# Oriented Measures* 

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A vector measure $\mu=\left(\mu_{1}, \ldots, \mu_{n}\right)$ defined on $[a, b]$ is oriented if for each $k$-tuple of disjoint measurable sets $\left(A_{1}, \ldots, A_{k}\right)$ such that $A_{1}<\cdots<A_{k}$ the determinant

$$
\left|\begin{array}{ccc}
\mu_{1}\left(A_{1}\right) & \cdots & \mu_{1}\left(A_{k}\right) \\
\vdots & \ddots & \vdots \\
\mu_{k}\left(A_{1}\right) & \cdots & \mu_{k}\left(A_{k}\right)
\end{array}\right|
$$

is positive. We study the range $\mathscr{R}$ of an oriented measure:

$$
\begin{aligned}
\mathscr{R} & =\left\{\mu(E): \chi_{E} \text { has } n \text { discontinuity points }\right\} \\
\partial \mathscr{R} & =\left\{\mu(E): \chi_{E} \text { has less than } n-1 \text { discontinuity points }\right\} .
\end{aligned}
$$

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## 1. INTRODUCTION

A theorem of Lyapunov states that the range $\mathscr{R}$ of a non-atomic vector measure $\mu$ on $[a, b]$

$$
\mathscr{R}=\{\mu(A): A \text { measurable subset of }[a, b]\}
$$

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coincides with the convex set

$$
\left\{\int_{a}^{b} \rho d \mu: 0 \leq \rho \leq 1\right\}
$$

However for a given $\rho, 0 \leq \rho \leq 1$, the usual proofs based on convexityextreme points arguments $[4 ; 5]$ do not give any information about the existence of a "nice" set $E$ such that

$$
\mu(E)=\int_{a}^{b} \rho d \mu
$$

Consider, for instance, the two-dimensional vector measure $\mu(A)=$ $(|A|,|A|+2|A \cap B|)$ where $B$ is a borelian subset of $[a, b]$ and $|\mid$ denotes the Lebesgue measure. For each set $E$, the equality $\mu(E)=\mu(B)$ implies $B=E$.

When the measure $\mu$ admits a density $f$, Halkin [3] showed that if for each vector $p \in \mathbb{R}^{n}$ the set

$$
\{t \in[a, b]: p \cdot f(t)>0\}
$$

(where $\cdot$ is the usual scalar product) is a finite (respectively countable) union of intervals, then there exists a set $E$ which is a finite (resp. countable) union of intervals.

In our paper [2], we introduced the stronger orientation condition $\Delta$ : we say that $n$ real functions $f_{1}, \ldots, f_{n}$ verify condition $\Delta$ on an interval $[a, b]$ if for each $k$ in $\{1, \ldots, n\}$ the determinant

$$
\left|\begin{array}{cccc}
f_{1}\left(x_{1}\right) & f_{1}\left(x_{2}\right) & \cdots & f_{1}\left(x_{k}\right) \\
f_{2}\left(x_{1}\right) & f_{2}\left(x_{2}\right) & \cdots & f_{2}\left(x_{k}\right) \\
\vdots & \vdots & \ddots & \vdots \\
f_{k}\left(x_{1}\right) & f_{k}\left(x_{2}\right) & \cdots & f_{k}\left(x_{k}\right)
\end{array}\right|
$$

is not equal to zero whenever the $x_{i}$ 's in $[a, b]$ are distinct and its sign is constant on the $k$-tuples $\left(x_{1}, \ldots, x_{k}\right)$ such that $a \leq x_{1}<x_{2}<\cdots<x_{k} \leq b$.

We showed that if a measure $\mu$ admits a density function whose components are continuous and satisfy the orientation condition $\Delta$ then the set $E$ may be built in such a way that its characteristic function has at most $n$ discontinuity points. Moreover, if $0<\rho<1$ there exist two such sets $E_{1}$ and $E_{2}$ whose characteristic functions $\chi_{E_{1}}$ and $\chi_{E_{2}}$ have exactly $n$ discontinuity points (one set is a neighbourhood of $a$ whereas the other is not).

Our proofs relied upon the fact that the map

$$
\left(\alpha_{1}, \ldots, \alpha_{n}\right) \mapsto \int_{\alpha_{1}}^{\alpha_{2}} f(x) d x+\int_{\alpha_{3}}^{\alpha_{4}} f(x) d x+\cdots
$$

is differentiable and has an invertible Jacobian whenever $a<\alpha_{1}<\cdots<$ $\alpha_{n}<b$.

We also showed that whenever a function $x$ satisfies $x(0)=\cdots=$ $x^{(n-2)}(0)=0$ and $x^{(n-1)}(0)=1$ then the $n$ functions $\left(x^{(n-1)}, \ldots, x^{\prime}, x\right)$ verify $\Delta$ on a neighbourhood of 0 . We applied these results to the study of reachable sets of constrained bang-bang solutions and to non-convex problems of the calculus of variations.

In this paper, we deal with measures which are not necessarily absolutely continuous with respect to the Lebesgue measure.

Oriented Measure. If $A_{1}, \ldots, A_{k}$ are $k$ measurable sets of $[a, b$ ], by $A_{1}<\cdots<A_{k}$ we mean that for all $k$-tuple $\left(x_{1}, \ldots, x_{k}\right)$ of $A_{1} \times \cdots \times A_{k}$ we have $x_{1}<\cdots<x_{k}$. A measure $\mu=\left(\mu_{1}, \ldots, \mu_{n}\right)$ is said to be oriented if for each $k$-tuple of measurable sets $A_{1}, \ldots, A_{k}$ such that $A_{1}<\cdots<A_{k}$ the determinant

$$
\left|\begin{array}{ccc}
\mu_{1}\left(A_{1}\right) & \cdots & \mu_{1}\left(A_{k}\right) \\
\vdots & \ddots & \vdots \\
\mu_{k}\left(A_{1}\right) & \cdots & \mu_{k}\left(A_{k}\right)
\end{array}\right|
$$

is positive.
In this more general framework, we give a new proof of the results stated in [2].

We carry out a deep study of the range $\mathscr{R}$ of the measure:

- for each point $q$ of its interior $\mathscr{R}$ there exist exactly two distinct "dual" sets $E_{1}, E_{2}$ whose characteristic functions have $n$ discontinuity points such that $\mu\left(E_{1}\right)=q=\mu\left(E_{2}\right)$;
- the set $\mathscr{\mathscr { R }}$ coincides with

$$
\left\{\int_{a}^{b} \rho d \mu: 0<\rho<1\right\}
$$

so that the above set is open;

- the set $\mathscr{R}$ is strictly convex;
- a point $\mu(E)$ belongs to the boundary $\mathscr{R}$ of $\mathscr{R}$ if and only if the characteristic function of $E$ has less than $n-1$ discontinuity points;
- finally, we give a recursive decomposition of the boundary $\partial \mathscr{R}$.


## 2. ORIENTED MEASURES

Throughout the paper, we will work with an interval $[a, b]$ equipped with the Lebesgue $\sigma$-field $\mathscr{L}$. Measurable will mean measurable with respect to this $\sigma$-field. A negligible set is a measurable set of Lebesgue measure zero. A vector measure on $[a, b]$ is a countably additive set function defined on the Lebesgue $\sigma$-field with values in $\mathbb{R}^{n}$ for some integer $n$.

Notation. If $A_{1}, \ldots, A_{k}$ are $k$ measurable sets of $[a, b]$, by $A_{1}<\cdots<$ $A_{k}$ we mean that $A_{1}, \ldots, A_{k}$ have non-zero Lebesgue measure and for all $k$-tuple $\left(x_{1}, \ldots, x_{k}\right)$ of $A_{1} \times \cdots \times A_{k}$ we have $x_{1}<\cdots<x_{k}$.

Let $\mu=\left(\mu_{1}, \ldots, \mu_{k}\right)$ be a vector measure. If $\rho$ belongs to $L_{\mu}^{1}([a, b])$, we note

$$
\mu_{i}(\rho)=\int_{a}^{b} \rho d \mu_{i}, \quad \mu(\rho)=\int_{a}^{b} \rho d \mu=\left(\mu_{1}(\rho), \ldots, \mu_{k}(\rho)\right)
$$

Definition 2.1. A vector measure $\mu=\left(\mu_{1}, \ldots, \mu_{n}\right)$ on $[a, b]$ is said to be oriented on $[a, b]$ if it is non-atomic and if for each $k$ in $\{1, \ldots, n\}$ and for each $k$-tuple of measurable sets $A_{1}, \ldots, A_{k}$ such that $A_{1}<\cdots<A_{k}$ the determinant

$$
\left|\begin{array}{ccc}
\mu_{1}\left(A_{1}\right) & \cdots & \mu_{1}\left(A_{k}\right) \\
\vdots & \ddots & \vdots \\
\mu_{k}\left(A_{1}\right) & \cdots & \mu_{k}\left(A_{k}\right)
\end{array}\right|
$$

is positive.
Remark. If $\mu$ is oriented, then $\mu_{1}$ is a positive measure which assigns positive values to sets of positive Lebesgue measure. In particular, if $I$ is a non-trivial interval, then $\mu(I)$ is non-zero.

Remark. If $\mu$ is oriented and $I_{1}, \ldots, I_{n}$ are $n$ disjoint non-trivial intervals, then the vectors $\mu\left(I_{1}\right), \ldots, \mu\left(I_{n}\right)$ form a basis of $\mathbb{R}^{n}$.

A very important fact concerning oriented measures is that their characteristic property carries on from sets to positive functions.

Notation. If $\rho$ is a function its support is the set supp $\rho=\{x: \rho(x) \neq 0\}$.
Theorem 2.2. Let $\mu=\left(\mu_{1}, \ldots, \mu_{n}\right)$ be an oriented measure. If $\rho_{1}, \ldots, \rho_{n}$ are $n \mu$-integrable non-negative functions such that supp $\rho_{1}<$ $\cdots<\operatorname{supp} \rho_{n}$, then the determinant

$$
\left|\begin{array}{ccc}
\mu_{1}\left(\rho_{1}\right) & \cdots & \mu_{1}\left(\rho_{n}\right) \\
\vdots & \ddots & \vdots \\
\mu_{n}\left(\rho_{1}\right) & \cdots & \mu_{n}\left(\rho_{n}\right)
\end{array}\right|
$$

is positive.

Let us first state a preparatory lemma.
Lemma 2.3. Let $\mu=\left(\mu_{1}, \ldots, \mu_{n}\right)$ be a vector measure and $\rho_{1}, \ldots, \rho_{n}$ be $n$ measurable $\mu$-integrable functions. The determinant

$$
\left|\begin{array}{ccc}
\int \rho_{1} d \mu_{1} & \cdots & \int \rho_{n} d \mu_{1} \\
\vdots & \ddots & \vdots \\
\int \rho_{1} d \mu_{n} & \cdots & \int \rho_{n} d \mu_{n}
\end{array}\right|
$$

is equal to

$$
\int \cdots \int \rho_{1}\left(s_{1}\right) \cdots \rho_{n}\left(s_{n}\right) d\left(\sum_{\sigma \in \Pi_{n}} \epsilon(\sigma) \mu_{\sigma(1)} \otimes \cdots \otimes \mu_{\sigma(n)}\right)\left(s_{1}, \ldots, s_{n}\right) .
$$

Proof of Lemma 2.3. The identity is obviously true whenever $\rho_{1}, \ldots, \rho_{n}$ are characteristic functions. The monotone class theorem yields the result.

Proof of Theorem 2.2. We apply the lemma. The domain of integration of the $n$-fold integral is reduced to supp $\rho_{1} \times \cdots \times \operatorname{supp} \rho_{n}$.

We first prove that the measure $\hat{\mu}=\sum_{\sigma \in \Pi_{n}} \epsilon(\sigma) \mu_{\sigma(1)} \otimes \cdots \otimes \mu_{\sigma(n)}$ is positive on the product space $\left(\operatorname{supp} \rho_{1}, \mathscr{L}\right) \times \cdots \times\left(\operatorname{supp} \rho_{n}, \mathscr{L}\right)$ equipped with the product $\sigma$-field (where $\mathscr{L}$ denotes the one-dimensional Lebesgue $\sigma$-field). Notice that the product $\sigma$-field $\mathscr{L}^{\otimes n}$ does not coincide in general with the $n$-dimensional Lebesgue $\sigma$-field (i.e., the completion of the $n$-dimensional Borel $\sigma$-field).

Consider first the case of a subset of supp $\rho_{1} \times \cdots \times \operatorname{supp} \rho_{n}$ which is a product set $A_{1} \times \cdots \times A_{n}$ (where the $A_{i}$ 's are measurable). Necessarily, each $A_{i}$ is a subset of supp $\rho_{i}$. If none of the $A_{i}$ 's is negligible, then we have $A_{1}<\cdots<A_{n}$ and $\hat{\mu}\left(A_{1} \times \cdots \times A_{n}\right)=\operatorname{det}\left[\mu\left(A_{1}\right), \ldots, \mu\left(A_{n}\right)\right]$ is positive by definition.

Suppose now some of the $A_{i}$ 's are negligible. For each index $i, 1 \leq i \leq n$, there exists a decreasing sequence $\left(B_{m}^{i}\right)_{m \in \mathbb{N}}$ of non-negligible measurable subsets of supp $\rho_{i}$ having an empty intersection (this is a consequence of the fact that $\operatorname{supp} \rho_{i}$ is not negligible). Now for each $m$ we have $A_{1} \cup B_{m}^{1}$ $<\cdots<A_{n} \cup B_{m}^{n}$ whence $\hat{\mu}\left(A_{1} \cup B_{m}^{1} \times \cdots \times A_{n} \cup B_{m}^{n}\right)$ is positive. By the continuity of the measure $\mu$, we have

$$
\hat{\mu}\left(A_{1} \times \cdots \times A_{n}\right)=\lim _{m \rightarrow \infty} \hat{\mu}\left(A_{1} \cup B_{m}^{1} \times \cdots \times A_{n} \cup B_{m}^{n}\right)
$$

so that $\hat{\mu}\left(A_{1} \ldots A_{n}\right)$ is non-negative. It follows that $\hat{\mu}$ is non-negative on the boolean algebra of the finite (disjoint) union of product sets: its unique
extension to the $\sigma$-field $\mathscr{L}^{\otimes n}$ generated by these products is also nonnegative.

The function $\left(s_{1}, \ldots, s_{n}\right) \mapsto \rho_{1}\left(s_{1}\right) \cdots \rho_{n}\left(s_{n}\right)$ is positive everywhere on this set and is measurable with respect to the $\sigma$-field $\mathscr{L}^{\otimes n}$ : thus the integral $\int \rho_{1}\left(s_{1}\right) \ldots \rho_{n}\left(s_{n}\right) d \hat{\mu}\left(s_{1}, \ldots, s_{n}\right)$ is positive.

Remark. If $\mu$ is absolutely continuous with respect to the Lebesgue measure, then Lyapunov theorem yields an alternative proof of Theorem 2.2. In fact,

$$
\forall k \in\{1, \ldots, n\} \exists E_{k} \subset \operatorname{supp} \rho_{k} \quad \mu\left(\rho_{k}\right)=\mu\left(E_{k}\right) .
$$

Necessarily $\mu\left(E_{k}\right)$ is non-zero for each $k$ (see remark after definition 2.1) and the absolute continuity hypothesis on $\mu$ implies that the $E_{k}$ 's are not negligible.

It follows that $E_{1}<\cdots<E_{n}$ and $\operatorname{det}\left[\mu\left(\rho_{1}\right), \ldots, \mu\left(\rho_{n}\right)\right]=$ $\operatorname{det}\left[\mu\left(E_{1}\right), \ldots, \mu\left(E_{n}\right)\right]>0$.

We shall denote by $\Gamma_{k}$ the subset

$$
\Gamma_{k}=\left\{\left(x_{1}, \ldots, x_{k}\right) \in \mathbb{R}^{k}: a \leq x_{1} \leq \cdots \leq x_{k} \leq b\right\}
$$

Definition 2.4. The measure $\mu$ is said to be locally oriented if for each $n$-tuple $x$ of $\Gamma_{n}$ there exists a neighbourhood $V=V_{1} \times \cdots \times V_{n}$ of $x$ such that for $k$-tuple of measurable sets $A_{1}<\cdots<A_{k}$ satisfying $A_{1} \times \cdots \times A_{k} \subset V_{1} \times \cdots \times V_{k}$, the determinant

$$
\left|\begin{array}{ccc}
\mu_{1}\left(A_{1}\right) & \cdots & \mu_{1}\left(A_{k}\right) \\
\vdots & \ddots & \vdots \\
\mu_{k}\left(A_{1}\right) & \cdots & \mu_{k}\left(A_{k}\right)
\end{array}\right|
$$

is positive.
As a curiosity, we prove the following:
Proposition 2.5. A locally oriented measure on $[a, b]$ is oriented on [ $a, b]$.

Proof. Let $\mu$ be a locally oriented measure. The compact set $\Gamma_{n}$ can be covered by a finite family of open sets $\left(V_{i}\right)_{i \in \Upsilon}$ where $V_{i}=I_{1}^{i} \times \cdots \times I_{n}^{i}$
 of measurable sets $A_{1}<\cdots<A_{k}$ satisfying $A_{1} \times \cdots \times A_{k} \subset V_{i}$ for some $i \in \Upsilon$, the determinant formed with the first $k$ components of the vectors $\mu\left(A_{1}\right), \ldots, \mu\left(A_{k}\right)$ is positive.

Let $\left(J_{l}\right)_{l \in \Sigma}$ be the finite family of the atoms of the algebra generated by the sets $\left(I_{k}^{i}, i \in \Upsilon, 1 \leq k \leq n\right)$ (thus the $J_{l}$ 's are exactly the sets of the form $\cap_{i, k: x \in I_{k} I_{k}^{i}}$ for some $\left.x \in[a, b]\right)$. Let us remark that for each
$\left(l_{1}, \ldots, l_{k}\right)$ in $\Sigma^{k}$, the product $J_{l_{1}} \times \cdots \times J_{l_{k}}$ is contained in some product $I_{1}^{i_{0}} \times \cdots \times I_{k}^{i_{0}}$. In fact,

$$
J_{l_{1}} \times \cdots \times J_{l_{k}} \subset \bigcup_{i \in \Upsilon} I_{1}^{i} \times \cdots \times I_{k}^{i}
$$

so that there exits $i_{0}$ such that $J_{l_{1}} \times \cdots \times J_{l_{k}} \cap I_{1}^{i_{0}} \times \cdots \times I_{k}^{i_{0}}$ is not empty. It follows that $J_{l_{1}} \cap I_{1}^{i_{0}} \neq \varnothing, \ldots, J_{l_{k}} \cap I_{k}^{i_{0}} \neq \varnothing$ and by the very construction of the sets $J_{l}^{\prime}$ 's we obtain $J_{l_{1}} \subset I_{1}^{i_{0}}, \ldots, J_{l_{k}} \subset I_{k}^{i_{0}}$. We denote by $\hat{\mu}_{k}$ the measure $\hat{\mu}_{k}=\sum_{\sigma \in \Pi_{k}} \epsilon(\sigma) \mu_{\sigma(1)} \otimes \cdots \otimes \mu_{\sigma(k)}$. Let $\left(A_{1}, \ldots, A_{k}\right)$ be a $k$-tuple of measurable sets such that $A_{1}<\cdots<A_{k}$. The product $A_{1} \times \cdots \times$ $A_{k}$ is the disjoint union of the sets $\left(A_{1} \times \cdots \times A_{k}\right) \cap\left(J_{l_{1}} \times \cdots \times J_{l_{k}}\right)$ when $\left(l_{1}, \ldots, l_{k}\right)$ varies in $\Sigma^{k}$. Let now $\left(l_{1}, \ldots, l_{k}\right)$ belong to $\Sigma^{k}$. Either $\left(A_{1} \times \cdots \times A_{k}\right) \cap\left(J_{l_{1}} \times \cdots \times J_{l_{k}}\right)$ is empty (and thus has a zero $\hat{\mu}_{k}$ measure) or it is not empty and necessarily, $J_{l_{1}}<\cdots<J_{l_{k}}$. Proceeding as in the proof of Theorem 2.2, we show that $\hat{\mu}_{k}$ is a positive measure on the set $\left(J_{l_{1}} \times \cdots \times J_{l_{k}}\right)$ whence $\hat{\mu}_{k}\left(\left(A_{1} \times \cdots \times A_{k}\right) \cap\left(J_{l_{1}} \times \cdots \times J_{l_{k}}\right)\right)$ is nonnegative. Since the set $A_{1} \times \cdots \times A_{k}$ is not negligible, at least one of these sets is not negligible. Let $\left(A_{1} \times \cdots \times A_{k}\right) \cap\left(J_{l_{1}} \times \cdots \times J_{l_{k}}\right)$ be such a set. It is a subset of one of the $V_{i}$ 's and, moreover, $\left(A_{1} \cap J_{l_{1}}\right)<\cdots<$ $\left(A_{k} \cap J_{l_{k}}\right)$ whence $\hat{\mu}_{k}\left(\left(A_{1} \cap J_{l_{1}}\right) \times \cdots \times\left(A_{k} \cap J_{l_{k}}\right)\right)$ is positive. Thus $\hat{\mu}_{k}\left(A_{1} \times \cdots \times A_{k}\right)$ is positive.

## 3. ORIENTED MEASURES WITH DENSITIES

Orientation Condition $\Delta$. We say that $n$ real functions $f_{1}, \ldots, f_{n}$ verify condition $\Delta$ on an interval $[a, b]$ if for each $k$ in $\{1, \ldots, n\}$, the determinant

$$
\left|\begin{array}{ccc}
f_{1}\left(x_{1}\right) & \cdots & f_{1}\left(x_{k}\right) \\
f_{2}\left(x_{1}\right) & \cdots & f_{2}\left(x_{k}\right) \\
\vdots & \ddots & \vdots \\
f_{k}\left(x_{1}\right) & \cdots & f_{k}\left(x_{k}\right)
\end{array}\right|
$$

is positive whenever the $x_{i}$ 's in $[a, b]$ are such that $a \leq x_{1}<x_{2}<\cdots<$ $x_{k} \leq b$.

Remark. In our previous paper [2], we did not impose the sign of the above determinant to be positive. When dealing with continuous functions, a connectedness argument shows immediately that the sign is constant on the set $\Gamma_{k}$. In our present framework (at the measure level), we find it convenient to work with this slightly more restrictive condition.

Examples. For $n=1$, condition $\Delta$ states that the function $f_{1}$ is positive; for $n=2$, the functions $f_{1}, f_{2}$ satisfy $\Delta$ if and only if $f_{1}$ is positive and $f_{2} / f_{1}$ is strictly increasing. The functions $f_{i}(t)=t^{i-1}(i \geq 1)$ satisfy condition $\Delta$ on $\mathbb{R}$ (the corresponding determinants are Vandermonde determinants).

Proposition 3.1. Let $f_{1}, \ldots, f_{n}$ be $n$ functions in $L^{1}([a, b])$ satisfying the orientation condition $\Delta$ on $[a, b]$. Let $\mu_{i}$ be the measure on $[a, b]$ whose density with respect to the Lebesgue measure is $f_{i}$. Then the measure $\mu=$ $\left(\mu_{1}, \ldots, \mu_{n}\right)$ is oriented.

Proof. Let $A_{1}<\cdots<A_{k}$ be $k$ measurable sets of $[a, b]$. Since the determinant is a multilinear continuous form, we can write

$$
\left|\begin{array}{ccc}
\int_{A_{1}} f_{1} & \cdots & \int_{A_{k}} f_{1} \\
\vdots & \ddots & \vdots \\
\int_{A_{1}} f_{k} & \cdots & \int_{A_{k}} f_{k}
\end{array}\right|=\int_{A_{1} \times \cdots \times A_{k}} \cdots \int_{\times}\left|\begin{array}{ccc}
f_{1}\left(s_{1}\right) & \cdots & f_{1}\left(s_{k}\right) \\
f_{2}\left(s_{1}\right) & \cdots & f_{2}\left(s_{k}\right) \\
\vdots & \ddots & \vdots \\
f_{k}\left(s_{1}\right) & \cdots & f_{k}\left(s_{k}\right)
\end{array}\right| d s_{1} \ldots d s_{k}
$$

By condition $\Delta$, the integrand is positive on $A_{1} \times \cdots \times A_{k}$.
If $f_{1}, \ldots, f_{k}$ are of class $\mathscr{C}^{k-1}$ on $[a, b]$, we will denote their Wronskian by $W\left(f_{1}, \ldots, f_{k}\right)$. The following operational criterion for the fulfilment of the orientation condition $\Delta$ has been used in [2].

Proposition 3.2. Let $f_{1}, \ldots, f_{n}$ in $\mathscr{E}^{n-1}([a, b])$ be such that

$$
\forall t \in[a, b] \quad W\left(f_{1}\right)(t)>0, \ldots, W\left(f_{1}, \ldots, f_{n}\right)(t)>0
$$

Then $f_{1}, \ldots, f_{n}$ satisfy the orientation condition $\Delta$ on $[a, b]$.

## 4. NOTATIONS AND PRELIMINARY LEMMAS

Let us introduce some notations.
If $u_{1}, \ldots, u_{n}$ are vectors of $\mathbb{R}^{n}$, their determinant is sometimes denoted by $\operatorname{det}\left[u_{1}, \ldots, u_{n}\right]$. Let $A$ be a $n \times n$ matrix with real coefficients; by $\operatorname{det} A$ or $|A|$, we denote its determinant. For each $i, j$ in $\{1, \ldots, n\}$, by $A_{i j}$ we mean the $(n-1) \times(n-1)$ matrix obtained by removing the $i$ th row and the $j$ th column from $A$. Surprisingly, the following simple algebraic trick will play an essential role in the existence part of the proof of theorem 5.1.

Lemma 4.1. Let $A=\left(a_{i j}\right)_{1 \leq i, j \leq n}$ be an $n \times n$ matrix with real coefficients. Let $x_{1}, \ldots, x_{n}$ be such that

$$
\left\{\begin{array}{lc}
a_{1,1} x_{1}+\cdots+a_{1, n-1} x_{n-1} & +a_{1, n} x_{n}=0 \\
a_{2,1} x_{1}+\cdots+a_{2, n-1} x_{n-1} & +a_{2, n} x_{n}=0 \\
\vdots & \vdots
\end{array} \vdots \vdots \vdots+a_{n-1, n-1} x_{n-1}+a_{n-1, n} x_{n}=0 ~ \$ ~ \$\right.
$$

If $\operatorname{det} A_{n n} \neq 0$, then

$$
a_{n 1} x_{1}+\cdots+a_{n n} x_{n}=\frac{|A|}{\left|A_{n n}\right|} x_{n}
$$

Proof. Cramer's rule applied to the above system yields

$$
\forall i \in\{1, \ldots, n-1\} \quad x_{i}=\frac{(-1)^{n+i}\left|A_{n i}\right|}{\left|A_{n n}\right|} x_{n}
$$

so that

$$
a_{n 1} x_{1}+\cdots+a_{n n} x_{n}=\frac{\sum_{i=1}^{n}(-1)^{n+i} a_{n i}\left|A_{n i}\right|}{\left|A_{n n}\right|} x_{n}=\frac{|A|}{\left|A_{n n}\right|} x_{n}
$$

since $\sum_{i=1}^{n}(-1)^{n+i} a_{n i}\left|A_{n i}\right|$ is the development of the determinant $|A|$ along the first row.

The next lemmas involve strongly the notion of oriented measure.
Lemma 4.2. Let $F$ and $G$ be two distinct subsets of $[a, b]$ which are the union of $l$ and $m$ disjoint closed intervals

$$
F=\bigcup_{i=1}^{l} I_{i}, \quad G=\bigcup_{j=1}^{m} J_{j}
$$

and let $\mu=\left(\mu_{1}, \ldots, \mu_{n}\right)$ be an oriented measure. Assume $\mu(F)=\mu(G)$.
Then $n<l+m$; moreover if $\partial F \cap \partial G \neq \varnothing$ then $n<l+m-1$.
Proof. Let us first remark that the symmetric difference

$$
\left(I_{1} \cup \cdots \cup I_{l}\right) \Delta\left(J_{1} \cup \cdots \cup J_{m}\right)=\left(\bigcup_{i, j}\left(I_{i} \cup J_{j}\right)\right) \quad\left(\bigcup_{i, j}\left(I_{i} \cap J_{j}\right)\right)
$$

is the union of at most $l+m$ non-trivial intervals and that whenever at least two intervals have a common boundary point then this number is
smaller than $l+m-1$. Since the intervals $I_{1}, \ldots, I_{l}$ are disjoint, as well as $J_{1}, \ldots, J_{m}$, we have

$$
\begin{aligned}
& \left(I_{1} \cup \cdots \cup I_{l}\right) \cup\left(J_{1} \cup \cdots \cup J_{m}\right) \backslash\left(I_{1} \cap J_{1}\right) \\
& \quad=\left(I_{1} \cup J_{1}\right) \backslash\left(I_{1} \cap J_{1}\right) \cup\left(I_{2} \cup \cdots \cup I_{l}\right) \cup\left(J_{2} \cup \cdots \cup J_{m}\right)
\end{aligned}
$$

Now, the set $\left(I_{2} \cup \cdots \cup I_{l}\right) \cup\left(J_{2} \cup \cdots \cup J_{m}\right)$ is a union of at most $l+$ $m-2$ disjoint intervals. Either $I_{1} \cap J_{1}=\varnothing$ or $I_{1} \cap J_{1} \neq \varnothing$ and $I_{1} \cup J_{1}$ is an interval. In both cases, $\left(I_{1} \cup J_{1}\right) \backslash\left(I_{1} \cap J_{1}\right)$ is the union of at most two intervals (at most one if $I_{1}$ and $J_{1}$ have a boundary point in common). A straightforward induction gives the result.

Since the sets $F$ and $G$ are distinct, $F \Delta G$ is not empty. Let $A_{1}<\cdots<$ $A_{p}$ be the connected components of $F \Delta G$. For $k$ in $\{1, \ldots, p\}$, we have

$$
A_{k}=\left(A_{k} \cap F\right) \cup\left(A_{k} \cap G\right)
$$

$\left(A_{k} \cap F\right) \cap\left(A_{k} \cap G\right) \subset A_{k} \cap(F \cap G) \subset(F \Delta G) \cap(F \cap G)=\varnothing ;$
the set $A_{k}$ being connected, either $A_{k} \subset F \backslash G$ or $A_{k} \subset G \backslash F$. Put

$$
\lambda_{k}= \begin{cases}+1 & \text { if } A_{k} \subset F \backslash G \\ -1 & \text { if } A_{k} \subset G \backslash F\end{cases}
$$

so that the equality $\mu(F)=\mu(G)$ can be rewritten as

$$
\left\{\begin{array}{ccc}
\lambda_{1} \mu_{1}\left(A_{1}\right) & +\cdots+ & \lambda_{p} \mu_{1}\left(A_{p}\right)=0 \\
\vdots & \ddots & \vdots \\
\lambda_{1} \mu_{1}\left(A_{1}\right) & +\cdots+ & \lambda_{p} \mu_{n}\left(A_{p}\right)=0
\end{array}\right.
$$

Suppose $n \geq p$; the first $p$ equations imply that the determinant

$$
\left|\begin{array}{ccc}
\mu_{1}\left(A_{1}\right) & \cdots & \mu_{1}\left(A_{p}\right) \\
\vdots & \ddots & \vdots \\
\mu_{p}\left(A_{1}\right) & \cdots & \mu_{p}\left(A_{p}\right)
\end{array}\right|
$$

vanishes, which contradicts the fact that $\mu$ is oriented.
The following notations will be used throughout the remainder of the paper.

Notations 4.3. We shall denote by $\Gamma_{k}$ the set

$$
\Gamma_{k}=\left\{\left(\gamma_{1}, \ldots, \gamma_{k}\right) \in \mathbb{R}^{k}: a \leq \gamma_{1} \leq \cdots \leq \gamma_{k} \leq b\right\}
$$

To each $k$-tuple $\gamma=\left(\gamma_{1}, \ldots, \gamma_{k}\right)$ belonging to $\Gamma_{k}$ we associate the two sets

$$
E_{\gamma}^{-}=\bigcup_{\substack{0 \leq i \leq k \\ i \text { odd }}}\left[\gamma_{i}, \gamma_{i+1}\right], \quad E_{\gamma}^{+}=\bigcup_{\substack{0 \leq i \leq k \\ i \text { even }}}\left[\gamma_{i}, \gamma_{i+1}\right]
$$

where by convention $\gamma_{0}=a, \gamma_{k+1}=b$.

Lemma 4.4 (Uniqueness). Let $\mu$ be an $n$-dimensional oriented measure on $[a, b]$. Assume the $n$-tuples $\gamma=\left(\gamma_{1}, \ldots, \gamma_{n}\right)$ and $\delta=\left(\delta_{1}, \ldots, \delta_{n}\right)$ of $\Gamma_{n}$ satisfy $\mu\left(E_{\gamma}^{-}\right)=\mu\left(E_{\delta}^{-}\right)$(respectively $\left.\mu\left(E_{\gamma}^{+}\right)=\mu\left(E_{\delta}^{+}\right)\right)$. Then $E_{\gamma}^{-}=E_{\delta}^{-}$ (resp. $E_{\gamma}^{+}=E_{\delta}^{+}$).

Proof. Assume $E_{\gamma}^{-}, E_{\delta}^{-}$are distinct and $\mu\left(E_{\gamma}^{-}\right)=\mu\left(E_{\delta}^{-}\right)$.
Now, two possible cases may occur according to the parity of $n$.

- If $n=2 r$, the sets $E_{\gamma}^{-}$and $E_{\delta}^{-}$are the union of at most $r$ intervals. Lemma 4.2 implies $n<r+r$, which is absurd.
- If $n=2 r+1$, the sets $E_{\gamma}^{-}$and $E_{\delta}^{-}$are the union of at most $r+1$ intervals. However, $b$ is a common boundary point. Lemma 4.2 implies $n<(r+1)+(r+1)-1$, which is absurd. The dual case $\mu\left(E_{\gamma}^{+}\right)=\mu\left(E_{\delta}^{+}\right)$ can be treated similarly.

The following essential lemma will be used repeatedly.
Lemma 4.5. Let $\mu=\left(\mu_{1}, \ldots, \mu_{n}\right)$ be an oriented measure on the interval [ $a, b$ ] and $I_{0}<I_{1}<\cdots<I_{n}$ be $n+1$ subintervals of $[a, b]$. Then, given $a$ positive $\epsilon$, there exist $n+1$ positive real numbers $\lambda_{0}, \ldots, \lambda_{n}$ such that

$$
\forall l \in\{0, \ldots, n\} \quad 0<\lambda_{l}<\epsilon \quad \text { and } \quad \sum_{l=0}^{n}(-1)^{l} \lambda_{l} \mu\left(I_{l}\right)=0 .
$$

Proof. Consider the $n \times n$ linear system

$$
\lambda_{0} \mu\left(I_{0}\right)-\lambda_{1} \mu\left(I_{1}\right)+\cdots+(-1)^{n-1} \lambda_{n-1} \mu\left(I_{n-1}\right)=(-1)^{n-1} \lambda_{n} \mu\left(I_{n}\right)
$$

where $\lambda_{n}$ is a parameter. The determinant of the system is

$$
\omega_{n}=(-1)^{n(n-1) / 2} \operatorname{det}\left[\mu\left(I_{0}\right), \ldots, \mu\left(I_{n-1}\right)\right]
$$

The measure $\mu$ being oriented, $\omega_{n}$ is not zero. Moreover, for each $i$ in $\{0, \ldots, n-1\}$,

$$
\omega_{i}=\left\lvert\, \begin{array}{cccc}
\mu_{1}\left(I_{0}\right) & \cdots & (-1)^{i-2} \mu_{1}\left(I_{i-2}\right) & (-1)^{n-1} \mu_{1}\left(I_{n}\right) \\
\mu_{2}\left(I_{0}\right) & \cdots & (-1)^{i-2} \mu_{2}\left(I_{i-2}\right) & (-1)^{n-1} \mu_{2}\left(I_{n}\right) \\
\vdots & \ddots & \vdots & \vdots \\
\mu_{n}\left(I_{0}\right) & \cdots & (-1)^{i-2} \mu_{n}\left(I_{i-2}\right) & (-1)^{n-1} \mu_{n}\left(I_{n}\right)
\end{array}\right.
$$

$$
\begin{array}{ccc}
(-1)^{i} \mu_{1}\left(I_{i}\right) & \cdots & (-1)^{n-1} \mu_{1}\left(I_{n-1}\right) \\
(-1)^{i} \mu_{2}\left(I_{i}\right) & \cdots & (-1)^{n-1} \mu_{2}\left(I_{n-1}\right) \\
\vdots & \ddots & \vdots \\
(-1)^{i} \mu_{n}\left(I_{i}\right) & \cdots & (-1)^{n-1} \mu_{n}\left(I_{n-1}\right)
\end{array}
$$

i.e.,

$$
\omega_{i}=(-1)^{n(n-1) / 2} \operatorname{det}\left[\mu\left(I_{0}\right), \ldots, \mu\left(I_{i-2}\right), \mu\left(I_{i}\right), \ldots, \mu\left(I_{n}\right)\right]
$$

By Cramer's formula, $\lambda_{i}$ equals $\lambda_{n} \omega_{i} / \omega_{n}$. The measure $\mu$ being oriented $\omega_{i}$ and $\omega_{n}$ have the same sign so that $\lambda_{i}$ is positive whenever $\lambda_{n}$ is positive. Choosing $\lambda_{n}$ such that

$$
0<\lambda_{n}<\min \left(\frac{\omega_{n}}{\omega_{0}} \epsilon, \ldots, \frac{\omega_{n}}{\omega_{n-1}} \epsilon, \epsilon\right)
$$

we obtain an $(n+1)$-tuple which solves the problem.

## 5. MAIN RESULT

The statement of the main result involves Notations 4.3.
Theorem 5.1. Let $\mu$ be an oriented measure on $[a, b]$ and let $\rho$ be a measurable function defined on $[a, b]$ with values in $[0,1]$.

There exist two $n$-tuples $\alpha=\left(\alpha_{1}, \ldots, \alpha_{n}\right)$ and $\beta=\left(\beta_{1}, \ldots, \beta_{n}\right)$ in $\Gamma_{n}$ such that

$$
\begin{equation*}
\mu\left(E_{\alpha}^{-}\right)=\int_{a}^{b} \rho d \mu=\mu\left(E_{\beta}^{+}\right) . \tag{*}
\end{equation*}
$$

If in addition $0<\rho<1$, then $\alpha$ and $\beta$ in $\Gamma_{n}$ satisfying (*) are unique and verify

$$
a<\alpha_{1}<\cdots<\alpha_{n}<b, \quad a<\beta_{1}<\cdots<\beta_{n}<b
$$

Remark. The measure $\mu$ being non-atomic we do not care about boundary points of intervals and we might write $\mu(\alpha, \beta)$ for the measure of the interval $\mu([\alpha, \beta])$.

Proof. We consider first the case $0<\rho<1$ and we prove the result by induction on $n$.

- $n=1$. The measure $\mu$ being oriented on $[a, b]$, the maps $\alpha \mapsto$ $\mu([\alpha, b])$ and $\beta \mapsto \mu([a, \beta])$ are continuous and strictly monotonic on $[a, b]$. It follows that there exist unique real numbers $\alpha_{1}$ and $\beta_{1}$ such that

$$
\mu\left(\left[\alpha_{1}, b\right]\right)=\int_{a}^{b} \rho d \mu=\mu\left(\left[a, \beta_{1}\right]\right)
$$

- Assume the result is true at rank $n-1$. We deal only with the $n$-tuple $\beta$ : existence of the $n$-tuple $\alpha$ corresponding to $\rho$ at rank $n$ follows from the fact that it coincides with the $n$-tuple $\beta$ corresponding to $1-\rho$.

Define for each $k$ in $\{1, \ldots, n\}$

$$
\mu_{k}(\rho)=\int_{a}^{b} \rho d \mu_{k}
$$

and for each $n$-tuple $\beta$ in $\Gamma_{n}$

$$
\theta_{k}(\beta)=\mu_{k}\left(E_{\beta}^{+}\right)
$$

The inductive assumption yields the existence of two ( $n-1$ )-tuples $\bar{\alpha}=$ $\left(\bar{\alpha}_{1}, \ldots, \bar{\alpha}_{n-1}\right)$ and $\bar{\beta}=\left(\bar{\beta}_{1}, \ldots, \bar{\beta}_{n-1}\right)$ such that

$$
a<\bar{\alpha}_{1}<\cdots<\bar{\alpha}_{n-1}<b, a<\bar{\beta}_{1}<\cdots<\bar{\beta}_{n-1}<b
$$

and for each $k$ in $\{1, \ldots, n-1\}$

$$
\begin{align*}
& \theta_{k}\left(a, \bar{\alpha}_{1}, \ldots, \bar{\alpha}_{n-1}\right)=\sum_{\substack{0 \leq i \leq n-1 \\
i \text { odd }}} \mu_{k}\left(\bar{\alpha}_{i}, \bar{\alpha}_{i+1}\right)=\mu_{k}(\rho), \\
& \theta_{k}\left(\bar{\beta}_{1}, \ldots, \bar{\beta}_{n-1}, b\right)=\sum_{\substack{0 \leq i \leq n-1 \\
i \text { even }}} \mu_{k}\left(\bar{\beta}_{i}, \bar{\beta}_{i+1}\right)=\mu_{k}(\rho) . \tag{**}
\end{align*}
$$

Put

$$
\begin{aligned}
\mathscr{S}=\{ & \beta=\left(\beta_{1}, \ldots, \beta_{n}\right) \in \Gamma_{n}: \beta_{1} \leq \bar{\beta}_{1}, \\
& \left.\forall k \in\{1, \ldots, n-1\} \theta_{k}(\beta)=\mu_{k}(\rho)\right\} .
\end{aligned}
$$

Since $\left(\bar{\beta}_{1}, \ldots, \bar{\beta}_{n-1}, b\right)$ and $\left(a, \bar{\alpha}_{1}, \ldots, \bar{\alpha}_{n-1}\right)$ belong to $\mathscr{S}$, the set $\mathscr{S}$ is not empty.

We show now that either

$$
\theta_{n}\left(\bar{\beta}_{1}, \ldots, \bar{\beta}_{n-1}, b\right)<\mu_{n}(\rho)<\theta_{n}\left(a, \bar{\alpha}_{1}, \ldots, \bar{\alpha}_{n-1}\right)
$$

or

$$
\theta_{n}\left(a, \bar{\alpha}_{1}, \ldots, \bar{\alpha}_{n-1}\right)<\mu_{n}(\rho)<\theta_{n}\left(\bar{\beta}_{1}, \ldots, \bar{\beta}_{n-1}, b\right)
$$

The equalities $(* *)$ yield for each $k$ in $\{1, \ldots, n-1\}$

$$
\sum_{\substack{0 \leq i \leq n-1 \\ i \text { even }}} \int_{\bar{\beta}_{i}}^{\bar{\beta}_{i+1}}(1-\rho) d \mu_{k}-\sum_{\substack{0 \leq i \leq n-1 \\ i \text { odd }}} \int_{\bar{\beta}_{i}}^{\bar{\beta}_{i+1}} \rho d \mu_{k}=0
$$

Put for $k, j$ in $\{1, \ldots, n\}$

$$
x_{j}^{\beta}=(-1)^{j+1}, \quad a_{k j}^{\beta}=\int_{\bar{\beta}_{j-1}}^{\bar{\beta}_{j}} \rho_{j}^{\beta} d \mu_{k}, \quad A^{\beta}=\left(a_{k j}^{\beta}\right)_{1 \leq k, j \leq n},
$$

where

$$
\rho_{j}^{\beta}=\left\{\begin{array}{cl}
\rho & \text { if } j \text { is even } \\
1-\rho & \text { if } j \text { is odd }
\end{array}\right.
$$

With these notations the above equalities become

$$
\forall k \in\{1, \ldots, n-1\} \quad \sum_{j=1}^{n} a_{k j}^{\beta} x_{j}^{\beta}=0 .
$$

Since the measure $\mu$ is oriented then the determinant $\left|A_{n n}^{\beta}\right|$ does not vanish by Theorem 2.2.

We are thus in the position to apply Lemma 4.1:

$$
\theta_{n}\left(\bar{\beta}_{1}, \ldots, \bar{\beta}_{n-1}, b\right)-\mu_{n}(\rho)=\sum_{j=1}^{n} a_{n j}^{\beta} x_{j}^{\beta}=\frac{\left|A^{\beta}\right|}{\left|A_{n n}^{\beta}\right|}(-1)^{n+1}
$$

Similarly, if we define for $k, j$ in $\{1, \ldots, n\}$

$$
x_{j}^{\alpha}=(-1)^{j}, \quad a_{k j}^{\alpha}=\int_{\bar{\alpha}_{j-1}}^{\bar{\alpha}_{j}} \rho_{j}^{\alpha} d \mu_{k}, \quad A^{\alpha}=\left(a_{k j}^{\alpha}\right)_{1 \leq k, j \leq n},
$$

where

$$
\rho_{j}^{\alpha}=\left\{\begin{array}{cc}
\rho & \text { if } j \text { is odd } \\
1-\rho & \text { if } j \text { is even }
\end{array}\right.
$$

we have

$$
\theta_{n}\left(a, \bar{\alpha}_{1}, \ldots, \bar{\alpha}_{n-1}\right)-\mu_{n}(\rho)=\frac{\left|A^{\alpha}\right|}{\left|A_{n n}^{\alpha}\right|}(-1)^{n}
$$

The measure $\mu$ being oriented, the determinants $\left|A^{\alpha}\right|$ and $\left|A^{\beta}\right|$ have the same sign, as do $\left|A_{n n}^{\alpha}\right|$ and $\left|A_{n n}^{\beta}\right|$. It follows that $\theta_{n}\left(\bar{\beta}_{1}, \ldots, \bar{\beta}_{n-1}, b\right)-$ $\mu_{n}(\rho)$ and $\theta_{n}\left(a, \bar{\alpha}_{1}, \ldots, \bar{\alpha}_{n-1}\right)-\mu_{n}(\rho)$ have opposite signs.

At this stage, we prove that the set $\mathscr{S}$ is the graph of a continuous function; this will imply that $\mathscr{S}$ is connected.

Let $\beta_{1}$ belong to $\left[a, \bar{\beta}_{1}\right]$. We are looking for an $(n-1)$-tuple $\left(\beta_{2}, \ldots, \beta_{n}\right)$ satisfying for each $k$ in $\{1, \ldots, n-1\}$

$$
\begin{aligned}
& \mu_{k}\left(a, \beta_{1}\right)+\sum_{\substack{2 \leq i \leq n \\
i \text { even }}} \mu_{k}\left(\beta_{i}, \beta_{i+1}\right) \\
& \quad=\mu_{k}(\rho)=\mu_{k}\left(a, \bar{\beta}_{1}\right)+\sum_{\substack{2 \leq i \leq n-1 \\
i \text { even }}} \mu_{k}\left(\bar{\beta}_{i}, \bar{\beta}_{i+1}\right)
\end{aligned}
$$

or equivalently

$$
\begin{aligned}
\forall k \in\{1, \ldots, n-1\} \quad \sum_{\substack{2 \leq i \leq n \\
i \text { even }}} \mu_{k}\left(\beta_{i}, \beta_{i+1}\right)= & \mu_{k}\left(\beta_{1}, \bar{\beta}_{1}\right) \\
& +\sum_{\substack{2 \leq i \leq n-1 \\
i \text { even }}} \mu_{k}\left(\bar{\beta}_{i}, \bar{\beta}_{i+1}\right) .
\end{aligned}
$$

Suppose first $\beta_{1}=\bar{\beta}_{1}$. The above equations become

$$
\forall k \in\{1, \ldots, n-1\} \quad \sum_{\substack{2 \leq i \leq n \\ i \text { even }}} \mu_{k}\left(\beta_{i}, \beta_{i+1}\right)=\sum_{\substack{2 \leq i \leq n-1 \\ i \text { even }}} \mu_{k}\left(\bar{\beta}_{i}, \bar{\beta}_{i+1}\right) .
$$

We put $\beta=\left(\beta_{2}, \ldots, \beta_{n-1}, \beta_{n}\right)$ and $\hat{\beta}=\left(\bar{\beta}_{2}, \ldots, \bar{\beta}_{n-1}, b\right)$.
If $n$ is odd, then

$$
E_{\beta}^{-}=\left[\beta_{2}, \beta_{3}\right] \cup \cdots \cup\left[\beta_{n-1}, \beta_{n}\right], \quad E_{\hat{\beta}}^{-}=\left[\bar{\beta}_{2}, \bar{\beta}_{3}\right] \cup \cdots \cup\left[\bar{\beta}_{n-1}, b\right]
$$

if $n$ is even, then

$$
E_{\beta}^{-}=\left[\beta_{2}, \beta_{3}\right] \cup \cdots \cup\left[\beta_{n}, b\right], \quad E_{\hat{\beta}}^{-}=\left[\bar{\beta}_{2}, \bar{\beta}_{3}\right] \cup \cdots \cup\left[\bar{\beta}_{n-2}, \bar{\beta}_{n-1}\right] .
$$

In both cases, the preceding formulae can be rewritten as

$$
\forall k \in\{1, \ldots, n-1\} \quad \mu_{k}\left(E_{\beta}^{-}\right)=\mu_{k}\left(E_{\hat{\beta}}^{-}\right)
$$

Lemma 4.4 implies that $E_{\beta}^{-}=E_{\hat{\beta}}$. Since in addition $\bar{\beta}_{2}<\cdots<\bar{\beta}_{n-1}<b$, then necessarily $\beta_{2}=\bar{\beta}_{2}, \ldots, \beta_{n-1}=\bar{\beta}_{n-1}, \beta_{n}=b$.

Suppose now $\beta_{1}<\bar{\beta}_{1}$. Since $\beta_{1}<\bar{\beta}_{1}<\cdots<\bar{\beta}_{n-1}<b$, then Lemma 4.5 yields the existence of $n$ real numbers $\lambda_{1}, \ldots, \lambda_{n}$ in ] $0,1 / 2$ [ such that for each $k$ in $\{1, \ldots, n-1\}$

$$
-\lambda_{1} \mu_{k}\left(\beta_{1}, \bar{\beta}_{1}\right)+\sum_{1 \leq i \leq n-1}(-1)^{i+1} \lambda_{i+1} \mu_{k}\left(\bar{\beta}_{i}, \bar{\beta}_{i+1}\right)=0
$$

The function

$$
\begin{aligned}
\tilde{\rho}= & \left(1-\lambda_{1}\right) \chi_{\left[\beta_{1}, \bar{\beta}_{1}\right]}+\sum_{\substack{1 \leq i \leq n-1 \\
i \text { odd }}} \lambda_{i+1} \chi_{\left[\bar{\beta}_{i}, \bar{\beta}_{i+1}\right]} \\
& +\sum_{\substack{2 \leq i \leq n-1 \\
i \text { even }}}\left(1-\lambda_{i+1}\right) \chi_{\left[\bar{\beta}_{i}, \bar{\beta}_{i+1}\right]}
\end{aligned}
$$

satisfies $0<\tilde{\rho}<1$ on $\left[\beta_{1}, b\right]$ and for each $k$ in $\{1, \ldots, n-1\}$

$$
\int_{\beta_{1}}^{b} \tilde{\rho} d \mu_{k}=\mu_{k}\left(\beta_{1}, \bar{\beta}_{1}\right)+\sum_{\substack{2 \leq i \leq n-1 \\ i \text { even }}} \mu_{k}\left(\bar{\beta}_{i}, \bar{\beta}_{i+1}\right)
$$

We are thus led to find an $(n-1)$-tuple $\left(\beta_{2}, \ldots, \beta_{n}\right)$ such that $\left(\beta_{1} \leq\right)$ $\beta_{2} \leq \cdots \leq \beta_{n}(\leq b)$ and for each $k$ in $\{1, \ldots, n-1\}$

$$
\sum_{\substack{2 \leq i \leq n \\ i \text { even }}} \mu_{k}\left(\beta_{i}, \beta_{i+1}\right)=\int_{\beta_{1}}^{b} \tilde{\rho} d \mu_{k},
$$

or equivalently, if we put $\tilde{\beta}=\left(\beta_{2}, \ldots, \beta_{n}\right)$,

$$
\forall k \in\{1, \ldots, n-1\} \quad \mu_{k}\left(E_{\tilde{\beta}}^{\bar{\beta}}\right)=\int_{\beta_{1}}^{b} \tilde{\rho} d \mu_{k}
$$

Existence and uniqueness of $\tilde{\beta}$ follow from the inductive assumption at rank $n-1$.

In addition, since $0<\tilde{\rho}<1$ on [ $\beta_{1}, b$ ], we have $\beta_{1}<\beta_{2}<\cdots<$ $\beta_{n}<b$.

We can thus define a map $\psi:\left[a, \bar{\beta}_{1}\right] \rightarrow \mathbb{R}^{n-1}$ such that for all $n$-tuple $\left(\beta_{1}, \ldots, \beta_{n}\right)$ in $\Gamma_{n}$

$$
\left(\beta_{1}, \ldots, \beta_{n}\right) \in \mathscr{S} \Leftrightarrow\left(\beta_{2}, \ldots, \beta_{n}\right)=\psi\left(\beta_{1}\right) .
$$

Thus $\mathscr{S}$ is the graph of $\psi$.
By the continuity of the measure $\mu$, the maps $\theta_{k}, 1 \leq k \leq n-1$, are continuous so that the set $\mathscr{S}$ is closed; moreover, the function $\psi$ takes its values in the compact set $[a, b]^{n-1}$. It follows that $\psi$ is continuous. Henceforth $\mathscr{S}$ is connected. As a consequence, the map $\theta_{n}$, being continuous on $\mathscr{S}$, reaches all the values between $\theta_{n}\left(\bar{\beta}_{1}, \ldots, \bar{\beta}_{n-1}, b\right)$ and $\theta_{n}\left(a, \bar{\alpha}_{1}, \ldots, \bar{\alpha}_{n-1}\right)$. In particular, there exists an $n$-tuple $\beta$ in $\mathscr{S}$ such that $\theta_{n}(\beta)=\mu_{n}(\rho)$. This $n$-tuple $\beta$ solves the problem.

Since $\theta_{n}\left(a, \bar{\alpha}_{1}, \ldots, \bar{\alpha}_{n-1}\right) \neq \mu_{n}(\rho)$ and $\theta_{n}\left(\bar{\beta}_{1}, \ldots, \bar{\beta}_{n-1}, b\right) \neq \mu_{n}(\rho)$ then $a<\beta_{1}<\bar{\beta}_{1}$ so that $a<\beta_{1}<\beta_{2}<\cdots<\beta_{n}<b$. Uniqueness of $\beta$ follows from Lemma 4.4.

Consider now the case $0 \leq \rho \leq 1$. Let $\left(\rho_{m}\right)_{m \in \mathbb{N}}$ be a sequence of measurable functions such that $0<\rho_{m}<1$ and $\rho_{m}$ converges to $\rho$ in $L_{\mu}^{1}([a, b])$. For each function $\rho_{m}$ there exists a unique $n$-tuple $\beta^{m}$ such that

$$
\mu\left(E_{\beta^{m}}^{+}\right)=\int_{a}^{b} \rho_{m} d \mu .
$$

By compactness, we may assume that $\beta^{m}$ converges to some $n$-tuple $\beta$ of $\Gamma_{n}$. Passing to the limit, we obtain $\mu\left(E_{\beta}^{+}\right)=\mu(\rho)$.

## 6. THE RANGE OF AN ORIENTED MEASURE

Let $\mu$ be an oriented measure on $[a, b]$. We denote by $\mathscr{R}$ the range of $\mu$, i.e.,

$$
\mathscr{R}=\{\mu(A): A \text { measurable subset of }[a, b]\}
$$

Lemma 6.1. Let $\bar{\rho}$ be a measurable function on $[a, b], 0 \leq \bar{\rho} \leq 1$. Suppose there exist a non-trivial interval I of $[a, b]$ and a positive real number $\epsilon$ such that $\epsilon \leq \bar{\rho} \leq 1-\epsilon$ on I. Then the set

$$
\left\{\int_{a}^{b} \rho d \mu: \rho=\nu \chi_{I}+\bar{\rho}, \nu \in L_{\mu}^{1}(I),|\nu|<\epsilon\right\}
$$

is a neighbourhood of $\int_{a}^{b} \bar{\rho} d \mu$ in $\mathbb{R}^{n}$.
Proof. Let $I_{1}<\cdots<I_{n}$ be $n$ non-trivial subintervals of $I$. The measure $\mu$ being oriented, the vectors $\mu\left(I_{1}\right), \ldots, \mu\left(I_{n}\right)$ form a basis of $\mathbb{R}^{n}$. The map

$$
\Lambda:\left(\lambda_{1}, \ldots, \lambda_{n}\right) \in \mathbb{R}^{n} \mapsto \sum_{1 \leq i \leq n} \lambda_{i} \mu\left(I_{i}\right) \in \mathbb{R}^{n}
$$

is a linear isomorphism and is thus open. Let

$$
V_{\epsilon}=\left\{\left(\lambda_{1}, \ldots, \lambda_{n}\right): \max _{1 \leq i \leq n}\left|\lambda_{i}\right|<\epsilon\right\} .
$$

Since $\Lambda\left(V_{\epsilon}\right)$ is a neighbourhood of the origin and is contained in the set

$$
\left\{\int_{I} \nu d \mu: \nu \in L_{\mu}^{1}(I),|\nu|<\epsilon\right\}
$$

the conclusion follows.
Remark. The hypothesis $\epsilon \leq \bar{\rho} \leq 1-\epsilon$ implies that $\mu(\bar{\rho})$ belongs to the interior of $\mathscr{R}$.

Remark. The conclusion of Lemma 6.1 does not hold for an arbitrary vector measure: consider, for instance, the $n$-dimensional Lebesgue measure.

Let $\theta: \Gamma_{n} \rightarrow \mathscr{R}$ be the function defined by $\theta(\gamma)=\mu\left(E_{\gamma}^{-}\right)$.
The interior of $\Gamma_{n}$ is the set $\stackrel{\circ}{\Gamma}_{n}=\left\{\left(\gamma_{1}, \ldots, \gamma_{n}\right) \in \mathbb{R}^{n}: a<\gamma_{1}<\cdots<\right.$ $\left.\gamma_{n}<b\right\}$.

Corollary 6.2. The set $\theta\left(\Gamma_{n}\right)$ is contained in $\mathscr{R}$.

Lemma 6.3. The set $\theta\left(\stackrel{\circ}{\Gamma}_{n}\right)$ coincides with the set

$$
F=\left\{\int_{a}^{b} \rho d \mu: 0<\rho<1\right\} .
$$

Proof. The existence part of Theorem 5.1 implies that $F$ is contained in $\theta\left(\Gamma_{n}\right)$.

Conversely, let $\gamma=\left(\gamma_{1}, \ldots, \gamma_{n}\right)$ belong to $\stackrel{\circ}{\Gamma}_{n}$; applying Lemma 4.5 to $\mu$, $\gamma$ and $\epsilon<1 / 2$, we obtain an $(n+1)$-tuple $\left(\lambda_{0}, \ldots, \lambda_{n}\right)$ such that

$$
\forall i \in\{0, \ldots, n\} \quad 0<\lambda_{i}<\epsilon \quad \text { and } \quad \sum_{i=0}^{n}(-1)^{i} \lambda_{i} \mu\left(\gamma_{i}, \gamma_{i+1}\right)=0
$$

Put

$$
\rho=\sum_{\substack{0 \leq i \leq n \\ i \text { even }}} \lambda_{i} \chi_{\left[\gamma_{i}, \gamma_{i+1}\right]}+\sum_{\substack{0 \leq i \leq n \\ i \text { odd }}}\left(1-\lambda_{i}\right) \chi_{\left[\gamma_{i}, \gamma_{i+1}\right]} .
$$

By construction, we have $0<\rho<1$ and

$$
\int_{a}^{b} \rho d \mu=\mu\left(E_{\gamma}^{-}\right)=\theta(\gamma)
$$

so that $\theta(\gamma)$ belongs to $F$.
We have the following:
Theorem 6.4. The range of $\theta$ coincides with $\mathscr{R}$; the map $\theta$ induces a homeomorphism from $\stackrel{\circ}{\Gamma}_{n}$ onto $\mathscr{\mathscr { R }}$ and maps $\partial \Gamma_{n}$ onto $\partial \mathscr{R}$.

Proof. The surjectivity of $\theta$ follows directly from Theorem 5.1. Injectivity of the restriction of $\theta$ of $\stackrel{\circ}{\Gamma}_{n}$ is a consequence of the uniqueness part of Theorem 5.1 together with Lemma 6.3. We claim that $\theta\left(\dot{\Gamma}_{n}\right)$ is open. Let $\gamma$ belong to $\stackrel{\circ}{\Gamma}_{n}$. Lemma 4.5 allows as usual to find a piecewise constant function $\bar{\rho}$ such that $0<\bar{\rho}<1$ and $\mu(\bar{\rho})=\theta(\gamma)$. Clearly there exist a positive $\epsilon$ and a subinterval $I$ of $[a, b]$ on which $\epsilon \leq \bar{\rho} \leq 1-\epsilon$. Put

$$
V_{\bar{\rho}}^{I, \epsilon}=\left\{\nu \chi_{I}+\bar{\rho}: \nu \in L_{\mu}^{1}(I),|\nu|<\epsilon\right\} .
$$

Lemma 6.1 implies that the set

$$
\mu\left(V_{\bar{\rho}}^{I, \epsilon}\right)=\left\{\int_{a}^{b} \rho d \mu: \rho \in V_{\bar{\rho}}^{I, \epsilon}\right\}
$$

is a neighbourhood of $\mu(\bar{\rho})$ in $\mathbb{R}^{n}$. Since each element $\rho$ of $V_{\bar{\rho}}^{I, \epsilon}$ satisfies $0<\rho<1$, then $\mu\left(V_{\bar{\rho}}^{I, \epsilon}\right)$ is entirely contained in $F$. Moreover, $F$ coincides with $\theta\left(\stackrel{\circ}{\Gamma}_{n}\right)$ and thus $\theta\left(\Gamma_{n}\right)$ is a neighbourhood of $\theta(\gamma)$.

Now each open convex set in $\mathbb{R}^{n}$ is the interior of its closure; by Lemma 6.3, the set $\theta\left(\Gamma_{n}\right)$ is convex and its closure is $\mathscr{R}$, whence $\theta\left(\stackrel{\circ}{\Gamma}_{n}\right)=\mathscr{R}$.

Finally, we show that the map $\theta$ is proper (i.e., that the inverse image of a compact subset is compact). Let $K$ be a compact subset of $F$ and $\left(\gamma^{m}\right)_{m \in \mathbb{N}}$ be a sequence in $\theta^{-1}(K)$ such that $\theta\left(\gamma^{m}\right)$ converges to $\mu(\rho)$ for some $\rho, 0<\rho<1$. Since the sequence $\left(\gamma^{m}\right)_{m \in \mathbb{N}}$ is contained in $\Gamma_{n}$, by compactness, we may assume that $\gamma^{m}$ converges to $\gamma$ in $\Gamma_{n}$. By the continuity of $\theta$, we have

$$
\theta(\gamma)=\mu\left(E_{\gamma}^{-}\right)=\int_{a}^{b} \rho d \mu
$$

The uniqueness part of Theorem 5.1 implies that $\gamma$ belongs to $\stackrel{\circ}{\Gamma}_{n}$.
The map $\theta$ is proper and thus closed. It follows that its inverse $\theta^{-1}$ is continuous.

The equality $\theta\left(\partial \Gamma_{n}\right)=\partial \mathscr{R}$ is a consequence of the inclusion $\theta\left(\Gamma_{n}\right) \subset \mathscr{R}$ and the fact that $\theta$ is one to one.

We refer to [7] for the definitions of classical notions associated with convex sets. We have the following:

Theorem 6.5. The range $\mathscr{R}$ of an oriented measure is strictly convex.
Proof. Let $\mu(E), \mu(F)$ be two distinct points of $\mathscr{R}$. By Theorem 5.1, we may assume that the sets $E$ and $F$ are finite unions of closed intervals. Let $\lambda \in] 0,1\left[\right.$ and put $\bar{\rho}=\lambda \chi_{E}+(1-\lambda) \chi_{F}$. Assume, for instance, $E \backslash F \neq \varnothing$. Then there exists a non-trivial interval $I$ such that

$$
\forall x \in I \quad \bar{\rho}(x)=\lambda \chi_{E}(x)+(1-\lambda) \chi_{F}(x)=\lambda
$$

Put $\epsilon=\min _{0}(\lambda, 1-\lambda)$. Lemma 6.1 applied to $\bar{\rho}, I, \epsilon$ shows that $\mu(\bar{\rho})$ belongs to $\mathscr{R}$.

Corollary 6.6. Let $E$ be a measurable subset of $[a, b]$. Then $\mu(E)$ belongs to the boundary of $\mathscr{R}$ if and only if there exists a set $F$ which is a finite union of intervals such that $\chi_{F}$ has less than $n-1$ discontinuity points and $E \Delta F$ is $\mu$-negligible (such a set has also a zero Lebesgue measure).

Proof. We first remark that the family of the sets which are a finite union of intervals and whose characteristic function has less than $n-1$ discontinuity points coincides with the family $\left\{E_{\gamma}^{-}: \gamma \in \partial \Gamma_{n}\right\} \cup\left\{E_{\gamma}^{+}: \gamma \in\right.$ $\left.\delta \Gamma_{n}\right\}$.

Theorem 6.4 shows that $\mu(F)$ belongs to $\partial \mathscr{R}$ whenever $F=E_{\gamma}^{-}$or $F=E_{\gamma}^{+}$for some $\gamma \in \partial \Gamma_{n}$.

Conversely let $E$ be such that $\mu(E)$ belongs to $\partial \mathscr{R}$. Theorem 6.4 yields the existence of an $n$-tuple $\gamma$ belonging to $\partial \Gamma_{n}$ such that $\mu\left(E_{\gamma}^{-}\right)=\mu(E)$; a consequence of Theorem 6.5 is that $\mu(E)$ is an extreme point of $\mathscr{R}$. The Olech Theorem [5, Th. 1] implies that $E \Delta E_{\gamma}^{-}$is $\mu$-negligible.

Our approach discloses the recursive structure of the boundary of the range of an oriented measure. For $k$ belonging to $\{0, \ldots, n\}$ let

$$
\mathscr{R}_{k}^{-}=\left\{\mu\left(E_{\gamma}^{-}\right): \gamma \in \Gamma_{k}\right\}, \quad \mathscr{R}_{k}^{+}=\left\{\mu\left(E_{\gamma}^{+}\right): \gamma \in \Gamma_{k}\right\} .
$$

Notice that $\Gamma_{0}=\varnothing, \mathscr{R}_{0}^{-}=\{0\}, \mathscr{R}_{0}^{+}=\{\mu(a, b)\}$.
PROPOSITION 6.7. The function $\gamma \in \stackrel{\circ}{\Gamma}_{k} \mapsto \mu\left(E_{\gamma}^{-}\right) \in \mathscr{R}_{k}^{-} \quad$ (resp. $\gamma \in$ $\left.\stackrel{\circ}{\Gamma}_{k} \mapsto \mu\left(E_{\gamma}^{+}\right) \in \mathscr{R}_{\mathrm{\circ}}^{+}\right)$is a homeomorphism from ${\stackrel{\circ}{\Gamma_{k}}}_{k}$ onto its range which coincides with $\mathscr{R}_{k}^{-}\left(\right.$resp. $\left.\mathscr{\mathscr { R }}_{k}^{+}\right)$.

Proof. Injectivity follows directly from Corollary 6.6. The rest of the proof uses the techniques of the proof of Theorem 6.4.

Remark. For each $k$ in $\{1, \ldots, n-1\}$, the set $\mathscr{R}_{k} \backslash \mathscr{R}_{k-1}$ is partitioned into two connected components $\mathscr{R}_{k}^{-}, \mathscr{R}_{k}^{+}$. However, for $k=n$, $\mathscr{R}_{n}^{-}=\mathscr{R}_{n}^{+}$ $=\mathscr{R}$.

These results yield the following:
Proposition 6.8. The boundary of the range $\mathscr{R}$ of an oriented $n$-dimensional measure admits the decomposition

$$
\partial \mathscr{R}=\mathscr{\mathscr { R }}_{n-1}^{-} \cup \cdots \cup \mathscr{R}_{1}^{-} \cup\{0\} \cup\{\mu(a, b)\} \cup \mathscr{R}_{1}^{+} \cup \cdots \cup \mathscr{\mathscr { R }}_{n-1}^{+} .
$$

Let $T$ be the symmetry with respect to $\mu(a, b) / 2$ (so that for each measurable subset $A$ of $[a, b], T(\mu(A))=\mu([a, b] \backslash A))$. Then for each $k$ belonging to $\{0, \ldots, n\}$, we have

$$
T\left(\mathscr{R}_{k}^{-}\right)=\mathscr{\mathscr { R }}_{k}^{+}, \quad T\left(\mathscr{R}_{k}\right)=\mathscr{R}_{k}
$$

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