# CONSISTENT EXTENSIONS OF LINEAR FUNCTIONALS AND OF PROBABILITY MEASURES 

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## 1. Introduction

Suppose we are given a family of fields $\mathscr{F}_{\alpha}, \alpha \in \mathscr{A}$, of subsets of a fixed set $X$, and for each $\alpha$ a finite measure $m_{\alpha}$ defined on $\mathscr{F}_{\alpha}$. Under what conditions will there exist a measure $m$, defined on a field $\mathscr{F}$ containing each $\mathscr{F}_{\alpha}$, and agreeing with each $m_{\alpha}$ on $\mathscr{F}_{\alpha}$ ? This problem is of some importance in probability theory (it arises in the theory of marginals), and has been studied, for instance, in [3], [5], [7], [14], [16]. The problem takes several forms: the measures considered can be positive or signed, dominated by a previously given measure or not, and countably or finitely additive. The countably additive cases are the most significant, but also the most difficult (see [1], [4], [8]) ; the most far reaching results here seem to be those of Kellerer [5], [6]. However, the finitely additive cases are also of interest (in fact, the case in which $X$ itself is finite is of significance; see [16]), and the present paper deals almost entirely with them.

Our main object is to give a unified, simple, general treatment of the finitely additive cases, allowing arbitrarily many measures $m_{\alpha}$ (in the literature, only finitely many are usually considered). From this we shall obtain one theorem providing a countably additive extension, under additional topological assumptions. We shall also deal with a second general problem (see [16]): under what conditions on the fields $\mathscr{F}_{\alpha}$ will every consistent set of measures have a (finitely additive) consistent extension?

It is convenient to reformulate the finitely additive problems in a slightly more general form, which reduces them to linear algebra. Let $L$ denote the set of all real valued functions on $X$; for each $\alpha \in \mathscr{A}$, let $L_{\alpha}$ be the subset of $L$ consisting of all $\mathscr{F}_{\alpha}$ measurable step functions on $X$, and for each $f \in L_{\alpha}$, let

$$
\begin{equation*}
\phi_{\alpha}(f)=\int_{X} f d m_{\alpha} \tag{1.1}
\end{equation*}
$$

Then the measures $m_{\alpha}$ have a common extension to a (finitely additive) measure $m$ if and only if the linear maps $\phi_{\alpha}: L_{\alpha} \rightarrow R$ have a common extension to a linear map $\phi$ of a linear space containing each $L_{\alpha}$ (which, as we shall see, may as well be taken to be $L$ itself). Thus, throughout most of what follows, we shall be concerned with the problem of extending the linear maps $\phi_{\alpha}$. Since it makes
things no harder, we permit the $L_{\alpha}$ to be arbitrary linear subspaces of an arbitrary linear space $L$, and allow the values of the linear maps $\phi_{\alpha}$ to be in an arbitrary linear space $M$. For simplicity, all linear spaces are understood to be over the real numbers as groundfield (though in Sections 2 through 5 any other field would do as well; in Sections 6 and 7 we would require an ordered field); and the letters $\lambda, \mu, \nu$ are reserved for real numbers, and $\delta, \varepsilon$ for positive real numbers.

The extension problem in its simplest form is dealt with in Section 2 (the content of which is essentially known). In Sections 3 through 5, we consider the second (and harder) problem in the form: under what conditions on the linear spaces $L_{\alpha}$ will every consistent system of linear maps $\phi_{\alpha}$ have a common extension? We show that our results here apply to marginals (in the finitely additive signed case). In Section 6, we return to the first problem for the case of positive measures (reformulated in terms of positive linear maps), and we also cover the case in which the (positive) measure is to be dominated by a previously given measure. In Section 7, we deal with the second problem for positive measures, and for positive dominated ones. Finally, in Section 8, we consider briefly the positive, countably additive case. Throughout, all measures considered will be finite.

## 2. Extension of consistent linear maps

Let $\mathscr{L}=\left\{L_{\alpha} \mid \alpha \in \mathscr{A}\right\}$ be a nonempty family of linear subspaces of a linear space $L$ and, for each $\alpha \in \mathscr{A}$, let $\phi_{\alpha}$ be a linear map of $L_{\alpha}$ in a fixed linear space $M$. We say that the maps $\phi_{\alpha}, \alpha \in \mathscr{A}$, are extendable (or, more accurately, form an extendable family) if there exists a linear map $\phi: L \rightarrow M$ such that, for each $\alpha \in \mathscr{A}$, the restriction $\phi \mid L_{\alpha}=\phi_{\alpha}$. Obviously, a necessary condition for this is that the maps $\phi_{\alpha}$ are consistent (that is, $\phi_{\alpha}\left|L_{\alpha} \cap L_{\beta}=\phi_{\beta}\right| L_{\alpha} \cap L_{\beta}$ for all $\alpha, \beta \in \mathscr{A}$ ) and we assume the consistency of the maps $\phi_{\alpha}$ throughout what follows. They then combine to give a well-defined (but in general not linear!) map $\phi_{U}: U \rightarrow M$, where

$$
\begin{equation*}
U=\bigcup \mathscr{L}=\bigcup\left\{L_{\alpha} \mid \alpha \in \mathscr{A}\right\} \tag{2.1}
\end{equation*}
$$

We write $S=S(U)$ for the linear subspace of $L$ spanned by $U$; generally, if $A \subset L, S(A)$ denotes the linear subspace of $L$ spanned by $A$.

The question of when the consistent family $\left\{\phi_{\alpha} \mid \alpha \in \mathscr{A}\right\}$ is extendable has the following simple and essentially known answer (see [5], Sätze 2.1 and 2.2, and [12]).

Theorem 2.1. The following statements are equivalent (for a consistent family of linear maps $\left.\phi_{\alpha}, \alpha \in \mathscr{A}\right)$ :
(i) the maps $\phi_{\alpha}, \alpha \in \mathscr{A}$, are extendable;
(ii) $\phi_{U}$ can be extended to a linear map $\phi: L \rightarrow M$;
(iii) $\phi_{U}$ can be extended to a linear map $\psi ; S \rightarrow M$;
(iv) whenever $x_{1}, x_{2}, \cdots, x_{n} \in U$ and $\Sigma_{i=1}^{n} x_{i}=0$, then $\Sigma_{i=1}^{n} \phi_{U}\left(x_{i}\right)=0$.

Proof. Obviously (i) and (ii) are equivalent, and (ii) $\rightarrow$ (iii) $\rightarrow$ (iv). It is an easy matter of routine to verify that (iv) $\rightarrow$ (iii), while, finally, the implication (iii) $\rightarrow$ (ii) is the "pre-Hahn-Banach" theorem (see, for instance, [15], p. 140).

Corollary 2.1. If $|\mathscr{A}|=2$, then every consistent pair of maps $\phi_{1}, \phi_{2}$ is extendable.

In fact, condition (iv) holds automatically in this case. For, suppose $x_{i} \in U=$ $L_{1} \cup L_{2}$ for $i=1,2, \cdots, n$, and that $x_{1}+x_{2}+\cdots+x_{n}=0$. Write

$$
\begin{equation*}
y_{1}=\sum\left\{x_{i} \mid x_{i} \in L_{1}\right\}, \quad y_{2}=\sum\left\{x_{i} \mid x_{i} \notin L_{1}\right\} \tag{2.2}
\end{equation*}
$$

(with the usual convention that empty sums are 0 ). Then $y_{j} \in L_{j}, j=1,2$, and $y_{1}+y_{2}=0$. Hence, $y_{2}=-y_{1} \in L_{1}$, so

$$
\begin{equation*}
\sum\left\{\phi_{U}\left(x_{i}\right) \mid i=1,2, \cdots, n\right\}=\phi_{1}\left(y_{1}\right)+\phi_{2}\left(y_{2}\right)=\phi_{1}\left(y_{1}+y_{2}\right)=0, \tag{2.3}
\end{equation*}
$$

as required.
However, when $|\mathscr{A}| \geqq 3$, not every consistent system is extendable; in other words, condition (iv) is not vacuous. This is shown by the following simple example.

Example 2.1. Let $L$ be two dimensional Euclidean space $R^{2}, M$ the real line $R$, and $\mathscr{L}=\left\{L_{1}, L_{2}, L_{3}\right\}$, where $L_{1}, L_{2}, L_{3}$ are the one dimensional subspaces of $L$ spanned by the vectors ( 1,0$),(0,1),(1,1)$, respectively. Define $\phi_{1}$ and $\phi_{2}$ to be the constant maps sending $L_{1}$ and $L_{2}$ to 0 , and define $\phi_{3}(x, x)=$ $x(x \in R)$. These maps are consistent, but clearly the only linear extension of $\phi_{1}$ and $\phi_{2}$ to $L$ is 0 , which disagrees with $\phi_{3}$.

In the next sections we consider the less trivial problem: under what conditions on $\mathscr{L}$ is every consistent system $\left\{\phi_{\alpha} \mid \alpha \in \mathscr{A}\right\}$ of linear maps extendable?

## 3. $\mathscr{L}$ equivalence

3.1. We say that the family $\mathscr{L}=\left\{L_{\alpha} \mid \alpha \in \mathscr{A}\right\}$ of linear subspaces of the linear space $L$ has the " $M$ extension property," where $M$ is a linear space, provided that every consistent system of linear maps $\phi_{\alpha}: L_{\alpha} \rightarrow M$ is extendable. It turns out that this property is independent of $M$ (providing $M \neq\{0\}$ ); once this is proved, we shall refer to it merely as the "extension property." In this section, we introduce and study an equivalence relation which we shall use to characterize the $M$ extension property. (For a completely different characterization in a special case, see [16].)

Notation. Let $\mathscr{F}$ denote the set of all finite sequences $\left(x_{1}, x_{2}, \cdots, x_{n}\right)$, where $x_{i} \in U=\cup \mathscr{L}$ for each $i=1,2, \cdots, n$ and where $n=1,2, \cdots$. If $F=$ $\left(x_{1}, \cdots, x_{n}\right) \in \mathscr{F}$ and $\lambda \in R$, then $\lambda F$ denotes the sequence $\left(\lambda x_{1}, \cdots, \lambda x_{n}\right)$. We abbreviate $(-1) F$ to $-F$ as usual. If $F=\left(x_{1}, \cdots, x_{n}\right)$ and $G=\left(y_{1}, \cdots, y_{m}\right)$, then $(F, G)$ denotes the sequence $\left(x_{1}, \cdots, x_{n}, y_{1}, \cdots, y_{m}\right)$. We define $\Sigma F$, where $F=\left(x_{1}, \cdots, x_{n}\right) \in \mathscr{F}$, to be the element $x_{1}+x_{2}+\cdots+x_{n}$ of $L$.

The elements $x_{1}, \cdots, x_{n}$ of $L$ are said to be " $\mathscr{L}$ related" provided that there exists some $L_{\alpha} \in \mathscr{L}$ which contains all of them. Finally, a map $\phi: A \rightarrow M$, where
$A \subset L$, is called " $\mathscr{L}$ linear" provided that: whenever $x_{1}, \cdots, x_{n}$ are $\mathscr{L}$ related elements of $A$ and $\lambda_{1}, \cdots, \lambda_{n}$ are real numbers such that $\Sigma_{i=1}^{n} \lambda_{i} x_{i}=0$, then $\Sigma_{i=1}^{n} \lambda_{i} \phi\left(x_{i}\right)=0$. Note that $\phi_{U}$, in Section 2, is automatically $\mathscr{L}$ linear; conversely, every $\mathscr{L}$ linear map $\phi: U \rightarrow M$, where $U=\cup \mathscr{L}$, is expressible in the form $\phi_{U}$ for a consistent family of linear maps $\phi_{\alpha}: L_{\alpha} \rightarrow M$ (namely, $\phi_{\alpha}=\phi \mid L_{\alpha}$ ). Thus, in what follows, the $\mathscr{L}$ linearity of $\phi_{U}$ enables us to forget about the individual linear maps $\phi_{\alpha}$.
3.2. Now consider the following relations on $\mathscr{F}: F R_{1} G$ means that $G$ is a permutation of $F ; F R_{2} G$ means that $F=\left(x_{1}, \cdots, x_{n-1}, x_{n}\right)$ and $G=\left(x_{1}, \cdots\right.$, $x_{n-1}, y_{1}, \cdots, y_{r}$ ), where $x_{n}, y_{1}, \cdots, y_{r}$ are $\mathscr{L}$ related and $x_{n}=y_{1}+\cdots+y_{r}$. Here $n, r$ are arbitrary positive integers. $F R_{3} G$ means that $G R_{2} F$.

We say that $F, G$ are $\mathscr{L}$ equivalent, and write $F \sim G$, provided there exists a finite sequence $H_{1}, H_{2}, \cdots, H_{m} \in \mathscr{F}$ (where $m \geqq 2$ ) such that $F=H_{1}, G=H_{m}$, and, for each $i=1, \cdots, m-1$, we have $H_{i} R_{j} H_{i+1}$ for some $j=1,2,3$. Clearly $\mathscr{L}$ equivalence is just the equivalence relation generated by $R_{1}$ and $R_{2}$.

We shall make frequent use of the following easily verified properties of $\mathscr{L}$ equivalence, where it is understood that $F, G, H \in \mathscr{F}$ and $\lambda \in R$ :
(a) if $F=\left(x_{1}, \cdots, x_{n}\right)$, and $x_{i_{1}}, x_{i_{2}}, \cdots, x_{i_{r}}$ are $\mathscr{L}$ related and $x_{j_{1}}, \cdots, x_{j_{s}}$ are the other elements of $F$, then $F \sim\left(x_{j_{1}}, \cdots, x_{j_{s}}, x_{i_{1}}+x_{i_{2}}+\cdots+x_{i_{r}}\right)$;
(b) if $F \sim G$ then $\lambda F \sim \lambda G$;
(c) $(F,(0,0, \cdots, 0)) \sim F$;
(d) $(F, G) \sim(G, F)$;
(e) if $F_{i} \sim G_{i}, i=1, \cdots, m$, then $\left(F_{1}, F_{2}, \cdots, F_{m}\right) \sim\left(G_{1}, G_{2}, \cdots, G_{m}\right)$;
(f) if $F \sim G$, then $\Sigma F=\Sigma G$ (the converse is not true; see Theorem 3.1 and Example 5.2);
(g) $F \sim G$ if and only if $(F,-G) \sim(0)$; in particular, $(H,-H) \sim(0)$ always;
(h) if $(F, H) \sim(G, H)$, then $F \sim G$ (and conversely, by (e)).

Of these assertions, only (g) and (h) need be proved here. We first establish the special case of $(\mathrm{g})$, that $(H,-H) \sim(0)$. Suppose $H=\left(h_{1}, \cdots, h_{n}\right)$; then

$$
\begin{align*}
(H,-H) & R_{1}\left(h_{1},-h_{1}, h_{2},-h_{2}, \cdots, h_{n},-h_{n}\right)  \tag{3.1}\\
& R_{3}\left(h_{1},-h_{1}, \cdots, h_{n-1},-h_{n-1}, 0\right) \\
& R_{3}\left(h_{1},-h_{1}, \cdots, h_{n-1},-h_{n-1}\right) \cdots R_{3}\left(h_{1},-h_{1}\right) R_{3}(0)
\end{align*}
$$

From (e), it now follows that
(i) $(F, H,-H) \sim F$.

Now if $F \sim G$ then, from (e), $(F,-G) \sim(G,-G) \sim(0)$, from the special case of $(\mathrm{g})$ just established. Conversely, if $(F,-G) \sim(0)$, then by (i) we have $F \sim(F,-G, G) \sim(0, G) \sim G$, by (e) and (c). This proves (g), and (h) follows easily from ( g ) and (i).

The following lemma permits us to express the $\mathscr{L}$ linearity of $\phi_{U}$ in terms of $\mathscr{L}$ equivalence.

Lemma 3.1. Let $\psi$ be a map from $U=\cup \mathscr{L}$ to $M$. The following statements are equivalent:
(i) $\psi$ is $\mathscr{L}$ linear ;
(ii) if $\left(\lambda_{1} x_{1}, \lambda_{2} x_{2}, \cdots, \lambda_{n} x_{n}\right) \sim(0)$, where $x_{i} \in U$ and $\lambda_{i} \in R, i=1,2, \cdots, n$, then $\Sigma_{i=1}^{n} \lambda_{i} \psi\left(x_{i}\right)=0$;
(iii) if $\left(\lambda_{1} x_{1}, \cdots, \lambda_{n} x_{n}\right) \sim\left(\mu_{1} y_{1}, \cdots, \mu_{m} y_{m}\right)$, where the $x$ and the $y$ are in $U$ and the $\lambda$ and the $\mu$ are real, then $\Sigma_{i=1}^{n} \lambda_{i} \psi\left(x_{i}\right)=\Sigma_{j=1}^{m} \mu_{j} \psi\left(y_{j}\right)$.

Suppose (ii) holds, and suppose that $x_{1}, \cdots, x_{n}$ are $\mathscr{L}$ related elements of $U$ and $\lambda_{1}, \cdots, \lambda_{n}$ are real numbers such that $\Sigma_{i} \lambda_{i} x_{i}=0$. Then $\lambda_{1} x_{1}, \cdots, \lambda_{n} x_{n}$ are also $\mathscr{L}$ related; hence, $\left(\lambda_{1} x_{1}, \cdots, \lambda_{n} x_{n}\right) R_{3}(0)$, so (ii) gives $\Sigma_{i} \lambda_{i} \psi\left(x_{i}\right)=0$, and (i) follows.

If (i) holds, we easily check for each of the three relations $R_{j}, j=1,2,3$, that if ( $\left.\lambda_{1} x_{1}, \cdots, \lambda_{n} x_{n}\right) R_{j}\left(\mu_{1} y_{1}, \cdots, \mu_{m} y_{m}\right)$, then $\Sigma_{i} \lambda_{i} \psi\left(x_{i}\right)=\Sigma_{j} \mu_{j} \psi\left(y_{j}\right)$; and from this (iii) follows.

Finally, the implication (iii) $\rightarrow$ (ii) is trivial.
We can now state the first main theorem.
Theorem 3.1. The following assertions are equivalent:
(i) $\mathscr{L}$ has the $M$ extension property, for some $M \neq\{0\}$;
(ii) $\mathscr{L}$ has the $M$ extension property for all $M$;
(iii) for each $F \in \mathscr{F}$ such that $\Sigma F=0$, we have $F \sim(0)$.

Remarks. The converse of (iii) always holds by (f), so (iii) could be restated: $\Sigma F=0$ is equivalent to $F \sim(0)$. It has already been shown (Example 2.1) that (ii) does not hold automatically, in general; thus, the requirement in (iii) is not vacuous. However, it follows from Corollary 2.1 that, when $|\mathscr{A}|=2$, then $\left\{L_{1}, L_{2}\right\}$ always has the $M$ extension property for all $M$, so that (iii) must hold automatically in this case.

Proof. Assume (iii) holds, and let $x_{1}, \cdots, x_{n}$ be elements of $U$ such that $x_{1}+\cdots+x_{n}=0$. Then $\left(x_{1}, \cdots, x_{n}\right) \sim(0)$, by (iii). Since $\phi_{U}$ is $\mathscr{L}$ linear, Lemma 3.1 (ii) (with all the $\lambda=1$ ) shows $\Sigma_{i} \phi_{U}\left(x_{i}\right)=0$. By Theorem 2.1, $\phi_{U}$ can be extended to a linear map $\phi: L \rightarrow M$, so (ii) follows.

The implication (ii) $\rightarrow$ (i) being trivial, all that remains is to prove (i) $\rightarrow$ (iii). This we do in the next section, after developing some further machinery.

## 4. $\mathscr{L}$ span and $\mathscr{L} S$ linearity

4.1. To complete the proof of Theorem 3.1, we introduce two further notions. First, if $A$ is a nonempty subset of $U=U \mathscr{L}$, its " $\mathscr{L}$ span" $\mathscr{L} S(A)$ is the set of all $y \in U$ such that, for some $x_{1}, \cdots, x_{n} \in A$ and $\lambda_{1}, \cdots, \lambda_{n} \in R$ (and some positive integer $n),(y) \sim\left(\lambda_{1} x_{1}, \cdots, \lambda_{n} x_{n}\right)$. The following properties of the $\mathscr{L}$ span are easily verified:
(j) $A \subset \mathscr{L} S(A) \subset S(A) \cap U($ from (f));
(k) if $y \in \mathscr{L} S(A)$ and $\lambda \in R$, then $\lambda y \in \mathscr{L} S(A)$ (from (b));
(l) if $x_{1}, \cdots, x_{n}$ are $\mathscr{L}$ related elements of $\mathscr{L} S(A)$, then $x_{1}+\cdots+x_{n} \in$ $\mathscr{L} S(A)$ (from (a) and (e)).

Together, (k) and (l) give:
(m) for each $\alpha \in \mathscr{A}, L_{\alpha} \cap \mathscr{L} S(A)$ is a linear subspace of $L_{\alpha}$;
(n) $\mathscr{L} S(\mathscr{L} S(A))=\mathscr{L} S(A)$ (from (j) and (e)).

Caution. $\mathscr{L} S(A)$ is not, in general, a linear space; it will be so only when it coincides with the span $S(A)$.
4.2. Next, with $A$ as above, suppose $\psi$ is a map of $A$ into the linear space $M$. We say that $\psi$ is $\mathscr{L} S$ linear (on $A$ ) if, whenever $x_{1}, \cdots, x_{n} \in A$ and $\lambda_{1}, \cdots, \lambda_{n} \in R$ are such that $\left(\lambda_{1} x_{1}, \cdots, \lambda_{n} x_{n}\right) \sim(0)$, then $\Sigma_{i=1}^{n} \lambda_{i} \psi\left(x_{i}\right)=0$. The following easily proved properties will be relevant:
(o) if $\psi$ is the restriction to $A$ of a linear map of $L$ into $M$, then $\psi$ is $\mathscr{L} S$ linear;
(p) if $\psi$ is $\mathscr{L} S$ linear, then $\psi$ is $\mathscr{L}$ linear (from (a)); the converse of ( p ) holds when $A=U$, from Lemma 3.1, but not in general;
(q) if $\psi$ and $\theta$ are $\mathscr{L} S$ linear on $A$, then so is $\lambda \psi+\mu \theta, \quad \lambda, \mu \in R$;
(r) if $0 \in A$ and $\psi$ is $\mathscr{L} S$ linear, then $\psi(0)=0$.

Lemma 4.1. If $\varnothing \neq A \subset U$, and if $\psi: A \rightarrow M$ is $\mathscr{L} S$ linear, then $\psi$ can be extended to an $\mathscr{L} S$ linear map $\Psi: \mathscr{L} S(A) \rightarrow M$.

If $y \in \mathscr{L} S(A)$, we have $(y) \sim\left(\lambda_{1} x_{1}, \cdots, \lambda_{n} x_{n}\right)$, where the $x$ are in $A$ (and the $\lambda$ in $R$ ). Define

$$
\begin{equation*}
\Psi(y)=\sum_{i=1}^{n} \lambda_{i} \psi\left(x_{i}\right) \tag{4.1}
\end{equation*}
$$

This is well defined, for if we also have $(y) \sim\left(\mu_{1} z_{1}, \cdots, \mu_{m} z_{m}\right)$, where the $z$ are in $A$, then by ( g ) we have ( $\lambda_{1} x_{1}, \cdots, \lambda_{n} x_{n},-\mu_{1} z_{1}, \cdots,-\mu_{m} z_{m}$ ) $\sim(0)$, so that the $\mathscr{L} S$ linearity of $\psi$ shows that $\Sigma_{i} \lambda_{i} \psi\left(x_{i}\right)=\Sigma_{j} \mu_{j} \psi\left(z_{j}\right)$.

To verify that $\Psi$ is $\mathscr{L} S$ linear on $\mathscr{L} S(A)$, suppose $y_{1}, \cdots, y_{k} \in \mathscr{L} S(A)$ and $\left(\mu_{1} y_{1}, \cdots, \mu_{k} y_{k}\right) \sim(0)$. For each $i=1, \cdots, k$, we have $\left(y_{i}\right) \sim\left(\lambda_{i, 1} x_{i, 1}, \cdots\right.$, $\lambda_{i, n_{i}} x_{i, n_{i}}$ ), where the $x$ are in $A$, and $\Sigma_{i} \mu_{i} \Psi\left(y_{i}\right)=\Sigma_{i, j} \mu_{i} \lambda_{i, j} \psi\left(x_{i, j}\right)$. But, by (e),

$$
\begin{align*}
&\left(\mu_{1} \lambda_{1,1} x_{1,1}, \cdots, \mu_{1} \lambda_{1, n_{1}} x_{1, n_{1}}, \mu_{2} \lambda_{2,1} x_{2,1}, \cdots, \mu_{k} \lambda_{k, n_{k}} x_{k, n_{k}}\right)  \tag{4.2}\\
& \sim\left(\mu_{1} y_{1}, \cdots, \mu_{k} y_{k}\right) \sim(0)
\end{align*}
$$

hence, the fact that $\psi$ is $\mathscr{L} S$ linear on $A$ shows that $\Sigma_{i} \mu_{i} \Psi\left(y_{i}\right)=0$, as required.
Remark. It is easily seen that the $\mathscr{L} S$ linear extension $\Psi$ here is unique.
Lemma 4.2. Let $A$ be a nonempty subset of $U$, and suppose $x_{0} \in U \backslash \mathscr{L} S(A)$, $m_{0} \in M$. Write $B=A \cup\left\{x_{0}\right\}$, and let $\psi$ be any $\mathscr{L} S$ linear map of $A$ into $M$. Then the map $\psi^{*}: B \rightarrow M$, defined to agree with $\psi$ on $A$ and to make $\psi^{*}\left(x_{0}\right)=m_{0}$, is $\mathscr{L} S$ linear on $B$.

Suppose $y_{1}, \cdots, y_{k} \in B$ and ( $\lambda_{1} y_{1}, \cdots, \lambda_{k} y_{k}$ ) $\sim(0)$; we must show that $\Sigma_{i} \lambda_{i} \psi^{*}\left(y_{i}\right)=0$. If none of the $y$ is $x_{0}$, this follows from the $\mathscr{L} S$ linearity of $\psi$. If all of them are $x_{0}$, we see from (f) that $\left(\lambda_{1}+\cdots+\lambda_{k}\right) x_{0}=0$; but $x_{0} \neq 0$ since $x_{0} \notin \mathscr{L} S(A)$; thus, $\Sigma_{i} \lambda_{i}=0$ and again the result follows. So we may assume, without loss of generality, that $y_{1}, \cdots, y_{r}$ are different from $x_{0}$ and that $y_{r+1}, \cdots, y_{k}$ all $=x_{0}$, where $0<r<k$. Then we have

$$
\begin{equation*}
(0) \sim\left(\lambda_{1} y_{1}, \cdots, \lambda_{k} y_{k}\right) \sim\left(\lambda_{1} y_{1}, \cdots, \lambda_{r} y_{r}, \mu x_{0}\right) \tag{4.3}
\end{equation*}
$$

where $\mu=\lambda_{r+1}+\cdots+\lambda_{k}$. If $\mu \neq 0$, an application of ( b ) and ( g ) would make $\left(x_{0}\right) \sim\left(v_{1} y_{1}, \cdots, v_{r} y_{r}\right)$, where $v_{i}=-\lambda_{i} / \mu$, so that $x_{0} \in \mathscr{L} S(A)$, a contradiction. So $\mu=0$, and it follows by (c) that ( $\lambda_{1} y_{1}, \cdots, \lambda_{r} y_{r}$ ) $\sim(0)$. Thus,

$$
\begin{equation*}
\lambda_{1} \psi^{*}\left(y_{1}\right)+\cdots+\lambda_{k} \psi^{*}\left(y_{k}\right)=\lambda_{1} \psi\left(y_{1}\right)+\cdots+\lambda_{r} \psi\left(y_{r}\right)+\mu m_{0}=0, \tag{4.4}
\end{equation*}
$$

as required.
Lemma 4.3. Let $\psi$ be an $\mathscr{L} S$ linear map of $A$ into $M$, where $\varnothing \neq A \subset U$. Then there exists an $\mathscr{L} S$ linear map $\Psi$ of $U$ into $M$ which extends $\psi$. Moreover, given $x_{0} \in U \backslash \mathscr{L} S(A)$ and $m_{0} \in M$, we may choose $\Psi$ so that $\Psi\left(x_{0}\right)=m_{0}$.

For, using Zorn's lemma, we obtain a maximal $\mathscr{L} S$ linear extension $\Psi$ of $\psi$, mapping, say, $B$ in $M$, where $A \subset B \subset U$ and $B$ is maximal. By Lemma 4.2, $\mathscr{L} S(B)=U$; by Lemma $4.1, B=\mathscr{L} S(B)$; hence, $B=U$ as asserted.

Proof of Theorem 3.1, Concluded. Suppose that Theorem 3.1 (iii) is false (we must show that Theorem 3.1 (i) is also false). There is then a smallest positive integer $n$ for which $n$ elements $x_{1}, \cdots, x_{n}$ of $U$ exist such that $x_{1}+\cdots+x_{n}=0$ but ( $x_{1}, \cdots, x_{n}$ ) is not $\sim(0)$. Clearly, $n \geqq 2$ here (it is not hard to see that in fact $n \geqq 3$ ). We first show:
(1) every $n-1$ of $x_{1}, \cdots, x_{n}$ are linearly independent.

For if, say, $x_{1}=\lambda_{2} x_{2}+\cdots+\lambda_{n-1} x_{n-1}$, write $y_{i}=\left(1+\lambda_{i}\right) x_{i}$ for $i=$ $2,3, \cdots, n-1$, and $y_{n}=x_{n}$. Then each $y_{i} \in U$ and $y_{2}+\cdots+y_{n}=0$. Because of the minimality of $n$, it follows that $\left(y_{2}, \cdots, y_{n}\right) \sim(0)$. Again, write $z_{1}=x_{1}$ and $z_{i}=-\lambda_{i} x_{i}$ for $i=2, \cdots, n-1$; a similar argument shows that $\left(z_{1}, \cdots, z_{n-1}\right) \sim(0)$. Thus, by (g), we have

$$
\begin{equation*}
\left(z_{1}, z_{2}, \cdots, z_{n-1}, y_{2}, \cdots, y_{n-1}, y_{n}\right) \sim(0) \tag{4.5}
\end{equation*}
$$

Now, for $i=2, \cdots, n-1$, we observe that $y_{i}$ and $z_{i}$ are $\mathscr{L}$ related and $y_{i}+z_{i}=x_{i}$. Thus, from the foregoing $\mathscr{L}$ equivalence we deduce $\left(x_{1}, \cdots\right.$, $\left.x_{n-1}, x_{n}\right) \sim(0)$, a contradiction; and (1) is established.

Next we show that, for all real $\lambda_{1}, \cdots, \lambda_{n-1}$,
(2) $\left(\lambda_{1} x_{1}, \cdots, \lambda_{n-1} x_{n-1}, x_{n}\right)$ is not $\sim(0)$.

For if (2) were false, (f) would give $\lambda_{1} x_{1}+\cdots+\lambda_{n-1} x_{n-1}+x_{n}=0$ and, hence, $\sum_{i=1}^{n=1}\left(\lambda_{i}-1\right) x_{i}=0$. This contradicts (1) unless $\lambda_{1}=\lambda_{2}=\cdots=$ $\lambda_{n-1}=1$, but in that case (2) holds by hypothesis.

Now put $A=\left\{x_{1}, x_{2}, \cdots, x_{n-1}\right\}$, a nonempty subset of $U$. From (2), $x_{n} \notin \mathscr{L} S(A)$. Choose any linear map $\phi: L \rightarrow M$ (for instance, the zero map), and any $m_{0} \in M$ different from $\phi\left(x_{n}\right)$ (this is where we use the hypothesis $M \neq$ $\{0\}$ ). Let $\psi$ be the restriction $\phi \mid A$; by ( $o$ ). $\psi$ is $\mathscr{L} S$ linear on $A$. By Lemma 4.3 there exists an $\mathscr{L} S$ linear map $\Psi: U \rightarrow M$ which extends $\psi$, and which is such that $\Psi\left(x_{n}\right)=m_{0}$. In particular, by $(p), \Psi$ is $\mathscr{L}$ linear on $U$, so that if Theorem 3.1 (i) holds, $\Psi$ has an extension to a linear map $\Psi^{*}: L \rightarrow M$. Since $x_{1}+\cdots+$ $x_{n}=0$, we see that

$$
\begin{equation*}
\phi\left(x_{n}\right)=-\sum_{i=1}^{n-1} \phi\left(x_{i}\right)=-\sum_{i=1}^{n-1} \Psi *\left(x_{i}\right)=\Psi^{*}\left(x_{n}\right)=\Psi\left(x_{n}\right)=m_{0}, \tag{4.6}
\end{equation*}
$$

a contradiction.

## 5. Simpler conditions

5.1. As Theorem 3.1 shows, if the family $\mathscr{L}$ of linear subspaces of the linear space $L$ has the $M$ extension property for one nontrivial linear space $M$, it has this property for all and, accordingly, we refer to it as "the extension property" in what follows. (Note also that, by Theorem 2.1, the extension property is in a sense independent of $L$, too, since $L$ can be replaced by $S(\cup \mathscr{L})$.) The characterization of the extension property provided by Theorem 3.1 is somewhat complicated (though, as we shall see, quite usable), and in the present section we consider two plausible simplifications of it. It turns out that one of these is necessary and the other is sufficient, and that no more is true in general. We conclude the section by applying the simpler sufficient condition to the problem of marginals (in the case of finitely additive signed measures).

Theorem 5.1. Consider the following statements.
(i) Whenever $L_{0}, L_{1}, \cdots, L_{k} \in \mathscr{L}$ (where $k$ is a positive integer), $L_{0} \cap S\left(L_{1} \cup \cdots \cup L_{k}\right)=S\left(\left(L_{0} \cap L_{1}\right) \cup \cdots \cup\left(L_{0} \cap L_{k}\right)\right)$.
(ii) $\mathscr{L}$ has the extension property.
(iii) For each $\alpha \in \mathscr{A}$,

$$
\begin{equation*}
L_{\alpha} \cap S\left(\bigcup\left\{L_{\beta} \mid \beta \in \mathscr{A} \backslash\{\alpha\}\right\}\right)=S\left(L_{\alpha} \cap \bigcup\left\{L_{\beta} \mid \beta \in \mathscr{A} \backslash\{\alpha\}\right\}\right) \tag{5.1}
\end{equation*}
$$

Then $(\mathrm{i}) \Rightarrow$ (ii) $\Rightarrow$ (iii); but in general neither implication is reversible.
Proof. For the proof that (i) $\Rightarrow$ (ii), suppose that (i) holds but (ii) fails; then, by Theorem 3.1, there exist $x_{0}, x_{1}, \cdots, x_{k} \in U$ such that $x_{0}+\cdots+x_{k}=0$ and $\left(x_{0}, \cdots, x_{k}\right)$ is not $\sim(0)$. We may suppose that $k$ here is as small as possible; note that $k \geqq 2$. We have $x_{i} \in L_{\alpha(i)}$ for some $\alpha(i) \in \mathscr{A}, i=0, \cdots, k$, and then

$$
\begin{equation*}
x_{0}=-\left(x_{1}+\cdots+x_{k}\right) \in L_{\alpha(0)} \cap S\left(L_{\alpha(1)} \cup \cdots \cup L_{\alpha(k)}\right) . \tag{5.2}
\end{equation*}
$$

From (i) we can, therefore, write $x_{0}=y_{1}+\cdots+y_{k}$, where $y_{i} \in L_{\alpha(0)} \cap L_{\alpha(i)}$, $i=1, \cdots, k$. Thus, $x_{0}, y_{1}, \cdots, y_{\mathrm{k}}$ are $\mathscr{L}$ related, and we have

$$
\text { (1) }\left(x_{0},-y_{1}, \cdots,-y_{k}\right) \sim(0) \text {. }
$$

Again, for each $i=1, \cdots, k$, we have that $x_{i}, y_{i}$ are $\mathscr{L}$ related, and $\left(x_{1}+y_{1}\right)+$ $\cdots+\left(x_{k}+y_{k}\right)=0$. The minimality of $k$ shows that $\left(x_{1}+y_{1}, \cdots, x_{k}+y_{k}\right) \sim$ (0), from which we obtain
(2) $\left(x_{1}, \cdots, x_{k}, y_{1}, \cdots, y_{k}\right) \sim(0)$.

By applying (c) and (e) to (1) and (2), we obtain

$$
\begin{equation*}
\left(x_{1}, \cdots, x_{k}, y_{1}, \cdots, y_{k},-y_{1}, \cdots,-y_{k}, x_{0}\right) \sim(0) \tag{5.3}
\end{equation*}
$$

and thence, $\left(x_{0}, x_{1}, \cdots, x_{n}\right) \sim(0)$, a contradiction.
For (ii) $\Rightarrow$ (iii), assume that Theorem 5.1 (ii) holds, and write

$$
\begin{align*}
M_{\alpha} & =L_{\alpha} \cap S\left(\bigcup\left\{L_{\beta} \mid \beta \in \mathscr{A} \backslash\{\alpha\}\right\}\right) \\
N_{\alpha} & =S\left(L_{\alpha} \cap \bigcup\left\{L_{\beta} \mid \beta \in \mathscr{A}-\{\alpha\}\right\}\right) \tag{5.4}
\end{align*}
$$

we must prove $M_{\alpha}=N_{\alpha}$, and since $N_{\alpha} \subset M_{\alpha}$, trivially, it is enough to prove $M_{\alpha} \subset N_{\alpha}$.

It will be convenient to introduce the following notation. For each $F=$ $\left(x_{1}, \cdots, x_{k}\right) \in \mathscr{F}$ and $\alpha \in \mathscr{A}$, write

$$
\begin{equation*}
\Sigma_{\alpha}(F)=\Sigma\left\{x_{i} \mid x_{i} \in L_{\alpha}, i=1,2, \cdots, k\right\} \tag{5.5}
\end{equation*}
$$

Clearly (for $F, G \in \mathscr{F}$ ), we have
(3) $\Sigma_{\alpha}(F,-G)=\Sigma_{\alpha}(F)-\Sigma_{\alpha}(G)$.

We shall show:
(4) if $F \sim G$, then $\Sigma_{\alpha}(F,-G) \in N_{\alpha}$.

Because of (3), it is enough to verify (4) in the three cases $F R_{j} G, j=1,2,3$, separately (see Section 3.2). The case $j=1$ (permutation) is trivial. If $F R_{2} G$, $\operatorname{say} F=\left(x_{1}, \cdots, x_{k}\right), G=\left(x_{1}, \cdots, x_{k-1}, x_{k+1}, \cdots, x_{k+r}\right)$, where $x_{k}=x_{k+1}+$ $\cdots+x_{k+r}$ and the elements $x_{k+i}, i=0,1, \cdots, r$, are $\mathscr{L}$ related; say they all belong to $L_{\beta}$. If $\beta=\alpha$, clearly, $\Sigma_{\alpha}(F,-G)=0 \in N_{\alpha}$; and if $\beta \neq \alpha$, it is also easy to see that $\Sigma_{\alpha}(F,-G) \in N_{\alpha}$. This verifies (4) in the $R_{2}$ case; and the $R_{3}$ case follows. Thus (4) is established.

Now suppose $x \in M_{\alpha}$; then $x \in L_{\alpha}$ and there exist $y_{1}, \cdots, y_{n}$ such that $x=y_{1}+\cdots+y_{n}, y_{i} \in L_{\beta(i)}$, and $\beta(i) \neq \alpha, i=1,2, \cdots, n$. We may suppose the $y$ so numbered that $y_{1}, \cdots y_{r} \notin L_{\alpha}$ and $y_{r+1}, \cdots, y_{n} \in L_{\alpha}$, where $0 \leqq r \leqq n$. Clearly, $y_{r+1}+\cdots+y_{n} \in N_{\alpha}$, so if $z=y_{1}+\cdots+y_{r}$ it will be enough to prove that $z \in N_{\alpha}$ also. Since $\mathscr{L}$ has the extension property, Theorem 3.1 (iii) shows that $\left(z,-y_{1}, \cdots,-y_{r}\right) \sim(0)$. Thus, by (4), $\Sigma_{\alpha}\left(z,-y_{1}, \cdots,-y_{r}\right)=0$. That is, $z=0$, and $x \in N_{\alpha}$, as required.

Example 5.1. To see that (ii) does not imply (i) in Theorem 5.1, let $\mathscr{L}$ be any system of linear subspaces $L_{\alpha}$ of a linear space $L$, which does not have the extension property (for instance, Example 2.1), and consider $\mathscr{L}^{*}=\mathscr{L} \cup\{L\}$. Trivially, $\mathscr{L}^{*}$ has the extension property (ii) (for $L \in \mathscr{L}^{*}$ ). But if (i) were true of $\mathscr{L}^{*}$, it would be true of $\mathscr{L}$ also, so that (from Theorem 5.1) $\mathscr{L}$ would have the extension property, a contradiction.

Example 5.2. To see that (iii) does not imply (ii) in Theorem 5.1, take $L=R^{3}$ and let $\mathscr{L}$ consist of the lines $L_{1}, L_{2}, L_{3}$ spanned by the vectors $(1,0,0)$, $(0,1,0),(1,1,0)$, respectively, together with the planes $L_{4}, L_{5}, L_{6}$, where $L_{i+3}$ is spanned by $L_{i}$ and ( $0,0,1$ ), $i=1,2,3$. One readily verifies that (iii) holds; but the extension property fails for the same reason as in Example 2.1.
5.2. Application. Consider a product space $X=\Pi\left\{X_{a} \mid a \in A\right\}$. For each $B \subset A$, we write $X_{B}=\Pi\left\{X_{a} \mid a \in B\right\}$ and $B^{\prime}=A \backslash B$; thus $X$ can be regarded as $X_{B} \times X_{B^{\prime}}$. For $x \in X$, we write $x_{B}$ for the projection of $x$ in $X_{B}$; thus, we can write $x=\left(x_{B}, x_{B^{\prime}}\right)$. (With the usual convention that points of $X$ are maps from $A$ to $\cup_{a} X_{a}, x_{B}=x \mid B$.) A function $f: X \rightarrow R$ is a "cylinder function on $B$ " if $f\left(x_{B}, x_{B^{\prime}}\right)$ is independent of $x_{B^{\prime}}$.

Now suppose that a nonempty family $\mathscr{T}$ of subsets of $A$ is given, and that for each $T \in \mathscr{T}$ a linear space $L_{T}$ of real valued functions on $X$ is given, in such a way that:
(s) $L_{T}$ consists of cylinder functions on $T, T \in \mathscr{T}$, and
(t) given $T, T_{1}, \cdots, T_{n} \in \mathscr{T}$, and $f_{i} \in L_{T_{i}}, i=1, \cdots, n$, there exists $y_{T^{\prime}} \in X_{T^{\prime}}$ such that the function $g_{i}: X \rightarrow R$, defined by $g_{i}(x)=f_{i}\left(x_{T}, y_{T^{\prime}}\right)$, is in $L_{T} \cap L_{T_{i}}$, $i=1, \cdots, n$. (Note that, in ( t ), $g_{i}$ is automatically a cylinder function on $T \cap T_{i}$.)

This situation arises, for example, in the theory of marginals (see [5]), where each $X_{a}$ has a given field $\mathscr{B}_{a}$ of measurable subsets specified; we put $\mathscr{B}_{\boldsymbol{T}}$ for the field (of subsets of $X$ ) generated by the cylinders on the sets in $\mathscr{B}_{a}, a \in T$, and $L_{T}$ consists of all the $\mathscr{B}_{T}$ measurable $T$ cylinder functions. (In this case, ( t ) is satisfied automatically by choosing $y_{T^{\prime}}$ arbitrarily in $X_{T^{\prime}}$.) The problem there is whether a consistent system of measures on the various spaces ( $X_{T}, \mathscr{B}_{T}$ ), $T \in \mathscr{T}$, can be extended to a measure on $X$. Kellerer ([5], Satz 2.2) shows that the answer is "yes" when we are dealing with a finite number of countably additive signed measures. The following theorem is a partial generalization; it allows the number $|\mathscr{T}|$ of fields to be arbitrary, but produces only a finitely additive (signed) measure.

Theorem 5.2. Under conditions (s) and ( t ), the family $\left\{L_{T} \mid T \in \mathscr{T}\right\}$ has the extension property.
(As remarked in 5.1 , it is not necessary to specify $L$ precisely; any linear space containing all the $L_{T}$ as subspaces would do, for instance, the space of all real functions on $X$.)

Proof. We show that Theorem 5.1 (i) applies. Let $f \in L_{T} \cap S\left(L_{T_{1}} \cup \cdots \cup\right.$ $L_{T_{k}}$ ), where $T, T_{i} \in \mathscr{T}, i=1, \cdots, k$. Then there exist functions $f_{i} \in L_{T_{i}}$ such that $f=f_{1}+\cdots+f_{k}$. Applying condition ( t ), we obtain functions $g_{i} \in L_{T} \cap L_{T_{i}}$, where $g_{i}(x)=f_{i}\left(x_{T}, y_{T^{\prime}}\right)$. Clearly, $f=g_{1}+\cdots+g_{k}$ (since $f$ is a cylinder function on $T)$. Thus $f \in S\left(L_{T} \cap L_{T_{1}} \cup \cdots \cup L_{T} \cap L_{T_{k}}\right)$. This proves one inclusion between the sets in Theorem 5.1 (i); the reverse inclusion holds in any case, so the proof is complete.

## 6. Extending positive maps

6.1. Positive extensions. In order to cover the case of positive (that is, nonnegative) finitely additive measures, we introduce the notion of (partial) order into the previous setup (Section 2). Accordingly, we make the following standing assumptions. We suppose that $L$ and $M$ are given p.o. (partially ordered) linear spaces over the reals (for simplicity; any other ordered field would do). For the definition and basic properties of p.o. linear spaces, see for instance [11], Chapter 1 . We assume that $M$ is order complete (that is, every nonempty subset which is bounded above has a least upper bound), and that $L$ has a given order unit $e$ (that is, $e>0$ and for each $x \in L$ there exists $\lambda \in R$ such that $\lambda e \geqq x$ ). (In the application to measures, as described in the introduction, where $L$ consists of all real functions on the space $X, e$ will be the constant function 1.) As before, we suppose given a family $\mathscr{L}=\left\{L_{\alpha} \mid \alpha \in \mathscr{A}\right\}$ of linear subspaces of $\mathscr{L}$, and for each $\alpha \in \mathscr{A}$ a linear map $\phi_{\alpha}: L_{\alpha} \rightarrow M$; but we now suppose that each $\phi_{\alpha}$ is positive in the sense: if $x \in L_{\alpha}$ and $x \geqq 0$, then $\phi_{\alpha}(x) \geqq 0$. We say that the maps
$\phi_{\alpha}$ are "positively extendable" if there exists a positive linear map $\phi: L \rightarrow M$ such that $\phi \mid L_{\alpha}=\phi_{\alpha}, \alpha \in \mathscr{A}$. As before, it is necessary for this that the maps $\phi_{\alpha}$ be consistent, and we assume the consistency of the maps $\phi_{\alpha}$ throughout what follows. Then they combine to give a well defined map $\phi_{U}: U \rightarrow M$, where $U=\cup \mathscr{L}$. Note that $\phi_{U}$ is now "positively $\mathscr{L}$ linear" in the sense: it is $\mathscr{L}$ linear (see Notation, Section 3.1), and its restriction to each $L_{\alpha}$ is positive.

The question of when the maps $\phi_{\alpha}$ are positively extendable is answered by the following analog of Theorem 2.1.

Theorem 6.1. The following statements are equivalent (for a consistent family of positive linear maps $\left.\phi_{\alpha}, \alpha \in \mathscr{A}\right)$ :
(i) the maps $\phi_{\alpha}, \alpha \in \mathscr{A}$, are positively extendable;
(ii) $\phi_{U}$ can be extended to a positive linear map $\phi: L \rightarrow M$;
(iii) $\phi_{U}$ can be extended to a positive linear map $\psi: S(U) \rightarrow M$;
(iv) whenever $x_{1}, x_{2}, \cdots, x_{n} \in U$ and $\Sigma_{i=1}^{n} x_{i} \geqq 0$, then $\Sigma_{i=1}^{n} \phi_{U}\left(x_{i}\right) \geqq 0$.

Remark. In the case of two maps $(|\mathscr{A}|=2)$, this theorem is due to Guy [3]. For related results (especially in the case of marginals) see [7], [14], Theorem 7.1, and [16], p. 156. The condition (iv) has been used, in a somewhat different connection, by von Neumann ([10], p. 90). It should be remarked that, even in the case of marginals, where the existence of an extension for signed measures is automatic (see Section 5.2), some condition is needed to ensure the existence of a positive extension for positive measures; see [16], p. 147 for an example.

Proof. Obviously (i) and (ii) are equivalent, and (ii) $\Rightarrow$ (iii) $\Rightarrow$ (iv). It is an easy matter of routine to verify that (iv) $\Rightarrow$ (iii) (note that (iv) here implies condition (iv) of Theorem 2.1). Finally, the implication (iii) $\Rightarrow$ (ii) follows easily from the following (essentially known) modification of the Hahn-Banach theorem.

Lemma 6.1. Suppose $X, Y$ are p.o. linear spaces, where $X$ has an order unit $e$ and $Y$ is order complete. Let $X^{\prime}$ be a linear subspace of $X$ containing $e$, and let $\phi^{\prime}$ be a positive linear map of $X^{\prime}$ in $Y$. Then $\phi^{\prime}$ can be extended to a positive linear map of $X$ in $Y$.

When $Y=R$ this is in [9], p. 119; the argument carries over to the general case without difficulty.

Example 6.1. We observe that, even when $|\mathscr{A}|=2$, not every consistent system of positive linear maps $\phi_{\alpha}, \alpha \in \mathscr{A}$, is positively extendable-in other words, condition. (iv) in Theorem 6.1 is not vacuous. Moreover, it can even happen that two positive linear maps $\phi_{1}, \phi_{2}$ are extendable without being positively extendable. To see this, take $L=R^{3}, M=R$, ordering $L$ so that $\left(x_{1}, x_{2}, x_{3}\right) \geqq 0$ if and only if all of $x_{1}, x_{2}, x_{3} \geqq 0$. Take $e=(1,1,1)$; let $L_{1}$ be the two dimensional subspace of $L$ spanned by $e$ and ( $1,1,0$ ); and let $L_{2}$ be the subspace spanned similarly by $e$ and $(1,0,1)$. Define $\phi_{1}$ on $L_{1}$ by $\phi_{1}(x, x, z)=8 x+4 z$, and $\phi_{2}$ on $L_{2}$ by $\phi_{2}(x, y, x)=3 x+9 y$. Then $\dot{\phi}_{1}$ and $\phi_{2}$ are consistent, and have a unique common linear extension $\phi$ to $L$, where $\phi(x, y, z)=-x+9 y+$ $4 z$. Clearly, $\phi$ is not positive.
6.2. Dominated extension. Suppose that, in the setup described in 6.1, we are given a positive linear map $\theta: L \rightarrow M$ which "dominates" each $\phi_{\alpha}$ in the sense: if $x \in L_{\alpha}$ and $x \geqq 0$, then $\theta(x) \geqq \phi_{\alpha}(x)$. Under what conditions will there exist a positive extension $\phi$ of the $\phi_{\alpha}$ which is likewise dominated by $\theta$ (that is, $\theta(x) \geqq \phi(x)$ for all positive $x \in L)$ ? This question is answered by the expected modification of Theorem 6.1.

Theorem 6.2. The following statements are equivalent:
(i) the maps $\phi_{\alpha}, \alpha \in \mathscr{A}$, are positively extendable to a map dominated by $\theta$;
(ii) $\phi_{U}$ can be extended to a positive linear map $\phi: L \rightarrow M$ which is dominated by 0 ;
(iii) $\phi_{U}$ can be extended to a positive linear map $\psi: S(U) \rightarrow M$ which is dominated by $\theta$;
(iv) whenever $x_{1}, x_{2}, \cdots, x_{n} \in U$ and $\Sigma_{i=1}^{n} x_{i} \geqq 0$, then $0 \leqq \Sigma_{i=1}^{n} \phi_{U}\left(x_{i}\right) \leqq$ $\sum_{i=1}^{n} \theta\left(x_{i}\right)$.

Remark. This is the analog of a characterization due to Kellerer ([5], Satz 4.3) in the case of marginals (with a finite number of countably additive measures).

The proof is essentially the same as that of Theorem 6.1. Instead of Lemma 6.1, we use the following variation of the Hahn-Banach theorem, proved by a slight sharpening of the proof of that lemma.

Lemma 6.2. Under the hypotheses of Lemma 6.1, suppose that a linear map $\theta: X \rightarrow Y$ is given such that $0 \leqq \phi^{\prime}(x) \leqq \theta(x)$ for each positive $x \in X^{\prime}$. Then $\phi^{\prime}$ can be extended to a linear map $\phi: X \rightarrow Y$ such that $0 \leqq \phi(x) \leqq \theta(x)$ for each positive $\boldsymbol{x} \in X$.

## 7. The positive extension property

7.1. With the same assumptions and notation as in 6.1 , we say that $\mathscr{L}$ has the positive $M$ extension property providing every consistent family of positive linear maps $\phi_{\alpha}: L_{\alpha} \rightarrow M$ is positively extendable. As with the $M$ extension property (3.1), it turns out that this is independent of $M$ (if $M \neq\{0\}$ ), as well as (in a sense) of $L$ (since $L$ could be replaced by $S(U)$, from Theorem 6.1). Once this is proved (Theorem 7.1) we refer to it simply as the "positive extension property". First we need some further notation. If $F=\left(x_{1}, \cdots, x_{n}\right) \in \mathscr{F}$, we say that $F$ is "positive" if there exists $G=\left(y_{1}, \cdots, y_{m}\right) \in \mathscr{F}$ such that $F \sim G$ and $y_{j} \geqq 0, j=1, \cdots, m$. We shall make use of the following obvious properties:
(u) if $x_{i}$ and $y_{i}$ are $\mathscr{L}$ related and $x_{i} \geqq y_{i}, i=1, \cdots, n$, and if ( $y_{1}, \cdots, y_{n}$ ) is positive, then so is $\left(x_{1}, \cdots, x_{n}\right)$;
(v) if $F_{i} \in \mathscr{F}$ and is positive, and if $\lambda_{i} \geqq 0, i=1, \cdots, n$, then $\left(\lambda_{1} F_{1}, \cdots, \lambda_{n} F_{n}\right)$ is positive;
(w) if $F \in \mathscr{F}$ is positive, then $\Sigma F \geqq 0$;
(x) if $F \sim G$ and $G$ is positive, then so is $F$.

Theorem 7.1. Under the assumptions of 6.1, the following statements are equivalent:
(i) $\mathscr{L}$ has the positive $M$ extension property, for some $M \neq\{0\}$;
(ii) $\mathscr{L}$ has the positive $M$ extension property for every (order complete) $M$;
(iii) if $F \in \mathscr{F}$ and $\Sigma F \geqq 0$, then ( $F, \varepsilon e$ ) is positive for all $\varepsilon>0$.

Remark. When $L$ is finite dimensional, a characterization along entirely different lines has been given by Vorob'ev [16].

Proof. Trivially (ii) $\Rightarrow$ (i). To prove (iii) $\Rightarrow$ (ii), suppose $\phi_{U}: U \rightarrow M$ is positively $\mathscr{L}$ linear (Section 6.1) and that (iii) holds. By Theorem 6.1 it is enough to prove that, given $x_{1}, \cdots, x_{n} \in U$ such that $\Sigma_{i} x_{i} \geqq 0$, then $\Sigma_{i} \phi_{U}\left(x_{i}\right) \geqq 0$. For each $\varepsilon>0,\left(x_{1}, \cdots, x_{n}, \varepsilon e\right)$ is positive (by (iii)); that is, there exist $y_{1}, \cdots, y_{m} \in U$ such that each $y_{j} \geqq 0$ and $\left(x_{1}, \cdots, x_{n}, \varepsilon e\right) \sim\left(y_{1}, \cdots, y_{m}\right)$. Lemma 3.1 shows that

$$
\begin{equation*}
\phi_{U}\left(x_{1}\right)+\cdots+\phi_{U}\left(x_{n}\right)+\varepsilon \phi_{U}(e)=\phi_{U}\left(y_{1}\right)+\cdots+\phi_{U}\left(y_{m}\right) \geqq 0 . \tag{7.1}
\end{equation*}
$$

Since this is true for all $\varepsilon>0$, it follows that $\phi_{U}\left(x_{1}\right)+\cdots+\phi_{U}\left(x_{n}\right) \geqq 0$, as required.

Before completing the proof by showing that (i) $\Rightarrow$ (iii), we need some further notation and some lemmas.
7.2. Suppose $\psi$ is a map of $U$ into $M$. It is easy to see that
(1) $\psi$ is positively $\mathscr{L}$ linear if and only if, whenever $x_{1}, \cdots, x_{n}$ are $\mathscr{L}$ related elements of $U$ and $\lambda_{1} x_{1}+\cdots+\lambda_{n} x_{n} \geqq 0$, then $\lambda_{1} \psi\left(x_{1}\right)+\cdots+\lambda_{n} \psi\left(x_{n}\right) \geqq 0$.

Now suppose, more generally, that $\psi$ is a map of $A$ into $M$, where $A$ is an arbitrary subset of $U$. We say that $\psi$ is positively $\mathscr{L} S$ linear (on $A$ ) if, whenever $x_{1}, \cdots, x_{n} \in A$ and $\left(\lambda_{1} x_{1}, \cdots, \lambda_{n} x_{n}\right)$ is positive, then $\lambda_{1} \psi\left(x_{1}\right)+\cdots+\lambda_{n} \psi\left(x_{n}\right) \geqq 0$.

When $A=U$, these notations are equivalent, as we state in Lemma 7.1.
Lemma 7.1. If $\psi: U \rightarrow M$ is positively $\mathscr{L} S$ linear (on $U$ ), then $\psi$ is positively $\mathscr{L}$ linear; and conversely.

For if $x_{1}, \cdots, x_{n}$ are as in (1), then ( $\left.\lambda_{1} x_{1}, \cdots, \lambda_{n} x_{n}\right) \sim\left(\lambda_{1} x_{1}+\cdots+\lambda_{n} x_{n}\right)$ and is, therefore, positive; thus, the conclusion of (1) follows. The converse is proved by an argument similar to the proof of Lemma 3.1.

Remark. If $\psi: A \rightarrow M$ is positively $\mathscr{L} S$ linear, then clearly $\psi$ is both $\mathscr{L} S$ linear (Section 4.2) and positive on $A$. The converse is false in general, though it is true if $A=U$.

Lemma 7.2. If $\varnothing \neq A \subset U$, and if $\psi: A \rightarrow M$ is positively $\mathscr{L} S$ linear, then $\psi$ can be extended to a positively $\mathscr{L} S$ linear map $\Psi: \mathscr{L} S(A) \rightarrow M$.

By Lemma 4.1, $\psi$ can be extended to an $\mathscr{L} S$ linear map $\Psi: \mathscr{L} S(A) \rightarrow M$. We verify that $\Psi$ is positively $\mathscr{L} S$ linear. Suppose, then, that $x_{1}, \cdots, x_{n} \in \mathscr{L} S(A)$, and that ( $\lambda_{1} x_{1}, \cdots, \lambda_{n} x_{n}$ ) is positive. Since $x_{i} \in \mathscr{L} S(A)$, we have $\left(x_{i}\right) \sim$ $\left(\mu_{i, 1} y_{i, 1}, \cdots, \mu_{i, m_{i}} y_{i, m_{i}}\right), i=1, \cdots, n$, where the $y$ are in $A$; hence

$$
\begin{equation*}
\left(\lambda_{1} x_{1}, \cdots, \lambda_{n} x_{n}\right) \sim\left(\lambda_{1} \mu_{1,1} y_{1,1}, \cdots, \lambda_{1} \mu_{1, m_{1}} y_{1, m_{1}}, \cdots, \lambda_{n} \mu_{n, m_{n}} y_{n, m_{n}}\right) \tag{7.2}
\end{equation*}
$$

which is, therefore, also positive. Because $\Psi$ is $\mathscr{L}$ linear by (p), Lemma 3.1 gives

$$
\begin{equation*}
\sum_{i} \lambda_{i} \Psi\left(x_{i}\right)=\sum_{i, j} \lambda_{i} \mu_{i, j} \Psi\left(y_{i, j}\right)=\sum_{i, j} \lambda_{i} \mu_{i, j} \psi\left(y_{i, j}\right) \geqq 0 \tag{7.3}
\end{equation*}
$$

since $\psi$ is positively $\mathscr{L} S$ linear, completing the proof.

Remark. As in Lemma 4.1, the map $\Psi$ here is unique.
Lemma 7.3. Suppose $e \in A \subset U$ and $x_{0} \in U \backslash A$, and let $\psi$ be a positively $\mathscr{L} S$ linear map of $A$ into $M$. Then $\psi$ can be extended to a positively $\mathscr{L} S$ linear map $\psi^{*}: B \rightarrow M$, where $B=A \cup\left\{x_{0}\right\}$.

Remark. In contrast to Lemma 4.2, the value of $\psi^{*}\left(x_{0}\right)$ here cannot be assigned arbitrarily, even if $x_{0} \notin \mathscr{L} S(A)$.

Proof. Because of Lemma 7.2, we may replace $A$ by $\mathscr{L} S(A)$; this allows us to use a somewhat simpler notation in view of $(\mathrm{k})$.

Let $\mathscr{G}$ be the set of all finite sequences $\left\{g_{1}, \cdots, g_{n}\right\}$ of elements of $A$ for which ( $g_{1}, \cdots, g_{n},-x_{0}$ ) is positive; similarly, let $\mathscr{H}$ be the set of all finite sequences $\left\{h_{1}, \cdots, h_{m}\right\}$ of elements of $A$ for which $\left(-h_{1}, \cdots,-h_{m}, x_{0}\right)$ is positive. We first note that neither $\mathscr{G}$ nor $\mathscr{H}$ is empty. In fact, if $N$ is a large enough integer, we have $N e \geqq x_{0}$; also $N e$ and $x_{0}$ are $\mathscr{L}$ related (because each $L_{\alpha}$ contains $e$ ); hence $\left(N e,-x_{0}\right) \sim\left(N e-x_{0}\right)$ and is positive. Thus, the one term sequence $\{N e\} \in \mathscr{G}$; and similarly $\{-N e\} \in \mathscr{H}$.

Next we show that, for all $\left\{g_{1}, \cdots, g_{n}\right\} \in \mathscr{G}$ and $\left\{h_{1}, \cdots, h_{m}\right\} \in \mathscr{H}$,
(2) $\psi\left(g_{1}\right)+\cdots+\psi\left(g_{n}\right) \geqq \psi\left(h_{1}\right)+\cdots+\psi\left(h_{m}\right)$.

In fact, ( v ) shows that ( $g_{1}, \cdots, g_{n},-x_{0},-h_{1}, \cdots,-h_{m}, x_{0}$ ) is positive, and thus ( $g_{1}, \cdots, g_{n},-h_{1}, \cdots,-h_{m}$ ) is also positive by ( x ). Our assumption on $\psi$ now gives

$$
\begin{equation*}
\psi\left(g_{1}\right)+\cdots+\psi\left(g_{n}\right)-\psi\left(h_{1}\right)-\cdots-\psi\left(h_{m}\right) \geqq 0, \tag{7.4}
\end{equation*}
$$

establishing (2). Define

$$
\begin{align*}
q & =\inf \left\{\psi\left(g_{1}\right)+\cdots+\psi\left(g_{n}\right) \mid\left\{g_{1}, \cdots, g_{n}\right\} \in \mathscr{G}\right\} \\
r & =\sup \left\{\psi\left(h_{1}\right)+\cdots+\psi\left(h_{m}\right) \mid\left\{h_{1}, \cdots, h_{m}\right\} \in \mathscr{H}\right\} \tag{7.5}
\end{align*}
$$

(using the fact that $M$ is order complete). From (2), we see that $q, r$ indeed exist and that $r \leqq q$. Choose any $m_{0} \in M$ such that $r \leqq m_{0} \leqq q$ (for instance, $m_{0}=q$ ), and extend $\psi$ to a map $\psi^{*}: B \rightarrow M$ by setting $\psi^{*}\left(x_{0}\right)=m_{0}$. We verify that $\psi^{*}$ is positively $\mathscr{L} S$ linear.

Suppose, then, that ( $\lambda_{1} f_{1}, \cdots, \lambda_{n} f_{n}$ ) is positive, where $f_{1}, \cdots, f_{n} \in B=A \cup$ $\left\{x_{0}\right\}$. We must prove
(3) $\lambda_{1} \psi^{*}\left(f_{1}\right)+\cdots+\lambda_{n} \psi^{*}\left(f_{n}\right) \geqq 0$.

If none of $f_{1}, \cdots, f_{n}$ is $x_{0}$, this follows from the assumption that $\psi$ is positively $\mathscr{L} S$ linear (on $A$ ). If all of them are $x_{0}$, we have from (w) that $\lambda x_{0} \geqq 0$, where $\lambda=\lambda_{1}+\cdots+\lambda_{n}$, and have to prove that $\lambda \psi^{*}\left(x_{0}\right) \geqq 0$. If $\lambda=0$ there is nothing to prove. If $\lambda>0$, we have $x_{0} \geqq 0$, so that $\left(0, x_{0}\right)$ is positive and $\{0\} \in \mathscr{H}$; thus $0 \leqq r \leqq m_{0}=\psi^{*}\left(x_{0}\right)$, as required. Finally, if $\lambda<0$ the argument is similar, using $\mathscr{G}$.

Thus, we may assume $f_{1}, \cdots, f_{k}$ different from $x_{0}$ and $f_{k+1}, \cdots, f_{n}$ all equal to $x_{0}$, where $1 \leqq k<n$. Write $\mu_{0}=\lambda_{k+1}+\cdots+\lambda_{n}$; then

$$
\begin{equation*}
\left(\lambda_{1} f_{1}, \cdots, \lambda_{n} f_{n}\right) \sim\left(\lambda_{1} f_{1}, \cdots, \lambda_{k} f_{k}, \mu_{0} x_{0}\right) \tag{7.6}
\end{equation*}
$$

and thus
(4) $\left(\lambda_{1} f_{1}, \cdots, \lambda_{k} f_{k}, \mu_{0} x_{0}\right)$ is positive.

Again there are three cases, depending now on $\mu_{0}$.
First, if $\mu_{0}=0$, (3) follows from (4) and the fact that $\psi$ is positively $\mathscr{L} S$ linear on $A$.

Second, if $\mu_{0}>0$, write $\mu_{i}=-\lambda_{i} / \mu_{0}, i=1, \cdots, k$; from (4), ( $-\mu_{1} f_{1}, \cdots$, $-\mu_{k} x_{k}, x_{0}$ ) is positive, and thus ( $\left.\mu_{1} f_{1}, \cdots, \mu_{k} f_{k}\right) \in \mathscr{H}$ (since each $\mu_{i} f_{i} \in A$; this is the point of replacing $A$ by $\mathscr{L} S(A)$ at the beginning). Hence, $\mu_{1} \psi\left(f_{1}\right)+\cdots+$ $\mu_{k} \psi\left(f_{k}\right) \leqq r \leqq m_{0}=\psi^{*}\left(x_{0}\right)$, from which (3) follows.

Third, if $\mu_{0}<0$, the argument is similar.
Lemma 7.4. Let A be a nonempty subset of $U$, and $\psi$ a positively $\mathscr{L} S$ linear map of $A$ into $M$. Then there exists a positively $\mathscr{L}$ linear map $\Psi$ of $U$ into $M$ which extends $\psi$.

This follows straightforwardly from the foregoing lemmas and Zorn's lemma.
7.3. Now we conclude the proof of Theorem 7.1. Assume that Theorem 7.1 (iii) fails; we show that Theorem 7.1 (i) fails with $M=R$.

For some positive integer $n$, there exist $f_{1}, \cdots, f_{n} \in U$ such that $f_{1}+\cdots+$ $f_{n} \geqq 0$, and $\varepsilon>0$, such that ( $f_{1}, \cdots, f_{n}, \varepsilon e$ ) is not positive. We take $n$ to be as small as possible; note that $n \geqq 2$ from (v). We are going to apply Lemma 7.4 to $A=\left\{f_{1}, \cdots, f_{n}, e\right\}$, but the definition of a suitable $\psi$ requires some preparation.

Consider the subset $K$ of $R^{n+1}$ consisting of all points $p=\left(\lambda_{1}, \cdots \lambda_{n+1}\right)$ for which ( $\lambda_{1} f_{1}, \cdots, \lambda_{n} f_{n}, \lambda_{n+1} e$ ) is positive. It is easy to see, using ( v ), that $K$ is a convex cone. Let $q$ be the point $(1, \cdots, 1,0) \in R^{n+1}$. We show that
(5) $q \notin \bar{K}$ (the closure of $K$ ).

For, take $N$ to be an integer large enough to make $-N e \leqq f_{i} \leqq N e$, $i=1, \cdots, n$, and put $\delta=\varepsilon /(N n+1)$. If (5) were false, there would exist $\mu_{1}, \cdots, \mu_{n+1}$ such that $\left|\mu_{i}\right|<\delta, i=1, \cdots, n+1$, and $\left(1+\mu_{1}, \cdots, 1+\mu_{n}\right.$, $\left.\mu_{n+1}\right) \in K$; that is,
(6) $F=\left(\left(1+\mu_{1}\right) f_{1}, \cdots,\left(1+\mu_{n}\right) f_{n}, \mu_{n+1} e\right)$ is positive.

Now $\delta N e-\mu_{i} f_{i} \geqq 0, i=1, \cdots, n$, from the choice of $N$ and the fact that $e>0$. Hence,
(7) $G=\left(\delta N e-\mu_{1} f_{1}, \cdots, \delta N e-\mu_{n} f_{n},\left(\delta-\mu_{n+1}\right) e\right)$ is positive.

Accordingly, $(F, G)$ is positive. But $e$ and $f_{i}$ are $\mathscr{L}$ related for each $i$; hence, we may perform an $\mathscr{L}$ equivalence on $(F, G)$ by combining each pair $\left(1+\mu_{i}\right) f_{i}$, $\delta N e-\mu_{i} f_{i}$, into the single term $f_{i}+\delta N e$. We can also replace the pair $\mu_{n+1} e$, $\left(\delta-\mu_{n+1}\right) e$ by the single term $\delta e$. Thus, $(F, G) \sim H=\left(f_{1}+\delta N e, \cdots, f_{n}+\right.$ $\delta N e, \delta e)$. Now we apply a further $\mathscr{L}$ equivalence to $H$, replacing the pair $f_{1}+\delta N e, \delta e$ first by their sum, and then this sum by the pair $f_{1}, \delta(N+1) e$. Applying this principle $n$ times, we obtain $H \sim\left(f_{1}, \cdots, f_{n}, \delta(N n+1) e\right)$, which therefore $\sim(F, G)$ and is positive. From (u) and the choice of $\delta$, it follows that $\left(f_{1}, \cdots, f_{n}, \varepsilon e\right)$ is positive, contrary to the hypothesis. Thus (5) is proved.

From (5), $q$ and $K$ can be separated by a hyperplane through the origin; that is, there exist real numbers $\theta_{1}, \cdots, \theta_{n+1}$ such that
(8) $\theta_{1}+\cdots+\theta_{n}<0 \leqq \theta_{1} \lambda_{1}+\cdots+\theta_{n+1} \lambda_{n+1}$ for all $\left(\lambda_{1}, \cdots, \lambda_{n+1}\right) \in K$.

Now, because of the minimality of $n$, no two of $f_{1}, \cdots, f_{n}$ can be $\mathscr{L}$ related, for if $f_{1}, f_{2}$ were $\mathscr{L}$ related we could have started with the $n-1$ elements $f_{1}+f_{2}, f_{3}, \cdots, f_{n}$ instead of $f_{1}, \cdots, f_{n}$. In particular, no two of them are equal, and (since $n \geqq 2$ ) none of them is $e$. Writing $f_{n+1}=e$, we have that $f_{1}, \cdots, f_{n+1}$ are all different. Thus, we can define a map $\psi: A \rightarrow R$, where $A=\left\{f_{1}, \cdots, f_{n+1}\right\}$, by: $\psi\left(f_{i}\right)=\theta_{i}, i=1, \cdots, n+1$.

We shall show that $\psi$ is positively $\mathscr{L} S$ linear. Suppose $x_{1}, \cdots, x_{m} \in A$ and that ( $\mu_{1} x_{1}, \cdots, \mu_{m} x_{m}$ ) is positive; we must show
(9) $\mu_{1} \psi\left(x_{1}\right)+\cdots+\mu_{m} \psi\left(x_{m}\right) \geqq 0$.

For each $i=1, \cdots, n+1$ write $J_{i}=\left\{j \mid l \leqq j \leqq m, x_{j}=f_{i}\right\}$; of course, $J_{i}$ may be empty. Let $\lambda_{i}=\Sigma\left\{\mu_{j} \mid j \in J_{i}\right\}$ (with the convention that an empty sum is 0 ). Then
(10) $\left(\mu_{1} x_{1}, \cdots, \mu_{m} x_{m}\right) \sim\left(\lambda_{1} f_{1}, \cdots, \lambda_{n+1} f_{n+1}\right)$,
because if $J_{i} \neq \varnothing$ we may apply an $\mathscr{L}$ equivalence (see (a)) to replace the terms $\mu_{j} x_{j}, j \in J_{i}, i$ fixed, by their sum, $\lambda_{i} f_{i}$. (If $J_{i}=\varnothing$, we appeal instead to (c).)

Thus, $\left(\lambda_{1} f_{1}, \cdots, \lambda_{n+1} f_{n+1}\right)$ is positive, that is, $\left(\lambda_{1}, \cdots, \lambda_{n+1}\right) \in K$. By (8), $\lambda_{1} \psi\left(f_{1}\right)+\cdots+\lambda_{n+1} \psi\left(f_{n+1}\right) \geqq 0$, and therefore,

$$
\begin{equation*}
\sum_{j} \mu_{j} \psi\left(x_{j}\right)=\sum_{i} \sum\left\{\mu_{j} \psi\left(f_{i}\right) \mid j \in J_{i}\right\}=\sum_{i} \lambda_{i} \psi\left(f_{i}\right) \geqq 0, \tag{7.7}
\end{equation*}
$$

as required.
By Lemma $7.4, \psi$ can be extended to a positively $\mathscr{L}$ linear map $\Psi: U \rightarrow R$. However, $\Psi$ can not be extended to a positive linear map $\Psi^{*}: L \rightarrow R$; for such an extension would make $\Psi^{*}\left(f_{1}+\cdots+f_{n}\right)=\theta_{1}+\cdots+\theta_{n}<0$, whereas $f_{1}+\cdots+f_{n} \geqq 0$. This completes the proof of Theorem 7.1.

Example 7.1. It is tempting to conjecture that, in Theorem 7.1, condition (iii) could be replaced by the simpler (and stronger) condition:
(iii') if $F \in \mathscr{F}$ and $\Sigma F \geqq 0$, then $F$ is positive.
However, this is not possible in general, even if $\mathscr{L}$ consists of only two linear spaces. This is shown by the following example. Take $L=m$, the space of all bounded sequences $x=\left\{\xi_{1}, \xi_{2}, \cdots\right\}$ of real numbers, ordered coordinatewise (that is, $x \geqq 0$ means $\xi_{n} \geqq 0$ for all $n=1,2, \cdots$ ). Put $e=\{1,1, \cdots\}$ (this in an order unit), $a=\{1,1 / 2, \cdots, 1 / n, \cdots\}, b=\left\{2,3 / 4, \cdots, 1 / n+1 / n^{2}, \cdots\right\}$, and let $L^{0}$ denote the subspace of $L$ consisting of the sequences which have only finitely many terms $\neq 0$. We define $L_{1}$ to be the $\operatorname{span} S\left(L^{0} \cup\{e\} \cup\{a\}\right)$, and $L_{2}$ to be $S\left(L^{0} \cup\{e\} \cup\{b\}\right) ; \mathscr{L}$ will consist of $L_{1}$ and $L_{2}$. The standing hypotheses in Section 6.1 are clearly satisfied.

We show that $\mathscr{L}$ has the positive extension property by verifying condition (iii) of Theorem 7.1. Suppose, then, that $x^{1}, x^{2}, \cdots, x^{n} \in U=L_{1} \cup L_{2}$, that $x^{1}+\cdots+x^{n} \geqq 0$, and that $\varepsilon>0$; we must show that ( $x^{1}, \cdots, x^{n}, \varepsilon e$ ) is positive. Writing $y^{1}=$ sum of the $x^{i}$ in $L_{1}, y^{2}=$ sum of the others, we see that this is equivalent to showing that $F=\left(y^{1}, y^{2}, \varepsilon e\right)$ is positive, where $y^{1} \in L_{1}$, $y^{2} \in L_{2}$, and $y^{1}+y^{2} \geqq 0$. We have $y^{1}=\lambda_{1} a+\mu_{1} e+p, y^{2}=\lambda_{2} b+\mu_{2} e+q$, where $p, q \in L^{0}$. Thus, $F \sim\left(\mu_{1} e, \lambda_{1} a, p, \mu_{2} e, \lambda_{2} b, q, \varepsilon e\right)$ and, therefore,
(1) $F \sim G=\left((\mu+\varepsilon) e, \lambda_{1} a, \lambda_{2} b, r\right)$, where $\mu=\mu_{1}+\mu_{2}$ and $r=p+q \in L^{0}$.

Take a positive integer $N$ so large that
(2) $\varepsilon N>4\left(\left|\lambda_{1}\right|+\left|\lambda_{2}\right|\right)$,
(3) the $n$th coordinate of $r$ is 0 for all $n \geqq N$.

Since $y^{1}+y^{2} \geqq 0$, we have
(4) $\mu+\lambda_{1} / n+\lambda_{2}\left(1 / n+1 / n^{2}\right) \geqq 0$ for all $n \geqq N$.

Now for each $x=\left\{\xi_{1}, \xi_{2}, \cdots\right\} \in L$, write

$$
\begin{align*}
x^{\prime} & =\left\{\xi_{1}, \cdots, \xi_{N}, 0,0, \cdots\right\} \\
x^{\prime \prime} & =x-x^{\prime}=\left\{0, \cdots, 0, \xi_{N+1}, \xi_{N+2}, \cdots\right\} \tag{7.8}
\end{align*}
$$

Note that if $x \in L_{j}, j=1,2$, then $x^{\prime} \in L^{0} \subset L_{1} \cap L_{2}, x^{\prime \prime} \in L_{j}$, so that $x, x^{\prime}, x^{\prime \prime}$ are $\mathscr{L}$ related for each $x \in U$. Thus, on writing

$$
\begin{equation*}
s=r+(\mu+\varepsilon) e^{\prime}+\lambda_{1} a^{\prime}+\lambda_{2} b^{\prime} \tag{7.9}
\end{equation*}
$$

we obtain
(5) $G \sim H=\left((\mu+\varepsilon) e^{\prime \prime}, \lambda_{1} a^{\prime \prime}, \lambda_{2} b^{\prime \prime}, s\right)$ where $s \in L^{0}$.

Further, from (3), we have
(6) the $n$th coordinate of $s$ is 0 for all $n \geqq N$.

Again, $e^{\prime \prime}$ and $a^{\prime \prime}$ are $\mathscr{L}$ related, as are $e^{\prime \prime}$ and $b^{\prime \prime}$. Thus, on writing $\delta=\varepsilon / 2$, we obtain
(7) $H \sim K=\left(\mu e^{\prime \prime}, \delta e^{\prime \prime}+\lambda_{1} a^{\prime \prime}, \delta e^{\prime \prime}+\lambda_{2} b^{\prime \prime}, s\right)$.

We assert that each of the terms of $K$ is $\geqq 0$, which then shows that $F$ is positive, as required. By making $n \rightarrow \infty$ in (4), we see that $\mu \geqq 0$ and, thus, $\mu e^{\prime \prime} \geqq 0$. From (2), it easily follows that $\delta e^{\prime \prime}+\lambda_{1} a^{\prime \prime}$ and $\delta e^{\prime \prime}+\lambda_{2} b^{\prime \prime}$ are $\geqq 0$. Finally, since $\Sigma K=\Sigma F \geqq 0$, we see that, in particular, each of the first $N$ coordinates of $s$ is $\geqq 0$. Thus, from (6), $s \geqq 0$, and $F$ is positive.

Nevertheless, (iii') does not hold. For consider $F=(-a, b)$. Clearly, $\Sigma F \geqq 0$, but we claim that $F$ is not positive. For otherwise we obtain, as before, $F \sim$ ( $y^{1}, y^{2}$ ), where $y^{j} \in L_{j}$ and $y^{j} \geqq 0, j=1,2$. As before, we write $y^{1}=\lambda_{1} a+$ $\mu_{1} e+p, y^{2}=\lambda_{2} b+\mu_{2} e+q$, where $p, q \in L^{0}$. Thus, the $n$th coordinate if $y^{1}+y^{2}$ is

$$
\begin{equation*}
\mu_{1}+\mu_{2}+\frac{\lambda_{1}+\lambda_{2}}{n}+O\left(n^{-2}\right) \tag{7.10}
\end{equation*}
$$

for large $n$. But $y^{1}+y^{2}=\Sigma F=b-a$, and its $n$th coordinate is $O\left(n^{-2}\right)$. This proves that $\mu_{1}+\mu_{2}=0=\lambda_{1}+\lambda_{2}$. However, the fact that $y^{1} \geqq 0$ shows that $\mu_{1} \geqq 0$, and $\lambda_{1} \geqq 0$ if $\mu_{1}=0$; similarly, $\mu_{2} \geqq 0$, and $\lambda_{2} \geqq 0$ if $\mu_{2}=0$. Thus, we have $\lambda_{1}=\lambda_{2}=\mu_{1}=\mu_{2}=0$, so that $b-a=p+q$. But the $n$th coordinate of $p+q$ is 0 if $n$ is large enough, while that of $b-a$ is not, giving the desired contradiction.

It would be interesting to know whether (iii') can replace (iii), in Theorem 7.1, when $L$ is finite dimensional.

Example 7.2. There is no implication between the extension property and the positive extension property, in general. That the extension property does not imply the positive extension property has been remarked in Theorem 6.1; see also Example 6.1 and the remark in Theorem 3.1. What is perhaps more surprising is that the positive extension property does not imply the extension property. This is shown by the following example.

Let $L, L^{0}, e$ have the same meaning as in Example 7.1, and select three elements $d^{1}, d^{2}, d^{3}$ of $L$, no two of them linearly dependent, such that their $n$th coordinates do not equal zero but tend to zero as $n \rightarrow \infty$, and $d^{1}+d^{2}+d^{3}=0$; this is easily done. Define $L_{j}=S\left(L^{0} \cup\{e\} \cup\left\{d^{j}\right\}\right), j=1,2,3 ; \mathscr{L}$ will consist of $L_{1}, L_{2}$ and $L_{3}$. Consider the following linear maps $\phi_{j}: L_{j} \rightarrow R, j=1,2,3$ : $\phi_{1}(x)=0$ for all $x \in L_{1}, \phi_{2}(x)=0$ for all $x \in L_{2}$, and, for each $p \in L^{0}$, $\phi_{3}\left(\lambda p+\mu e+v d^{3}\right)=v$.

It is easy to see that $L_{1} \cap L_{2}=L_{2} \cap L_{3}=L_{3} \cap L_{1}=S\left(L^{0} \cup\{e\}\right)$, so that the maps $\phi_{1}, \phi_{2}, \phi_{3}$ form a consistent system. However, they have no common linear extension because $d^{1}+d^{2}+d^{3}=0$, whereas $\phi_{1}\left(d^{1}\right)+\phi_{2}\left(d^{2}\right)+$ $\phi_{3}\left(d^{3}\right)=1$. Thus, the extension property fails here. Nevertheless, the positive extension property holds, as we now show by verifying condition (iii) of Theorem 7.1.

Suppose $H \in \mathscr{F}$ and $\Sigma H \geqq 0$. As before we may assume $H=\left(x^{1}, x^{2}, x^{3}\right)$, where $x^{j} \in L_{j}$; thus, $x^{j}=\lambda_{j} p^{j}+\mu_{j} e+v_{j} d^{j}$, where $p^{j} \in L^{0}, j=1,2,3$. Since $x^{1}+x^{2}+x^{3} \geqq 0$, we see (by considering the $n$th coordinates and making $n \rightarrow \infty)$ that $\mu_{1}+\mu_{2}+\mu_{3}=\mu$, say, $\geqq 0$. Given $\varepsilon>0$, put $\delta=\varepsilon / 3$, and pick $N$ large enough. We use the same decomposition $x=x^{\prime}+x^{\prime \prime}$ as in Example $7.1 ;$ thus $p^{j \prime \prime}=0, j=1,2,3$. The same reasoning as in Example 7.1 will show that

$$
\begin{equation*}
(H, \varepsilon e) \sim\left(\mu e^{\prime \prime}, \delta e^{\prime \prime}+\lambda_{1} d^{1 \prime \prime}, \delta e^{\prime \prime}+\lambda_{2} d^{2 \prime \prime}, \delta e^{\prime \prime}+\lambda_{3} d^{3 \prime \prime}, s\right), \tag{7.11}
\end{equation*}
$$

where each of the terms appearing on the right is $\geqq 0$, as required.

### 7.4. Dominated positive extension.

Theorem 7.2. If $\mathscr{L}$. has the positive extension property, then $\mathscr{L}$ has the dominated positive extension property.

By this we mean that if linear maps $\phi_{\alpha}: L_{\alpha} \rightarrow M$ and $\theta: L \rightarrow M$ are given such that $0 \leqq \phi(x) \leqq \theta(x)$ for each positive $x \in L_{\alpha}, \alpha \in \mathscr{A}$, and $0 \leqq \theta(x)$ for each positive $x \in L$, then the $\phi_{\alpha}$ can be extended to a linear map $\phi: L \rightarrow M$ such that $0 \leqq \phi(x) \leqq \theta(x)$ for each positive $x \in L$.

By Theorem 6.2 it is enough to prove that, if $x_{1}, \cdots, x_{n} \in U$ and $\Sigma_{i} x_{i} \geqq 0$ then $0 \leqq \Sigma_{i} \phi_{U}\left(x_{i}\right) \leqq \Sigma_{i} \theta\left(x_{i}\right)$. By Theorem 7.1, we know that ( $x_{1}, \cdots, x_{n}, \varepsilon e$ ) is positive for each $\varepsilon>0$, say $\left(x_{1}, \cdots, x_{n}, \varepsilon e\right) \sim\left(y_{1}, \cdots, y_{m}\right)$, where $y_{j} \in U$ and $y_{j} \geqq 0, j=1, \cdots, m$. By Lemma 3.1, we have

$$
\begin{equation*}
\sum_{i} \phi_{U}\left(x_{i}\right)+\varepsilon \phi_{U}(e)=\sum_{j} \phi_{U}\left(y_{j}\right) \geqq 0 \tag{7.12}
\end{equation*}
$$

and also

$$
\begin{equation*}
\sum_{i} \phi_{U}\left(x_{i}\right)+\varepsilon \phi_{U}(e) \leqq \sum_{j} \theta\left(y_{j}\right)=\sum_{i} \theta\left(x_{i}\right)+\varepsilon \theta(e) \tag{7.13}
\end{equation*}
$$

On making $\varepsilon \rightarrow 0$, the result follows.

Corollary 7.1. Suppose that $L$ has the positive extension property, and that linear maps $\theta_{*}, \theta^{*}: L \rightarrow M$ such that $\theta_{*}(x) \leqq \phi_{\alpha}(x) \leqq \theta^{*}(x)$ for each positive $x \in L_{\alpha}, \alpha \in \mathscr{A}$, are given. Then there exists a linear map $\phi: L \rightarrow M$, extending the $\phi_{a}$, such that $\theta_{*}(x) \leqq \phi(x) \leqq \theta^{*}(x)$ for each positive $x \in L$.

It would be interesting to know whether, conversely, the dominated positive extension property implies the positive extension property.

## 8. Countably additive measures

8.1. In this section we turn to the difficult but important problem of extending a consistent family of countably additive measures to a countably additive measure. We consider only positive finite measures.

There is a sense in which the preceding considerations suffice here, too, provided we are willing to enlarge the space $X$ on which the measure is to be defined. We have obtained criteria which are necessary and sufficient for there to exist a finitely additive common extension, say $m$, on a field $\mathscr{F}$ of measurable sets. Now there always exists a space $\tilde{X} \supset X$, with a $\sigma$-field $\mathscr{B}$ of subsets of $\tilde{X}$ and a countably additive measure $\tilde{m}$ on $\mathscr{B}$, such that each $F \in \mathscr{F}$ is of the form $X \cap B$ for some $B \in \mathscr{B}$, and for each such $B$ we have $\tilde{m}(B)=m(F)$. To see this, we take $\tilde{X}$ to be the Stone representation space of $F$ (see [13], p. 24); the points of $\tilde{X}$ are the ultrafilters of $F$, and we identify $x \in X$ with the ultrafilter consisting of those elements of $\mathscr{F}$ which contain $x$. Each $F \in \mathscr{F}$ determines a set $\tilde{F} \subset \tilde{X}$ (consisting of those ultrafilters on $\mathscr{F}$ to which $F$ belongs); the family of all such $\widetilde{F}$ is a field $\mathscr{F} *$ of subsets of $X$ which is isomorphic to $\mathscr{F}$ (under the correspondence $F \rightarrow \tilde{F}$; note that $F=\tilde{F} \cap X)$. We take $\mathscr{B}$ to be the Borel field generated by $\mathscr{F} *$, and we define $\tilde{m}$ on $\mathscr{F}^{*}$ by: $\tilde{m}(\tilde{F})=m(F), F \in \mathscr{F}$. It is easily seen (by using the Stone topology on $\tilde{X}$, which makes $X$ compact and $\mathscr{F}^{*}$ consist of the open-closed sets) that if $\widetilde{F}_{1}, \tilde{F}_{2}, \cdots$, are disjoint elements of $\mathscr{F}^{*}$ whose union is also in $\mathscr{F}^{*}$, then $\tilde{m}\left(U_{i} \tilde{F}_{i}\right)=\Sigma_{i} \tilde{m}\left(\tilde{F}_{i}\right) ;$ thus, by a standard theorem, $\tilde{m}$ has a countably additive extension to all of $\mathscr{B}$, as asserted. Under reasonable topological conditions on $X$, one can even embed $X$ topologically in $\tilde{X}$; see [2].
8.2. But it is more interesting to try to obtain a countably additive extension without enlarging the space. In this direction we have the following fairly simple theorem.

Theorem 8.1. Let $m_{\alpha}, \alpha \in \mathscr{A}$, be a positive, finite, finitely additive measure defined on a field $\mathscr{F}_{\alpha}$ of subsets of a fixed Hausdorff space $X$, and assume that each $m_{a}$ is "regular" in the sense that, for each $F \in \mathscr{F}_{a}$,

$$
\begin{equation*}
m_{\alpha}(F)=\sup \left\{m_{\alpha}(K) \mid K \in \mathscr{F}_{\alpha}, K \subset F, K \text { compact }\right\} . \tag{8.1}
\end{equation*}
$$

Suppose further that, whenever $f_{i}$ is an $m_{\alpha_{i}}$ measurable real valued step function on $X, i=1, \cdots, n$, such that $\Sigma_{i=1}^{n} f_{i}(x) \geqq 0$ for all $x \in X$, then

$$
\begin{equation*}
\sum_{i=1}^{n} \int_{X} f_{i} d m_{\alpha_{i}} \geqq 0 \tag{8.2}
\end{equation*}
$$

Then there exists a countably additive (positive, finite) measure $\mu$, defined on the Borel field generated by $\cup_{\alpha} F_{\alpha}$, which extends all the $m_{\alpha}$ simultaneously.

Remark. Though we have not explicitly required the $m_{\alpha}$ to be countably additive, or consistent, it follows easily from (8.2) that they are consistent, and from (8.1) that they have countably additive extensions.

Proof. By Theorem 6.1, the measures $m_{\alpha}$ have a common extension to a finitely additive (positive, finite) measure $\mu$ which, for convenience, we restrict to the field $\mathscr{F}$ generated by $\bigcup_{\alpha} \mathscr{F}_{\alpha}$. We need only show that $\mu$ can be extended to a countably additive measure on $\mathscr{B}$, the $\sigma$-field generated by $\mathscr{F}$. The first step is to show that, given $F \in \mathscr{F}$ and $\varepsilon>0$, there exist a compact $K \subset F$ and open $G \supset F$ such that
(1) $K, G \in \mathscr{F}$ and $\mu(K)+\varepsilon>\mu(F)>\mu(G)-\varepsilon$.

Suppose first that $F$ is of the form $A_{1} \cap A_{2} \cap \cdots \cap A_{n}$, where $A_{i} \in \mathscr{F}_{\alpha_{i}}$, $i=1, \cdots, n$. Using (8.1) (applied both to $A_{i}$ and to its complement), we find compact $K_{i} \in \mathscr{F}_{\alpha_{i}}$ and open $G_{i} \in \mathscr{F}_{\alpha_{i}}$ such that

$$
\begin{equation*}
\mu\left(K_{i}\right)+\frac{\varepsilon}{n}>\mu\left(A_{i}\right)>\mu\left(G_{i}\right)-\frac{\varepsilon}{n} . \tag{8.3}
\end{equation*}
$$

Put $K=\cap_{i} K_{i}, G=\cap_{i} G_{i}$; it is easy to see that (1) holds in this case.
But if $F$ is an arbitrary member of $\mathscr{F}$, it is well known that $F$ is expressible as a union of disjoint sets $F_{1}, \cdots, F_{m}$ each of the form just considered. Since (1) holds for each $F_{j}$, it follows that (1) holds for $F$.

The rest is standard, but for completeness we sketch the argument. It is enough to prove that if $F_{1}, F_{2}, \cdots$ are disjoint members of $\mathscr{F}$ whose union $F$ is also in $\mathscr{F}$, then $\mu(F)=\Sigma_{n} \mu\left(F_{n}\right)$. One inequality is trivial. For the other, given $\varepsilon>0$, take open $G_{n} \in \mathscr{F}$ such that $G_{n} \supset F_{n}$ and $\mu\left(G_{n} \backslash F_{n}\right)<\varepsilon / 2^{n+1}$, and take compact $K \in \mathscr{F}$ such that $K \subset F$ and $\mu(F \backslash K)<\varepsilon / 2$. Then $K$ is covered by the sets $G_{1}, G_{2}, \cdots$, and, hence, $K \subset G_{1} \cup \cdots \cup G_{n}$ for some $n$. It follows that $\mu(K) \leqq$ $\Sigma_{i=1}^{n} \mu\left(G_{i}\right)$, whence $\mu(F) \leqq \Sigma_{n} \mu\left(F_{n}\right)+\varepsilon$, for all $\varepsilon>0$.
8.3. We remark, in conclusion, that in the case of marginals, and with only finitely many measures involved, Kellerer ([5], Satz 2.2) has shown (by explicit construction) that a consistent family of countably additive signed measures always has a countably additive common extension. It would be desirable to have an abstract theorem which includes this.

## REFERENCES

[1] E. S. Andersen and B. Jessen, "On the introduction of measures in infinite product sets," Danske Vid. Selsk. Mat.-Fys. Medd., Vol. 25 (1948), pp. 3-7.
[2] D. N. Dudin, "Completion of an additive measure to a $\sigma$-additive measure by an extension of the space," Math. Notes, Vol. 3 (1968), pp. 42-44.
[3] D. L. Guy, "Common extension of finitely additive probability measures," Portugal Math., Vol. 20 (1961), pp. 1-5.
[4] C. Ionescu Tulcea, "Mesures dans les espaces produits," Atti Accad. Naz. Lincei Rend. Cl. Sci. Fis. Mat. Natur., Vol. 7 (1947), pp. 208-211.
[5] H. G. Kellerer, "Masstheoretische Marginalprobleme," Math. Ann., Vol. 153 (1964), pp. 168-198.
[6] —_, "Marginalprobleme für Funktionen," Math. Ann., Vol. 154 (1964), pp. 147-156.
[7] E. Marczewski, "Measures in almost independent fields," Fund. Math., Vol. 38 (1951), pp. 217-229.
[8] E. Marczewski and C. Ryll-Nardzewski, "Remarks on compactness and direct products of measure spaces," Fund. Math., Vol. 40 (1953), pp. 165-170.
[9] L. Nachbin, Topology and Order, Princeton, Van Nostrand, 1965.
[10] J. von Neumann, "Zur allgemeinen Theorie des Masses," Fund. Math., Vol. 13 (1929), pp. 73-116.
[11] A. L Peressini, Ordered Topological Spaces, New York, Harper and Row, 1967.
[12] V. Pták, "Simultaneous extensions of two functionals," Czechoslovak Math.J., Vol. 3 (1969), pp. 3-19.
[13] R. Sikorski, Boolean Algebras, Berlin. Springer Verlag, 1964, (2nd ed.).
[14] V. Strassen, '"The existence of probability measures with given marginals," Ann. Math. Statist., Vol. 36 (1965), pp. 423-439.
[15] A. E. Taylor, Functional Analysis, New York, Wiley, 1958.
[16] N. N. Vorob'ev, "Consistent families of measures and their extensions," Theor. Probability Appl., Vol. 7 (1962), pp. 147-162.

