m} \\
= & \frac{a_{1} \nu_{1}}{\nu_{1}} d \xi_{2} \wedge \ldots \wedge d \xi_{m}+\sum_{i=2}^{m}(-1)^{i} \frac{a_{i} \nu_{i}}{\nu_{1}} d \xi_{i} \wedge d \xi_{2} \wedge \ldots \widehat{d \xi}_{i} \ldots \wedge d \xi_{m} \\
= & \frac{v \cdot n}{\nu_{1}} d \xi_{2} \wedge \ldots \wedge d \xi_{m}
\end{aligned}
$$

and similarly, $\mu_{n}=\left(1 / \nu_{1}\right) d \xi_{2} \wedge \ldots \wedge d \xi_{m}$. The lemma follows.
Throughout, we orient the linear submanifold $N_{j}=\varepsilon_{j}\left(\mathbf{R}^{m-1}\right)$ of $\mathbf{R}^{m}$ by the form $\mu_{n_{j}}, j=0, \ldots, m$.
7.2. Lemma. Let $C, S$, and $T$ be, respectively, a countable, slight, and thin subset of $\Delta$, and let $\omega$ be a bounded $(m-1)$-form on $\Delta$ which is continuous in $\Delta-S$ and differentiable in $\Delta^{\circ}-C \cup T$. Then $\omega \upharpoonright \Delta_{j}$ is $\mathscr{M}^{N_{j}}$ measurable for $j=0, \ldots, m, d \omega$ is integrable in $\Delta$, and

$$
\int_{\Delta} d \omega=\sum_{j=0}^{m}(L) \int_{\Delta_{j}} \omega
$$

Proof. Let $\omega=\mu_{v}$ where $v=\left(a_{1}, \ldots, a_{m}\right)$. Then

$$
d \omega(x)=\sum_{i=1}^{m}(-1)^{i+1} \partial_{i} a_{i}(x) d \xi_{i} \wedge d \xi_{1} \wedge \ldots \widehat{d \xi}_{i} \ldots \wedge d \xi_{m}=[\nabla \cdot v(x)] \mu
$$

for each $x \in \Delta^{\circ}-C \cup T$. It follows from Theorem 5.6 that $v \upharpoonright \Delta^{\circ}$ is $\bar{\sigma}_{\Delta}$-measurable, $d \omega$ is integrable in $\Delta$, and $\int_{\Delta} d \omega=(L) \int_{\Delta} \cdot v \cdot n_{\Delta} d \bar{\sigma}_{\Delta}$.

It is easy to check that $\varepsilon_{0}^{*} \mu_{n_{0}}=m^{1 / 2} \mu^{m-1}$ and $\varepsilon_{j}^{*} \mu_{n_{i}}=(-1)^{j} \mu^{m-1}$ for $j=$ $1, \ldots, m$. Thus ( $N_{j},(-1)^{j} \varepsilon_{j}^{-1}$ ) is a positively oriented chart of $N_{j}, j=0, \ldots, m$. If $\chi$ is a characteristic function of $\Delta_{0}$, then in the notation of Lemma 5.1, we have

$$
a_{i}\left(\zeta_{1}, \ldots, \zeta_{m-1}, 1-\sum_{k=1}^{m-1} \zeta_{k}\right)=\left(a_{i} \chi\right)^{*}\left(\zeta_{1}, \ldots, \zeta_{m-1}\right), \quad i=1, \ldots, m
$$

It follows from Lemma 5.1(iv) that $v \circ \varepsilon_{0}$ is $\lambda_{m-1}$-measurable in $\Delta^{m-1}$, and similarly we establish that so are $v \circ \varepsilon_{1}, \ldots, v \circ \varepsilon_{m}$. Thus $\omega \upharpoonright \Delta_{j}$ is $\mathscr{M}^{N_{j} \text {-measurable }}$ for $j=0, \ldots, m$.

By Lemma 7.1,

$$
\begin{aligned}
\left\langle\varepsilon_{j}^{*}[\omega(x)] \mid u_{1}, \ldots, u_{m-1}\right\rangle & =\left\langle\mu_{v}(x) \mid \varepsilon_{j *} u_{1}, \ldots, \varepsilon_{j *} u_{m-1}\right\rangle \\
& =\left[v(x) \cdot n_{j}\right] \cdot\left\langle\mu_{n_{j}} \mid \varepsilon_{j *} u_{1}, \ldots, \varepsilon_{j *} u_{m-1}\right\rangle \\
& =\left[v(x) \cdot n_{j}\right] \cdot\left\langle\varepsilon_{j}^{*} \mu_{n} \mid u_{1}, \ldots, u_{m-1}\right\rangle
\end{aligned}
$$

for each $x \in \Delta_{j}, j=0, \ldots, m$, and each $u_{1}, \ldots, u_{m-1}$ in $\mathbf{R}^{m-1}$. Thus on $\Delta^{m-1}$, we have $\varepsilon_{0}^{*} \omega=m^{1 / 2}\left(v \circ \varepsilon_{0} \cdot n_{0}\right) \mu^{m-1}$, and $(-1)^{j} \varepsilon_{j}^{*} \omega=\left(v \circ \varepsilon_{j} \cdot n_{j}\right) \mu^{m-1}$ for $j=$ $1, \ldots, m$. Consequently,

$$
\begin{aligned}
\sum_{j=0}^{m}(L) \int_{\Delta_{j}} \omega & =m^{1 / 2}(L) \int_{\Delta^{m-1}} v \circ \varepsilon_{0} \cdot n_{0} d \lambda_{m-1}+\sum_{i=1}^{m}(L) \int_{\Delta^{m-1}} v \circ \varepsilon_{i} \cdot n_{i} d \lambda_{m-1} \\
& =\sum_{i=1}^{m}(L) \int_{\Delta^{m-1}}\left(a_{i} \circ \varepsilon_{0}-a_{i} \circ \varepsilon_{i}\right) d \lambda_{m-1}
\end{aligned}
$$

If $n_{\Delta}=\left(\nu_{1}, \ldots, \nu_{m}\right)$, then by Lemma 5.1(iv),
$(L) \int_{\Delta^{m-1}}\left(a_{m} \circ \varepsilon_{0}-a_{m} \circ \varepsilon_{m}\right) d \lambda_{m-1}$

$$
\begin{aligned}
= & (L) \int_{\Delta^{m-1}} a_{m}\left(1-\sum_{k=1}^{m-1} \zeta_{k}, \zeta_{1}, \ldots, \zeta_{m-1}\right) d \lambda_{m-1}\left(\zeta_{1}, \ldots, \zeta_{m-1}\right) \\
& -(L) \int_{\Delta^{m-1}} a_{m}(z, 0) d \lambda_{m-1}(z)
\end{aligned}
$$

$$
=(L) \int_{\Delta^{m-1}} a_{m}\left(\zeta_{1}, \ldots, \zeta_{m-1}, 1-\sum_{k=1}^{m-1} \zeta_{k}\right) d \lambda_{m-1}\left(\zeta_{1}, \ldots, \zeta_{m-1}\right)
$$

$$
-(L) \int_{\Delta^{m-1}} a_{m}(z, 0) d \lambda_{m-1}(z)
$$

$$
=(L) \int_{\Delta^{*}}\left(a_{m}\right)^{*} d \lambda_{m-1}
$$

$$
=(L) \int_{\Delta^{*}} a_{m} \nu_{m} d \bar{\sigma}_{\Delta}
$$

here $\Delta^{*}$ and $\left(a_{m}\right)^{*}$ are the symbols introduced in Lemma 5.1. Now by symmetry,

$$
(L) \int_{\Delta^{m-1}}\left(a_{i} \circ \varepsilon_{0}-a_{i} \circ \varepsilon_{i}\right) d \lambda_{m-1}=(L) \int_{\Delta^{*}} a_{i} \nu_{i} d \bar{\sigma}_{\Delta}
$$

for $i=1, \ldots, m-1$. It follows that

$$
\sum_{j=0}^{m}(L) \int_{\Delta_{j}} \omega=(L) \int_{\Delta^{*}} v \cdot n_{\Delta} d \bar{\sigma}_{\Delta}
$$

and the lemma is proved.
Now we can prove the Stokes theorem.
7.3. Theorem. Let $M$ be an m-dimensional compact oriented manifold, and let $C$, $S$, and $T$ be, respectively, a countable, slight and thin subset of $M$. Let $\omega$ be a bounded ( $m-1$ )-form on $M$ which is continuous in $M-S$ and differentiable in $M-(\partial M \cup C \cup T)$. Then $\omega \upharpoonright \partial M$ is $\mathscr{M}^{\partial M}$-measurable, d $\omega$ is integrable in $M$, and $\int_{M} d \omega=(L) \int_{\partial M} \omega$.

Proof. By [21, Chapter IV, Theorem 12A], the manifold $M$ has a finite $C^{1}$ triangulation $\mathscr{T}$. For each $A \in \mathscr{T}$, there is a positively oriented chart $\left(U_{A}, \Phi_{A}\right)$ such that $A \subset U_{A}$ and $\Phi_{A}(A)=\Delta$. Given $A \in \mathscr{T}$ and $j=0, \ldots, m$, we denote by $A_{j}$ the $(m-1)$-dimensional submanifold $\Phi_{A}^{-1}\left(\Delta_{j}^{\circ}\right)$ of $M$ oriented by the form $\Phi_{A}^{*} \mu_{n_{j}}$. In particular, if $A_{j} \subset \partial M$ the the orientation of $A_{j}$ is induced by that of $\partial M$. If $A \in \mathscr{T}$, then $A$ is an elementary admissible set, and it follows from


$$
\int_{A} d \omega=\int_{\Delta}\left(\Phi_{A}^{-1}\right)^{*} d \omega=\int_{\Delta} d\left(\Phi_{A}^{-1}\right)^{*} \omega=\sum_{j=0}^{m}(L) \int_{\Delta_{j}}\left(\Phi_{A}^{-1}\right)^{*} \omega=\sum_{j=0}^{m}(L) \int_{A_{j}} \omega .
$$

If $A$ and $B$ are in $\mathscr{T}$ and $\Phi_{A}^{-1}\left(\Delta_{j}\right)=\Phi_{B}^{-1}\left(\Delta_{i}\right)$ for some integers $j$ and $i$, then $(L) \int_{\mathcal{A}_{j}} \omega=-(L) \int_{B_{i}} \omega$; for the oriented manifolds $A_{j}$ and $B_{i}$ have opposite orientations. Since $M=\bigcup \mathscr{T}$ and $\partial M=\bigcup\left\{A_{j}^{-}: A \in \mathscr{T}, A_{j} \subset \partial M\right\}$, we see that $d \omega$ is integrable in $M$ and $\omega \upharpoonright \partial M$ is $\mathscr{M}^{\partial M}$-measurable. As the integrals of $d \omega$ and $\omega$ vanish on the overlaps, we have

$$
\int_{M} d \omega=\sum_{A \in \mathscr{F}} \int_{A} d \omega=\sum_{A \in \mathscr{F}} \sum_{j=0}^{m}(L) \int_{A_{j}} \omega=\sum_{A \in \mathscr{F}} \sum_{A_{j} \subset \partial M}(L) \int_{A_{j}} \omega=(L) \int_{\partial M} \omega,
$$

and the theorem is proved.
Letting $C=S=T=\varnothing$ in Theorem 7.3, we obtain the following corollary.
7.4. Corollary. Let $M$ be an m-dimensional compact oriented manifold, and let $\omega$ be a continuous ( $m-1$ )-form on $M$ which is differentiable in $M-\partial M$. Then $d \omega$ is integrable in $M$ and $\int_{M} d \omega=(L) \int_{\partial M} \omega$.
7.5. Remark. At this point it appears appropriate to mention that our choice of admissible sets is only one of many. For example, replacing $\mathscr{A}$ by a smaller family of all diffeomorphic images of simplices or convex linear cells, leads to an integration process very similar to ours. Whether there is an optimal family of admissible sets is unclear.

## 8. The one-dimensional case

Throughout this section, we assume that $m=1$. Our aim is to compare the integral of Definition 3.1 to the classical Denjoy-Perron integral (see [19, Chapter VI, Section 6, and Chapter VIII]), or equivalently to the Henstock-Kurzweil integral defined in [10]. For a simple proof of the fact that the Denjoy-Perron and Henstock-Kurzweil integrals are indeed equivalent we refer to [7] or [15, Appen$\operatorname{dix} B]$.

### 8.1. Lemma. Slight sets are empty, and thin sets are finite.

Proof. The first part of the lemma is obvious, because $|U(E, \eta)| \geqslant 2 \eta$ for each nonempty set $E \subset \mathbf{R}$ and each $\eta>0$. Let $E \subset \mathbf{R}$ be infinite, and let $\left\{x_{1}, x_{2}, \ldots\right\}$ be an infinite sequence of distinct points from $E$. Given an integer $n \geqslant 1$, find $\eta>0$ so that $2 \eta<\left|x_{i}-x_{j}\right|$ for all distinct $i, j=1, \ldots, n$. Then $|U(E, \eta)| \geqslant 2 n \eta$, and it follows that $E$ is not thin.
8.2. Corollary. Modulo finite sets, each admissible set is a finite disjoint union of nondegenerate compact intervals.

For $A, B \subset R$, we write $A \sim B$ whenever the symmetric difference $(A-B) \cup$ ( $B-A$ ) is finite.

Let $A \in \mathscr{A}$. By Corollary 8.2, there is an integer $k>0$ and real numbers $a_{1}<b_{1}<\cdots<a_{k}<b_{k}$ such that $A \sim \bigcup_{i=1}^{k}\left(a_{i}, b_{i}\right)$. It follows from Lemma 5.1 that $\sigma_{A}$ is the counting measure concentrated on $\left\{a_{1}, \ldots, a_{k}, b_{1}, \ldots, b_{k}\right\}$, and that $n_{A}\left(b_{i}\right)=-n_{A}\left(a_{i}\right)=1$ for $i=1, \ldots, k$. If $F$ is a function on $A$; then extending the notation of Lemma 5.1, we let $F^{*}(A)=\sum_{i=1}^{k}\left[F\left(b_{i}\right)-F\left(a_{i}\right)\right]$. As $F$ is also a vector field on $A^{*}$, we have $F^{*}(A)=(L) \int_{A} \cdot F \cdot n_{A} d \sigma_{A}$. Now if $F$ is a function on $A^{-}$, then the map $F^{*}$ given by $B \mapsto F^{*}(B)$ for $B \in \mathscr{A}(A)$ is an additive function in $A$. Moreover, $F^{*}$ is continuous if and only if $F$ is; for the continuity of $F^{*}$ clearly implies that of $F$, and the converse follows from Lemma 5.4.

If $f$ is an integrable function in $[a, b]$, we write $\int_{a}^{b} f$ instead of $\int_{[a, b]} f$. From Theorem 5.6, we obtain immediately the following proposition.
8.3. Proposition. Let $f$ and $F$ be functions on $[a, b]$ such that $F$ is continuous and $F^{\prime}(x)=f(x)$ for all but countably many $x \in(a, b)$. Then $f$ is integrable in $[a, b]$ and $\int_{a}^{b} f=F(b)-F(a)$.

Next we show that the integral is closed with respect to the formation of improper integrals.
8.4. Proposition. Let $f$ be a function on $[a, b]$ which is integrable in $[a, x]$ for each $x \in[a, b)$, and let a finite $\lim _{x \rightarrow b-} \int_{a}^{x} f=I$ exist. Then $f$ is integrable in $[a, b]$ and $\int_{a}^{b} f=I$.

Proof. If $F(x)=\int_{a}^{b} f$ for $x \in[a, b)$ and $F(b)=I$, then $F$ is continuous in [a,b], and by Lemma 5.4, so is $F^{*}$. Choose an $\varepsilon>0$ and a sequence $\left\{b_{n}\right\}$ such that $a=b_{1}<b_{2}<\cdots<b$ and $\lim b_{n}=b$. Now $F^{*}=\int f$ in each $\left[b_{n}, b_{n+1}\right]$, $n=1,2, \ldots$ By extending the sequence $\left\{b_{n}\right\}$, we may assume that $\left\{b_{n}, b_{n+1}\right\}$ is a thin subset of $\left[b_{n}, b_{n+1}\right]$ associated to $F^{*}$ by Definition 3.1. Choose an $\varepsilon>0$, and for $n=1,2, \ldots$, find $\delta_{n}:\left[b_{n}, b_{n+1}\right] \rightarrow \mathbf{R}_{+}$so that $\sum_{i=1}^{q}\left|f\left(y_{i}\right)\right| B_{i}\left|-F^{*}\left(B_{i}\right)\right|<$ $\varepsilon / 2^{n+1}$ for each $\left\{\left(B_{1}, y_{1}\right), \ldots,\left(B_{q}, y_{q}\right)\right\}$ in $\mathscr{P}\left(\left[b_{n}, b_{n+1}\right], \varepsilon ; \delta_{n}\right)$. There is a $\delta:[a, b]$ $\rightarrow \mathbf{R}_{+}$such that for $n=1,2, \ldots$ the following conditions are satisfied:
(i) $\delta(x)=\min \left(\delta_{n}(x), x-b_{n}, b_{n+1}-x\right)$ for each $x \in\left(b_{n}, b_{n+1}\right)$;
(ii) $\left|f\left(b_{n}\right)\right| \delta\left(b_{n}\right)<\varepsilon / 3 \cdot 2^{n+1}$, and $\left|F^{*}(B)\right|<\varepsilon / 3 \cdot 2^{n+1}$ for each admissible set $B \subset[a, b] \cap U\left[\left\{b_{n}\right\}, \delta\left(b_{n}\right)\right]$ with $\|B\|<1 / \varepsilon$.

Let $P=\left\{\left(A_{1}, x_{1}\right), \ldots,\left(A_{p}, x_{p}\right)\right\}$ be in $\mathscr{P}([a, b], \varepsilon ; \delta)$. By (i), we see that $\left\{\left(A_{1}, x_{i}\right) \in P: x_{i} \in\left(b_{n}, b_{n+1}\right)\right\}$ belongs to $\mathscr{P}\left(\left[b_{n}, b_{n+1}\right], \varepsilon ; \delta_{n}\right), n=1,2, \ldots$. Since

$$
\left\|A_{i}\right\|=\frac{\left|A_{i}\right|}{d\left(A_{i}\right) r\left(A_{i}\right)} \leqslant \frac{1}{\varepsilon}
$$

we obtain from (ii) that $\left|f\left(x_{i}\right)\right| A_{i}\left|-F^{*}\left(A_{i}\right)\right|<\varepsilon / 3 \cdot 2^{n}$ whenever $x_{i}=b_{n}$ for some $n=2,3, \ldots$; of course, no $x_{i}$ equals $b_{1}=a$. As the map $i \mapsto x_{i}$ is at most three-to-one, we have

$$
\begin{aligned}
& \sum_{i=1}^{p}\left|f\left(x_{i}\right)\right| A_{i}\left|-F^{*}\left(A_{i}\right)\right| \\
& \quad=\sum_{n=2}^{\infty}\left[\sum_{x_{i}=b_{n}}\left|f\left(x_{i}\right)\right| B_{i}\left|-F\left(B_{i}\right)\right|+\sum_{x_{i} \in\left(b_{n-1}, b_{n}\right)}\left|f\left(x_{i}\right)\right| A_{i}\left|-F^{*}\left(A_{i}\right)\right|\right] \\
& \quad<\varepsilon \sum_{n=2}^{\infty} 2^{-n+1}=\varepsilon
\end{aligned}
$$

and the proposition is proved.
The Denjoy-Perron integral over $[a, b]$ of a Denjoy-Perron integrable function $f$ in $[a, b]$ is denoted by $(D P) \int_{a}^{b} f$.
8.5. Proposition. If $f$ is an integrable function in $[a, b]$, then $f$ is Denjoy-Perron integrable in $[a, b]$ and $(D P) \int_{a}^{b} f=\int_{a}^{b} f$.

Proof. Let $f$ be an integrable function in $[a, b]$. We may assume that $\{a, b\}$ is a thin set $T_{F}$ associated to $F=\int f$ by Definition 3.1 ; for if $T_{F} \subset\left\{z_{0}, \ldots, z_{n}\right\}$ where $a=z_{0}<\cdots<z_{n}=b$, we can repeat the argument for each interval $\left[z_{i-1}, z_{i}\right], i=1, \ldots, n$. Given $\varepsilon>0$, there is a $\delta:[a, b] \rightarrow \mathbf{R}_{+}$satisfying the following conditions:
(i) $\delta(x) \leqslant \min (x-a, b-x)$ for each $x \in(a, b)$;
(ii) $\delta(a)=\delta(b)=\eta$ where $\max (|f(a)|,|f(b)|)<\varepsilon / 8 \eta$;
(iii) $|F(B)|<\varepsilon / 8$ for each admissible set $B \subset[a, b] \cap U(\{a, b\}, \eta)$ with $\|B\|$ $\leqslant 2$;
(iv) $\sum_{i=1}^{p}\left|f\left(x_{i}\right)\right| A_{i}\left|-F\left(A_{i}\right)\right|<\varepsilon / 2$ for each $\left\{\left(A_{1}, x_{1}\right), \ldots,\left(A_{p}, x_{p}\right)\right\}$ in $\mathscr{P}([a, b], 1 / 2 ; \delta)$.

Now let $a=t_{0}<\cdots<t_{k}=b$ and $x_{1}, \ldots, x_{k}$ be such that $t_{i-1} \leqslant x_{i} \leqslant t_{i}$ and $t_{i}-t_{i-1}<\delta\left(x_{i}\right), i=1, \ldots, k$. Set $A_{i}=\left[t_{i-1}, t_{i}\right)$ for $i=1, \ldots, k-1$, and $A_{k}=$ $\left[t_{k-1}, t_{k}\right]$. By (i), $x_{1}=a, x_{k}=b$, and $\left\{\left(A_{2}, x_{2}\right), \ldots,\left(A_{k-1}, x_{k-1}\right)\right\}$ belongs to $\mathscr{P}([a, b], 1 / 2 ; \delta)$. Thus by (ii)-(iv), we have

$$
\begin{gathered}
\left|\sum_{i=1}^{k} f\left(x_{i}\right)\left(t_{i}-t_{i-1}\right)-\int_{a}^{b} f\right| \leqslant|f(a)| \cdot\left|A_{1}\right|+|f(b)| \cdot\left|A_{k}\right|+\left|F\left(A_{1}\right)\right| \\
+\left|F\left(A_{2}\right)\right|+\sum_{i=2}^{k-1}\left|f\left(x_{i}\right)\right| A_{i}\left|-F\left(A_{i}\right)\right|<\varepsilon
\end{gathered}
$$

and the proposition follows from the aforementioned equivalence of the Henstock-Kurzwẹil and Denjoy-Perron integrals.

The next example shows that the converse of Proposition 8.5 is false. In view of Propositions 4.1, 8.3, and 8.4, this is somewhat surprising.
8.6. Example. If $J=(a, b) \subset \mathbf{R}$ is nonempty, let
$J_{+}^{n}=\left(a+2^{-2 n}|J|, a+2^{-2 n+1}|J|\right), \quad J_{-}^{n}=\left(a+2^{-2 n+1}|J|, a+2^{-2 n+2}|J|\right)$, $n=1,2, \ldots$ For $x \in \mathbf{R}$, set $f_{J}(x)=2^{2 n} / n$ if $x \in J_{+}^{n}, f_{J}(x)=2^{2 n-1} / n$ if $x \in$ $J_{-}^{n}$, and $f_{J}(x)=0$ otherwise. Since $\int_{J_{+}^{n}} f_{J}=-\int_{J_{-}^{n}} f_{J}=|J| / n, n=1,2, \ldots$, it follows from Proposition 8.4 that $f_{J}$ is integrable in $J$ with $\int_{a}^{b} f_{J}=0$, and that

$$
\sup \left\{\left|\int_{\alpha}^{\beta} f_{J}\right|: a \leqslant \alpha \leqslant \beta \leqslant b\right\}=|J|
$$

Let $D$ be the Cantor ternary set in $[0,1]$, and let $\mathscr{J}$ be the collection of all components of $[0,1]-D$. Letting $f=\sum_{J \in g} f_{J}$, it is easy to see from [19, Chapter VIII, Theorem (5.1)] that $f$ is Denjoy-Perron integrable in $[0,1]$ with ( $D P$ ) $\int_{0}^{1} f=$ 0 . We show next that the integral $\int_{0}^{1} f$ does not exist.

Proceeding towards contradiction, suppose that $f$ is integrable in [ 0,1 ], and choose a $\delta:[0,1] \rightarrow \mathbf{R}_{+}$. Using the Baire category theorem in $D$, we find an open set $U \subset \mathbf{R}$ with $D \cap U \neq \varnothing$, and an $\eta>0$ such that the set $E=\{x \in D \cap$ $U: \delta(x) \geqslant \eta)$ is dense in $D \cap U$. Select a $J=(a, b)$ in $\mathscr{J}$ with $J^{-} \subset U$, and construct a sequence $\left\{x_{n}\right\}$ in $E$ so that $x_{1}<x_{2}<\cdots<a$, and $a-x_{n}<2^{-2 n}|J|$ for $n=1,2, \ldots$. If $A_{n}=J_{+}^{n} \cup\left\{x_{n}\right\}$, then $\left\|A_{n}\right\|=2,\left|A_{n}\right|=2^{-2 n}|J|$, and $d\left(A_{n}\right)$ $\leqslant 3 \cdot 2^{-2 n}|J|, n=1,2, \ldots$. There is an $n_{0}$ with $d\left(A_{n_{0}}\right)<\eta$, and we see that $\left\{\left(A_{n_{0}}, x_{n_{0}}\right), \ldots,\left(A_{n_{0}+p}, x_{n_{0}+p}\right)\right\}$ belongs to $\mathscr{P}([0,1], 1 / 6 ; \eta)$ for each $p=0,1, \ldots$ Since thin sets are finite and

$$
\sum_{k=n_{0}}^{n_{0}+p}\left|f\left(x_{k}\right)\right| A_{k}\left|-\int_{A_{k}} f\right|=\sum_{k=n_{0}}^{n_{0}+p} \int_{J_{+}^{\prime}} f_{J}=|J| \sum_{k=n_{0}}^{n_{0}+p} \frac{1}{k},
$$

the contradiction follows.

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