

Diameter-preserving maps on various classes of function spaces

by

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Abstract. Under some mild assumptions, non-linear diameter-preserving bijections between (vector-valued) function spaces are characterized with the help of a well-known theorem of Ulam and Mazur. A necessary and sufficient condition for the existence of a diameter-preserving bijection between function spaces in the complex scalar case is derived, and a complete description of such maps is given in several important cases.

Introduction. Several papers on diameter-preserving linear bijections of function spaces have appeared in recent years (see, for example, [GM], [GU], [S], [RR]). The present paper is a contribution to this circle of ideas, its principal motivation coming from an attempt to clarify and extend some of the results in [RR] in response to an interesting question raised there.

To give a summary in this introduction to our main results, let us first explain briefly the notations and terminology that we use (unexplained terms and notations will be found in [RR]). K , L , and S will generally denote compact convex sets in a Hausdorff locally convex topological vector space X over \mathbb{C} or \mathbb{R} , as the context will make clear. The (non-empty) set of extreme points of K will be denoted by $\text{ext}(K)$, and $A(K)$ will stand for the space of complex-valued continuous affine functions on K endowed with the sup-norm. Our general references for facts concerning compact convex sets are [A] and [AE].

Let Q be a compact Hausdorff space, and X a Banach space. Then the diameter of $f \in C(Q, X)$ is defined by $d(f) = \sup_{x,y \in Q} \|f(x) - f(y)\|$. If $A_1 \subseteq C(Q_1, X_1)$ and $A_2 \subseteq C(Q_2, X_2)$ are function spaces, i.e. sup-norm closed subspaces, containing the constant functions, and separating points, then a (not necessarily linear) map $T : A_1 \rightarrow A_2$ is called a *diameter-preserving* (*d-preserving*) *bijection* if T is a 1-1 map of A_1 onto A_2 , and $d(a) = d(Ta)$ for all $a \in A_1$.

We begin, in Section 1, by discussing non-linear d -preserving bijections between (vector-valued) function spaces A_1 and A_2 of the kind described in the last paragraph and show, with the help of a well-known theorem of Mazur and Ulam (see [MU] and [B]), that under certain mild assumptions, T can be characterized as $\bar{T} + \varphi$ where $\bar{T} : A_1 \rightarrow A_2$ is a *linear* d -preserving bijection, and $\varphi : A_1 \rightarrow X_2$ is a function, generally non-linear, with well-defined properties. We give several examples of such maps in Section 2.

In [RR, Theorem 1, p. 5], it was found that all linear d -preserving bijections between the spaces $A(K)$ and $A(L)$ are essentially induced by affine homeomorphisms between K and L when K, L are compact convex sets with all points of $\text{ext}(K)$ and $\text{ext}(L)$ split. It was asked there whether the same conclusion could be derived by assuming the “splittability” property only for $\text{ext}(K)$. We exhibit in Section 3 a very simple two-dimensional counterexample to this question involving a simplex K and a hexagon S for which a linear d -preserving bijection exists between $A_{\mathbb{R}}(K)$ and $A_{\mathbb{R}}(S)$. [Here the subscript \mathbb{R} signifies real-valued functions.] The problem therefore naturally arises of characterizing compact convex sets K_1 and K_2 for which linear d -preserving bijections exist between $A(K_1)$ and $A(K_2)$. We consider the question more generally for function spaces $A_i \subseteq C(Q_i)$, $i = 1, 2$, and show that the existence of an affine homeomorphism between certain compact convex sets in A_1^* and A_2^* —these sets being associated in a very natural manner with the state spaces of A_1 and A_2 —is necessary and sufficient for this purpose. (For example, for the spaces $A_{\mathbb{R}}(K)$ and $A_{\mathbb{R}}(S)$ as above, the condition is that $K - K$ must be affinely homeomorphic to $S - S$.) We pursue this further by describing in Theorem 9 (in Section 4) the precise form of these d -preserving bijections from $A(K)$ to $A(S)$ where all the points in $\text{ext}(K)$ are split and S is any compact convex set when the aforementioned condition is satisfied.

The discussion in Section 3 leads naturally to the question of solvability in S (a compact convex set) of the equation $S - S = K - K$ where K is a given compact convex set with all points of $\text{ext}(K)$ split. The question in general appears too difficult, but we can exhibit, in \mathbb{R}^n , the class of polytopes S for which $S - S = K - K$ when K is a given simplex in \mathbb{R}^n . This is the content of Section 5.

1. d -Preserving maps on function spaces. Assume that E is a normed linear space. Let $C(X, E)$ be the space of all continuous E -valued functions on X . For $v \in E$, we also denote by v the constant function, $v(x) = v$ for all $x \in X$. Thus, vectors and constant functions are denoted in the same way; the meaning can be determined from the context. Also, we use E to denote the closed subspace of $C(X, E)$ consisting of all the constant functions.

A subspace $A \subseteq C(X, E)$ is a (vector-valued) *function space* if A is closed in the sup-norm, $E \subseteq A$, and A separates points of X . For $a \in A$, let $[a] = a + E$ be the residue class of a in the quotient space A/E . We remark that A is linearly isomorphic to $(A/E) \oplus E$. For let $\psi : A \rightarrow (A/E) \oplus E$ be defined by $\psi(a) = [a] \oplus a(x_0)$ where $x_0 \in X$ is fixed. Then ψ is 1-1 since if $a, b \in A$ with $\psi(a) = \psi(b)$, then $[a] = [b]$, so a and b differ by a constant, and also $a(x_0) = b(x_0)$. Thus, $a = b$. The linear map ψ maps A onto $(A/E) \oplus E$, since given $[a] \oplus v$, setting $c = a + (v - a(x_0))$, we have $\psi(c) = [a] \oplus v$.

The *diameter* of $a \in A$ is defined by $d(a) \equiv \sup_{x,y \in X} \|a(x) - a(y)\|$.

Note that $d(a) = 0$ if, and only if, a is a constant function. For $[a] \in A/E$, define $d[a] = d(a)$. Then $[a] \mapsto d[a]$ is a norm on A/E which we call the *d-norm*.

PROPOSITION 1. Assume that $A_1 \subseteq C(X_1, E_1)$ and $A_2 \subseteq C(X_2, E_2)$ are function spaces. Make the linear identifications $A_1 \approx (A_1/E_1) \oplus E_1$ and $A_2 \approx (A_2/E_2) \oplus E_2$ as above, using the maps ψ_1 and ψ_2 , $\psi_k(a_k) = [a_k] \oplus a_k(x_k)$ where $x_k \in X_k$ are fixed points, $a_k \in A_k$ for $k = 1, 2$. Assume that $D : A_1/E_1 \rightarrow A_2/E_2$ and $J : E_1 \rightarrow E_2$ are linear bijections, and that D is an isometry with respect to the d-norms. Define $D \oplus J : (A_1/E_1) \oplus E_1 \rightarrow (A_2/E_2) \oplus E_2$ by $(D \oplus J)([a] \oplus v) = D[a] \oplus J(v)$. Finally, define $\bar{D} : A_1 \rightarrow A_2$ by $\bar{D}(a) = \psi_2^{-1}(D \oplus J)\psi_1(a)$. Then \bar{D} is a linear bijection which is d-preserving and has the property that the induced map $\tilde{\bar{D}} : A_1/E_1 \rightarrow A_2/E_2$ defined by $\tilde{\bar{D}}([a]) = [\bar{D}(a)]$ is D .

Proof. That $\bar{D}(a)$ is a linear bijection is clear, since ψ_2^{-1} , $D \oplus J$ and ψ_1 are all linear bijections. Let $a \in A_1$. Then by a straightforward computation,

$$\bar{D}(a) = b + -b(x_2) + J(a(x_1)) \quad \text{where } [b] = D[a].$$

Thus, $d(\bar{D}(a)) = d(b) = d[b] = d[D[a]] = d[a] = d(a)$, so \bar{D} is d-preserving. ■

REMARK. When $J : E_1 \rightarrow E_2$ is a bijection, but *not* linear, and $D \oplus J$ and \bar{D} are defined as above, then \bar{D} is still a bijection and a d-preserving function which is not linear.

Let A_1 and A_2 be function spaces as above, and assume that $T : A_1 \rightarrow A_2$ is a function (not necessarily linear). Now we prove a result which shows that with fairly weak assumptions on T plus the assumption that E_1 and E_2 are linearly isomorphic, there exists a linear bijection $\bar{T} : A_1 \rightarrow A_2$ with \bar{T} d-preserving. The key tool here is the *Ulam-Mazur Theorem* [MU, B] which we state for the convenience of the reader: Assume that $(E_1, \| \cdot \|_1)$ and $(E_2, \| \cdot \|_2)$ are normed linear spaces and that $D : E_1 \rightarrow E_2$ is a bijection with the properties: (a) $\|D(x) - D(y)\|_2 = \|x - y\|_1$ for all $x, y \in E_1$, (b) $D(0) = 0$, and in the case of complex scalars, (c) $D(ia) = iD(a)$ for all $a \in E_1$. Then D is linear.

THEOREM 2. Assume that $T : A_1 \rightarrow A_2$ is a function with the properties:

- (1) T is a bijection;
- (2) T has the property $d(T(a) - T(b)) = d(a - b)$ for all $a, b \in A_1$;
- (3) $T(0) = 0$;
- (4) [in the case of complex scalars] $T(ia) = iT(a)$ for all $a \in A_1$.

Also, assume that there exists $J : E_1 \rightarrow E_2$ such that J is a linear bijection. Then $T(a) = \bar{T}(a) + \varphi(a)$ where $\bar{T} : A_1 \rightarrow A_2$ is a d -preserving linear bijection with $\bar{T}(v) = J(v)$, $v \in E_1$, and $\varphi : A_1 \rightarrow E_2$ is a function with the properties:

- (i) $\varphi(0) = 0$;
- (ii) [in the complex scalar case] for all $a \in A$, $\varphi(ia) = i\varphi(a)$;
- (iii) the map $a \mapsto a + J^{-1}(\varphi(a))$ from A_1 to A_1 is 1-1;
- (iv) given $a \in A_1$, the equation $J(x) + \varphi(a+x) = 0$ is solvable for $x \in E_1$.

Conversely, assume $\bar{T} : A_1 \rightarrow A_2$ is a d -preserving linear bijection, and define $J : E_1 \rightarrow E_2$ by $J(v) = \bar{T}(v)$, $v \in E_1$. Further, assume that $\varphi : A_1 \rightarrow E_2$ is a function with properties (i)–(iv) above. Then $T(a) = \bar{T}(a) + \varphi(a)$ is a (possibly non-linear) map from A_1 onto A_2 having the properties (1)–(4) listed in the theorem.

Proof. Assume that $T : A_1 \rightarrow A_2$ has properties (1)–(4). Define $\tilde{T} : A_1/E_1 \rightarrow A_2/E_2$ in the obvious way: $\tilde{T}([a]) = [Ta]$. From properties (1) and (2), \tilde{T} is a bijection and has property (a) above (in the statement of the Ulam–Mazur Theorem). Also, \tilde{T} inherits properties (3) and (4). Therefore the Ulam–Mazur Theorem applies, so \tilde{T} is linear. Now using the construction in Proposition 1, define $\bar{T}(a) = \psi_2^{-1}(\tilde{T} \oplus J)\psi_1(a)$. By Proposition 1, \bar{T} is a linear bijection which is d -preserving. For $a \in A_1$, $\bar{T}(a) = T(a) + c$ for some $c \in E_2$. Define $\varphi(a) = T(a) - \bar{T}(a) \in E_2$. Thus, for all $a \in A_1$, $T(a) = \bar{T}(a) + \varphi(a)$. It follows from the definition of \bar{T} that $\bar{T}(a) = J(a)$ when $a \in E_1$. That φ satisfies (i) and (ii) is clear. We verify that (iii) and (iv) hold:

(iii) The map $a \mapsto a + J^{-1}(\varphi(a))$ from A_1 to A_1 is 1-1: Suppose $a + J^{-1}(\varphi(a)) = b + J^{-1}(\varphi(b))$. Then $v = b - a \in E_1$, and $v = J^{-1}(\varphi(a) - \varphi(b))$. Thus, $\bar{T}(v) = J(v) = \varphi(a) - \varphi(b)$. Therefore, $\bar{T}(b) - \bar{T}(a) = \varphi(a) - \varphi(b)$, so $T(a) = T(b)$. Then since T is 1-1, $a = b$.

(iv) Given $a \in A_1$, the equation $J(x) + \varphi(a+x) = 0$ is solvable for $x \in E_1$: Choose b such that $T(b) = \bar{T}(a)$ (T is surjective). Then $0 = d(T(a) - T(b)) = d(a - b)$, so $b = a + v$ for some $v \in E_1$. Thus, $\bar{T}(a) = T(b) = T(a + v) = \bar{T}(a + v) + \varphi(a + v) = \bar{T}(a) + \bar{T}(v) + \varphi(a + v) = \bar{T}(a) + J(v) + \varphi(a + v)$. Therefore, $J(v) + \varphi(a + v) = 0$.

Now we do the converse. Assume that \bar{T} , J , and φ are as stated in the last paragraph of the theorem. Also, define $T : A_1 \rightarrow A_1$ as above: $T(a) =$

$\overline{T}(a) + \varphi(a)$. That T has properties (2), (3), and (4) follows immediately. We prove that T is bijective.

T is 1-1: Assume $T(a) = T(b)$, so

$$(1) \quad \overline{T}(a) + \varphi(a) = \overline{T}(b) + \varphi(b).$$

Then $\overline{T}(a - b) = \varphi(b) - \varphi(a)$. Therefore, $a - b = x \in E_1$. Then $\overline{T}(a) = \overline{T}(b) + J(x)$ (since $J = \overline{T}$ on E_1). Substituting this equality into (1), we have $\overline{T}(b) + J(x) + \varphi(a) = \overline{T}(b) + \varphi(b)$, so $J(x) = \varphi(b) - \varphi(a)$. Thus, $J(x) + \varphi(a) = \varphi(b) = \varphi(a - x)$. Then $a = b + x = b + J^{-1}(\varphi(b) - \varphi(a))$. Finally, this implies $a + J^{-1}(\varphi(a)) = b + J^{-1}(\varphi(b))$. Applying (iii), it follows that $a = b$.

T maps onto A_2 : Assume $c \in A_2$. We want to find $a \in A_1$ such that $\overline{T}(a) + \varphi(a) = c$. There exists $b \in A_1$ such that $c = \overline{T}(b)$. Now $\overline{T}(a) + \varphi(a) = \overline{T}(b)$ implies $\overline{T}(b - a) = \varphi(a)$. Therefore, $b - a \in E_1$. Setting $x = b - a$, we have $\overline{T}(b) = \overline{T}(a) + J(x)$, so $\overline{T}(a) + \varphi(a) = \overline{T}(a) + J(x)$. Thus, $\varphi(a) = J(x)$. Then $\varphi(b - x) = J(x)$ and $J(-x) + \varphi(b - x) = 0$. Finally, applying (iv), we can solve this last equation for x . Therefore, $a = b - x$ will be a solution of $T(a) = \overline{T}(a) + \varphi(a) = c$. ■

REMARK. \overline{T} , the linear part of T , has been characterized in [RR, Theorem 2 and Proposition 2] for some vector-valued function spaces.

The assumption in Theorem 2 of the existence of the linear bijection $J : E_1 \rightarrow E_2$ is not very restrictive. First, it is not assumed that J is continuous, so the existence of J is equivalent to E_1 and E_2 having (algebraic) bases of the same cardinality.

We use $\text{card}(S)$ to denote the cardinality of a set S , and set $c = \text{card}(\mathbb{R})$. Also, the dimension of a linear space E is denoted by $\text{dim}(E)$. Here is a useful known result:

$$\text{If } \text{dim}(E) \geq c, \text{ then } \text{card}(E) = \text{dim}(E) \quad [\text{LT, Problem 2, p. 43}].$$

Note that because of the hypotheses in Theorem 2, $T(E_1) = E_2$, and T is 1-1. It follows that $\text{card}(E_1) = \text{card}(E_2)$. Thus, if $\text{dim}(E_1) \geq c$ and $\text{dim}(E_2) \geq c$, then

$$\text{dim}(E_1) = \text{card}(E_1) = \text{card}(E_2) = \text{dim}(E_2).$$

CONCLUSION. *If $\text{dim}(E_1) \geq c$ and $\text{dim}(E_2) \geq c$, then there exists a linear bijection of E_1 onto E_2 .*

Our main concern in this paper being with function spaces whose members take values in Banach spaces, it is pertinent to point out that, as a consequence of the Baire Category Theorem, the algebraic dimension of a Banach space is either finite or uncountable ($\geq c$).

Now assume the dimensions of E_1 and E_2 are both finite, $E_1 \approx \mathbb{R}^m$ and $E_2 \approx \mathbb{R}^n$. Assume that T and T^{-1} are continuous, so that T is a homeo-

morphism of $E_1 \approx \mathbb{R}^m$ onto $E_2 \approx \mathbb{R}^n$. Then the Invariance of Dimension Theorem [D, p. 359] implies that $m = n$.

CONCLUSION. *If E_1 and E_2 are both finite-dimensional and T is a homeomorphism, then there exists a linear bijection of E_1 onto E_2 .*

2. Examples of non-linear d-preserving maps. In this section we present several examples of non-linear d-preserving maps T which satisfy the hypotheses (1)–(4) of Theorem 2. Of course, by that theorem, T must be the sum of a linear d-preserving bijection and a non-linear part.

First we give an example when $E_1 = E_2 = \mathbb{C}$. Let A be a function space, $A \subseteq C(X)$. Note the relations for a function $a \in A$: $a = \operatorname{Re}(a) + i \operatorname{Im}(a)$; $ia = i \operatorname{Re}(a) - \operatorname{Im}(a)$; $\operatorname{Re}(ia) = -\operatorname{Im}(a)$; $\operatorname{Im}(ia) = \operatorname{Re}(a)$. Now set

$$\varphi(a) = \sup(\operatorname{Re}(a)) + i \sup(\operatorname{Im}(a)) - \sup(-\operatorname{Re}(a)) - i \sup(-\operatorname{Im}(a)).$$

For $a \in A$, $\varphi(ia) = i\varphi(a)$. Proof:

$$\begin{aligned} \varphi(ia) &= \sup(\operatorname{Re}(ia)) + i \sup(\operatorname{Im}(ia)) - \sup(-\operatorname{Re}(ia)) - i \sup(-\operatorname{Im}(ia)) \\ &= \sup(-\operatorname{Im}(a)) + i \sup(\operatorname{Re}(a)) - \sup(\operatorname{Im}(a)) - i \sup(-\operatorname{Re}(a)) \\ &= i[\sup(\operatorname{Re}(a)) - \sup(-\operatorname{Re}(a)) + i \sup(\operatorname{Im}(a)) - i \sup(-\operatorname{Im}(a))] \\ &= i\varphi(a). \end{aligned}$$

NOTE. For $b \in A$, $\mu \in \mathbb{C}$, $\varphi(b + \mu) = \varphi(b) + 2\mu$. Proof:

$$\begin{aligned} \varphi(b + \mu) &= [\sup(\operatorname{Re}(b)) + \operatorname{Re}(\mu)] + i[\sup(\operatorname{Im}(b)) + \operatorname{Im}(\mu)] \\ &\quad - [\sup(-\operatorname{Re}(b)) - \operatorname{Re}(\mu)] - i[\sup(-\operatorname{Im}(b)) - \operatorname{Im}(\mu)] \\ &= \varphi(b) + 2\mu. \end{aligned}$$

Now define $T(a) = a + \varphi(a)$. Then T maps onto A . Proof: Given $b \in A$, we want $a \in A$ such that $b = a + \varphi(a)$. Set $a = b + \mu$ where $\mu \in \mathbb{C}$. Then we want that $b = (b + \mu) + \varphi(b + \mu) = b + \mu + \varphi(b) + 2\mu$ (from Note). Solving, we see that $3\mu = -\varphi(b)$. Letting $\mu = -\frac{1}{3}\varphi(b)$, we find that $T(a) = b$.

T is 1-1. Proof: Suppose that $a, b \in A$ with $a + \varphi(a) = b + \varphi(b)$, so $a - b = \varphi(b) - \varphi(a)$. Then $a = b + \mu$ where $\mu \in \mathbb{C}$. Therefore, $\mu = \varphi(b) - \varphi(a) = \varphi(a - \mu) - \varphi(a) =$ (from Note) $\varphi(a) - 2\mu - \varphi(a) = -2\mu$, so $\mu = 0$.

Thus, $T(a) = a + \varphi(a)$ is non-linear, but satisfies the hypotheses of Theorem 2.

When the scalar field is \mathbb{R} , there is a much simpler example. Assume that A is a function space of \mathbb{R} -valued functions. The reader can check that $T(a) = a + \sup(a)$ or $T(a) = a + \max[\sup(a), 0]$, $a \in A$, is a non-linear bijection of A onto A that satisfies the hypotheses of Theorem 2. The latter definition of T shows that the function $\varphi(x)$ in the statement of Theorem 2 is not generally surjective.

Now we give an example in the vector-valued case. Let Y be a compact Hausdorff space, and set $E_1 = E_2 = C(Y)$. We construct a function J with the properties:

- (i) J is a bijection of $C(Y)$ onto $C(Y)$;
- (ii) $J(0) = 0$ and $J(if) = iJ(f)$ for all $f \in C(Y)$;
- (iii) J is not linear.

For $f \in C(Y)$, define

$$J(f) = [\operatorname{Re}(f)]^3 - [\operatorname{Im}(f)]^3 - i[\operatorname{Re}(if)]^3 + i[\operatorname{Im}(if)]^3.$$

[In the real scalar case, $J(f) = f^3$ will work.] (ii) follows from a straightforward computation, and (iii) is clear. We verify (i).

J is 1-1. Proof: Suppose $f_k \in C(Y)$, $f_k = u_k + iv_k$, where u_k, v_k are \mathbb{R} -valued functions in $C(Y)$, $k = 1, 2$. Note that $if_k = -v_k + iu_k$, $k = 1, 2$. Suppose that $J(f_1) = J(f_2)$. Then $u_1^3 - v_1^3 = u_2^3 - v_2^3$, and $u_1^3 + v_1^3 = u_2^3 + v_2^3$. Therefore, $2u_1^3 = 2u_2^3$, so $u_1 = u_2$. Then $v_1^3 = v_2^3$, so $v_1 = v_2$.

J maps onto $C(Y)$. Proof: Assume that $g = u + iv \in C(Y)$, u, v \mathbb{R} -valued. Set

$$f = \left[\frac{u+v}{2} \right]^{1/3} + i \left[\frac{v-u}{2} \right]^{1/3}, \quad \text{so} \quad if = i \left[\frac{u+v}{2} \right]^{1/3} - \left[\frac{v-u}{2} \right]^{1/3}.$$

Then

$$J(f) = \left[\frac{u+v}{2} \right] - \left[\frac{v-u}{2} \right] + i \left[\frac{v-u}{2} \right] + i \left[\frac{u+v}{2} \right] = u + iv = g.$$

Now let $A \subseteq C(X, C(Y))$ be a function space. Let I denote the identity map on A . Use Proposition 1 to define $T = \psi^{-1}(I \oplus J)\psi$ where ψ is the map defined in the discussion just prior to Proposition 1. Then T is a non-linear bijection with properties (1)–(4) in Theorem 2.

3. d-Preserving maps on complex-valued function spaces. The real-scalar version of Theorem 1 in [RR] is:

Let K and S be compact convex sets, both of which have the property that every extreme point is a split face. If $D : A(K) \rightarrow A(S)$ is a d -preserving linear bijection, then there exists an affine homeomorphism $\tau : S \rightarrow K$ and a functional α defined on $A(K)$ such that for all $a \in A(K)$,

$$D(a) = c(a \circ \tau) + \alpha(a), \quad \text{where } c = \pm 1, \text{ and } \alpha(1) \neq -c.$$

A question raised in [RR, Remark 1] is: Does the same result hold if the hypothesis, “every extreme point is a split face”, is assumed only for K ? We now give an example which answers this question in the negative.

In \mathbb{R}^2 , let $K = \operatorname{co}\{(1, 0), (-1, 0), (0, 1)\}$. The three extreme points of K are split faces in the sense defined in [A] because K is a *simplex*, and for each extreme point x of K , $\{x\}'$ (\equiv the complementary set of $\{x\}$) is a face,

and every point $p \in K$ can be written uniquely as $p = \alpha x + (1 - \alpha)y$, where $y \in \{x\}'$ and $0 \leq \alpha \leq 1$. [Strictly speaking, to use the definition given in [A], one should regard $K \subseteq \{\varphi \in A(K)^* : \varphi(1) = 1\}$. The analysis could be done in $A(K)^*$. But this seems an unnecessary and technical approach to the elementary example under consideration.]

Now let $S = \frac{1}{2}(K - K)$. It is easily checked that S is a hexagon; in fact $S = \text{co}\{x_1, x_2, x_3, -x_1, -x_2, -x_3\}$ where $x_1 = (1, 0)$, $x_2 = (1/2, 1/2)$, and $x_3 = (-1/2, 1/2)$.

Thus, there is no affine homeomorphism of S onto K . Using the identity

$$\frac{1}{2}(z - w) - \frac{1}{2}(x - y) = \left(\frac{z + y}{2}\right) - \left(\frac{w + x}{2}\right),$$

it is easy to check that $K - K = S - S$. Then it follows from Corollary 7 of this paper that there exists a d-preserving linear bijection of $A(K)$ onto $A(S)$.

Now we investigate function spaces $A_i \subseteq C(Q_i)$, $i = 1, 2$, to determine necessary and sufficient conditions for the existence of a linear d-preserving $T : A_1 \rightarrow A_2$. First we give some relevant definitions and results which will be needed for this analysis.

Let S be a compact convex set which is symmetric ($s \in S \Rightarrow -s \in S$). When the scalar field is \mathbb{R} , define $A_0(S) = \{f \in A(S) : f(0) = 0\}$. When the scalar field is \mathbb{C} , assume $(s \in S, \alpha \in \mathbb{C}, |\alpha| = 1) \Rightarrow \alpha s \in S$. In this case define $A_0(S) = \{f \in A(S) : f(is) = if(s) \text{ for all } s \in S\}$. Note that if $f \in A_0(S)$, then $f(i0) = if(0)$, so $f(0) = 0$.

PROPOSITION 3. *Let X be a Banach space. Assume that $f \in A_0(X_1^*)$ (here X_1^* is the closed unit ball of the dual of X , and the topology on X_1^* is the w^* -topology). Extend f to \tilde{f} by*

$$\tilde{f}(\varphi) = \|\varphi\|f(\varphi/\|\varphi\|), \quad \varphi \in X^*.$$

Then \tilde{f} is a w^ -continuous linear functional on X^* .*

Proof. We do the complex scalar case, so it is assumed that $f(i\varphi) = if(\varphi)$ for all $\varphi \in X_1^*$. Also, note that f has the properties: $f(-\varphi) = -f(\varphi)$; and $(0 \leq \alpha \leq 1, \varphi \in X_1^*) \Rightarrow \alpha f(\varphi) = f(\alpha\varphi)$. To be proved:

- (a) $\tilde{f}(\varphi + \psi) = \tilde{f}(\varphi) + \tilde{f}(\psi)$ for all $\varphi, \psi \in X^*$;
- (b) $\tilde{f}(\alpha\varphi) = \alpha\tilde{f}(\varphi)$ for all $\alpha \in \mathbb{C}, \varphi \in X^*$.

Assume that both (a) and (b) hold. Now $\ker(\tilde{f}) \cap X_1^* = \{\text{the zero set of } f\}$, which is w^* -closed. Therefore by the Kreĭn-Shmul'yan Theorem (see [DS] or [LT]), \tilde{f} is w^* -continuous.

Now we prove (a). Let $\varphi, \psi \in X^* \setminus \{0\}$. Since $\|\varphi + \psi\|/(\|\varphi\| + \|\psi\|) \leq 1$,

$$\left(\frac{\|\varphi + \psi\|}{\|\varphi\| + \|\psi\|}\right) f\left(\frac{\varphi + \psi}{\|\varphi + \psi\|}\right) = f\left(\frac{\varphi + \psi}{\|\varphi\| + \|\psi\|}\right).$$

Therefore

$$\begin{aligned} \tilde{f}(\varphi + \psi) &= \|\varphi + \psi\| f\left(\frac{\varphi + \psi}{\|\varphi + \psi\|}\right) = (\|\varphi\| + \|\psi\|) f\left(\frac{\varphi + \psi}{\|\varphi\| + \|\psi\|}\right) \\ &= (\|\varphi\| + \|\psi\|) \left[\frac{\|\varphi\|}{\|\varphi\| + \|\psi\|} f\left(\frac{\varphi}{\|\varphi\|}\right) + \frac{\|\psi\|}{\|\varphi\| + \|\psi\|} f\left(\frac{\psi}{\|\psi\|}\right) \right] \\ &= \|\varphi\| f\left(\frac{\varphi}{\|\varphi\|}\right) + \|\psi\| f\left(\frac{\psi}{\|\psi\|}\right) = \tilde{f}(\varphi) + \tilde{f}(\psi). \end{aligned}$$

To prove (b), first note that

$$\tilde{f}(-\varphi) = \|\varphi\| f\left(\frac{-\varphi}{\|\varphi\|}\right) = -\|\varphi\| f\left(\frac{\varphi}{\|\varphi\|}\right) = -\tilde{f}(\varphi).$$

Now suppose $\alpha \in \mathbb{R}$ and $\alpha > 0$. Then $\tilde{f}(\alpha\varphi) = \|\alpha\varphi\| f(\alpha\varphi/\|\alpha\varphi\|) = \alpha\|\varphi\| f(\varphi/\|\varphi\|) = \alpha\tilde{f}(\varphi)$. The same equality for $\alpha \in \mathbb{R}$ and $\alpha < 0$ follows from this and the fact that $\tilde{f}(-\varphi) = -\tilde{f}(\varphi)$. Also, $\tilde{f}(i\varphi) = \|i\varphi\| f(i\varphi/\|i\varphi\|) = \|\varphi\| f(i\varphi/\|\varphi\|) = i\|\varphi\| f(\varphi/\|\varphi\|) = i\tilde{f}(\varphi)$.

Finally, assume that $\alpha = \beta + i\delta$, $\beta, \delta \in \mathbb{R}$. Then

$$\tilde{f}(\alpha\varphi) = \tilde{f}(\beta\varphi + i\delta\varphi) = \tilde{f}(\beta\varphi) + \tilde{f}(i\delta\varphi) = \beta\tilde{f}(\varphi) + i\delta\tilde{f}(\varphi). \blacksquare$$

Let $A \subseteq C(Q)$ where Q is a compact Hausdorff space, be a function space equipped with the usual sup-norm $\|a\|_\infty$. We work in the complex scalar case. For $[a] \in A/\mathbb{C}$, define

$$\|[a]\|_\infty = \inf\{\|a + \lambda\|_\infty : \lambda \in \mathbb{C}\},$$

the usual quotient norm.

NOTE 4. *The d -norm on A/\mathbb{C} is equivalent to the quotient norm.*

Proof. For $a \in A$, fix $x, y \in X$ such that $d(a) = |a(x) - a(y)|$. Note that it is clear that $d(a) \leq 2\|a\|_\infty$, so for all $\lambda \in \mathbb{C}$, $d(a) = d(a + \lambda) \leq 2\|a + \lambda\|_\infty$. It follows that $d[a] = d(a) \leq 2\|[a]\|_\infty$. Also, $\|[a]\|_\infty \leq \|a - a(y)\|_\infty = \sup_{z \in X} |a(z) - a(y)| = |a(x) - a(y)| = d(a) = d[a]$. \blacksquare

Note that by the Hahn–Banach Theorem, $(A/\mathbb{C})^*$ is isometrically isomorphic to $\{\varphi \in A^* : \varphi(1) = 0\}$. For $a \in A$, $\varphi \in (A/\mathbb{C})^*$, let $\widehat{[a]}(\varphi) = \varphi([a])$. Define $\Gamma = \{\alpha \in \mathbb{C} : |\alpha| = 1\}$, and $T = \{\alpha(q - r) : \alpha \in \Gamma, q, r \in Q\} \subseteq (A/\mathbb{C})^*$ (here we identify Q as a subset of A^* via the evaluation map $q \mapsto e_q$ where

$e_q \in A^*$ is defined by $e_q(a) = a(q)$, $a \in A$, $q \in Q$). Then

$$\begin{aligned} T^0 &= \{[a] : \operatorname{Re}(\alpha(a(q) - a(r))) \leq 1 \text{ for all } \alpha \in \Gamma, \text{ and all } q, r \in Q\} \\ &= \{[a] : d[a] = d(a) \leq 1\} \\ &= \{\text{the closed unit ball in } A/\mathbb{C} \text{ with respect to the d-norm}\}. \end{aligned}$$

$$\begin{aligned} T^{00} &= \{\varphi \in (A/\mathbb{C})^* : \operatorname{Re}(\widehat{[a]}(\varphi)) \leq 1 \text{ for all } [a] \in T^0\} \\ &= \{\text{the closed unit ball in } (A/\mathbb{C})^* \text{ with respect to the dual d-norm}\}, \\ &= \overline{\operatorname{co}}(T) \quad (\text{by the Bipolar Theorem [LT, Thm. 7.3, p. 162]}) \\ &= \overline{\operatorname{co}}(\Gamma(S_A - S_A)) = \overline{\operatorname{aco}}(S_A - S_A) \end{aligned}$$

where $S_A = \{\varphi \in A^* : \|\varphi\| = \varphi(1) = 1\}$ is the *state space* of A , and $\overline{\operatorname{aco}}(S)$ is the *absolute convex hull* of a set $S (\subseteq A^*)$ where the closure is taken with respect to the w^* -topology in A^* .

We use this notation in the theorem below.

THEOREM 5. *The map $[a] \mapsto \widehat{[a]}$ is a linear bijection of A/\mathbb{C} onto $A_0(T^{00})$. Also, it is an isometry of $(A/\mathbb{C}, d\text{-norm})$ onto $(A_0(T^{00}), \|\cdot\|_\infty)$ where for $b \in A_0(T^{00})$,*

$$\|b\|_\infty = \sup\{|b(\varphi)| : \varphi \in T^{00}\}.$$

Proof. Since $T^{00} = \{\text{the closed unit ball in } (A/\mathbb{C})^* \text{ with respect to the dual d-norm}\}$, for $a \in A$, $d[a] = \sup\{|\varphi(a)| : \varphi \in T^{00}\} = \|\widehat{[a]}\|_\infty$. Thus, the map $[a] \mapsto \widehat{[a]}$ is an isometry. This map is clearly linear and 1-1. Now assume that $b \in A_0(T^{00})$. Proposition 3 applies with $X = (A/\mathbb{C}, d\text{-norm})$. Therefore b has a w^* -continuous extension \tilde{b} in $(AC)^{**}$. It follows that there exists $[a] \in A/\mathbb{C}$ such that $\widehat{[a]} = b$. ■

Let J and K be convex circled subsets of a complex linear space. We say that an affine map $\tau : J \rightarrow K$ is a *complex affine map* if $\tau(ix) = i\tau(x)$ for all $x \in J$. [Note that the map $z \mapsto \bar{z}$ on the closed unit disk in the complex plane is affine, but not complex affine.]

Now let $A_k \subseteq C(Q_k)$, $k = 1, 2$, be function spaces on compact Hausdorff spaces Q_1 and Q_2 . Let D be a linear isometry of $(A_1/\mathbb{C}, d\text{-norm})$ onto $(A_2/\mathbb{C}, d\text{-norm})$. Letting $T_k = \{\alpha(q - r) : \alpha \in \Gamma, q, r \in Q_k\}$, $k = 1, 2$, we deduce by the discussion prior to Theorem 5 that T_k^{00} is the closed unit ball in $(A_k\mathbb{C})^*$. Thus, $D^*(T_2^{00}) = T_1^{00}$, where D^* is the adjoint of D . For $\varphi \in T_2^{00}$, define $\tau(\varphi) = D^*(\varphi)$. Then τ is a complex affine homeomorphism (w^* -topology) of T_2^{00} onto T_1^{00} . Also, for $\varphi \in T_2^{00}$, $[a] \in A_1/\mathbb{C}$, $\widehat{D[a]}(\varphi) = \widehat{[a]}(D^*(\varphi)) = \widehat{[a]}(\tau(\varphi))$, so $\widehat{D[a]} = \widehat{[a]} \circ \tau$, $[a] \in A_1/\mathbb{C}$.

We summarize this discussion in the following theorem.

THEOREM 6. *There exists a d -preserving linear bijection of A_1 onto A_2 if, and only if, there exists a complex affine homeomorphism of the set $\overline{\text{aco}}(S_{A_2} - S_{A_2})$ onto $\overline{\text{aco}}(S_{A_1} - S_{A_1})$.*

Proof. First note that T_k is compact as it is the continuous image of the compact set $\Gamma \times Q_k \times Q_k$ under the map $(\alpha, q, r) \mapsto \alpha(q - r)$. Thus, $T_k^{00} = \overline{\text{aco}}(S_{A_k} - S_{A_k})$ is compact. Suppose that $\tau : \overline{\text{aco}}(S_{A_2} - S_{A_2}) \rightarrow \overline{\text{aco}}(S_{A_1} - S_{A_1})$ is a complex affine homeomorphism. Then for $a \in A_1$, $D[\widehat{a}] = [\widehat{a}] \circ \tau$ is a linear bijection of $A_0(T_1^{00})$ onto $A_0(T_2^{00})$ which is an isometry with respect to the sup-norm. By Theorem 5, this implies the existence of a linear bijection \widetilde{D} which is an isometry of $(A_1/\mathbb{C}, d\text{-norm})$ onto $(A_2/\mathbb{C}, d\text{-norm})$. Then by Proposition 1, \widetilde{D} lifts to a linear bijection of A_1 onto A_2 which is d -preserving.

Conversely, assume that \overline{D} is a linear bijection of A_1 onto A_2 which is d -preserving. Define $D : A_1/\mathbb{C} \rightarrow A_2/\mathbb{C}$, as usual, by $D[a] = [\overline{D}(a)]$. Then D is a linear bijection which is an isometry with respect to the d -norm. Then as argued in the discussion before the theorem, D^* is a complex affine homeomorphism of $T_2^{00} = \overline{\text{aco}}(S_{A_2} - S_{A_2})$ onto $T_1^{00} = \overline{\text{aco}}(S_{A_1} - S_{A_1})$. ■

When A_k is the space $A(K_k)$, i.e., the space of continuous affine functions on a compact convex set K_k with the sup-norm, then $S_{A_k} = K_k$, $k = 1, 2$. In the real scalar case, we see that $T_k^{00} = \overline{\text{co}}(K_k - K_k) = K_k - K_k$.

Also note that if $\tau : K_2 - K_2 \rightarrow K_1 - K_1$ is an affine homeomorphism, then τ carries a point of symmetry to a point of symmetry, and 0 is the only point of symmetry for both the above sets. Therefore, we must have $\tau(0) = 0$.

Using the same notation as in Theorem 6, we have the following corollary:

COROLLARY 7. *In the case where the scalar field is \mathbb{R} , there exists a d -preserving linear bijection of $A(K_1)$ onto $A(K_2)$ if, and only if, there exists an affine homeomorphism of $K_2 - K_2$ onto $K_1 - K_1$.*

The proof of Theorem 6 applies verbatim to Corollary 7, *except* that $T_k^{00} = K_k - K_k$, $k = 1, 2$, as we noted above.

Corollary 7 raises the natural question: When K_1 and K_2 are compact convex sets, under what conditions are $K_1 - K_1$ and $K_2 - K_2$ affinely homeomorphic? This question seems too difficult to answer in general, although in some cases conditions can be found. For example, the results in [RR] show that (in the real scalar case), when K_1 and K_2 both have the property that all their extreme points are split faces, then $K_1 - K_1$ and $K_2 - K_2$ are affinely homeomorphic if, and only if, K_1 and K_2 are affinely homeomorphic.

Here is an especially simple situation. Suppose that K is a compact convex set which is symmetric. Then clearly $K + K = K - K$. Also, $K + K = 2K$, since for all $x, y \in K$, $x + y = 2(\frac{x+y}{2})$. Thus, $K - K = 2K$. It follows that,

when both K_1 and K_2 are symmetric, then again $K_1 - K_1$ and $K_2 - K_2$ are affinely homeomorphic if, and only if, K_1 and K_2 are affinely homeomorphic.

We derive more information concerning this question in the last section.

4. A characterization of some linear d-preserving maps. Assume that L and S are compact convex sets, that $0 \in L$, $0 \in S$, and $\overline{\text{aco}}(L - L) = \overline{\text{aco}}(S - S)$ [in the case of real scalars, the assumption is $L - L = S - S$]. For $a \in A(L)$, the function $\widehat{a} = [\widehat{a}]$ is in $A_0(\overline{\text{aco}}(L - L))$. Now for $a \in A(L)$, define $a_S \in A(S)$ by $a_S = \widehat{a}|_S + a(0)$ on S . Since the hypotheses are the same for L and S , for $a \in A(S)$, we define a_L in the same way. Note that $a_S \in A(S)$ and $a_L \in A(L)$. Also, $a_S(0) = a_L(0) = a(0)$. We use this notation in the next result.

PROPOSITION 8. *Assume that L and S are compact convex sets with the properties above.*

- (1) For $a \in A(L)$, $\widehat{a}_S = \widehat{a}$; for $a \in A(S)$, $\widehat{a}_L = \widehat{a}$.
- (2) For $a \in A(L)$, $(a_S)_L = a$; for $a \in A(S)$, $(a_L)_S = a$.
- (3) For $a \in A(L)$, $d(a) = d(a_S)$; for $a \in A(S)$, $d(a) = d(a_L)$.

Also, if $\lambda \in A(L)$ is a constant function, then $\lambda_S = \lambda$ (and the same statement with L and S interchanged).

Proof. We do the proof in the complex scalar case.

First we prove (1) when $a \in A(L)$. It is enough to verify that $\widehat{a}_S(\varphi) = \widehat{a}(\varphi)$ for all $\varphi \in \overline{\text{aco}}(S - S)$ of the form $\varphi = t(s_1 - s_2)$, $|t| = 1$, $s_1, s_2 \in S$, since these generate $\overline{\text{aco}}(S - S)$. Assume that $\varphi = t(s_1 - s_2)$ as above, and $a \in A(L)$. Now $\frac{1}{2}(\varphi + ts_2) = \frac{1}{2}ts_1$, so $\frac{1}{2}\widehat{a}(\varphi) + \frac{1}{2}t\widehat{a}(s_2) = \frac{1}{2}t\widehat{a}(s_1)$. Therefore, $\widehat{a}_S(\varphi) = t(a_S(s_1) - a_S(s_2)) = t(\widehat{a}(s_1) - \widehat{a}(s_2)) = \widehat{a}(\varphi)$.

Now we prove (2) for $a \in A(L)$. By definition $(a_S)_L = \widehat{a}_S|_L + a_S(0)$, so by (1), $(a_S)_L = \widehat{a}|_L + a(0)$. For $l \in L$, $l = l - 0$, so $\widehat{a}(l) = a(l) - a(0)$. Thus, $a(l) = \widehat{a}(l) + a(0)$. Then $(a_S)_L(l) = \widehat{a}(l) + a(0) = a(l)$. This establishes (2).

Assume that $a \in A(L)$. By (1), $\widehat{a} = \widehat{a}_S$. Then $d(a) = \|\widehat{a}\|_\infty$ and $d(a_S) = \|\widehat{a}_S\|_\infty$ (Theorem 5), so (3) follows from these equalities. ■

We can now describe the general form of the d-preserving linear bijection raised in the question in [RR] that we mentioned at the beginning of Section 3.

THEOREM 9. *Let $D : A(K) \rightarrow A(S)$ be a d-preserving linear bijection where K, S are compact convex sets with the former having the property that all the points of $\text{ext}(K)$ are split. [In particular, K could be a Choquet simplex.] We assume, as we may by translation in $A(S)^*$, that $0 \in S$. Then there exist a compact convex set $L \subseteq \overline{\text{aco}}(S - S)$ [$S - S$ in the real scalar case], affinely homeomorphic to K , such that $\overline{\text{aco}}(S - S) = \overline{\text{aco}}(L - L)$, $0 \in L$,*

and an affine homeomorphism $\tau : L \rightarrow K$, and $\alpha \in A(L)'$ such that for all $a \in A(K)$,

$$D(a) = c(a \circ \tau)_S + \alpha(a), \quad \text{where } |c| = 1, \alpha(1) \neq -c.$$

Proof. We do the proof in the complex scalar case. First we construct L . Define $\tilde{D} : A(K)/\mathbb{C} \rightarrow A(S)/\mathbb{C}$ in the usual way: $\tilde{D}[a] = [Da]$. Then as seen in the proof of Theorem 6, \tilde{D}^* maps $\overline{\text{aco}}(S - S)$ onto $\overline{\text{aco}}(K - K)$. Fix $x_0 \in \text{ext}(K)$. Define $L = (\tilde{D}^*)^{-1}(K - \{x_0\})$. Then for $x \in L$, $x \mapsto \tilde{D}^*(x) + x_0 \in K$, and this map is an affine homeomorphism of L onto K . Note $0 \in L$. Observe that $\text{ext}(L) - \text{ext}(L) = (\tilde{D}^*)^{-1}(\text{ext}(K) - \text{ext}(K))$, and consequently, $\overline{\text{aco}}(L - L) = (\tilde{D}^*)^{-1}(\overline{\text{aco}}(K - K)) = \overline{\text{aco}}(S - S)$. Thus, $S \subseteq S - S \subseteq \overline{\text{aco}}(S - S)$, $0 \in S$, and $L \subseteq L - L \subseteq \overline{\text{aco}}(L - L)$, $0 \in L$. Therefore Proposition 8 applies.

Now for $a \in A(K)$, $D(a) \in A(S)$ and $(D(a))_L \in A(L)$. By Proposition 8, $a \mapsto (D(a))_L$ is a d-preserving linear bijection of $A(K)$ onto $A(L)$. Applying [RR, Theorem 1], we have $(D(a))_L = c(a \circ \tau) + \alpha(a)$, where $|c| = 1$, $\alpha \in A(L)'$, $\tau : L \rightarrow K$ is an affine homeomorphism, and $\alpha(1) \neq -c$. Using Proposition 8 again, we have

$$D(a) = ((D(a))_L)_S = (c(a \circ \tau) + \alpha(a))_S = c(a \circ \tau)_S + \alpha(a). \quad \blacksquare$$

5. A geometric problem involving $K - K$, K a simplex. Let K be a simplex. We assume that K is embedded in $A(K)^*$ as the base of a cone \tilde{K} which generates $A(K)^*$; see [P, p. 59]. It is important to keep in mind that *distances between points will be computed in the dual norm on $A(K)^*$.*

First, let K be the simplex, $K = \text{co}\{x_1, x_2, x_3\}$. A simple observation using the linear independence of the vectors $x_1 - x_2$ and $x_2 - x_3$ is that

$$K - K = \{\alpha(x_1 - x_2) + \beta(x_2 - x_3) : \alpha, \beta \in \mathbb{R}, |\alpha| + |\alpha - \beta| + |\beta| \leq 2\}.$$

Also, representation of points in $K - K$ is unique. These remarks will be useful in what follows.

The problem is to find polytopes S with the property that $K - K = S - S$.

We may assume by translating that $S \subseteq \tilde{K}$. Also, we assume that

$$S = \text{co}\{s_1, \tilde{s}_1, s_2, \tilde{s}_2, s_3, \tilde{s}_3\} \subseteq A(K)^*$$

is such that

$$s_1 - \tilde{s}_1 = x_1 - x_2 \quad (= y_1, \text{ say}), \quad s_2 - \tilde{s}_2 = x_2 - x_3 = y_2,$$

$$s_3 - \tilde{s}_3 = x_1 - x_3 = y_3,$$

where $\text{ext}(S) = \{s_1, \tilde{s}_1, s_2, \tilde{s}_2, s_3, \tilde{s}_3\}$. Therefore, $S - S = K - K$. Note that $S - S$ will have only the extreme points $\pm y_1, \pm y_2, \pm y_3$.

By the decomposition property of vector lattices [P, p. 61], one sees easily that:

$$\begin{aligned}
 (\$) \quad & s_1 = x_1 + ax, & s_2 = x_2 + by, & s_3 = x_1 + cz, \\
 & \tilde{s}_1 = x_2 + ax, & \tilde{s}_2 = x_3 + by, & \tilde{s}_3 = x_3 + cz,
 \end{aligned}$$

where $a, b, c \geq 0, x, y, z \in K$.

We first consider the case where the vectors $s_1 - \tilde{s}_1, s_2 - \tilde{s}_2$, and $s_3 - \tilde{s}_3$ intersect in *distinct* points $P, P',$ and P'' ; see Fig. 1.

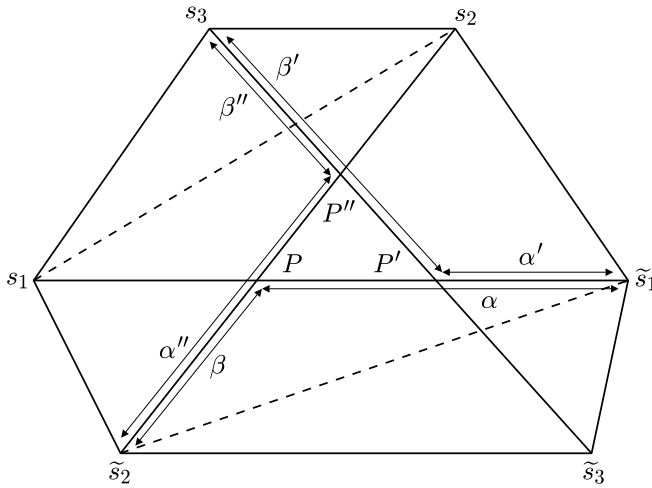


Fig. 1

Consider the point P as a typical case; refer to the quadrilateral $\{s_1, s_2, \tilde{s}_1, \tilde{s}_2\}$ in Fig. 1. Here

$$P = \alpha s_1 + (1 - \alpha)\tilde{s}_1 = \beta s_2 + (1 - \beta)\tilde{s}_2$$

(it is assumed that $1/2 < \alpha < 1$).

Recall that $\|s_1 - \tilde{s}_1\| = \|x_1 - x_2\| = 1 + 1 = 2$ by the splittability of extreme points for a simplex. Thus, we have $\alpha = \|P - \tilde{s}_1\|/\|s_1 - \tilde{s}_1\| = \frac{1}{2}\|P - \tilde{s}_1\|$. Now $P = \alpha(s_1 - \tilde{s}_1) + \tilde{s}_1 = \beta(s_2 - \tilde{s}_2) + \tilde{s}_2$. Also, $\tilde{s}_1 - \tilde{s}_2 = \beta(x_2 - x_3) - \alpha(x_1 - x_2) = (\alpha + \beta)x_2 - \alpha x_1 - \beta x_3$. Therefore, $\|\tilde{s}_1 - \tilde{s}_2\| = \alpha + \beta + \alpha + \beta = 2(\alpha + \beta) \leq 2$, so $\alpha + \beta \leq 1$. Also,

$$\begin{aligned}
 s_1 - s_2 &= (P - s_2) + (s_1 - P) \\
 &= [\beta s_2 + (1 - \beta)\tilde{s}_2 - s_2] + [s_1 - \alpha s_1 - (1 - \alpha)\tilde{s}_1] \\
 &= (1 - \beta)(\tilde{s}_2 - s_2) + (1 - \alpha)(s_1 - \tilde{s}_1) \\
 &= (1 - \beta)(x_3 - x_2) + (1 - \alpha)(x_1 - x_2) \\
 &= (1 - \alpha)x_1 - (2 - \alpha - \beta)x_2 + (1 - \beta)x_3.
 \end{aligned}$$

This implies that $\|s_1 - s_2\| = 4 - 2(\alpha + \beta) \leq 2$, so $\alpha + \beta \geq 1$. Thus, $\alpha + \beta = 1$. By referring to Fig. 1, we have similarly:

$$\begin{aligned} P' &= \alpha' s_1 + (1 - \alpha') \tilde{s}_1 \\ &= \beta' \tilde{s}_3 + (1 - \beta') s_3, \quad \alpha' + \beta' = 1, \text{ where } \alpha' = \frac{1}{2} \|P' - \tilde{s}_1\|, \\ P'' &= \alpha'' s_2 + (1 - \alpha'') \tilde{s}_2 \\ &= \beta'' \tilde{s}_3 + (1 - \beta'') s_3, \quad \alpha'' + \beta'' = 1, \text{ where } \alpha'' = \frac{1}{2} \|P'' - \tilde{s}_2\|. \end{aligned}$$

To find the relations among $\alpha, \beta, \alpha', \beta', \alpha'', \beta''$, argue as follows:

$$\begin{aligned} P - s_3 &= (P - P'') + (P'' - s_3) = (1 - \beta - \beta'')(\tilde{s}_2 - s_2) + \beta''(\tilde{s}_3 - s_3) \\ &= (\alpha - \beta'')(\tilde{s}_2 - s_2) + \beta''(\tilde{s}_3 - s_3) \\ &= (\beta'' - \alpha)(x_2 - x_3) - \beta''(x_1 - x_3) \\ &= -\beta''(x_1 - x_2) - \alpha(x_2 - x_3). \end{aligned}$$

Also,

$$\begin{aligned} P - s_3 &= (P - P') + (P' - s_3) = (1 - \alpha' - \beta)(s_1 - \tilde{s}_1) + \beta'(\tilde{s}_3 - s_3) \\ &= (\beta' - \beta)(x_1 - x_2) - \beta'(x_1 - x_3) = -\beta(x_1 - x_2) - \beta'(x_2 - x_3). \end{aligned}$$

It follows by the uniqueness of representation of points in $K - K$ that $\beta = \beta''$ and $\alpha = \beta'$. Then $1 = \alpha + \beta = \beta' + \beta''$, so $\alpha'' = \beta'$ and $\alpha' = \beta''$. Therefore, $\alpha = \beta' = \alpha''$ and $\beta = \alpha' = \beta''$.

Note that now the various distances can be computed. For example,

$$s_2 - \tilde{s}_1 = (s_2 - P) + (P - \tilde{s}_1) = \alpha(s_2 - \tilde{s}_2) + \alpha(s_1 - \tilde{s}_1) = \alpha(s_3 - \tilde{s}_3).$$

Therefore, $\|s_2 - \tilde{s}_1\| = 2\alpha$. This also shows that $s_2 - \tilde{s}_1$ is *not extreme* in $S - S$. The same kind of argument shows that $\{s_i - \tilde{s}_i\}_{i=1,2,3}$ are the only extreme points in $S - S$.

Returning to equation (§), we can explicitly write down the form of $s_1, \tilde{s}_1, s_2, \tilde{s}_2, s_3, \tilde{s}_3$ as follows (from Fig. 1):

$$\begin{aligned} P &= \alpha s_1 + \beta \tilde{s}_1 = \alpha(x_1 + ax) + \beta(x_2 + ax) = \alpha x_1 + \beta x_2 + ax \\ &= \alpha \tilde{s}_2 + \beta s_2 = \alpha(x_3 + by) + \beta(x_2 + by) = \alpha x_3 + \beta x_2 + by. \end{aligned}$$

This implies $\alpha(x_1 - x_3) = by - ax$, so $b = a$ and $by = \alpha(x_1 - x_3) + ax$. Similarly, $c = a$ and $cz = \alpha(x_2 - x_3) + ax$. Thus,

$$\begin{aligned} s_1 &= x_1 + ax, & \tilde{s}_1 &= x_2 + ax, \\ (\%) \quad s_2 &= x_2 + \alpha(x_1 - x_3) + ax, & \tilde{s}_2 &= x_3 + \alpha(x_1 - x_3) + ax, \\ s_3 &= x_1 + \alpha(x_2 - x_3) + ax, & \tilde{s}_3 &= x_3 + \alpha(x_2 - x_3) + ax. \end{aligned}$$

REMARKS. (1) The above is the general solution when $\alpha > 1/2$. When $\alpha \rightarrow 1/2$, we check that $P_0 = \frac{1}{2}(s_1 + \tilde{s}_1) = \frac{1}{2}(s_2 + \tilde{s}_2) = \frac{1}{2}(s_3 + \tilde{s}_3) = \frac{1}{2}(x_1 + x_2) + ax$, and we get the symmetric solution $S = \frac{1}{2}(K - K) + P_0$.

(2) When $\alpha \rightarrow 1$, we get

$$\begin{aligned} s_1 &= x_1 + ax, & s_2 &= x_1 + x_2 - x_3 + ax, & s_3 &= x_1 + x_2 - x_3 + ax, \\ \tilde{s}_1 &= x_2 + ax, & \tilde{s}_2 &= x_1 + ax, & \tilde{s}_3 &= x_2 + ax. \end{aligned}$$

Hence, $s_1 = \tilde{s}_2$, $\tilde{s}_1 = \tilde{s}_3$ and $s_2 = s_3$, and we get the simplex $\text{co}\{x_1, x_1 + x_2 - x_3, x_2\}$. Translating by x_3 , we have the simplex $\text{co}\{x_1 + x_3, x_1 + x_2, x_2 + x_3\}$, and the latter is obtained by translating $\text{co}\{-x_1, -x_2, -x_3\}$ by $x_1 + x_2 + x_3$.

(3) “Uniqueness” of solutions. To show that there cannot be points in S other than those specified, take s as a typical point outside $\text{co}\{s_1, \tilde{s}_1, s_2, \tilde{s}_2, s_3, \tilde{s}_3\}$ of the form

$$s = \tilde{s}_1 + t_0(\tilde{s}_1 - s_1) + t(s_2 - \tilde{s}_1), \quad t_0 > 0, \quad 0 < t < 1.$$

Then from (%) we obtain $s = x_2 + t_0(x_2 - x_1) + t\alpha(x_1 - x_3) + ax$ and $\tilde{s}_2 = x_3 + \alpha(x_1 - x_3) + ax$. Hence,

$$\begin{aligned} s - \tilde{s}_2 &= x_2 - x_3 + \alpha(t - 1)(x_1 - x_3) - t_0(x_1 - x_2) \\ &= -[\alpha(1 - t) + t_0](x_1 - x_2) + [1 - \alpha(1 - t)](x_2 - x_3) \\ &= a(x_1 - x_2) + b(x_2 - x_3). \end{aligned}$$

Then $|a| + |b| + |a - b| = \alpha(1 - t) + t_0 + 1 - \alpha(1 - t) + |1 + t_0| = 2(1 + t_0) > 2$. It follows that $s - \tilde{s}_2 \notin K - K$.

The problem in \mathbb{R}^3 . Consider the simplex $K = \text{co}\{x_1, x_2, x_3, x_4\}$ in \mathbb{R}^3 . Let $S = \{s_1, \tilde{s}_1, s_2, \tilde{s}_2, \dots, s_6, \tilde{s}_6\}$. Again, we want to find conditions under which $S - S = K - K$. For this purpose we use the equations in (%). By (%), for $\text{co}\{x_1, x_2, x_3\}$, with $1/2 < \alpha_1 < 1$, we have:

$$\begin{aligned} s_1 &= x_1 + ax, & \tilde{s}_1 &= x_2 + ax, \\ \text{(A)} \quad s_2 &= x_2 + \alpha_1(x_1 - x_3) + ax, & \tilde{s}_2 &= x_3 + \alpha_1(x_1 - x_3) + ax, \\ s_3 &= x_1 + \alpha_1(x_2 - x_3) + ax, & \tilde{s}_3 &= x_3 + \alpha_1(x_2 - x_3) + ax. \end{aligned}$$

Similarly, for $\text{co}\{x_1, x_2, x_4\}$, with $1/2 < \alpha_1 < 1, a_1\bar{x}$, we have:

$$\begin{aligned} s_1 &= x_1 + a_1\bar{x}, & \tilde{s}_1 &= x_2 + a_1\bar{x}, \\ \text{(B)} \quad s_5 &= x_2 + \alpha_2(x_1 - x_4) + a_1\bar{x}, & \tilde{s}_5 &= x_4 + \alpha_2(x_1 - x_4) + a_1\bar{x}, \\ s_6 &= x_1 + \alpha_2(x_2 - x_4) + a_1\bar{x}, & \tilde{s}_6 &= x_4 + \alpha_2(x_2 - x_4) + a_1\bar{x}. \end{aligned}$$

Note that from the first two equations in the systems (A) and (B), it is clear that $ax = a_1\bar{x}$. For $\text{co}\{x_1, x_3, x_4\}$, with $\frac{1}{2} < \alpha_3 < 1, a_2\bar{\bar{x}}$, we have:

$$\begin{aligned} s_3 &= x_1 + \alpha_3(x_4 - x_1) + a_2\bar{\bar{x}}, & \tilde{s}_3 &= x_3 + \alpha_3(x_4 - x_1) + a_2\bar{\bar{x}}, \\ \text{(C)} \quad s_4 &= x_4 + a_2\bar{\bar{x}}, & \tilde{s}_4 &= x_3 + a_2\bar{\bar{x}}, \\ s_6 &= x_1 + \alpha_3(x_3 - x_1) + a_2\bar{\bar{x}}, & \tilde{s}_6 &= x_4 + \alpha_3(x_3 - x_1) + a_2\bar{\bar{x}}. \end{aligned}$$

Equating s_6 from (B) and (C), we have

$$s_6 = x_1 + \alpha_2(x_2 - x_4) + a_1\bar{x} = s_6 = x_1 + \alpha_3(x_3 - x_1) + a_2\bar{x},$$

and this implies $a = a_2$ and $a_2\bar{x} = \alpha_2(x_2 - x_4) + \alpha_3(x_1 - x_3) + ax$. Then using the fifth equation in (C), we obtain

$$s_3 = x_1 + \alpha_3(x_4 - x_1) + \alpha_2(x_2 - x_4) + \alpha_3(x_1 - x_3) + ax = x_1 + \alpha_1(x_2 - x_3) + ax$$

(by the fifth equation in (A)). It follows that

$$\alpha_3(x_4 - x_1) + \alpha_2(x_2 - x_4) + \alpha_3(x_1 - x_3) = \alpha_1(x_2 - x_3),$$

or

$$\alpha_3(x_4 - x_1) - \alpha_2(x_4 - x_2) + \alpha_3(x_1 - x_4 + x_4 - x_3) = \alpha_1(x_2 - x_4 + x_4 - x_3),$$

or

$$(\alpha_1 - \alpha_2)(x_2 - x_4) + (\alpha_1 - \alpha_3)(x_4 - x_3) = 0.$$

By the linear independence of the vectors $\{x_1 - x_4, x_2 - x_4, x_3 - x_4\}$, we must have $\alpha_1 - \alpha_2 = 0 = \alpha_1 - \alpha_3$, which implies $\alpha_1 = \alpha_2 = \alpha_3$ ($= \alpha$, say), and the solution can now be written

$s_1 = x_1 + ax,$	$\tilde{s}_1 = x_2 + ax,$
$s_2 = x_2 + \alpha(x_1 - x_3) + ax,$	$\tilde{s}_2 = x_3 + \alpha(x_1 - x_3) + ax,$
$s_3 = x_1 + \alpha(x_2 - x_3) + ax,$	$\tilde{s}_3 = x_3 + \alpha(x_2 - x_3) + ax,$
$s_4 = x_4 + a(x_1 + x_2 - x_3 - x_4) + ax,$	$\tilde{s}_4 = x_3 + a(x_1 + x_2 - x_3 - x_4) + ax,$
$s_5 = x_2 + \alpha(x_1 - x_4) + ax,$	$\tilde{s}_5 = x_4 + \alpha(x_1 - x_4) + ax,$
$s_6 = x_1 + \alpha(x_2 - x_4) + ax,$	$\tilde{s}_6 = x_4 + \alpha(x_2 - x_4) + ax.$

The form of this solution simplifies if one translates by $\alpha(x_3 + x_4)$:

$s_1 = x_1 + \alpha(x_3 + x_4) + ax,$	$\tilde{s}_1 = x_2 + \alpha(x_3 + x_4) + ax,$
$s_2 = x_2 + \alpha(x_1 + x_4) + ax,$	$\tilde{s}_2 = x_3 + \alpha(x_1 + x_4) + ax,$
$s_3 = x_1 + \alpha(x_2 + x_4) + ax,$	$\tilde{s}_3 = x_3 + \alpha(x_2 + x_4) + ax,$
$s_4 = x_4 + a(x_1 + x_2) + ax,$	$\tilde{s}_4 = x_3 + a(x_1 + x_2) + ax,$
$s_5 = x_2 + \alpha(x_1 + x_3) + ax,$	$\tilde{s}_5 = x_4 + \alpha(x_1 + x_3) + ax,$
$s_6 = x_1 + \alpha(x_2 + x_3) + ax,$	$\tilde{s}_6 = x_4 + \alpha(x_2 + x_3) + ax.$

As before, one checks easily that $s_i - s_j$, $s_i - \tilde{s}_j$, and $\tilde{s}_i - \tilde{s}_j$ are not extreme for $i \neq j$.

REMARKS. (1) The points $\{s_i - \tilde{s}_i\}_{i=1}^6$ are the extreme points of $S - S$. It is a curious fact that the points in this set are equidistant from $x_1 + x_2 + x_3 + x_4$ [in the $A(K)^*$ metric].

(2) It is more or less clear that the set S is non-symmetric, but here is a formal proof of this fact: If P_0 were the centre of symmetry of S , then for each extreme point (say, s_1), there exists $t_1 \in S$ such that $(s_1 + t_1)/2 = P_0$.

This implies $s_1 - P_0 = P_0 - t_1$. Similarly, $\tilde{s}_1 - P_0 = P_0 - t_2$ for some $t_2 \in S$. Then we must have $s_1 - \tilde{s}_1 = t_2 - t_1$, and so $s_1 = t_2$ and $\tilde{s}_1 = t_1$ (by the uniqueness of the expression of an extreme point in $S - S$). Thus, $P_0 = (s_1 + \tilde{s}_1)/2 = \frac{1}{2}(x_1 + x_2) + ax$. But $(s_2 + \tilde{s}_2)/2$ is something different, which is a contradiction (unless $\alpha = 1/2$).

(3) $\alpha \rightarrow 1/2$ gives the symmetric solution as before with centre of symmetry $\frac{1}{2}(x_1 + x_2)$.

(4) $\alpha \rightarrow 1$ gives a translate of K as a solution.

Uniqueness. Since the solution for the simplex $K = co\{x_1, x_2, x_3, x_4\}$ in \mathbb{R}^3 was obtained by solving for each face, the solution should therefore be unique, modulo translation and the ordering of the points $s_1, \tilde{s}_1, s_2, \tilde{s}_2, \dots, s_6, \tilde{s}_6$.

It is now apparent what the solution for $K = co\{x_1, x_2, x_3, x_4, x_5\}$ in \mathbb{R}^4 will be like:

$$\begin{aligned} \{x_1, x_2\} &\rightarrow s_1 = x_1 + \alpha(x_3 + x_4 + x_5) + ax, \\ &\quad \tilde{s}_1 = x_2 + \alpha(x_3 + x_4 + x_5) + ax, \\ \{x_1, x_3\} &\rightarrow s_2 = x_1 + \alpha(x_1 + x_4 + x_5) + ax, \\ &\quad \tilde{s}_2 = x_3 + \alpha(x_1 + x_4 + x_5) + ax, \\ \{x_1, x_4\} &\rightarrow \\ \{x_1, x_5\} &\rightarrow \\ &\quad \vdots \\ \{x_4, x_5\} &\rightarrow s_{10} = x_4 + \alpha(x_1 + x_2 + x_3) + ax, \\ &\quad \tilde{s}_{10} = x_5 + \alpha(x_1 + x_2 + x_3) + ax. \end{aligned}$$

It is now clear that a solution S for a simplex $K = co\{x_1, x_2, \dots, x_{n+1}\}$ in \mathbb{R}^n can be written down by inspection.

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References

- [A] E. Alfsen, *Compact Convex Sets and Boundary Integrals*, Springer, Berlin, 1971.
- [AE] L. Asimow and A. J. Ellis, *Convexity Theory and Its Applications in Functional Analysis*, London Math. Soc. Monographs 16, Academic Press, 1980.
- [B] S. Banach, *Théorie des Opérations Linéaires*, Monografje Mat. 1, Warszawa, 1932.

- [D] J. Dugundji, *Topology*, Allyn and Bacon, Boston, 1966.
- [DS] N. Dunford and J. T. Schwartz, *Linear Operators, Vol. I*, Interscience, New York, 1958.
- [GU] F. González and V. Uspenskij, *On homomorphisms of groups of integer-valued functions*, Extracta Math. 14 (1999), 19–29.
- [GM] M. Györy and L. Molnár, *Diameter preserving linear bijections of $C(X)$* , Arch. Math. (Basel) 71 (1998), 301–310.
- [LT] D. Lay and A. Taylor, *Functional Analysis*, Wiley, 1990.
- [MU] S. Mazur et S. Ulam, *Sur les transformations isométriques d'espaces vectoriels normés*, C. R. Acad. Sci. Paris 194 (1932), 946–948.
- [P] R. R. Phelps, *Lectures on Choquet's Theorem*, Van Nostrand, Princeton, 1965.
- [RR] T. S. S. R. K. Rao and A. K. Roy, *Diameter-preserving linear bijections of function spaces*, J. Austral. Math. Soc. Ser. A 70 (2001), 323–335.
- [R] W. Rudin, *Functional Analysis*, McGraw-Hill, 1991.
- [S] F. C. Sanchez, *Diameter preserving linear maps and isometries (II)*, Proc. Indian Acad. Sci. Math. Sci. 110 (2000), 205–211.

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