Metrics for more than two points at once

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The conventional definition of a topological metric over a space specifies properties that must be obeyed by any measure of "how separated" two points in that space are. Here it is shown how to extend that definition, and in particular the triangle inequality, to concern arbitrary numbers of points. Such a measure of how separated the points within a collection are can be bootstrapped, to measure "how separated" from each other are two (or more) collections. The measure presented here also allows fractional membership of an element in a collection. This means it directly concerns measures of "how spread out" a probability distribution over a space is. When such a measure is bootstrapped to compare two collections, it allows us to measure how separated two probability distributions are, or more generally, how separated a distribution of distributions is.

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I. INTRODUCTION

The conventional definition of a topological metric formalizes the concept of distance. It specifies properties required of any function that purports to measure "how separated" two elements of a space are. However often one wants to measure "how separated" the members of a collection of more than two elements is. The conventional way to do this is to combine the pair-wise metric values for all pairs of elements in the collection, into an aggregate measure. This is ad hoc however.

As an alternative, here the formal definition of a topological metric is extended to apply to collections of more than two elements. In particular, the triangle inequality is extended to concern such collections. The measure presented here applies even to collections with duplicate elements (i.e., to bags). It also applies to collections with "fractional" numbers of elements, i.e., to probability distributions.

This measure can be directly incorporated into many domains where ad hoc combinations of pair-wise metrics are currently used. In addition, when applied to different projections of a high-dimensional data set, it provides a novel type of vector-valued characterization of that data set.

This new measure can be bootstrapped in a natural way, to measure "how separated" from each other two collections are. In other words, given a measure \( \rho \) of how separated from each other the elements in an arbitrary collection \( \xi \) are, one can define a measure of how separated from each other two collections \( \xi_1 \) and \( \xi_2 \) are. (Intuitively, the idea is to subtract the sum of the measure's values for each of the two separate collections \( \xi_1 \) and \( \xi_2 \) from the value of the measure for the union of the collections.) More generally, one can measure how separated a collection of such collections is. Indeed, with fractional memberships, such bootstrapping allows us to measure how separated a distribution of distributions is.

In the next section the definition of a multi-argument metric (multimetric, for short) is presented. Also in that section is an extensive set of examples and a list of some elementary properties. For instance, it is shown that the standard deviation of a probability distribution across \( \mathbb{R}^n \) is a multimetric, whereas the variance of that distribution is not.

The following section presents a way to bootstrap from a multimetric for elements within a collection to a multimetric over collections. Some examples and elementary properties of this bootstrapped measure are also in that section.

A short concluding section considers some of the possible uses of multimetrics.

II. MULTIMETRICS

Collections of elements from a space \( X \) are represented as vectors of counts, i.e., functions from \( x \in X \rightarrow \{0,1,2,\ldots\} \). So for example, if \( X = \{A,B,C\} \), and we have the collection of three A's, no B's, and one C, we represent that as \( (3,0,1) \). This use of the support of a vector to indicate points is analogous to how wave functions are interpreted in quantum mechanics. Like in quantum mechanics, here it is natural to extend the representation; for current purposes we only need extend it to include functions from \( x \in X \) to \( \mathbb{R} \). In particular, doing this allows us to represent probability distributions (or density functions, depending on the cardinality of \( X \)) over \( X \). Accordingly, our formalization of multimetrics will provide a measure for how how spread out a distribution over distributions is.[11] Given \( X \), the associated space of all functions from \( X \) to \( \mathbb{R} \) is written as \( \mathbb{R}^X \). The subspace of functions that are nowhere-negative is written as \((\mathbb{R}^+)^X\).

As a notational comment, integrals are written with the measure implicitly set by the associated space. In particular, for a finite space, the point-mass measure is implied, and the integral symbol indicates a sum. In addition \( \delta_x \) is used to indicate the appropriate type of delta function (Dirac, Kronecker, etc.) about \( x \). Other shorthand is \( R^X \equiv (\mathbb{R}^+)^X - \{0\} \) and \( ||v|| \) means \( \int dx \, v(x) \).

In this representation of collections of elements from \( X \), any conventional metric taking two arguments in \( X \) can be written as a function \( \rho \) over a subset of the vectors
in \( \mathcal{R}^X \). That subset consists of all vectors that either have two of their components equal to 1 and all others 0, or one component equal to 2 and all others 0. For example, for \( X = \{A, B, C\} \), the metric distance between \( A \) and \( B \) is \( \rho(1,0,1) \), and from \( A \) to itself is \( \rho(2,0,0) \).

Generalizing this, a multimetric for \( T(X) \subseteq \mathcal{R}^X \) is defined as a real-valued function \( \rho \) over \( \mathcal{R}^X \) such that

\[
1) u, v, w \in T(X) \Rightarrow \rho(u + v) \leq \rho(u + w) + \rho(v + w).
\]

\[
2) \rho(u) \geq 0, \rho(k \delta_x) = 0 \forall x, k > 0.
\]

\[
3) \rho(u) = 0 \Rightarrow u = k \delta_x \text{ for some } k, x.
\]

Consider the requirement for a (two-argument) metric that it be symmetric under permutation of the points its considering. The analogous requirement that a multimetric be symmetric under permutations of the elements in the collection is satisfied automatically, simply due to bow that collection is represented. Next consider the case where collections, if only one \( x \in X \) is in a collection (perhaps occurring more than once), then only one component of \( u \) is non-zero. Accordingly, conditions (2) and (3) are extensions of the usual condition defining a metric that it be non-negative and equal 0 if its arguments are the same. Condition (1) is an extension of the triangle inequality, to both allow repeats of elements from \( X \) and/or more than two elements from \( X \) to be in the collection. Note that condition (1) involves sums in its argument rather than (as in a conventional norm-based metric for a Euclidean space) differences. Intuitively, \( T(X) \) is that subset of \( \mathcal{R}^X \) over which the generalized version of the triangle inequality holds.[12]

Condition (1) implies that multimetrics obey a second triangle inequality, just as conventional metrics do:

\[
\rho(u + v) \geq |\rho(u + w) - \rho(v + w)|.
\]

This follows by rewriting condition (1) as \( \rho(u + w) \geq \rho(u + v) - \rho(v + w) \), and then relabeling twice.

**Example 1:** Set \( X = \mathbb{R}^N \). Take \( T(X) \) to be those elements of \( \mathcal{R}^X \) whose norm equals 1, i.e., the probability density functions over \( \mathbb{R}^N \). Then have \( \rho(s) \) for any \( s \in \mathcal{R}^X \) (whether in \( T(X) \) or not) be the standard deviation of the distribution \( \tilde{s} = \{s\} \), i.e.,

\[
\rho(s) = \sqrt{\frac{1}{2} \int dx dx' s(x) s(x') (x - x')^2}.
\]

Conditions (2) and (3) are immediate. To understand condition (1), first, as an example, say that all three of \( u, v \) and \( w \) are separate single delta functions over \( X \). Then condition (1) reduces to the conventional triangle inequality over \( \mathcal{R}^X \), here relating the points (in the supports of \( u, v, w \)). This example also demonstrates that the variance (i.e., the square of our \( \rho \)) is not a multimetric.

For a vector \( s \) that involves multiple delta functions, \( \rho(s) \) measures the square root of the sum of the squares of the Euclidean distances between the points (in the support of \( s \)). In this sense it tells us how "spread out" those points are. Condition (1) even holds for vectors that are not sums of delta functions however (see appendix).

**Example 2:** As a variant of Ex. 1, have \( X = \text{the unit simplex in } \mathbb{R}^N \), and use the same \( \rho \) as in Ex. 1. In this case any element of \( X \) is a probability distribution over a variable with \( N \) possible values. So any element of \( T(X) \) is a probability density function over such probability distributions. In particular, say \( s \) is a sum of some delta functions for such an \( X \). Then \( \rho(s) \) measures how spread out the probability distributions in (the support of) \( s \) are. If those probability distributions are themselves sums of delta functions, they just constitute subsets of our \( N \) values, and \( \rho(s) \) measures how spread out from one another those subsets are.

**Example 3:** As another variant of Ex. 1, for any \( X \), take \( T(X) = \mathcal{R}^X \). Define the tensor contraction \( \langle s \mid t \rangle \equiv \int dx dx' s(x) t(x') F(x, x') \) where \( F \) is symmetric and nowhere-negative, and where \( F(x, x') = 0 \iff x = x' \). Then \( \rho(s) = \sqrt{\langle s \mid s \rangle} \) obeys conditions (2) and (3) by inspection. It also obeys condition (1) (see appendix).

Note that the \( \langle \cdot, \cdot \rangle \) operator is not an inner product over \( \mathcal{R}^X \), the extension of \( T(X) \) to a full vector space. When components of \( s \) can be negative, \( \langle s, s \rangle \) may be as well. Note also that there is a natural differential geometric interpretation of this \( \rho \) when \( X \) consists of \( N \) values. Say we have a curve on an \( N \)-dimensional manifold with metric tensor \( F \) at a particular point on the curve, and that at that point the tangent vector to the curve is \( s \). Then \( \rho(s) \) is the derivative of arc length along that curve, evaluated at that point.

This suggest an extension of this multimetric, in which rather than a tensor contraction between two vectors, we form the tensor contraction of \( n \) vectors: \( 
\int \left[ dx ight] \left[ dx' \right] \ldots \left[ dx^{n-1} \right] \left[ dx^n \right] s(x) \ldots s^{(n)}(x^n) F(x_1, \ldots, x^n), \)
where \( F \) is invariant under permutation of its arguments, nowhere-negative, and equals 0 if and only if all its arguments have the same value. Any \( \rho(s) \) that is a monotonically increasing function of \( \langle s, s, \ldots, s \rangle^{1/2} \) automatically obeys conditions (2) and (3).

It is worth collecting a few elementary results concerning multimetrics:

**Lemma 1:**

1. Let \( \{\rho_1\} \) be a set of functions that obey conditions (2) and (3), and \( \{a_i\} \) a set of non-negative real numbers at least one of which is non-zero. Then \( \sum_i a_i \rho_i \) also obeys conditions (2) and (3).

2. Let \( \{\rho_1\} \) be a set of functions that obey condition (1), and \( \{\rho_2\} \) a set of non-negative real numbers at
least one of which is non-zero. Then $\sum a_i \rho_i$ also obeys condition (1).

3. Let $f : \mathbb{R} \to \mathbb{R}^+ \to \mathbb{R}^+$ be a monotonically increasing concave function that equals 0 when its argument does. Then if $\rho$ is a multimetric for some $T(X)$, $f(\rho)$ is also a multimetric for $T(X)$ (see appendix).

4. Let $f : X \to Y$ be invertible, and let $\rho_Y$ be a multimetric over $Y$. Define the operator $B_f : \mathcal{R}^X \to \mathcal{R}^Y$ by $[B_f(s)](y) \equiv s(f^{-1}(y))$ if $f^{-1}(y)$ exists, 0 otherwise. $B_f$ is a linear operator. This means $\rho_X(s) \equiv \rho_Y(A_f(s))$ is a multimetric [13].

Example 4: Take $X = \mathbb{R}^N$ again, and let $T(X)$ be all of $\mathcal{R}^X$ with bounded support. Then by Lemma 1, the width along $x_1$ of (the support of) $s \in T(X)$ is a multimetric function of $s$ (see appendix).

This means that the average of the width in $x_1$ over all possible rotations of $X$ is also a multimetric. Similarly, consider the smallest axis-parallel box enclosing the (support of the) Euclidean points in $s$. Then the sum of the lengths of the edges of that box is a multimetric function of $s$.

On the other hand, while the volume of that box obeys conditions (2) and (3), in general it can violate condition (1). Similarly, the volume of the convex hull of the (support of) the points in $s$ obeys conditions (2) and (3) but can violate (1). (In general, multimetrics have the dimension of a length, so volumes have to be raised to the appropriate power to make them be multimetrics.)

It is worth comparing the sum-of-edge-lengths multimetric to the standard deviation multimetric of Ex. 1 for the case where all arguments $s$ are finite sums of delta functions (i.e., “consist of a finite number of points”). For such an $s$ we can write the sum-of-edge-lengths multimetric as a sum over all $N$ dimensions $i$ of $\max_j s_{ij}^2 - \min_j s_{ij}^2$, where $s_{ij}$ is the $j$th point in $s$. In contrast, the (square of the) standard deviation multimetric is also a sum over all $i$, but of (the square of) the standard deviation of the $i$th components of the points in $s$. Another difference is that the standard deviation multimetric is a continuous function of its argument, unlike the sum-of-edge-lengths multimetric.

Example 5: Let $X$ be countable and have $T(X) = \mathcal{R}^X$. Then $\rho(s) = \int dx \Theta(s(x)) - 1$ where $\Theta$ is the Heaviside function is a multimetric (see appendix). This is the volume of the support of $s$, minus 1.

Example 6: Let $X$ be countable and have $T(X) = \mathcal{R}^X$. Then $\rho(s) = ||s|| - \max_s s(x)$ obeys conditions (2) and (3), by canceling. For this $\rho$ condition (1) holds iff $\max_x u(x) + v(x)) \geq \max_x u(x) + w(x) + \max_x v(x) + w(x) - 2||u||$. This is not true in general, for example when $||u|| = 0$ and the supports of $u$ and $v$ are disjoint. However if we take $T(X)$ to be the unit simplex in $\mathcal{R}^X$, then condition (1) is obeyed, and $\rho$ is a multimetric (see appendix).

Example 7: Let $X$ have a finite number of elements and set $T(X) = \mathcal{R}^X$. Say that $\rho(s) = 0$ along all of the axes, and that everywhere else, $k \leq \rho(s) \leq 2k$ for some fixed $k > 0$. Then $\rho$ is a multimetric.

A. Vector-valued multimetrics

It is straightforward to extend the definition of a multimetric to have range $\mathbb{R}^M$ rather than $\mathbb{R}$, so long as one has a linear ordering over $\mathbb{R}^M$ to specify the appropriate extension of condition (1). For example, consider the component-wise ordering: $\forall a, b, \mathbb{R} \leq b \Leftrightarrow a_i \leq b_i \forall i \in \{1, 2, \ldots, M\}$. Say we have a set of $M$ scalar multimetrics.

Then the $M$-fold Cartesian product of those multimetrics is an $M$-dimensional multimetric, when component-wise ordering defines the inequality in condition (1).

More generally, say we have chosen such a linear ordering over $\mathbb{R}^M$, and have an $M$-dimensional function with domain $\mathcal{R}^X$. Say this function obeys conditions (1) and (2) of an $M$-dimensional multimetric for our linear ordering. Then this function can be used as a low-dimensional characterization of an element of $\mathcal{R}^X$. In general, such characterizations may have $M$ is less than $|\mathcal{R}^X|$, the dimension of the space in which $\mathcal{R}^X$ is embedded, and may violate condition (1). The following examples illustrates this:

Example 8: Consider again Ex. 1. To define our vector-valued multimetric for the $X$ of Ex. 1, say we have a scalar multimetric $\rho_X$ for the subspace, $X' = T'(X') = \mathbb{R}$. Let $\{v^1, v^2, \ldots, v^M\}$ be a set of $M$ unit norm vectors living in $X$. Then we can define our $M$-dimensional multimetric by

$$\rho_X(u) \equiv \rho_X[f_{u,v'}(t)];$$

$$f_{u,v'}(t) = \int dx \frac{u(x)dx}{\delta(v' \cdot x - t)}.$$

To illustrate this, take $M = N$ and have the $\{v^i\}$ be the unit normals along the $N$ axes of $X$. Let $u$ be a sum of delta functions; $u = \delta_x + \delta_2$. Let $\rho_X$ be the standard deviation multimetric of Ex. 1 for one-dimensional probability density functions. So each component $\rho_i(u)$ is just the $i$th component of the difference $x^i - x^i$. Accordingly, $u$ can be reconstructed from the vector $\rho(u)$.

In this illustration conditions (1) and (2) are immediate. If $M = N$, then condition (3) also holds for $u$'s like the one considered here that are sums of two delta functions, but not more generally. Now modify this illustration by having $M < N$ and the $\{v^i\}$ not all point along the axes of $X$. Then for general $u$, the components $\rho_i(u)$ are the projections of $u$ along the different vectors $\{v^i\}$. As in techniques like Principal Components Analysis [2], those projections provide a low-dimensional characterization of $u$. 
III. CONCAVITY GAPS AND DISPERSIONS

In Ex. 1, $\rho$ can be used to tell us how spread out a distribution over $\mathbb{R}^N$ is. One would like to be able to use that $\rho$ to construct a measure of how spread out a collection of multiple distributions over $\mathbb{R}^N$ is. Intuitively, we want a way to construct a metric for a space of sets (generalized to be able to work sets with duplicates, fractional memberships, etc.) from a metric for the members within a single set. This would allow us to directly incorporate the distance relation governing $X$ into the distance relation for $\mathcal{R}^X$.

To do this, first let $\{Y, S(Y)\}$ be any pair of a subset of a vector space together with a subset of $\mathbb{R}^Y$ such that $\forall g \in S(Y), \frac{\partial g(y)}{\|g\|} \in Y$. (As an example, we could take $Y$ to be any convex subspace of a vector space, with $S(Y)$ any subset of $Y$.) Then the associated concavity gap operator $C : S(Y) \rightarrow \mathbb{R}^{S(Y)}$ is defined by

$$ (C\sigma)(g) = \sigma(\frac{\partial g(y)}{\|g\|}) - \frac{\partial g(y)}{\|g\|}\sigma(y) $$

where $y \in Y$, and both $\sigma$ and $g$ are arbitrary elements in $S(Y)$. So the concavity gap operator takes any single, member of the space $S(Y)$ (namely $\sigma$) and uses it to generate a function (namely, $C\sigma$) over all of $S(Y)$.[14]

In particular, say $Y = T(X)$ for some space $X$. Say we are given a multimetric $\sigma$ measuring the (X-space) spread specified by any element of $Y$. Say we are also given a $g$ which is a normalized distribution over $Y$. Then $C\sigma(g)$ is a measure of how spread out the distribution $g$ is. Note that in this example space $S(Y)$ is both the space of multimetrics over $Y$ and the space of distributions for $Y$, exemplified by $\sigma$ and $g$, respectively.

We can rewrite the definition of the concavity gap in several ways:

$$ C\sigma(g) = \sigma(E_g(y)) - E_g(\sigma) $$

where $E_g$ means expected value evaluated according to the probability distribution $\frac{\partial g(y)}{\|g\|}$, and in the last expression $y$ is the (infinite-dimensional) matrix whose $y$th column is the vector $y_i$, and the inner products are over the vector space $S(Y)$. Taken together, these equations say that the concavity gap of $\sigma$, applied to the distribution $g$, is given by evaluating the function $\sigma$ at the center of mass of the distribution $g$, and then subtracting the inner product between $\sigma$ and $g$.

**Example 9:** Let $Y = \mathbb{R}^N$, and choose $S(Y)$ to be the set of nowhere-negative functions of $Y$ with non-zero magnitude. Choose $\sigma(y) = 1 - \sum_{i=1}^{N} y_i^2$. Then $C\sigma(g) = \text{Var}(\frac{\partial g(y)}{\|g\|})$.

**Example 10:** Say $X$ has $N$ values, with $T(X) = \mathbb{R}^X$. Consider a $u \in T(X)$ whose components are all either 0 or some particular constant $a$ such that $\int dx u(x) = 1$. So $u$ is a point on the unit hypercube in $T(X)$, projected down to the unit simplex. Let $T$ be the set of all such points $u$. In the usual way, the support of each element of $T$ specifies a set of elements of $X$.

Let $Y = T(X)$, and have $S(Y) = \mathbb{R}^Y$. Have $g$ be a uniform average of a countable set of delta functions, each of which is centered on a member of $T$. So each of the delta functions making up $g$ specifies a set of elements of $X$; $g$ is a specification of a collection of such $X$-sets.

In this scenario $\sigma(E_g(y))$ is $\sigma$ applied to the union (over $X$) of all the $X$-sets specified in $g$. In contrast, $E_g(\sigma)$ is the average value you get when you apply $\sigma$ to one of the $X$-sets specified in $g$. $C\sigma(g)$ is the difference between these two values. Intuitively, it reflects how much overlap there is among the $X$-sets specified in $g$.

**Example 11:** Say $X$ has $N$ values, with $T(X) = \mathbb{R}^X$. Have $Y = T(X)$, and $S(Y) = \mathbb{R}^Y$, i.e., the set of all nowhere-negative non-zero functions over those points in $\mathbb{R}^N$ with no negative components. Choose $\sigma(y) = H(y) \forall y \in Y$, where $H(.) = -\int dy y(x)ln[y(x)]$, the Shannon entropy function extended to non-normalized $y$. This $\sigma$ is a natural choice to measure how "spread out" any point in $Y$ with magnitude 1 is.

Have $g$ be a sum of a set of delta functions, about the distributions over $\mathbb{R}$, $\{\delta_1, \delta_2, \ldots\}$. Then $C\sigma(g)$ is a measure of how "spread out" those distributions are. In the special case where $g = \delta_{a_1} + \delta_{a_2}$, $C\sigma(g)$ is the Jensen-Shannon divergence between $\delta_1$ and $\delta_2$ [6, 8]. More generally, if $g$ is a probability density function across the space of all distributions over $\mathbb{R}$, $C\sigma(g)$ is a measure of how "spread out" that density function is.

There are several elementary properties of concavity gaps worth mentioning:

**Lemma 2:**

1. C is linear.
2. C is linear $\Leftrightarrow$ it equals 0 everywhere $\Leftrightarrow$ $\sigma$ is linear.
3. C is continuous $\Leftrightarrow$ $\sigma$ is continuous.
4. $C\sigma(g) = 0$ if $g \propto \delta_{y'}$ for some $y' \in Y$.
5. Giving $C\sigma$ and the values of $\sigma$ at $1 + |Y|$ distinct points in $Y$ fixes the value of $\sigma$ across all $Y$. ($|Y|$ is the dimension of $Y$.)
6. The equivalence class of all $\sigma'$ having a particular concavity gap $C\sigma$ is the set of functions of $g \in Y$ having the form $\{\sigma(y) + b \cdot y + a : a \in \mathbb{R}, b \in Y, \sigma(y) + b \cdot y + a \in S(Y)\}$.

**Proof:** (2.1) and (2.4) are immediate. The first iff in (2.2) follows from the fact that $C\sigma(g) = C\sigma(\alpha g)$ $\forall \alpha \in \mathbb{R}$.
Lemma 3: \( g = \delta y / 2 \) and \( h = \delta y / 2 \) and expand \( C \sigma (g + h) = \sigma (y + y') - [\sigma (y) + \sigma (y')]. \) To see the forward direction of (2.3), choose \( g \) to have its center of mass infinitesimally to one side of the discontinuity in \( \sigma \), and then move it infinitesimally to the other side to get a discontinuity in the associated values of \( (C(\sigma))(g) \).

To prove (2.5), consider the case where \( Y \) is one-dimensional for simplicity. Say I give you \( \sigma \) at \( A \) and \( B > A \), and also give you \( C \sigma \). Then for every \( C > B \), choose \( g = \frac{B}{2} \delta A + \frac{B-C}{2} \delta C \). Evaluating the associated value \( C \sigma (g) = \sigma (B) - \frac{B}{2} \delta A + \frac{B-C}{2} \sigma (C) \) allows us to solve for \( \sigma (C) \). Similar reasoning holds for \( C < A \) and \( C \in (A,B) \). For higher dimensions we need the value of \( \sigma \) at one extra point for each extra dimension of \( Y \). This completes the proof.

To prove (2.6), first note that all members of that set do indeed have the same concavity gap, \( C \sigma \). To complete the proof we must show that there are no other \( \sigma \) with that concavity gap. Let \( \sigma' \) be any element of \( S(Y) \) with the same concavity gap as \( \sigma \). By (2.5), if we know the value of \( \sigma' \) at a total of \( 1 + |Y| \) points in \( Y \), then we know \( \sigma' \) in toto. In turn, for any such set of \( 1 + |Y| \) values, we can always find an \( a \) and \( b \) such that \( a + b \cdot y + \sigma (y) \) lies in \( S(Y) \) and has those values. This means that \( \sigma' \) is identical to \( a + b \cdot y + \sigma (y) \).

Proof: (3.1) arises from the fact that a dispersion is both concave and nowhere negative. To establish (3.2), first consider any two vectors \( u,v \) in the interior of \( T(X) \) that differ in only one component, \( i \). Since no component of \( u \) equals \( 0 \), there must be an \( s \in T(X) \) such that \( ||u|| = ||v|| - s \). Otherwise \( s = 0 \), so \( s_j = u_j - v_j \) for \( j \neq i \). This means that if \( \sigma \) is independent of the magnitude of its argument, \( \sigma (v) \leq \sigma (u) \). Since the reverse argument also holds, we have \( \sigma (u) = \sigma (v) \). Now repeat this reasoning to equate \( \sigma (v) \) with \( \sigma (w) \) for some \( w \) that differs from \( v \) in only one component, but differs from \( u \) in two components. Continuing in this way, we equate \( \sigma (u) \) with \( \sigma (z) \) for any \( z \) that differs in an arbitrary number of components from \( u \).

Example 12: Take \( X = \{1,2\} \), with \( T(X) = \mathbb{R}^X = \mathbb{R}^2 \). Define \( \sigma (u) \) to equal \( 0 \) if \( u \) equals \( 0 \) and \( \ln (1 + u_1) + \ln (1 + u_2) \) otherwise. Then \( \sigma \) is a (not everywhere continuous) strict dispersion.

Example 13: The \( X, \mathbb{R}^X \), and \( \rho \) of Ex. 3 form a strict dispersion (see appendix).

Example 14: The \( X, \mathbb{R}^X \), and \( \sigma \) of Ex. 5 form a strict dispersion.

Example 15: The \( X, \mathbb{R}^X \), and \( \sigma \) of Ex. 11 form a strict dispersion.

There are several relations between concavity gaps and dispersions:

**Lemma 3:** Let \( T(X) = \mathbb{R}^X \).

1. \( \sigma \) is a dispersion over \( T(X) \) \( \Rightarrow \) \( \sigma \) is nowhere-decreasing over \( T(X) \).

2. \( \sigma \) is a dispersion over \( T(X) \) and \( \sigma (s) \) is independent of \( |s| \) \( \forall s \neq 0 \in T(X) \) \( \Rightarrow \) \( \sigma \) is constant over the interior of \( T(X) \).

3. \( \sigma \) is (strictly) concave over \( T(X) \) \( \Leftrightarrow \) \( C \sigma \) obeys condition (2) in full (and condition (3) over \( T(X) \).

4. Say that \( \sigma \) is continuous over \( T(X) \). Then \( C \sigma \) is separately (strictly) concave over each simplex in \( T(X) \) if \( \sigma \) is (strictly) concave over \( T(X) \).

**Proof:** (3.1) arises from the fact that a dispersion is both concave and nowhere negative. To establish (3.2), first consider any two vectors \( u,v \) in the interior of \( T(X) \) that differ in only one component, \( i \). Since no component of \( u \) equals \( 0 \), there must be an \( s \in T(X) \) such that \( ||u|| = ||v|| - s \). Otherwise \( s = 0 \), so \( s_j = u_j - v_j \) for \( j \neq i \). This means that if \( \sigma \) is independent of the magnitude of its argument, \( \sigma (v) \leq \sigma (u) \). Since the reverse argument also holds, we have \( \sigma (u) = \sigma (v) \). Now repeat this reasoning to equate \( \sigma (v) \) with \( \sigma (w) \) for some \( w \) that differs from \( v \) in only one component, but differs from \( u \) in two components. Continuing in this way, we equate \( \sigma (u) \) with \( \sigma (z) \) for any \( z \) that differs in an arbitrary number of components from \( u \).

(3.3) is immediate from the definition of concavity and Jensen's inequality. To derive (3.4), first expand \( C \sigma (s_{a+b}) - C \sigma (s_a) + C \sigma (s_b) = \frac{1}{2} (E_u(s_a) + E_u(s_b)) - \frac{1}{2} (E_u(s_a) + \sigma (E_u(s_b))) \) when \( ||a|| = ||b|| \) (\( y \) being a generic argument of \( T(X) \)). In other words, this equality holds when \( a \) and \( b \) are on the same (not necessarily unit) simplex. Next invoke (2.3) to allow us to apply Jensen's inequality. QED.

Let \( f : \mathbb{R} \to \mathbb{R} \) be monotonically increasing and strictly concave. Then by Lemma 3.3, if \( \sigma \) is strictly concave, \( f(C \sigma) \) obeys conditions (2) and (3). For example, this is the case for \( \sqrt{C \sigma} \). In other words, so long as \( \sigma \) is a strict dispersion, \( \sqrt{C \sigma} \) obeys those conditions.

On the other hand, Lemma 3.2 means that any non-trivial \( \sigma \) that normalizes its argument (so that it is a probability distribution) and then evaluates a function of that normalized argument cannot be a dispersion. So for example, if a concavity gap is a dispersion, it must be constant.

Fortunately it is not the case that if \( f(C \sigma) \) is a multimetric it must be constant. In particular, often for a strictly concave \( \sigma \), \( \sqrt{C \sigma} \) for space \( \{Y,S(Y)\} \) is a multimetric for an appropriate \( T(Y) \subseteq S(Y) \).

Example 16: Choose \( \{Y,S(Y)\} \) as in Ex. 11, and take \( T(Y) \) to be all elements of \( S(Y) \) which are sums of two delta functions. This \( \sigma \) is strictly concave, so we know conditions (2) and (3) are obeyed by \( \sqrt{C \sigma} \). Furthermore, for this choice of \( T(Y) \), obeying condition (1) reduces to obeying the conventional triangle inequality of two-argument metrics, and it is known that the square root of the Jensen Shannon divergence obeys that inequality [3, 8]. Therefore all three conditions are met.
Example 17: Choose \( \{ \sigma, Y, S(Y) \} \) as in Ex. 9. As in Ex. 16, this \( \sigma \) is strictly concave, and therefore \( \sqrt{\sigma} \) automatically obeys conditions (2) and (3). Now take \( T(Y) = S(Y) \). Write \( C\sigma(g) \) as \( \langle g, g \rangle \) for the tensor contraction of Ex. 3, where \( F(y, y') = \frac{(y-y')^2}{2} \). So by that example, we know that \( \sqrt{\sigma} \) is a multimetric.

IV. POTENTIAL USES OF MULTIMETRICS

In addition to their intrinsic mathematical interest, multimetrics have numerous potential applications. One of them is to allow more nuanced complexity measures for physical systems, as described in [9]. Following is a list of some other from machine learning [2]:

1. Mixture of Gaussians density estimation: In density estimation one is given a data set of vectors \( \{x^i\} \) that were generated by IID sampling an unknown distribution over \( X \) and wants to infer that distribution. Say we have a probability distribution across \( X \) given by a linear combination of \( n \) Gaussian distributions over \( X \), centered at the \( n \) points \( \mu^i \). Such a distribution induces a probability of the set \( \{x^i\} \). Accordingly, one way to estimate the distribution that generated \( \{x^i\} \) is to search for the linear combination of \( n \) Gaussians that maximizes the associated probability of \( \{x^i\} \).

One shortcoming of this procedure is that even though there are a total of \( n \) points used to parameterize the distribution, that distribution is based solely on metric values for pairs of points (namely the distances between \( x \) and each of the \( \mu^i \)). If we have a multimetric \( \rho \) though, we have several ways to avoid this. For example, we could model the probability of each \( x^i \) as a Gaussian of \( \rho(\delta_{x^i} + \sum_j \delta_{x^j}) \). We would then take the probability of \( \{x^i\} \) to be the product of the probabilities of the \( x^i \), as in conventional Gaussian mixtures modeling. We could even model the probability of \( \{x^i\} \) given the \( n \) points \( \mu^i \) as a single Gaussian, with argument \( \rho(\sum_i \delta_{x^i} + \sum_j \delta_{x^j}) \).

2. Kernel density estimation: In kernel density estimation, one does not estimate the distribution over \( x \) as a linear combination of \( n \) kernel functions (e.g., Gaussians) that are free to be centered anywhere in \( X \), and then search for which such linear combination maximizes the probability of one’s data. Instead one centers a kernel function on each of the data points, and searches for the optimal parameters of those kernels functions. Conventionally such kernel functions only take two arguments. However exactly as in application 1, if one has a multimetric over \( X \), one can use kernel functions whose argument involves more than two points at once.

3. Classification can always be done via density estimation and Bayes’ theorem. So with applications 1 and 2, we have new ways of doing classification.

4. Kernel machines are a recent advance in machine learning in which data is first mapped non-linearly into a feature space where standard algorithms (like linear regression, linear discriminant analysis, PCA, etc) are applied [1]. Because of the non-linear mapping such methods work even when relationships in the data are highly non-linear. All that is required for such methods is a positive definite kernel function, \( k(x, x') \), giving inner products in the feature space. Multimetrics are not positive definite functions but can easily be made so by taking \( k(x, x') = \exp\{-\rho(\delta_x + \delta_{x'})\} \) as the kernel. So any of the multimetrics discussed above can be used for statistical analysis with kernel-based learning algorithms. In particular, this is the case for either supervised or unsupervised learning with kernel machines. In particular, we can use exponentials of multimetrics for regression by using them instead of the conventional kernels of kernel machines, with the multiplicative coefficients of each kernel (in the linear combination that gives our fit to the data) set to minimize some appropriate quadratic objective function.

V. APPENDIX

A. Proof of Lemma 1.3

That \( f(\rho) \) obeys conditions (2) and (3) when \( \rho \) does is immediate. To prove that condition (1) is obeyed, consider any \( u, v, w \in T(X) \) such that \( \rho(u + v) \leq \rho(u + w) + \rho(v + w) \). First assume that \( \rho(u + v) = \max[\rho(u + w), \rho(v + w)] \). Then since \( f \) is increasing, \( f(\rho(u + v)) \leq \max[f(\rho(u + w)), f(\rho(v + w))] \). Since in turn \( \max[f(\rho(u + w)), f(\rho(v + w))] \leq f(\rho(u + w)) + f(\rho(v + w)) \), condition (1) is obeyed.

Now consider the other case, where \( \rho(u + v) > \max[\rho(u + w), \rho(v + w)] \). In this situation, because \( f \) is concave, we know that \( f \) increases \( \rho(u + v) \) less than it increases both \( \rho(u + w) \) and/or \( \rho(v + w) \). So again condition (1) is obeyed. QED.

B. Proof of claim in Ex. 1

Consider any \( u, v \) and \( w \) whose norms equal 1. Then squaring both sides of condition (1) for our \( \rho \) implies that

\[
\frac{\text{Var}(\frac{u + v}{2})}{\text{Var}(\frac{u + w}{2})} \leq \frac{\text{Var}(\frac{u + v}{2}) + \text{Var}(\frac{u + w}{2}) + 2\sqrt{\text{Var}(\frac{u + v}{2})\text{Var}(\frac{u + w}{2})}}{\text{Var}(\frac{u + w}{2})}\]

Use the expansion \( \text{Var}(\frac{x + t}{2}) = \frac{\text{Var}(x) + \text{Var}(t) + (E_x - E_x)^2}{2} \) and cancel terms. The hardest case for
the resultant inequality to hold is where our three vari-
ances all equal 0. Setting them to 0, we see that condition
(1) holds if for any three real numbers \(a, b, c\),
\[
|a - b| \leq |a - c| + |b - c|.
\]
This is just the conventional triangle inequality though.
So condition (1) always holds. QED.

C. Proof of claim in Ex. 3

First note that for any \(s \in T(X), (s | s) \geq 0\), since all
all components of \(F\) are non-negative. Furthermore,
all \(s, t \in T(X), (s | t) \geq 0\), since all components of those
vectors are non-negative as are all components of \(F\). In
addition, we can use the properties of \(F\) to prove that our
tensor contraction obeys the Cauchy-Schwarz inequal-
ity: \(\langle u | v \rangle^2 \leq \langle u | u \rangle \langle v | v \rangle \forall u, v \in T(X)\). (Examped
\(s, s) \geq 0\) for \(s = u - \alpha v\). Solve for the \(\alpha\) minimizing the
lethand side (which is quadratic in \(\alpha\)), and plug that in.
Collecting terms establishes the desired inequality.

Now to check condition (1) for our \(\rho\), square both sides
of it and cancel terms. So the lefthand side is just \((u | v)\).
Since all inner products are non-negative, the right-hand
side is bounded below by \(\sqrt{(u | u)(v | v)}\). Plugging in the
Cauchy-Schwarz inequality establishes that condition (1)
does indeed hold. QED.

D. Proof of claim in Ex. 4

Conditions (2) and (3) are immediate. To prove con-
dition (1), first note that it holds for \(\rho_1(s) = \max(x_1 : s(x_1 \neq 0)\). Then note that it holds for \(\rho_1(s) = \min(x_1 : s(x_1 \neq 0)\), and invoke Lemma 1.2, to see that the width
in \(x_1\) of the support obeys condition (1). QED.

E. Proof of claim in Ex. 5

Conditions (2) and (3) are immediate. Condition (1)
also holds if (the supports of) \(u\) and \(v\) overlap, since any
non-zero volume must equal at least 1, and that overlap
volume gets counted twice in the sum \(\rho(u + w) + \rho(v + w)\), regardless of \(w\). If (the supports of) \(u\) and \(v\) do
not overlap, then (the support of) \(w\) must either extend
outside of (the support) of \(u\) or of \(v\). This means that
condition (1) must hold in this case as well. QED.

F. Proof of claim in Ex. 6

Define \(\text{argmax}_w u(x) \equiv a, \text{argmax}_w v(x) \equiv b, \max_w(u(x) + v(x)) \equiv M, \text{and } \max_w(u(x) + w(x)) + \max_w(v(x) + w(x)) - 2||w|| \equiv N\); we want to prove
that \(M > N\). To that end, note that if the support
of \(w(x)\) is restricted to \(a\) and \(b\), then \(N\) becomes
\(u(a) + v(b) - ||w|| = u(a) + v(b) - 1 \leq u(a)\). On the
other hand, \(M\) is bounded below by \(u(a)\). So condition
(1) holds for this situation.

We now consider the situation where \(w\)'s support is
not restricted to \(a\) and \(b\). It will be useful to define
\(\text{argmax}_w(u(x) + w(x)) = d\) and \(\text{argmax}_w(v(x) + w(x)) = e\).
First consider the case where \(d \neq e\). Then it is immedi-
ate that by transferring any \(w(x) \not\in \{d, e\}\) to \(w(d)\) and/or
\(w(e)\), we do not decrease \(N\) (since \(||w||\) doesn't change).
We can then transfer \(w(d)\) to \(w(a)\) and \(w(e)\) to \(w(b)\),
again not decreasing \(N\). After doing this for all such
points \(c\), we recover the case where the support of \(w(x)\)
is restricted to \(a\) and \(b\), as in the preceding paragraph.
So we can conclude that condition (1) is obeyed for this
case.

The remaining case to consider is where \(d = e\). For this
case we can transfer all \(w(c \neq d)\) to \(w(d)\), and in doing so
increase \(N\). Doing this for all such \(c\) restricts \(w\)'s support
to \(d\). After having done this, \(\max_w(u(x) + w(x)) = w(d) +
||w|| \leq u(a) + 1\), and similarly \(\max_w(v(x) + w(x)) \leq
v(b) + 1\). So we again get \(N \leq u(a) = v(b) - 1\), which
means \(M > N\). There are no more cases to consider.
QED.

G. Proof of claim in Ex. 13

First note that we have already established that our
\(\rho\) obeys conditions (2) and (3) of being a multimetric,
and therefore only need to establish that it is strictly
concave. That will be the case if \(\rho(\alpha u + (1 - \alpha)v) \leq
\alpha \rho(u) + (1 - \alpha) \rho(v) \forall u \in T(X), v \in T(X), \alpha \in [0, 1]\).
Square both sides of this inequality and cancel terms.
Then exploit the Cauchy Schwarz inequality for \(\langle \cdot, \cdot \rangle\),
established in the proof of the claim in Ex. 3. QED.

VI. ACKNOWLEDGEMENTS

I would like to thank Bill Macready and Creon Levit
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[1] Nello Christianini and John Shawe-Taylor. An introd-
uction to support vector machines and other kernel-based
[12] A natural extension of this analysis allows $\mathcal{R}^X$ to be any subset of $\mathbb{R}^X$ — potentially including vectors with negative components — so long as $k\delta_x \in \mathcal{R}^X \forall x \in X, k > 0$. Indeed, with care, one can even take the range of $\rho$ to include complex numbers, so that $\rho$ includes Hilbert vectors like those that arise in quantum mechanics. In the interests of clarity, we do not pursue such extensions here.
[13] Such mappings are analogous to the mapping from a “data space” to a “feature space” underlying kernel machines [1].
[14] Equivalently, it can be viewed as a non-symmetric function from $S(Y) \times S(Y) \rightarrow \mathbb{R}$, although we will not exploit that perspective here.