AN OPERATOR THEORETIC APPROACH TO NONPARAMETRIC MIXTURE MODELS∗

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When estimating finite mixture models, it is common to make assumptions on the mixture components, such as parametric assumptions. In this work, we make no distributional assumptions on the mixture components and instead assume that observations from the mixture model are grouped, such that observations in the same group are known to be drawn from the same mixture component. We precisely characterize the number of observations $n$ per group needed for the mixture model to be identifiable, as a function of the number $m$ of mixture components. In addition to our assumption-free analysis, we also study the settings where the mixture components are either linearly independent or jointly irreducible. Furthermore, our analysis considers two kinds of identifiability – where the mixture model is the simplest one explaining the data, and where it is the only one. As an application of these results, we precisely characterize identifiability of multinomial mixture models. Our analysis relies on an operator-theoretic framework that associates mixture models in the grouped-sample setting with certain infinite-dimensional tensors. Based on this framework, we introduce a general spectral algorithm for recovering the mixture components.

1. Introduction. A finite mixture model $\mathcal{P}$ is a probability measure over a space of probability measures where $\mathcal{P}(\{\mu_i\}) = w_i > 0$ for some finite collection of probability measures $\mu_1, \ldots, \mu_m$ and $\sum_{i=1}^{m} w_i = 1$. A realization from this mixture model first randomly selects some mixture component $\mu \sim \mathcal{P}$ and then draws from $\mu$. Mixture models have seen extensive use in statistics and machine learning.

A central theoretical question concerning mixture models is that of identifiability. A mixture model is said to be identifiable if there is no other mixture model that defines the same distribution over the observed data. Classically mixture models were concerned with the case where the observed data

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\(X_1, X_2, \ldots\) are iid with \(X_i\) distributed according to some unobserved random measure \(\mu_i\) with \(\mu_i \overset{iid}{\sim} \mathcal{P}\). This situation is equivalent to \(X_i \overset{iid}{\sim} \sum_{j=1}^{m} w_j \mu_j\). If we impose no restrictions on the mixture components \(\mu_1, \ldots, \mu_m\) one could easily concoct many choices of \(\mu_j\) and \(w_j\) which yield an identical distribution on \(X_i\). Because of this, most previous work on identifiability assumes some sort of structure on \(\mu_1, \ldots, \mu_m\), such as Gaussianity \([3, 9, 25]\). In this work we consider an alternative scenario where we make no assumptions on \(\mu_1, \ldots, \mu_m\) and instead have access to groups of samples that are known to come from the same component. We will call these groups of samples “random groups.” Mathematically a random group is a random element \(X_i\) where \(X_i = (X_{i,1}, \ldots, X_{i,n})\) with \(X_{i,1}, \ldots, X_{i,n} \overset{iid}{\sim} \mu_i\) and \(\mu_i \overset{iid}{\sim} \mathcal{P}\).

In this setting identifiability is now concerned with the distribution over \(X_i\) and the value of \(n\), the number of samples in each random group. We call a mixture of measures \(\mathcal{P}\) \(n\)-identifiable if it is the simplest mixture model (in terms of number of mixture components) that yields the observed distribution on \(X_i\). We also introduce a concept which is stronger than identifiability. We call \(\mathcal{P}\) \(n\)-determined if it is the only mixture model that yields the observed distribution on \(X_i\).

In this paper we show that every mixture model with \(m\) components is \((2m-1)\)-identifiable and \(2m\)-determined. Furthermore we show that any mixture model with linearly independent components is 3-identifiable and 4-determined, and any mixture model with jointly irreducible components is 2-determined. These results, presented in Section 4, hold for any mixture model over any space and cannot be improved. The operator theoretic framework underlying our analysis is presented in Section 5, and selected proofs of our main results appear in Section 6, with the rest appearing in Supplemental Material. In Section 7, we apply our main results to demonstrate some new and old results on the identifiability of multinomial mixture models. Section 8 describes a spectral algorithm for the recovery of the mixture components and weights, and experimental results on simulated data are presented in Section 9. Related work, the problem formulation, and a concluding discussion are offered in Sections 2, 3, and 10, respectively.

To keep the paper length reasonable, many of the proofs have been omitted and can be found in Supplemental Material. Supplemental Material also contains an in-depth description of the application of our spectral algorithm to categorical data (including a consistency proof) and additional technical details regarding the experiments in Section 9.

2. Previous Work. In classical mixture model theory identifiability is achieved by making assumptions about the mixture components. Some
assumptions which yield identifiability are Gaussian or binomial mixture components \cite{9, 24}. If one makes no assumptions on the mixture components then one must leverage some other type of structure in order to achieve identifiability. An example of such structure exists in the context of multiview models. In a multiview model samples have the form $X_i = (X_{i,1}, \ldots, X_{i,n})$ and the distribution of $X_i$ is defined by $\sum_{i=1}^m w_i \prod_{j=1}^n \mu_j^i$. In \cite{1} it was shown that if $\mu_j^i$ are probability distributions on $\mathbb{R}$ with $\mu_1^j, \ldots, \mu_m^j$ linearly independent for all $j$ and $n \geq 3$, then the model is identifiable. In \cite{7} the authors perform a smoothed analysis of tensor decompositions. They demonstrate that, with high probability, a tensor’s components are both identifiable and can be recovered using a polynomial time algorithm, provided the component dimensionality is sufficiently high. In that paper the authors go on to apply the result to multiview models, demonstrating bounds on identifiability.

The setting which we investigate is a special case of the multiview model where $\mu_j^i = \mu_j^{i'}$ for all $i, j, j'$. If the sample space of the $\mu_i$ is finite then this problem is exactly the topic modelling problem with a finite number of topics and one topic for each document. In topic modelling each $\mu_i$ is a “topic” and the sample space is a finite collection of words. This setting is well studied and it has been shown that one can recover the true topics provided certain assumptions on the topics are satisfied \cite{1, 2, 4, 5}. This problem was studied for arbitrary topics in \cite{22}. In this paper the authors introduce an algorithm that recovers any mixture of $m$ topics provided $2^m - 1$ words per document. They also show, in a result analogous to our own, that this $2^m - 1$ value cannot be improved. Our proof techniques are quite different than those used in \cite{22}, hold for arbitrary sample spaces, and are less complex.

In Lemma 7.1 we show that, when restricted to finite sample spaces, the grouped sample setting introduced in this paper is equivalent to a multinomial mixture model. Fundamental bounds on the identifiability of multinomial mixture models can be found in \cite{17, 12}. We will reproduce these results (and develop some new results) using techniques developed in this paper. Additional connections to previous work are given later.

3. Problem Setup. We treat this problem in a general setting. For any measurable space we define $\delta_x$ as the Dirac measure at $x$. For $\mathcal{Y}$ a set, $\sigma$-algebra, or measure, we denote $\mathcal{Y}^{\times a}$ to be the standard $a$-fold product associated with that object. Let $\mathbb{N}$ be the set of integers greater than or equal to zero and $\mathbb{N}_+$ be the integers strictly greater than 0. For $k \in \mathbb{N}_+$, we define $[k]$ to be those elements in $\mathbb{N}_+$ which are less than or equal to $k$. Let $\Omega$ be a set containing more than one element. This set is the sample
space of our data. Let $\mathcal{F}$ be a $\sigma$-algebra over $\Omega$. Assume $\mathcal{F} \neq \{\emptyset, \Omega\}$, i.e. $\mathcal{F}$ contains nontrivial events. We denote the space of probability measures over a measurable space $(\Psi, \mathcal{G})$ as $\mathcal{D}(\Psi, \mathcal{G})$. The space $\mathcal{D}(\Omega, \mathcal{F})$ will be shortened to $\mathcal{D}$ for brevity. We equip $\mathcal{D}$ with the $\sigma$-algebra $2^\mathcal{D}$ so that each Dirac measure over $\mathcal{D}$ is unique. Define $\Delta(\mathcal{D}) \triangleq \text{span}(\{\delta_x : x \in \mathcal{D}\})$. This is the ambient space where our mixtures of probability measures live.

Let $P = \sum_{i=1}^{m} w_i \delta_{\mu_i}$ be a probability measure in $\Delta(\mathcal{D})$. Let $\mu \sim P$ and $X_1, \ldots, X_n \overset{iid}{\sim} \mu$ and denote $X = (X_1, \ldots, X_n)$. Here $X$ is a random group sample, which was described in the introduction.

We now derive the probability distribution of $X$. Let $A \in \mathcal{F}^\times n$. Letting $P$ reflect both the draw of $\mu \sim P$ and $X_1, \ldots, X_n \overset{iid}{\sim} \mu$, we have

\[ P(X \in A) = \sum_{i=1}^{m} P(X \in A | \mu = \mu_i) P(\mu = \mu_i) \]

\[ = \sum_{i=1}^{m} w_i \mu_i^{\times n}(A). \]

The second equality follows from Lemma 3.10 in [16]. So the probability distribution of $X$ is

\[ \sum_{i=1}^{m} w_i \mu_i^{\times n}. \]

We want to view the probability distribution of $X$ as a function of $P$ in a mathematically rigorous way, which requires a bit of technical buildup. Let $\mathcal{Q} \in \Delta(\mathcal{D})$. From the definition of $\Delta(\mathcal{D})$ it follows that $\mathcal{Q}$ admits a representation $\mathcal{Q} = \sum_{i=1}^{r} \alpha_i \delta_{\nu_i}$. From the well-ordering principle there must exist some representation with minimal $r$ and we define this $r$ as the order of $\mathcal{Q}$. We can show that the minimal representation of any $\mathcal{Q} \in \Delta(\mathcal{D})$ is unique up to permutation of its indices.

**Lemma 3.1.** Let $\mathcal{Q} \in \Delta(\mathcal{D})$ and admit minimal representations $\mathcal{Q} = \sum_{i=1}^{r} \alpha_i \delta_{\nu_i} = \sum_{j=1}^{r} \alpha'_j \delta_{\nu'_j}$. There exists some permutation $\psi : [r] \to [r]$ such that $\nu_{\psi(i)} = \nu'_i$ and $\alpha_{\psi(i)} = \alpha'_i$ for all $i$.

Proofs of most of the lemmas in this paper are omitted and can be found in Supplemental Material. The only lemma proved in this paper is Lemma
6.5 since it is nontrivial and quite crucial for showing some of our bounds are tight.

Henceforth when we define an element of \( \Delta(D) \) with a summation we will assume that the summation is a minimal representation.

**Definition 3.1.** We call \( \mathcal{P} = \sum_{i=1}^{m} w_i \delta_{\mu_i} \) a mixture of measures if it is a probability measure in \( \Delta(D) \). The elements \( \mu_1, \ldots, \mu_m \), are called mixture components.

Any minimal representation of a mixture of measures \( \mathcal{P} \) with \( m \) components satisfies \( \mathcal{P} = \sum_{i=1}^{m} w_i \delta_{\mu_i} \) with \( w_i > 0 \) for all \( i \) and \( \sum_{i=1}^{m} w_i = 1 \). Hence any mixture of measures is a convex combination of Dirac measures at elements in \( D \).

For a measurable space \((\Psi, \mathcal{G})\) we define \( \mathcal{M}(\Psi, \mathcal{G}) \) as the space of all finite signed measures over \((\Psi, \mathcal{G})\). We can now introduce the operator \( V_n : \Delta(D) \to \mathcal{M}(\Omega^x, \mathcal{F}^x) \). For a minimal representation \( \mathcal{Q} = \sum_{i=1}^{r} \alpha_i \delta_{\nu_i} \), we define \( V_n \), with \( n \in \mathbb{N}_+ \), as

\[
V_n(\mathcal{Q}) = \sum_{i=1}^{r} \alpha_i \nu_i^x.
\]

(4)

This mapping is well defined as a consequence of Lemma 3.1. From this definition we have that \( V_n(\mathcal{P}) \) is simply the distribution of \( X \) which we derived earlier. In the following definitions, two mixtures of measures are considered equal if they define the same measure.

**Definition 3.2.** We call a mixture of measures, \( \mathcal{P} \), \( n \)-identifiable if there does not exist a different mixture of measures \( \mathcal{P}' \), with order no greater than the order of \( \mathcal{P} \), such that \( V_n(\mathcal{P}) = V_n(\mathcal{P}') \).

**Definition 3.3.** We call a mixture of measures, \( \mathcal{P} \), \( n \)-determined if there exists no other mixture of measures \( \mathcal{P}' \) such that \( V_n(\mathcal{P}) = V_n(\mathcal{P}') \).

Definitions 3.2 and 3.3 are central objects of interest in this paper. Given a mixture of measures, \( \mathcal{P} = \sum_{i=1}^{m} w_i \delta_{\mu_i} \), then \( V_n(\mathcal{P}) \) is equal to \( \sum_{i=1}^{m} w_i \mu_i^x \), the measure from which \( X \) is drawn. If \( \mathcal{P} \) is not \( n \)-identifiable then we know that there exists a different mixture of measures that is no more complex (in terms of number of mixture components) than \( \mathcal{P} \) which induces the same distribution on \( X \). Practically speaking this means we need more samples in each random group \( X \) in order for the full richness of \( \mathcal{P} \) to be manifested in \( X \). A stronger version of \( n \)-identifiability is \( n \)-determinedness where we
enforce the requirement that our mixture of measures be the only mixture of measures (of any order) that admits the distribution on $X$.

A quick note on terminology. We use the term “mixture of measures” rather than “mixture model” to emphasize that a mixture of measures should be interpreted a bit differently than a typical mixture model. A “mixture model” connotes a probability measure on the sample space of observed data $\Omega$, whereas a “mixture of measures” connotes a probability measure on the sample space of the unobserved latent measures $\mathcal{D}$.

4. Main Results. The first result is a bound on the $n$-identifiability of all mixtures of measures with $m$ or fewer components. This bound cannot be uniformly improved.

**Theorem 4.1.** Let $(\Omega, \mathcal{F})$ be a measurable space. Mixtures of measures with $m$ components are $(2m - 1)$-identifiable.

**Theorem 4.2.** Let $(\Omega, \mathcal{F})$ be a measurable space with $\mathcal{F} \neq \{\emptyset, \Omega\}$. For all $m \geq 2$, there exists a mixture of measures with $m$ components that is not $(2m - 2)$-identifiable.

We mention again that the previous two theorems had been previously found in [22] for finite sample spaces, using techniques different from our own. To be explicit, a “finite sample space” in our problem setting is the assumption that $|\Omega| < \infty$ and $\mathcal{F} = 2^\Omega$, which implies that the mixture components are categorical distributions. The following lemmas convey the unsurprising fact that $n$-identifiability is, in some sense, monotonic.

**Lemma 4.1.** If a mixture of measures is $n$-identifiable then it is $q$-identifiable for all $q > n$.

**Lemma 4.2.** If a mixture of measures is not $n$-identifiable then it is not $q$-identifiable for any $q < n$.

Viewed alternatively these results say that $n = 2m - 1$ is the smallest value for which $V_n$ is injective over the set of mixtures of measures with $m$ or fewer components.

We also present an analogous bound for $n$-determinedness. This bound also cannot be improved.

**Theorem 4.3.** Let $(\Omega, \mathcal{F})$ be a measurable space. Mixtures of measures with $m$ components are $2m$-determined.
Theorem 4.4. Let $(\Omega, \mathcal{F})$ be a measurable space with $\mathcal{F} \neq \{\emptyset, \Omega\}$. For all $m$, there exists a mixture of measures with $m$ components that is not $(2m - 1)$-determined.

Again $n$-determinedness is monotonic in the number of samples per group.

Lemma 4.3. If a mixture of measures is $n$-determined then it is $q$-determined for all $q > n$.

Lemma 4.4. If a mixture of measures is not $n$-determined then it is not $q$-determined for any $q < n$.

This collection of results can be interpreted in an alternative way. Consider some pair of mixtures of measures $\mathcal{P}, \mathcal{P}'$. If $n \geq 2m$ and either mixture of measures is of order $m$ or less, then $V_n(\mathcal{P}) = V_n(\mathcal{P}')$ implies $\mathcal{P} = \mathcal{P}'$. Furthermore $n = 2m$ is the smallest value of $n$ for which the previous statement is true for all pairs of mixtures of measures.

Our definitions of $n$-identifiability, $n$-determinedness, and their relation to previous works on identifiability deserve a bit of discussion. Some previous works on identifiability contain results related to what we call “identifiability” and others contain results related what we call “determinedness.” Both of these are simply called “identifiability” in these works. For example in [25] it is shown that different finite mixtures of multivariate Gaussian distributions will always yield different distributions, a result which we could call “determinedness.” Alternatively, in [24] it is demonstrated that mixtures of binomial distributions, with a fixed number of trials $n$ for every mixture component, are identifiable provided we only consider mixtures with $m$ mixture components and $n \geq 2m - 1$. In this result allowing for more mixture components may destroy identifiability and thus this is what we would call an “identifiability” result. The fact that the value $2m - 1$ occurs in both the previous binomial mixture model result and Theorem 4.1 is not a coincidence. We will demonstrate a new determinedness result for multinomial mixture models later in the paper, under the assumption that $n \geq 2m$. We will prove these results using Theorems 4.1 and 4.3. To our knowledge our work is the first to consider both identifiability and determinedness.

Finally we also include results that are analogous to previously shown results for the finite sample space setting. We note that our proof techniques are markedly different than the previous proofs for the finite sample space case.

Theorem 4.5. If $\mathcal{P} = \sum_{i=1}^{m} w_i \delta_{\mu_i}$ is a mixture of measures where $\mu_1, \ldots, \mu_m$ are linearly independent then $\mathcal{P}$ is $3$-identifiable.
This bound is tight as a consequence of Theorem 4.2 with $m = 2$ since any pair of distinct measures must be linearly independent.

A version of this theorem was first proven in [1] by making use of Kruskal’s Theorem [18]. Kruskal’s Theorem demonstrates that order 3 tensors over $\mathbb{R}^d$ admit unique decompositions (up to scaling and permutation) given certain linear independence assumptions. The linear independence assumption in Theorem 4.5 is stronger than that contained in Kruskal’s Theorem and thus yields a simple proof which does not invoke Kruskal’s Theorem. An efficient algorithm for recovering linearly independent mixture components for finite sample spaces with 3 samples per random group is described in [2]. Interestingly, with one more sample per group, these mixtures of measures become determined.

**Theorem 4.6.** If $\mathcal{P} = \sum_{i=1}^{m} w_i \delta_{\mu_i}$ is a mixture of measures where $\mu_1, \ldots, \mu_m$ are linearly independent then $\mathcal{P}$ is 4-determined.

This bound is tight as a result of Theorem 4.4 with $m = 2$.

Our final result is related to the “separability condition” found in [11]. The separability condition in the finite sample space setting requires that, for each mixture component $\mu_i$, there exists $B_i \in \mathcal{F}$ such that $\mu_i(B_i) > 0$ and $\mu_j(B_i) = 0$ for all $i \neq j$. There exists a generalization of the separability condition, known as joint irreducibility.

**Definition 4.1.** A collection of probability measures $\mu_1, \ldots, \mu_m$ are said to be jointly irreducible if $\sum_{i=1}^{m} w_i \mu_i$ being a probability measure implies $w_i \geq 0$.

In other words, any probability measure in the span of $\mu_1, \ldots, \mu_m$ must be a convex combination of those measures. It was shown in [8] that separability implies joint irreducibility, but not visa-versa. In that paper it was also shown that joint irreducibility implies linear independence, but the converse does not hold.

**Theorem 4.7.** If $\mathcal{P} = \sum_{i=1}^{m} w_i \delta_{\mu_i}$ is a mixture of measures where $\mu_1, \ldots, \mu_m$ are jointly irreducible then $\mathcal{P}$ is 2-determined.

A straightforward consequence of the corollary of Theorem 1 in [11] is that any mixture of measures on a finite sample space with jointly irreducible components is 2-identifiable. The result in [11] is concerned with the uniqueness of nonnegative matrix factorizations and Theorem 4.7, when applied to a finite sample space, can be posed as a special case of the result.
in [11]. In the context of nonnegative matrix factorization the result in [11] is significantly more general than our result. In another sense our result is more general since it applies to spaces where joint irreducibility and the separability condition are not equivalent. Furthermore [11] only implies that the mixture of measures in Theorem 4.7 are identifiable. The determinedness result is, as far as we know, totally new. It is worth mentioning that, in finite sample spaces, the separability condition assumption yields efficient algorithms for nonnegative matrix factorization [4, 5], whereas we are not aware of analogous algorithms which are applicable to the more general joint irreducibility setting.

5. Tensor Products of Hilbert Spaces. Our proofs will rely heavily on the geometry of tensor products of Hilbert spaces which we will introduce in this section.

5.1. Overview of Tensor Products. First we introduce tensor products of Hilbert spaces. Instead of a rigorous primer to the subject, we will simply state some basic facts about tensor products of Hilbert spaces and hopefully instill some intuition for the uninitiated by way of example. A thorough treatment of tensor products of Hilbert spaces can be found in [15].

Let $H$ and $H'$ be Hilbert spaces. From these two Hilbert spaces the “simple tensors” are elements of the form $h \otimes h'$ with $h \in H$ and $h' \in H'$. We can define an inner product on the simple tensors by setting

$$\langle h_1 \otimes h'_1, h_2 \otimes h'_2 \rangle = \langle h_1, h_2 \rangle \langle h'_1, h'_2 \rangle.$$  

Let $H_0$ be the inner product space spanned by the simple tensors. The tensor product of $H$ and $H'$ is the completion of $H_0$ and is denoted $H \otimes H'$. To avoid potential confusion we note that the notation just described is standard in operator theory literature. In some literature our definition of $H_0$ is denoted as $H \hat{\otimes} H'$ and our definition of $H \otimes H'$ is denoted $H \tilde{\otimes} H'$.

As an illustrative example we consider the tensor product $L^2(\mathbb{R}) \otimes L^2(\mathbb{R})$. It can be shown that there exists an isomorphism between $L^2(\mathbb{R}) \otimes L^2(\mathbb{R})$ and $L^2(\mathbb{R}^2)$ that maps the simple tensors to separable functions [15], $f \otimes f' \mapsto f(\cdot)f'(\cdot)$. We can demonstrate this isomorphism with a simple example. Let $f, g, f', g' \in L^2(\mathbb{R})$. Taking the $L^2(\mathbb{R}^2)$ inner product of $f(\cdot)f'(\cdot)$ and $g(\cdot)g'(\cdot)$ gives us

$$\int \int (f(x)f'(y))(g(x)g'(y))dxdy = \int f(x)g(x)dx \int f'(y)g'(y)dy = \langle f, g \rangle \langle f', g' \rangle \quad \langle f \otimes f', g \otimes g' \rangle.$$  

In [11].
Beyond tensor product we will need to define tensor power. To begin we will first show that tensor products are, in a certain sense, associative. Let $H_1, H_2, H_3$ be Hilbert spaces. Proposition 2.6.5 in [15] states that there is a unique unitary operator, $U : (H_1 \otimes H_2) \otimes H_3 \to H_1 \otimes (H_2 \otimes H_3)$, that satisfies the following for all $h_1 \in H_1, h_2 \in H_2, h_3 \in H_3$,

\[(9) \quad U ((h_1 \otimes h_2) \otimes h_3) = h_1 \otimes (h_2 \otimes h_3).\]

This implies that for any collection of Hilbert spaces, $H_1, \ldots, H_n$, the Hilbert space $H_1 \otimes \cdots \otimes H_n$ is defined unambiguously regardless of how we decide to associate the products. In the space $H_1 \otimes \cdots \otimes H_n$ we define a simple tensor as a vector of the form $h_1 \otimes \cdots \otimes h_n$ with $h_i \in H_i$. In [15] it is shown that $H_1 \otimes \cdots \otimes H_n$ is the closure of the span of these simple tensors. To conclude this primer on tensor products we introduce the following notation. For a Hilbert space $H$ we denote $H \otimes^n = H \otimes H \otimes \cdots \otimes H$ and for $h \in H$,

\[h \otimes^n = h \otimes h \otimes \cdots \otimes h.\]

5.2. Tensor Rank. A tool we will use frequently in our proofs is tensor rank, which is similar to matrix rank.

**Definition 5.1.** Let $h \in H \otimes^n$ where $H$ is a Hilbert space. The rank of $h$ is the smallest natural number $r$ such that $h = \sum_{i=1}^{r} h_i$ where $h_i$ are simple tensors.

In an infinite dimensional Hilbert space it is possible for a tensor to have infinite rank. We will only be concerned with finite rank tensors.

5.3. Some Results for Tensor Product Spaces. We present some technical results concerning tensor product spaces that will be useful for the rest of the paper. These lemmas are similar to or are straightforward extensions of previous results which we needed to modify for our particular purposes. The following lemma is used in the proof of Lemma 5.2 (Supplemental Material) and Lemma 6.5.

**Lemma 5.1.** Let $H_1, \ldots, H_n, H'_1, \ldots, H'_n$ be a collection of Hilbert spaces and $U_1, \ldots, U_n$ a collection of unitary operators with $U_i : H_i \to H'_i$ for all $i$. There exists a unitary operator $U : H_1 \otimes \cdots \otimes H_n \to H'_1 \otimes \cdots \otimes H'_n$ satisfying $U (h_1 \otimes \cdots \otimes h_n) = U_1(h_1) \otimes \cdots \otimes U_n(h_n)$ for all $h_1 \in H_1, \ldots, h_n \in H_n$.

Let $(\Psi, \mathcal{G}, \gamma)$ be a $\sigma$-finite measure space. We have the following lemma that connects tensor power of a $L^2$ space to the $L^2$ space of the product measure.
Lemma 5.2. There exists a unitary transform $U : L^2(\Psi, \mathcal{G}, \gamma) \rightarrow L^2(\Psi \times \mathcal{G}^n, \gamma \times n)$ such that, for all $f_1, \ldots, f_n \in L^2(\Psi, \mathcal{G}, \gamma)$,

\begin{equation}
U(f_1 \otimes \cdots \otimes f_n) = f_1(\cdot) \cdots f_n(\cdot).
\end{equation}

A statement of the following lemma for $\mathbb{R}^d$ can be found in [10]. We present our own proof for the Hilbert space setting in Supplemental Material.

Lemma 5.3. Let $n > 1$ and let $h_1, \ldots, h_n$ be elements of a Hilbert space such that no elements are zero and no pairs of elements are collinear. Then $h_1 \otimes n - 1, \ldots, h_n \otimes n - 1$ are linearly independent.

The following lemma is a Hilbert space version of a well known property for positive semi-definite matrices.

Lemma 5.4. Let $h_1, \ldots, h_m$ be elements of a Hilbert space. The rank of $\sum_{i=1}^m h_i \otimes 2$ is the dimension of $\text{span} \{h_1, \ldots, h_m\}$.

6. Proofs of Theorems. With the tools developed in the previous sections we will now prove a few, selected, theorems. Due to space constraints we only prove Theorems 4.1 to 4.4 in this document. These proofs give a good overview of the general techniques used to prove the other identifiability and determinedness results, which can be found in Supplemental Material. These theorems are proved for general measure spaces which introduces a fair amount of mathematical overhead. There are two basic components to these proofs, transforming the measure problem into and out of a tensor framework, and using the tensor framework as a means to manipulate these objects geometrically. For concreteness it can be helpful to consider the situation where $\Omega = \{1, 2, \ldots, d\}$, i.e. a finite sample space. In this situation a mixture component $\mu$ can be directly associated with a probability vector $p \in \mathbb{R}^d$ where $[p]_i = \mu(\{i\})$ and the tensor $p^\otimes m$ represents the density of $m$ iid samples of $\mu$:

\begin{align}
[p^\otimes m]_{i_1, \ldots, i_m} &= [p]_{i_1} [p]_{i_2} \cdots [p]_{i_m} \\
&= \mu(\{i_1\}) \mu(\{i_2\}) \cdots \mu(\{i_m\}) \\
&= \mu^\times m (\{(i_1, i_2, \ldots, i_m)\}).
\end{align}

In essence all of our results directly parallel this setting once transformed into the tensor space.

Before we begin our proofs we need to introduce one additional piece of notation. For a function $f$ on a domain $\mathcal{X}$ we define $f^{\times k}$ as simply the
product of the function \( k \) times on the domain \( \mathcal{X}^\times k \),

\[
f^\times k = f(\cdot) \cdots f(\cdot).\]

For a set, \( \sigma \)-algebra, or measure the notation continues to denote the standard \( k \)-fold product.

In these proofs we will be making extensive use of various \( L^2 \) spaces. These spaces will be equivalence classes of functions which are equal almost everywhere with respect to the measure associated with that space. When considering elements of these spaces, equality will always mean almost everywhere equality with respect to the measure associated with that space. When performing integrals or other manipulations of elements in \( L^2 \) spaces, we will be performing operations that do not depend on the representative of the equivalence class. The following lemma will be quite useful.

**Lemma 6.1.** Let \( \gamma_1, \ldots, \gamma_m, \pi_1, \ldots, \pi_l \) be probability measures on a measurable space \( (\Psi, \mathcal{G}) \), \( a_1, \ldots, a_m, b_1, \ldots, b_l \in \mathbb{R} \), and \( n \in \mathbb{N}_+ \). If

\[
\sum_{i=1}^{m} a_i \gamma_i^\times n = \sum_{j=1}^{l} b_j \pi_j^\times n
\]

then for all \( n' \in \mathbb{N}_+ \) with \( n' \leq n \) we have that

\[
\sum_{i=1}^{m} a_i \gamma_i^\times n' = \sum_{j=1}^{l} b_j \pi_j^\times n'.
\]

**Proof of Theorem 4.1.** We proceed by contradiction. Suppose there exist \( m, l \in \mathbb{N}_+ \) with \( l \leq m \) such that there two different mixtures of measures \( \mathcal{P} = \sum_{i=1}^{m} a_i \delta_{\mu_i} \neq \mathcal{P}' = \sum_{j=1}^{l} b_j \delta_{\nu_j} \), and

\[
\sum_{i=1}^{m} a_i \mu_i^\times 2m-1 = \sum_{j=1}^{l} b_j \nu_j^\times 2m-1.
\]

Clearly \( m > 1 \) otherwise we immediately arrive at a contradiction. By the well-ordering principle there exists a minimal \( m \) such that the previous statement holds. For that minimal \( m \) there exists a minimal \( l \) such that the previous statement holds. We will assume that the \( m \) and \( l \) are both minimal in this way. This assumption implies that \( \mu_i \neq \nu_j \) for all \( i, j \). To prove this we will assume that there exists \( i, j \) such that \( \mu_i = \nu_j \), and show that this
assumption leads to a contradiction. Without loss of generality we will assume that $\mu_m = \nu_l$. We will consider the three cases where $a_m = b_l$, $a_m > b_l$, and $a_m < b_l$.

**Case 1.** If $a_m = b_l$ then we have that

$$
\sum_{i=1}^{m-1} \frac{a_i}{1-a_m} \mu_i^{x2m-1} = \sum_{j=1}^{l-1} \frac{b_j}{1-b_l} \nu_j^{x2m-1}
$$

and from Lemma 6.1 we have

$$
\sum_{i=1}^{m-1} \frac{a_i}{1-a_m} \mu_i^{x2(m-1)-1} = \sum_{j=1}^{l-1} \frac{b_j}{1-b_l} \nu_j^{x2(m-1)-1}.
$$

Setting $P = \sum_{i=1}^{m-1} \frac{a_i}{1-a_m} \delta_{\mu_i}$ and $P' = \sum_{j=1}^{l-1} \frac{b_j}{1-b_l} \delta_{\nu_j}$, we have that $V_2^{x2(m-1)-1}(P) = V_2^{x2(m-1)-1}(P')$ which contradicts the minimality of $m$.

**Case 2.** If $a_m > b_l$ then we have

$$
\sum_{i=1}^{m-1} \frac{a_i}{1-b_l} \mu_i^{x2m-1} + \frac{a_m - b_l}{1-b_l} \mu_m^{x2m-1} = \sum_{j=1}^{l-1} \frac{b_j}{1-b_l} \nu_j^{x2m-1}
$$

which contradicts the minimality of $l$ by an argument similar to that in Case 1.

**Case 3.** If $a_m < b_l$ we have that

$$
\sum_{i=1}^{m-1} \frac{a_i}{1-a_m} \mu_i^{x2m-1} = \sum_{j=1}^{l-1} \frac{b_j}{1-a_m} \nu_j^{x2m-1} + \frac{b_l - a_m}{1-a_m} \nu_l^{x2m-1}.
$$

Again we will use arguments similar to the one used in Case 1. If $l = m$ then swapping the mixtures associated with $m$ and $l$ gives us a pair of mixtures of measures which violates the minimality of $l$. If $l < m$ then from Lemma 6.1 we have that

$$
\sum_{i=1}^{m-1} \frac{a_i}{1-a_m} \mu_i^{x2(m-1)-1} = \sum_{j=1}^{l-1} \frac{b_j}{1-a_m} \nu_j^{x2(m-1)-1} + \frac{b_l - a_m}{1-a_m} \nu_l^{x2(m-1)-1},
$$

which violates the minimality of $m$, thus completing Case 3.

We have now established that $\mu_i \neq \nu_j$, for all $i,j$. We will use the following lemma to embed the mixture components in a Hilbert space.

- **Lemma.**
LEMMA 6.2. Let $\gamma_1, \ldots, \gamma_n$ be finite measures on a measurable space $(\Psi, \mathcal{G})$. There exists a finite measure $\pi$ and non-negative functions $f_1, \ldots, f_n \in L^1(\Psi, \mathcal{G}, \pi) \cap L^2(\Psi, \mathcal{G}, \pi)$ such that, for all $i$ and all $B \in \mathcal{G}$

$$\gamma_i(B) = \int_B f_i d\pi.$$ (23)

From Lemma 6.2 there exists a finite measure $\xi$ and non-negative functions $p_1, \ldots, p_m, q_1, \ldots, q_l \in L^1(\Omega, \mathcal{F}, \xi) \cap L^2(\Omega, \mathcal{F}, \xi)$ such that, for all $B \in \mathcal{F}$, $\mu_i(B) = \int_B p_i d\xi$ and $\nu_j(B) = \int_B q_j d\xi$ for all $i, j$. Clearly no two of these functions are equal (in the $\xi$-almost everywhere sense). If one of the functions were a scalar multiple of another, for example $p_1 = \alpha p_2$ for some $\alpha \neq 1$, it would imply

$$\mu_1(\Omega) = \int p_1 d\xi = \int \alpha p_2 d\xi = \alpha.$$ (24)

This is not true so no pair of these functions are collinear.

We can use the following lemma to extend this new representation to a product measure.

LEMMA 6.3. Let $(\Psi, \mathcal{G})$ be a measurable space, $\gamma$ and $\pi$ a pair of finite measures on that space, and $f$ a non-negative function in $L^1(\Psi, \mathcal{G}, \pi)$ such that, for all $A \in \mathcal{G}$, $\gamma(A) = \int_A f d\pi$. Then for all $n$, for all $B \in \mathcal{G}^n$ we have

$$\gamma^\times n(B) = \int_B f^\times n d\pi^\times n.$$ (25)

Thus for any $R \in \mathcal{F}^\times 2m-1$ we have

$$\int_R \sum_{i=1}^m a_i p_i^\times 2m-1 d\xi^\times 2m-1 = \sum_{i=1}^m a_i \mu_i^\times 2m-1(R)$$ (26)

$$= \sum_{j=1}^l b_j \nu_j^\times 2m-1(R)$$ (27)

$$= \int_R \sum_{j=1}^l b_j q_j^\times 2m-1 d\xi^\times 2m-1.$$ (28)

The following lemma is a well known result in real analysis (Proposition 2.23 in [13]), but it is worth mentioning explicitly.

LEMMA 6.4. Let $(\Psi, \mathcal{G}, \gamma)$ be a measure space and $f, g \in L^1(\Psi, \mathcal{G}, \gamma)$. Then $f = g$ $\gamma$-almost everywhere iff, for all $A \in \mathcal{G}$, $\int_A f d\gamma = \int_A g d\gamma$. 

From this lemma it follows that

\[ \sum_{i=1}^{m} a_i p_i^{2m-1} = \sum_{j=1}^{l} b_j q_j^{2m-1}. \]  

(29)

Applying the \( U^{-1} \) operator from Lemma 5.2 to the previous equation yields

\[ \sum_{i=1}^{m} a_i p_i^{\otimes 2m-1} = \sum_{j=1}^{l} b_j q_j^{\otimes 2m-1}. \]  

(30)

Since \( l + m \leq 2m \) Lemma 5.3 states that

\[ p_1^{\otimes 2m-1}, \ldots, p_m^{\otimes 2m-1}, q_1^{\otimes 2m-1}, \ldots, q_l^{\otimes 2m-1} \]

are all linearly independent and thus \( a_i = 0 \) and \( b_j = 0 \) for all \( i, j \), a contradiction.

**Proof of Theorem 4.2.** To prove this theorem we will construct a pair of mixture of measures, \( \mathcal{P} \neq \mathcal{P}' \) which both contain \( m \) components and satisfy \( V_{2m-2}(\mathcal{P}) = V_{2m-2}(\mathcal{P}') \). From our definition of \((\Omega, \mathcal{F})\) we know there exists \( F \in \mathcal{F} \) such that \( F \) and \( F^C \) are nonempty. Let \( x \in F \) and \( x' \in F^C \). It follows that \( \delta_x \) and \( \delta_{x'} \) are different probability measures on \((\Omega, \mathcal{F})\). The theorem follows from the next lemma. We will prove the lemma after the theorem proof.

**Lemma 6.5.** Let \((\Psi, \mathcal{G})\) be a measurable space and \( \gamma, \gamma' \) be distinct probability measures on that space. Let \( \varepsilon_1, \ldots, \varepsilon_t \) be \( t \geq 3 \) distinct values in \([0, 1]\). Then there exist \( \beta_1, \ldots, \beta_t \), a permutation \( \sigma : [t] \rightarrow [t] \), and \( l \in \mathbb{N}_+ \) such that

\[ \sum_{i=1}^{l} \beta_i \left( \varepsilon_{\sigma(i)} \gamma + (1 - \varepsilon_{\sigma(i)}) \gamma' \right)^{t-2} = \sum_{j=l+1}^{t} \beta_j \left( \varepsilon_{\sigma(j)} \gamma + (1 - \varepsilon_{\sigma(j)}) \gamma' \right)^{t-2} \]

(32)

where \( \beta_i > 0 \) for all \( i \), \( \sum_{i=1}^{l} \beta_i = \sum_{j=l+1}^{t} \beta_j = 1 \), and \( l, t - l \geq \left\lceil \frac{t}{2} \right\rceil \).

Let \( \varepsilon_1, \ldots, \varepsilon_{2m} \in [0, 1] \) be distinct and let \( \mu_i = \varepsilon_i \delta_x + (1 - \varepsilon_i) \delta_{x'} \) for \( i \in [2m] \). From Lemma 6.5 with \( t = 2m \) there exists a permutation \( \sigma : [2m] \rightarrow [2m] \) and \( \beta_1, \ldots, \beta_{2m} \) such that

\[ \sum_{i=1}^{m} \beta_i \mu_{\sigma(i)}^{2m-2} = \sum_{j=m+1}^{2m} \beta_j \mu_{\sigma(j)}^{2m-2}, \]

(33)
Lemma 5.1 there exists a unitary operator $U$ spanned by $f$. $U$ is isomorphic to $\mathbb{R}$ if $r$ follows that $\mathbb{R}$.

An order $(34)$ for all $\mathbb{R}$ and non-negative functions $\pi_i \in \mathbb{R}$.

Let $\mathbb{R}$ and $\mathbb{C}$ be the real component when applied to an element of $\mathbb{C}$, and the

Proof of Lemma 6.5. From Lemma 6.2, there exists a finite measure $\pi$ and non-negative functions $f, f' \in L^1(\Psi, \mathcal{G}, \pi) \cap L^2(\Psi, \mathcal{G}, \pi)$ such that, for all $A \in \mathcal{G}$, $\gamma(A) = \int_A f d\pi$ and $\gamma'(A) = \int_A f' d\pi$.

Let $H_2$ be the Hilbert space associated with the subspace in $L^2(\Psi, \mathcal{G}, \pi)$ spanned by $f$ and $f'$. Let $(f_i)_{i=1}^m$ be non-negative functions in $L^1(\Psi, \mathcal{G}, \pi) \cap L^2(\Psi, \mathcal{G}, \pi)$ with $f_i = \varepsilon_i f + (1 - \varepsilon_i) f'$. Clearly $f_i$ is a pdf over $\pi$ for all $i$ and there are no pair in this collection which are collinear. Since $H_2$ is isomorphic to $\mathbb{R}^2$ there exists a unitary operator $U : H_2 \rightarrow \mathbb{R}^2$. From Lemma 5.1 there exists a unitary operator $U_{t-2} : H_{2}^{\otimes t-2} \rightarrow \mathbb{R}^{2^{t-2}}$, with $U_{t-2}(h_1 \otimes \cdots \otimes h_{t-2}) = U(h_1) \otimes \cdots \otimes U(h_{t-2})$. Because $U$ is unitary it follows that

\begin{equation}
U_{t-2}\left(\text{span}\left\{h^{\otimes t-2} : h \in H_2\right\}\right) = \text{span}\left\{x^{\otimes t-2} : x \in \mathbb{R}^2\right\}.
\end{equation}

An order $r$ tensor, $A_{i_1, \ldots, i_r}$, is symmetric if $A_{i_1, \ldots, i_r} = A_{\psi(i_1), \ldots, \psi(i_r)}$ for any $i_1, \ldots, i_r$ and permutation $\psi : [r] \rightarrow [r]$. A consequence of Lemma 4.2 in [10] is that span $\left\{x^{\otimes t-2} : x \in \mathbb{R}^2\right\} \subset S^{t-2}(\mathbb{C}^2)$, the space of all symmetric order $t-2$ tensors over $\mathbb{C}^2$. Complex symmetric tensor spaces will always be viewed as a vector space over the complex numbers and real symmetric tensor spaces will be always be viewed as a vector space over the real numbers.

From Proposition 3.4 in [10] it follows that the dimension of $S^{t-2}(\mathbb{C}^2)$ is $\binom{2 + t - 2 - 1}{t - 2} = t - 1$. From this it follows that $\dim S^{t-2}(\mathbb{R}^2) \leq t - 1$, where $S^{t-2}(\mathbb{R}^2)$ is the space of all symmetric order $t-2$ tensors over $\mathbb{R}^2$.

To see this consider some set of linearly dependent tensors $x_1, \ldots, x_r \in S^{t-2}(\mathbb{C}^2)$ each containing only real valued entries, i.e. the tensors are in $S^{t-2}(\mathbb{R}^2)$. Then it follows that there exists $c_1, \ldots, c_r \in \mathbb{C}$ such that

\begin{equation}
\sum_{i=1}^r c_i x_i = 0.
\end{equation}

Let $\Re$ denote the real component when applied to an element of $\mathbb{C}$, and the
real component applied entrywise when applied to a tensor. We have that

\[ 0 = \Re \left( \sum_{i=1}^{r} c_i x_i \right) = \sum_{i=1}^{r} \Re (c_i x_i) = \sum_{i=1}^{r} \Re (c_i) x_i. \]  

Thus it follows that \( x_1, \ldots, x_r \) are linearly dependent in \( S^{t-2} (\mathbb{R}^2) \) and thus the dimensionality bound holds, \( \dim S^{t-2} (\mathbb{R}^2) \leq t - 1 \).

From this we get that

\[ \dim \left( \text{span} \left( \{ h^{\otimes t-2} : h \in H_2 \} \right) \right) \leq t - 1. \]  

(37)

The bound on the dimension of \( \text{span} \left( \{ h^{\otimes t-2} : h \in H_2 \} \right) \) implies that \( (f_i^{\otimes t-2})_{i=1}^{t} \) are linearly dependent. Conversely Lemma 5.3 implies that removing a single vector from \( (f_i^{\otimes t-2})_{i=1}^{t} \) yields a set of vectors which are linearly independent. It follows that there exists \( (\alpha_i)_{i=1}^{t} \) with \( \alpha_i \neq 0 \) for all \( i \) and

\[ \sum_{i=1}^{t} \alpha_i f_i^{\otimes t-2} = 0. \]  

(38)

There exists a permutation \( \sigma : [t] \to [t] \) such that \( \alpha_{\sigma(i)} < 0 \) for all \( i \in [l] \) and \( \alpha_{\sigma(j)} > 0 \) for all \( j > l \) with \( l \leq \left\lfloor \frac{t}{2} \right\rfloor \) (ensuring that \( l \leq \left\lfloor \frac{t}{2} \right\rfloor \) may also require multiplying (38) by \(-1\)). This \( \sigma \) appears in the lemma statement, but for the remainder of the proof we will simply assume without loss of generality that \( \alpha_i < 0 \) for \( i \in [l] \) with \( l \leq \left\lfloor \frac{t}{2} \right\rfloor \).

From this we have

\[ \sum_{i=1}^{t} -\alpha_i f_i^{\otimes t-2} = \sum_{j=l+1}^{t} \alpha_j f_j^{\otimes t-2}. \]  

(39)

From Lemma 5.2 we have

\[ \sum_{i=1}^{l} -\alpha_i f_i^{\times t-2} = \sum_{j=l+1}^{t} \alpha_j f_j^{\times t-2} \]  

(40)

and thus

\[ \int \sum_{i=1}^{l} -\alpha_i f_i^{\times t-2} d\pi^{\times t-2} = \int \sum_{j=l+1}^{t} \alpha_j f_j^{\times t-2} d\pi^{\times t-2} \]  

(41)

\[ \Rightarrow \sum_{i=1}^{l} -\alpha_i = \sum_{j=l+1}^{t} \alpha_j. \]  

(42)
Let $r = \sum_{i=1}^{l} -\alpha_i$. We know $r > 0$ so dividing both sides of (39) by $r$ gives us

$$
\sum_{i=1}^{l} -\frac{\alpha_i}{r} f_i^\otimes t - 2 = \sum_{j=l+1}^{t} \frac{\alpha_j}{r} f_j^\otimes t - 2
$$

where the left and the right side are convex combinations. Let $(\beta_i)_{i=1}^{l}$ be positive numbers with $\beta_i = -\frac{\alpha_i}{r}$ for $i \in [l]$ and $\beta_j = \frac{\alpha_j}{r}$ for $j \in [t] \setminus [l]$. This gives us

$$
\sum_{i=1}^{l} \beta_i f_i^\otimes t - 2 = \sum_{j=l+1}^{t} \beta_j f_j^\otimes t - 2.
$$

We will now consider 3 cases for the value of $t$.

**Case 1.** If $t = 3$ then $l = 1$ and $l, t - l \geq \lfloor \frac{t}{2} \rfloor$ is satisfied.

**Case 2.** If $t$ is divisible by two then we can do the following,

$$
\sum_{i=1}^{l} \beta_i f_i^\otimes \frac{t}{2} - 1 \otimes f_i^\otimes \frac{t}{2} - 1 = \sum_{j=l+1}^{t} \beta_j f_j^\otimes \frac{t}{2} - 1 \otimes f_j^\otimes \frac{t}{2} - 1.
$$

Consider the elements in the last equation as order two tensors in $L^2(\Psi, G, \pi) \otimes^\frac{t}{2} - 1 \otimes L^2(\Psi, G, \pi) \otimes^\frac{t}{2} - 1$. From Lemma 5.3 and Lemma 5.4 we have that the RHS of the previous equation has rank at least $\lfloor \frac{t}{2} \rfloor$ and since $l \leq \frac{t}{2}$ it follows that $l = \frac{t}{2}$. Again we have that $l, t - l \geq \lfloor \frac{t}{2} \rfloor$.

**Case 3.** If $t$ is greater than 3 and not divisible by 2 then we can apply Lemma 5.2 to get

$$
\int_{\Psi} \sum_{i=1}^{l} \beta_i f_i^\times t - 3 f_i(x) d\pi(x) = \int_{\Psi} \sum_{j=l+1}^{t} \beta_j f_j^\times t - 3 f_j(y) d\pi(y)
$$

$$
\Rightarrow \sum_{i=1}^{l} \beta_i f_i^\times t - 3 = \sum_{j=l+1}^{t} \beta_j f_j^\times t - 3.
$$

Applying Lemma 5.2 again we get

$$
\sum_{i=1}^{l} \beta_i f_i^\otimes t - 3 = \sum_{j=l+1}^{t} \beta_j f_j^\otimes t - 3
$$

$$
\Rightarrow \sum_{i=1}^{l} \beta_i f_i^\otimes \frac{t}{2} - 1 \otimes f_i^\otimes \frac{t}{2} - 1 = \sum_{j=l+1}^{t} \beta_j f_j^\otimes \frac{t}{2} - 1 \otimes f_j^\otimes \frac{t}{2} - 1.
$$
Recall that \(\left\lfloor \frac{t}{2} \right\rfloor \geq l\) so we also have that

\[
\left\lfloor \frac{t}{2} \right\rfloor - l \geq 0
\]

\[
\Rightarrow \frac{t}{2} - l \geq \frac{1}{2}
\]

\[
\Rightarrow t - l \geq \frac{t - 1}{2}.
\]

From Lemma 5.3 and Lemma 5.4 we have that the RHS of (49) has rank at least \(t - \frac{1}{2}\) and thus \(l \geq t - \frac{1}{2}\). From this we have that \(t - l, l \geq \left\lfloor \frac{t}{2} \right\rfloor\) once again, which finishes Case 3.

So \(l, t - l \geq \left\lfloor \frac{t}{2} \right\rfloor\) for any \(t \geq 3\). Applying Lemma 5.2 to (44) we have

\[
\sum_{i=1}^{t} \beta_i f_i^{x_{t-2}} = \sum_{j=l+1}^{t} \beta_j f_j^{x_{t-2}}.
\]

From Lemma 6.3 we have

\[
\sum_{i=1}^{t} \beta_i (\varepsilon_i \gamma + (1 - \varepsilon_i) \gamma')^{x_{t-2}} = \sum_{j=l+1}^{t} \beta_j (\varepsilon_j \gamma + (1 - \varepsilon_j) \gamma')^{x_{t-2}}.
\]

**Proof of Theorem 4.3.** Let \(\mathcal{P} = \sum_{i=1}^{m} a_i \delta_{\mu_i}\) and \(\mathcal{P}' = \sum_{j=1}^{l} b_j \delta_{\nu_j}\) be mixtures of measures such that \(\mathcal{P}' \neq \mathcal{P}\). We will proceed by contradiction. Suppose that \(\sum_{i=1}^{m} a_i \mu_i^{x_{2m}} = \sum_{j=1}^{l} b_j \nu_j^{x_{2m}}\). From Theorem 4.1 we know that \(\mathcal{P}\) is \(2m - 1\)-identifiable and therefore \(2m\)-identifiable by Lemma 4.1. It follows that \(l > m\). From Lemma 6.2 there exists a finite measure \(\xi\) and non-negative functions \(p_1, \ldots, p_m, q_1, \ldots, q_l \in L^1(\Omega, \mathcal{F}, \xi)\cap L^2(\Omega, \mathcal{F}, \xi)\) such that, for all \(B \in \mathcal{F}\), \(\mu_i(B) = \int_B p_i d\xi\) and \(\nu_j(B) = \int_B q_j d\xi\) for all \(i, j\). Using Lemmas 6.3 and 6.4 we have

\[
\sum_{i=1}^{m} a_i p_i^{x_{2m}} = \sum_{j=1}^{l} b_j q_j^{x_{2m}}.
\]

By Lemma 5.2 we have

\[
\sum_{i=1}^{m} a_i p_i^{x_{2m}} = \sum_{j=1}^{l} b_j q_j^{x_{2m}},
\]
and therefore
\begin{equation}
\sum_{i=1}^{m} a_i p_i^{\otimes m} \otimes p_i^{\otimes m} = \sum_{j=1}^{l} b_j q_j^{\otimes m} \otimes q_j^{\otimes m}.
\end{equation}

Consider the elements in the last equation as tensors in $L^2(\Omega, \mathcal{F}, \xi)^{\otimes m} \otimes L^2(\Omega, \mathcal{F}, \xi)^{\otimes m}$. Since no pair of vectors in $p_1, \ldots, p_m$ are collinear, from Lemma 5.3 and Lemma 5.4 we know that the LHS has rank $m$. On the other hand, no pair of vectors $q_1, \ldots, q_l$ are collinear either, so Lemma 5.3 says that there is a subset of $\{q_1^{\otimes m}, \ldots, q_l^{\otimes m}\}$ which contains at least $m + 1$ linearly independent elements. By Lemma 5.4 it follows that the RHS has rank at least $m + 1$, a contradiction.

**Proof of Theorem 4.4.** To prove this theorem we will construct a pair of mixture of measures, $\mathcal{P} \neq \mathcal{P}'$ which contain $m$ and $m + 1$ components respectively and satisfy $V_{2m-1}(\mathcal{P}) = V_{2m-1}(\mathcal{P}')$. From our definition of $(\Omega, \mathcal{F})$ we know there exists $F \in \mathcal{F}$ such that $F, F^C$ are nonempty. Let $x \in F$ and $x' \in F^C$. It follows that $\delta_x$ and $\delta_{x'}$ are different probability measures on $(\Omega, \mathcal{F})$. Let $\varepsilon_1, \ldots, \varepsilon_{2m+1}$ be distinct values in $[0, 1]$. Applying Lemma 6.5 with $t = 2m + 1$ and letting $\mu_i = \varepsilon_i \delta_x + (1 - \varepsilon_i) \delta_{x'}$, there exists a permutation $\sigma : [2m + 1] \to [2m + 1]$ and $\beta_1, \ldots, \beta_{2m+1}$, with $\beta_i > 0$ for all $i$ and $\sum_{i=1}^{m} \beta_i = \sum_{j=m+1}^{2m+1} \beta_j = 1$, such that
\begin{equation}
\sum_{i=1}^{m} \beta_i \mu_{\sigma(i)}^{\times 2m-1} = \sum_{j=m+1}^{2m+1} \beta_j \mu_{\sigma(j)}^{\times 2m-1}.
\end{equation}

If we let $\mathcal{P} = \sum_{i=1}^{m} \beta_i \delta_{\mu_{\sigma(i)}}$ and $\mathcal{P}' = \sum_{j=m+1}^{2m+1} \beta_j \delta_{\mu_{\sigma(j)}}$, we have that $V_{2m-1}(\mathcal{P}) = V_{2m-1}(\mathcal{P}')$. □

7. Identifiability and Determinedness of Mixtures of Multinomial Distributions. Using the previous results we can show analogous identifiability and determinedness results for mixtures of multinomial distributions. The identifiability of mixtures of multinomial distributions was originally studied in [17] which contains a proof of Corollary 7.1 from this paper. An alternative proof of this corollary can be found in [12]. These results are analogous to identifiability results presented in this paper. Our proofs (see Supplemental Material) use techniques which are very different from those used in [17, 12]. Our techniques can also be used to prove a determinedness style result, Corollary 7.2, which we have not seen addressed elsewhere in the multinomial mixture model literature.
Central to the results in this section is Lemma 7.1 which establishes an equivalence between the grouped sample setting and multinomial mixture models. A sample from a multinomial distribution can be viewed as totalling the outcomes from an iid sampling of a categorical distribution. Consider some probability measure $\mu$ over a finite sample space and let $X = (X_1, \ldots, X_m)$ be a collection of $m$ iid samples from $\mu$. Here $X$ has the form of what we would call a “random group.” Because $X$ contains iid samples no useful statistical information is contained in the order of the samples. It follows that we can simply tally the number of results for each outcome and not lose any useful statistical information. Lemma 7.1 formalizes this intuition so that we can apply tools developed earlier in this paper to the multinomial mixture model setting.

Before stating our results we must first introduce some definitions and notation. Any multinomial distribution is completely characterized by positive integers $n$ and $q$ and a probability vector in $\mathbb{R}^q$, $p = [p_1, \ldots, p_q]^T$. The value $q$ represents the number of possible outcomes of a trial, $p$ is the likelihood of each outcome on a trial, and $n$ is the number of trials. For whole numbers $k, l$ we define
\[ C_{k,l} = \left\{ x \in \mathbb{N}^l : \sum_{i=1}^l x_i = k \right\} \]
These are vectors of the form $[x_1, \ldots, x_l]^T$ where $\sum_{i=1}^l x_i = k$. Using the values $n$ and $q$ above, the multinomial distribution is a probability measure over $C_{n,q}$. If $Q$ is a multinomial distribution with parameters $n, p, q$ as defined above then its probability mass function is
\[
Q\left(\left\{ [x_1, \ldots, x_q]^T \right\} \right) = \frac{n!}{x_1! \cdots x_q!} p_1^{x_1} \cdots p_q^{x_q}
\]
for $x \in C_{n,q}$. We will denote this measure as $Q_{n,p,q}$. Let
\[
\mathcal{M}(n, q) \triangleq \{ Q_{n,p,q} : p \text{ is a probability vector in } \mathbb{R}^q \},
\]
i.e. the space of all multinomial distributions with $n$ and $q$ fixed.

At the heart of our multinomial mixture model identifiability and determinedness results is the construction of a linear operator $T_{n,q}$ from $\text{span}(\mathcal{D}(C_{n,q}, 2^{C_{n,q}}))$ to $\text{span}(\mathcal{D}([q]^{\times n}, 2^{[q]^{\times n}}))$ and its use to show that non-identifiable mixtures of multinomial distributions yield non-identifiable mixtures of measures and non-determined mixtures of multinomial distributions yields non-determined mixtures of measures.

Since $C_{n,q}$ is a finite set, the vector space of finite signed measures on $(C_{n,q}, 2^{C_{n,q}})$ is a finite dimensional space and the set $\{ \delta_x : x \in C_{n,q} \}$ is a basis for this space. Note that $\{ \delta_x : x \in C_{n,q} \}$ is the set of all point masses on $C_{n,q}$, not vectors in the ambient space of $C_{n,q}$. Thus, to completely define
the operator $T_{n,q}$, we need only define $T_{n,q}(\delta_x)$ for all $x \in C_{n,q}$. To this end let $x \in C_{n,q}$. We define the function $F_{n,q} : C_{n,q} \to [q]^{\times n}$ as $F_{n,q}(x) = 1^{\times x_1} \times \cdots \times q^{\times x_q}$, where the exponents represent Cartesian powers. The definition of $F_{n,q}$ is a bit dense so we will do a simple example. Suppose $n = 6, q = 4$ and $x = [1, 0, 3, 2]^T$ then $F_{n,q}(x) = [1, 3, 3, 3, 4, 4]^T$. Intuitively the $F_{n,q}$ operator undoes the totalling which transforms a collection of trials from a categorical distribution into a draw from a multinomial distribution; $F_{n,q}$ returns these trials in nondecreasing order. Let $S_n$ be the symmetric group on $n$ symbols. We define our linear operator as follows

\begin{equation}
T_{n,q}(\delta_x) = \frac{1}{n!} \sum_{\sigma \in S_n} \delta_{\sigma(F_{n,q}(x))},
\end{equation}

where $\sigma$ is permuting the entries of $F_{n,q}(x)$. This operator is similar to the projection operator onto the set of order $n$ symmetric tensors [10]. The following lemma makes the crucial connection between the space of multinomial distributions and the probability measures of grouped samples.

**Lemma 7.1.** Let $Q_{n,p,q} \in \mathcal{M}(n, q)$, then

\begin{equation}
T_{n,q}(Q_{n,p,q}) = V_n \left( \delta_{\sum_{i=1}^p \delta_i} \right).
\end{equation}

This lemma allows us to make some assertions about the identifiability of mixtures of multinomial distributions.

In the following we will assume that all multinomial mixture models under consideration have only nonzero summands and distinct components. In the context of multinomial mixture models, a multinomial mixture model $\sum_{i=1}^m a_i Q_{n,p_i,q}$ is identifiable if it being equal to a different multinomial mixture model,

\begin{equation}
\sum_{i=1}^m a_i Q_{n,p_i,q} = \sum_{j=1}^s b_j Q_{n,r_j,q},
\end{equation}

with $s \leq m$ implies that $s = m$ and there exists some permutation $\sigma$ such that $a_i = b_{\sigma(i)}$ and $Q_{n,p_i,q} = Q_{n,r_{\sigma(i)},q}$ for all $i$. The mixture model is determined if the previous statement holds without the restriction $s \leq m$.

Multinomial mixture models are identifiable if the number of components $m$ and the number of trials in each component $n$ satisfy $n \geq 2m - 1$.

**Corollary 7.1.** Let $m \in \mathbb{N}_+$, $n \geq 2m - 1$, and fix $q \in \mathbb{N}_+$. Let $Q_{n,p_1,q}, \ldots, Q_{n,p_m,q}, Q_{n,r_1,q}, \ldots, Q_{n,r_s,q} \in \mathcal{M}(n, q)$ with $Q_{n,p_1,q}, \ldots, Q_{n,p_m,q}$
distinct, \(Q_{n,r_1,q}, \ldots, Q_{n,r_s,q}\) distinct, and \(s \leq m\). If

\[
\sum_{i=1}^{m} a_i Q_{n,p_i,q} = \sum_{j=1}^{s} b_j Q_{n,r_j,q}
\]

with \(a_i > 0, b_i > 0\) for all \(i\) and \(\sum_{i=1}^{m} a_i = \sum_{j=1}^{s} b_j = 1\), then \(s = m\) and there exists some permutation \(\sigma\) such that \(a_i = b_{\sigma(i)}\) and \(p_i = r_{\sigma(i)}\).

Alternatively, this corollary says that, given two different finite mixtures with components in \(\mathcal{M}(n,q)\), one mixture with \(m\) components and the other with \(s\) components, if \(n \geq 2m - 1\) and \(n \geq 2s - 1\) then the mixtures induce different measures. Additionally, multinomial mixture models are determined if the number of components \(m\) and the number of trials in each component \(n\) satisfy \(n \geq 2m\).

**Corollary 7.2.** Let \(n \geq 2m\) and fix \(q \in \mathbb{N}\). Let \(Q_{n,p_1,q}, \ldots, Q_{n,p_m,q}\) and \(Q_{n,r_1,q}, \ldots, Q_{n,r_s,q}\) be elements of \(\mathcal{M}(n,q)\) with \(Q_{n,p_1,q}, \ldots, Q_{n,p_m,q}\) distinct and \(Q_{n,r_1,q}, \ldots, Q_{n,r_s,q}\) distinct. If

\[
\sum_{i=1}^{m} a_i Q_{n,p_i,q} = \sum_{j=1}^{s} b_j Q_{n,r_j,q}
\]

with \(a_i > 0, b_i > 0\) for all \(i\) and \(\sum_{i=1}^{m} a_i = \sum_{j=1}^{m} b_j = 1\), then \(m = s\) and there exists some permutation \(\sigma\) such that \(a_i = b_{\sigma(i)}\) and \(p_i = r_{\sigma(i)}\).

Using the proof techniques employed in the proofs of these corollaries (see Supplemental Material) one could establish additional identifiability/determinedness style results for multinomial mixture models along the lines of Theorems 4.5, 4.6, and 4.7. Furthermore, it seems likely that one could use the algorithm described in the next section or from [2, 4, 22] to recover these components, using the transform \(T_{n,q}\).

**8. Algorithm.** Here we present an algorithm for the recovery of mixture components and proportions from data. The algorithm is quite general and can be applied to any measurable space. Supplemental Material contains a detailed description and analysis of the algorithm applied to categorical data, including a consistency proof.

Let \(\sum_{i=1}^{m} w_i \delta_{\mu_i}\) be an arbitrary mixture of measures on some measurable space \((\Omega, \mathcal{F})\) which we are interested in recovering. Let \(p_1, \ldots, p_m\) be square integrable densities with respect to a dominating measure \(\xi\), with \(\int_A p_i d\xi = \mu_i(A)\) for all \(i \in [m]\) and \(A \in \mathcal{F}\). A measure \(\xi\) and densities
$p_1, \ldots, p_m$ satisfying these properties are guaranteed to exist as a consequence of Lemma 6.2.

We will consider the situation where we have $2m - 1$ samples per random group and have access to the tensors $\sum_{i=1}^{m} w_i p_i \otimes 2m - 1$ and $\sum_{i=1}^{m} w_i p_i \otimes 2m - 2$. In a finite sample space, estimating these tensors is equivalent to estimating moment tensors of order $2m - 1$ and $2m - 2$. For measures over $\mathbb{R}^d$ dominated by the Lebesgue measure, one could estimate these tensors using a kernel density estimator in $\mathbb{R}^d(2m - 1)$ and $\mathbb{R}^d(2m - 2)$ using each sample group as a kernel center. We will also assume that $p_1, \ldots, p_m$ have distinct norms. Note that it is still possible to recover the mixture components if they do not have distinct norms. One way to do this is to choose $\xi$ so that the norms are distinct. In Supplemental Material we describe a method which is guaranteed to do this when $(\Omega, \mathcal{F})$ is a finite sample space by choosing $\xi$ randomly. We term this method “random dominating measure” in the experiments and supplement. Alternatively, if one is capable of choosing an element in span ($\{p_1, \ldots, p_m\}$) in an appropriate random way, one could recover the mixture components using a variation of Jenrich’s algorithm.

To describe the algorithm we will need to make use of bounded linear operators on Hilbert spaces. Given a pair of Hilbert spaces $H, H'$ we define $\mathcal{L}(H, H')$ as the space of bounded linear operators from $H$ to $H'$ and $\mathcal{L}(H) \triangleq \mathcal{L}(H, H)$. An operator, $T$, is in this space if there exists a nonnegative number $C$ such that $\|Tx\|_{H'} \leq C \|x\|_H$ for all $x \in H$. The space of bounded linear operators is a Banach space when equipped with the norm

$$\|T\| \triangleq \sup_{x \neq 0} \frac{\|Tx\|}{\|x\|_H}.$$  \hspace{1cm} (66)$$

In addition we will need to make use of tensor products of bounded linear operators. The following lemma is exactly Proposition 2.6.12 from [15].

**Lemma 8.1.** Let $H_1, \ldots, H_n, H'_1, \ldots, H'_n$ be Hilbert spaces and let $U_i \in \mathcal{L}(H_i, H'_i)$ for all $i \in [n]$. There exists a unique

$$U \in \mathcal{L}(H_1 \otimes \cdots \otimes H_n, H'_1 \otimes \cdots \otimes H'_n),$$

such that $U (h_1 \otimes \cdots \otimes h_n) = U_1 (h_1) \otimes \cdots \otimes U_n (h_n)$ for all $h_1 \in H_1, \ldots, h_n \in H_n$.

**Definition 8.1.** The operator constructed in Lemma 8.1 is called the tensor product of $U_1, \ldots, U_n$ and is denoted $U_1 \otimes \cdots \otimes U_n$. Finally we will need to employ Hilbert-Schmidt operators which are a subspace of the bounded linear operators.
Definition 8.2. Let $H, H'$ be Hilbert spaces and $T \in \mathcal{L}(H, H')$. $T$ is called a Hilbert-Schmidt operator if $\sum_{x \in J} \|Tx\|^2 < \infty$ for an orthonormal basis $J \subset H$. We denote the set of Hilbert-Schmidt operators in $\mathcal{L}(H, H')$ by $\mathcal{HS}(H, H')$.

This definition does not depend on the choice of orthonormal basis: the sum $\sum_{x \in J} \|T(x)\|^2$ will always yield the same value regardless of the choice of orthonormal basis $J$.

Finally we will also need to utilize the equivalence between tensor products and linear operators ([15] Proposition 2.6.9).

Lemma 8.2. Let $H, H'$ be Hilbert spaces. There exists a unitary operator $U : H \otimes H' \to \mathcal{HS}(H, H')$ such that, for any simple tensor $h \otimes h' \in H \otimes H'$, $U(h \otimes h') = \langle h, \cdot \rangle h'$.

Before we introduce the algorithm we will discuss an important point regarding computational implementation and Lemmas 8.2 and 8.1. For the remainder of this paragraph we will assume that Euclidean spaces are equipped with the standard inner product. Vectors in a space of tensor products of Euclidean spaces, for example $\mathbb{R}^{d_1} \otimes \cdots \otimes \mathbb{R}^{d_s}$ are easily represented on computers as elements of $\mathbb{R}^{d_1 \times \cdots \times d_s}$ [10]. Linear operators from some Euclidean tensor space to another can also be easily represented. Furthermore the transformation in Lemma 8.2 and the construction of new operators from Lemma 8.1 can be implemented in computers by “unfolding” the tensors into matrices, applying common linear algebraic manipulations and “folding” them back into tensors. The inner workings of these manipulations are beyond the scope of this paper and we refer the reader to [14] for details. Practically speaking this means the manipulations mentioned in Lemmas 8.2 and 8.1 are straightforward to implement with a bit of tensor programming know-how. Implementation may also be streamlined by using programming libraries that assist with these tensor manipulations such as the NumPy library for Python.

Because of the points mentioned in the previous paragraph, the following algorithm is readily implementable for estimating categorical distributions, where the measures can be represented as probability vectors on a Euclidean space. Similarly, we expect that these techniques could be extended to probability densities on Euclidean space using kernel density estimators with a kernel function with easily computable $L^2$ inner products (for example Gaussian kernels) although we suspect that implementation of such an algorithm may be significantly more involved.
To begin our description of the abstract algorithm we will apply the transform from Lemma 8.2 to $\sum_{i=1}^{m} w_i p_i^{\otimes m-2}$ to get the operator

$$C = \sum_{i=1}^{m} w_i p_i^{\otimes m-1} \langle p_i^{\otimes m-1} , \cdot \rangle = \sum_{i=1}^{m} \sqrt{w_i} p_i^{\otimes m-1} \langle \sqrt{w_i} p_i^{\otimes m-1} , \cdot \rangle.$$  

(68)

Here $C$ is a positive semi-definite (PSD) operator in $L \left( L^2 (\Omega, \mathcal{F}, \xi) \otimes^{m-1} \right)$. Let $C^\dagger$ be the (Moore-Penrose) pseudoinverse of $C$ and $W = \sqrt{C^\dagger}$. Now $W$ is an operator that whitens $\sqrt{w_1} p_1^{\otimes m-1}, \ldots, \sqrt{w_m} p_m^{\otimes m-1}$. That is, $W\sqrt{w_1} p_1^{\otimes m-1}, \ldots, W\sqrt{w_m} p_m^{\otimes m-1}$ are orthonormal vectors. Using the operator construction from Lemma 8.1 we can construct $I \otimes W \otimes W$ where, for all simple tensors in $L^2 (\Omega, \mathcal{F}, \xi) \otimes^{2m-1}$ we have,

$$\left( I \otimes W \otimes W \right) (x_1 \otimes \cdots \otimes x_{2m-1}) = x_1 \otimes W (x_2 \otimes \cdots \otimes x_m) \otimes W (x_{m+1} \otimes \cdots \otimes x_{2m-1}).$$  

(69)

Applying $I \otimes W \otimes W$ to $\sum_{i=1}^{m} w_i p_i^{\otimes m-1}$ yields

$$A \triangleq \sum_{i=1}^{m} w_i p_i \otimes W (p_i^{\otimes m-1}) \otimes W (p_i^{\otimes m-1})$$  

(70)

$$= \sum_{i=1}^{m} p_i \otimes W (\sqrt{w_i} p_i^{\otimes m-1}) \otimes W (\sqrt{w_i} p_i^{\otimes m-1}).$$  

(71)

From Lemma 8.2 we can transform the tensor $A$ into the operator $T$,

$$T = \sum_{i=1}^{m} p_i \otimes W (\sqrt{w_i} p_i^{\otimes m-1}) \langle \sqrt{w_i} p_i^{\otimes m-1}, \cdot \rangle.$$  

(72)

Because $W$ is a whitening operator, the operator $TT^H$ is

$$TT^H = \sum_{i=1}^{m} p_i \otimes W (\sqrt{w_i} p_i^{\otimes m-1}) \left\langle W (\sqrt{w_i} p_i^{\otimes m-1}), \cdot \right\rangle \sum_{j=1}^{m} W (\sqrt{w_j} p_j^{\otimes m-1}) \left\langle p_j \otimes W (\sqrt{w_j} p_j^{\otimes m-1}), \cdot \right\rangle$$  

(73)

$$= \sum_{i=1}^{m} p_i \otimes W (\sqrt{w_i} p_i^{\otimes m-1}) \langle p_i \otimes W (\sqrt{w_i} p_i^{\otimes m-1}), \cdot \rangle$$  

(74)

which is a PSD operator. We set $S \triangleq TT^H$. 


If \( p \) is the unique spectral decomposition of \( S \), and apply this operator to some vector \( W \). Also note that

\[
\begin{align*}
(75) & \quad \langle p_i \otimes W \sqrt{w_i p_i^{\otimes m-1}}, p_j \otimes W \sqrt{w_j p_j^{\otimes m-1}} \rangle \\
(76) & \quad = \langle p_i, p_j \rangle \langle W \sqrt{w_i p_i^{\otimes m-1}}, W \sqrt{w_j p_j^{\otimes m-1}} \rangle \\
(77) & \quad = \langle p_i, p_j \rangle 0 \\
(78) & \quad = 0.
\end{align*}
\]

Also note that

\[
\begin{align*}
(79) & \quad \| p_i \otimes W \sqrt{w_i p_i^{\otimes m-1}} \|^2 = \langle p_i \otimes W \sqrt{w_i p_i^{\otimes m-1}}, p_i \otimes W \sqrt{w_i p_i^{\otimes m-1}} \rangle \\
(80) & \quad = \langle p_i, p_i \rangle \langle W \sqrt{w_i p_i^{\otimes m-1}}, W \sqrt{w_i p_i^{\otimes m-1}} \rangle \\
(81) & \quad = \| p_i \|^2.
\end{align*}
\]

If \( p_1, \ldots, p_m \) have distinct norms then it follows that

\[
(82) \quad \sum_{i=1}^{m} p_i \otimes W \sqrt{w_i p_i^{\otimes m-1}} \langle p_i \otimes W \sqrt{w_i p_i^{\otimes m-1}}, \cdot \rangle
\]

is the unique spectral decomposition of \( S \) since the vectors \( p_1 \otimes W \sqrt{w_1 p_1^{\otimes m-1}}, \ldots, p_m \otimes W \sqrt{w_m p_m^{\otimes m-1}} \) are orthogonal, have distinct norms, and thus distinct positive eigenvalues. Given an eigenvector of \( S \), \( p_i \otimes W \sqrt{w_i p_i^{\otimes m-1}} \), we need only view it as a linear operator \( p_i \langle W \sqrt{w_i p_i^{\otimes m-1}}, \cdot \rangle \) and apply this operator to some vector \( z \) which is not orthogonal to \( W \sqrt{w_i p_i^{\otimes m-1}} \), thus yielding \( p_i \cdot z \). Were the norms of \( p_1, \ldots, p_m \) not distinct, then there would not be a spectral gap between some of the eigenvalues in \( S \), and a spectral decomposition of \( S \) may contain some eigenvectors that are not \( p_1 \otimes W \sqrt{w_1 p_1^{\otimes m-1}}, \ldots, p_m \otimes W \sqrt{w_m p_m^{\otimes m-1}} \), but are instead linear combinations of these vectors.

The following is a concise summary of the main points of the full algorithm:

1. Let \( C = \sum_{i=1}^{m} w_i p_i^{\otimes m-1} \langle p_i^{\otimes m-1}, \cdot \rangle \) by transforming \( \sum_{i=1}^{m} w_i p_i^{\otimes 2m-2} \).
2. Let \( W = \sqrt{C} \).
3. Let \( A = I \otimes W \otimes W \left( \sum_{i=1}^{m} w_i p_i^{\otimes 2m-1} \right) \). Note that

\[
I \otimes W \otimes W \left( \sum_{i=1}^{m} w_i p_i^{\otimes 2m-1} \right) = \sum_{i=1}^{m} p_i \otimes W \left( \sqrt{w_i p_i^{\otimes m-1}} \right) \otimes W \left( \sqrt{w_i p_i^{\otimes m-1}} \right)
\]

by direct evaluation and rearrangement of coefficients.
4. Let \( T = \sum_{i=1}^{m} p_i \otimes W \left( \sqrt{w_i} p_i^{\otimes m-1} \right) \langle W \left( \sqrt{w_i} p_i^{\otimes m-1} \right), \cdot \rangle \) by transforming \( A \).

5. Performing spectral decomposition on \( T T^H \) gives us eigenvectors \( \{ p_i \otimes W \left( \sqrt{w_i} p_i^{\otimes m-1} \right) \}_{i=1}^{m} \), up to scaling.

6. For all \( i \), let \( \tilde{p}_i = p_i \langle W \left( \sqrt{w_i} p_i^{\otimes m-1} \right), z \rangle \) by transforming the eigenvectors into linear operators and selecting \( z \) to be any vector such that the inner product does not evaluate to 0. Now \( \tilde{p}_i \) is a scaled version of \( p_i \).

7. Normalize \( \tilde{p}_i \) to get \( p_i \).

Once the mixture components \( p_1, \ldots, p_m \) are recovered from the spectral decomposition we can calculate the mixture proportions. From these mixture components we can construct the tensors \( p_1^{\otimes 2m-2}, \ldots, p_m^{\otimes 2m-2} \). These tensors are linearly independent by Lemma 5.3. The tensor \( \sum_{i=1}^{m} w_i p_i^{\otimes 2m-2} \) is known. By the linear independence of the components there is exactly one solution for \( a_1, \ldots, a_m \) in the equation

\[
\sum_{i=1}^{m} w_i p_i^{\otimes 2m-2} = \sum_{j=1}^{m} a_j p_j^{\otimes 2m-2},
\]

so simply minimizing \( \| \sum_{i=1}^{m} w_i p_i^{\otimes 2m-2} - \sum_{j=1}^{m} a_j p_j^{\otimes 2m-2} \| \) over \( a_1, \ldots, a_m \) will give us the mixture proportions.

In Supplemental Material we study this algorithm applied to finite sample spaces in further detail. In the supplement we demonstrate how to recover mixture components without the spectral gap assumption, how to construct the estimator given data (which we evaluate experimentally in Section 9), and prove that it is consistent.

Taking inspiration from [2] and [23] we can suggest yet another algorithm. The previous papers demonstrate algorithms for recovering mixture components which are measures on finite sample spaces and \( \mathbb{R}^d \), from random groups of size 3, provided the mixture components are linearly independent. Given a mixture of measures \( \mathcal{P} = \sum_{i=1}^{m} w_i \delta_{\mu_i} \) with density functions \( p_1, \ldots, p_m \), the tensors \( p_1^{\otimes m-1}, \ldots, p_m^{\otimes m-1} \) are linearly independent. Thus, with \( 3m - 3 \) samples per random group, we can estimate the tensors \( \sum_{i=1}^{m} w_i p_i^{\otimes 3m-3} \) and we can use the algorithms from the previous papers to recover \( p_1^{\otimes m-1}, \ldots, p_m^{\otimes m-1} \) from which it is straightforward to recover \( p_1, \ldots, p_m \).

9. Experiments. Here we will present some experimental results of our algorithm applied to a simple synthetic dataset. The sample space for the experiments is \( \Omega = \{0, 1, 2\} \). The mixture components of our dataset
are $\mu_1, \mu_2, \mu_3$ with $\mu_1$ distributed according to a binomial distribution with $n = 2$ and $p = 0.2$, $\mu_2$ is similar with $p = 0.8$ and $\mu_3 = \frac{1}{3}\mu_1 + \frac{2}{3}\mu_2$. The component weights are $w_1 = 0.5, w_2 = 0.3, w_3 = 0.2$. We chose these mixture components so that they are not particularly nice. Specifically, the mixture components are not linearly independent, and when considered as vectors in $\mathbb{R}^3$, $\mu_1$ and $\mu_2$ have the same norm. Our mixture of measures is $\mathcal{P} = \sum_{i=1}^{3} w_i \delta_{\mu_i}$ and our samples come from either $V_5(\mathcal{P})$ or $V_6(\mathcal{P})$ depending on the algorithm used.

We construct our own performance measure which allows us to judge the performance of the estimated components jointly. Let $\hat{\mu}_1, \hat{\mu}_2, \hat{\mu}_3$ be the three estimates for the mixture components from some algorithm. We will view these estimates as vectors in $\mathbb{R}^3$. Our performance measure is

$$\min_{\sigma \in S_3} \frac{1}{3} \sum_{i=1}^{3} \left\| \mu_i - \hat{\sigma}(i) \right\|_{\ell_1(\mathbb{R}^3)}.$$

That is, we take the average of total variations of the best matching of the estimated mixture components to the true components.

9.1. Proposed Algorithms. We include two different implementations of our proposed algorithm, one where we use what we call a “random dominating measure” and in the other we use what we call a “fixed dominating measure.” In the following we describe the two implementations and the rationale for presenting both of them.

In the description of our algorithm in Section 8 we make the assumption that the mixture components, when represented as square integrable densities over some dominating measure, have distinct $L^2$ norms. This is necessary to ensure that (82) admits a unique spectral decomposition. Because $\mu_1$ and $\mu_2$ have the same norm when considered as vectors in $\mathbb{R}^3$ this assumption does not hold for the experiments we present here. We use the aforementioned “random dominating measure” technique (details in Supplemental Material) which transforms the measure space so that the mixture components have distinct norms. To do this we choose a dominating measure randomly so that, with probability one, the mixture components have different norms when represented as densities over this measure space. We theoretically demonstrate that this technique works in Supplemental Material. In this paper we present experimental evidence that this technique also works in practice.

The purpose of the random dominating measure is to create a spectral gap between the mixture components. Intuitively, it seems reasonable that if we choose the dominating measure “well” then we will end up with large spectral gaps without making any of the component norms so diminutive as to become unnoticeable. In the interest of exploring this idea we tested
different dominating measures until we found one that improved algorithmic performance significantly and include these experimental results as well. We found that the dominating measure $\xi$ with $\xi(\{0\}) = 3^2, \xi(\{1\}) = 2^2$ and $\xi(\{2\}) = 1$ improved performance significantly and we refer to this as the “fixed dominating measure” implementation. These experimental results indicate the possibility for significant improvements to our algorithm by choosing the dominating measure intelligently. Additional specifics for these proposed implementations can be found in Supplemental Material.

Both of these implementations were run on two experimental scenarios, one with 50,000 random groups and the other with 10,000,000 random groups, with all groups drawn from $V_5(\mathcal{P})$. We repeated each experiment 20 times and report relevant statistics.

9.2. Competing Algorithms. As a baseline we compare our algorithm against simply choosing 3 measures uniformly at random from the probabilistic simplex. The randomly selected components algorithm was repeated 1000 times. We also compare our algorithm to a modified version of the algorithm introduced in [2]. The algorithm in [2] is designed to work on random groups with three samples and a mixture of measures with linearly independent components. Because of this we apply the algorithm in [2] to random groups from $V_6(\mathcal{P})$ rather than $V_5(\mathcal{P})$, with the adaptation described at the end of Section 8. This algorithm was also run on experimental scenarios with 50,000 and 10,000,000 random groups. Again these experiments were repeated 20 times.

9.3. Results. The results are summarized in Table 1. Our algorithm demonstrates a clear improvement as the number of random groups increases. Our modification of the algorithm in [2] performs noticeably better than the other algorithms, likely owing to the fact that it has more information per group and/or the fact that it does not depend on the "random dominating measure" trick. Using the fixed dominating measure narrows this gap considerably, and it seems likely that this gap could be further improved with a better choice of dominating measure.

10. Discussion. In closing, we offer the following observations related to our results.

10.1. Potential Statistical Test and Estimator. The results on determinedness suggest the possibility of a goodness of fit test. Suppose we have grouped samples from some mixture of measures $\mathcal{P}' = \sum_{i=1}^{m'} w_i' \delta_{\mu_i'}$. Further
Table 1

Experimental Results

<table>
<thead>
<tr>
<th>Method</th>
<th>Performance</th>
</tr>
</thead>
<tbody>
<tr>
<td>Random Dominating Measure, 50,000 samples</td>
<td>Mean: 0.1407, Variance: 0.0169</td>
</tr>
<tr>
<td>Fixed Dominating Measure, 50,000 samples</td>
<td>Mean: 0.0524, Variance: 0.0011</td>
</tr>
<tr>
<td>Anandkumar, et al. [2], 50,000 samples</td>
<td>Mean: 0.0503, Variance: 0.0145</td>
</tr>
<tr>
<td>Random Dominating Measure, 10,000,000 samples</td>
<td>Mean: 0.0433, Variance: 0.0062</td>
</tr>
<tr>
<td>Fixed Dominating Measure, 10,000,000 samples</td>
<td>Mean: 0.0037, Variance: (4e^{-6})</td>
</tr>
<tr>
<td>Anandkumar, et al. [2], 10,000,000 samples</td>
<td>Mean: 0.0026, Variance: (4e^{-6})</td>
</tr>
<tr>
<td>Randomly Selected Measures</td>
<td>Mean: 0.5323, Variance: 0.0203</td>
</tr>
</tbody>
</table>

suppose some null hypothesis

\[ H_0 : \mathcal{P}' = \mathcal{P} \triangleq \sum_{i=1}^{m} w_i \delta_{\mu_i}. \]

Given data from \(V_{2m}(\mathcal{P}')\) we may be able to reject the null hypothesis provided we have some way of estimating \(M \triangleq \sum_{i=1}^{m} w_i \mu_i^{2m}\) from the groups of samples. We will call such an estimator \(\hat{M}\). If \(M\) does not converge to \(M\) then we can reject the null hypothesis. The implementation and analysis of such an estimator would depend on the setting and is outside the scope of this paper.

One interesting observation from the proof of Theorem 4.3 is that, if \(\mathcal{P} = \sum_{i=1}^{m} w_i \delta_{\mu_i}\) is a mixture of measures, \(p_i\) is a pdf for \(\mu_i\) for all \(i\), and \(n > m\), then the rank of \(\sum_{i=1}^{m} a_i \mu_i^{\otimes n} \otimes p_i^{\otimes n}\) will be exactly \(m\). This suggests a statistical estimator for the number of mixture components. The form of this tensor is amenable to spectral methods since it is a positive semi-definite tensor of order 2, which is akin to a positive semi-definite matrix. Embedding the data with the kernel mean mapping, using a universal kernel [19], seems like a promising approach to constructing such a test or estimator.

10.2. Identifiability and the Value \(2n - 1\). The value \(2n - 1\) seems to carry some significance for identifiability beyond the setting we proposed. This value can also be found in results concerning metrics on trees [20], hidden Markov models [21], and frame theory, with applications to signal processing [6]. All of these results are related to identifiability of an object or the injectivity of an operator. We can offer no further insight as to why this value recurs, but it appears to be an algebraic phenomenon.

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SUPPLEMENTARY MATERIAL


References.


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