Abstract

We study the problem of learning mixtures of distributions, a natural formalization of clustering. A mixture of distributions is a collection of distributions $\mathcal{D} = \{D_1, \ldots, D_T\}$, and mixing weights $\{w_1, \ldots, w_T\}$ such that $\sum_i w_i = 1$. A sample from a mixture is generated by choosing $i$ with probability $w_i$ and then choosing a sample from distribution $D_i$.

The problem of learning the mixture is that of finding the parameters of the distributions comprising $\mathcal{D}$, given only the ability to sample from the mixture. In this paper, we restrict ourselves to learning mixtures of product distributions.

The key to learning the mixtures is to find a few vectors, such that points from different distributions are sharply separated upon projection onto these vectors. Previous techniques use the vectors corresponding to the top few directions of highest variance of the mixture. Unfortunately, these directions may be directions of high noise and not directions along which the distributions are separated. Further, skewed mixing weights amplify the effects of noise, and as a result, previous techniques only work when the separation between the input distributions is large relative to the imbalance in the mixing weights.

In this paper, we show an algorithm which successfully learns mixtures of distributions with a separation condition that depends only logarithmically on the skewed mixing weights. In particular, it succeeds for a separation between the centers that is $\Theta(\sqrt{T \log \Lambda})$, where $\sigma$ is the maximum directional standard deviation of any distribution in the mixture, $T$ is the number of distributions, and $\Lambda$ is polynomial in $T$, $\sigma$, $\log n$ and the imbalance in the mixing weights. For our algorithm to succeed, we require a spreading condition, that the distance between the centers be spread across $\Theta(T \log \Lambda)$ coordinates. Additionally, with arbitrarily small separation, i.e., even when the separation is not enough for clustering, with enough samples, we can approximate the subspace containing the centers. Previous techniques failed to do so in polynomial time for non-spherical distributions regardless of the number of samples, unless the separation was large with respect to the maximum directional variance $\sigma$ and polynomially large with respect to the imbalance of mixing weights. Our algorithm works for Binary Product Distributions and Axis-Aligned Gaussians. The spreading condition above is implied by the separation condition for binary product distributions, and is necessary for algorithms that rely on linear correlations.

Finally, when a stronger version of our spreading condition holds, our algorithm performs successful clustering when the separation between the centers is only $\Theta(\sigma^* \sqrt{T \log \Lambda})$, where $\sigma^*$ is the maximum directional standard deviation in the subspace containing the centers of the distributions.
1 Introduction

Clustering, the problem of grouping together data points in high dimensional space using a similarity measure, is a fundamental problem of statistics with numerous applications in a wide variety of fields. A natural model for clustering is that of learning mixtures of distributions. A mixture of distributions is a collection of distributions \( \mathcal{D} = \{D_1, \ldots, D_T\} \), and mixing weights, \( \{w_1, \ldots, w_T\} \) such that \( \sum_i w_i = 1 \). A sample from a mixture is generated by choosing \( i \) with probability \( w_i \) and choosing a sample from distribution \( D_i \). The problem of learning the mixture is that of finding the parameters of the distributions comprising \( D \), given only the ability to sample from the mixture.

If the distributions \( D_1, \ldots, D_T \) are very close to each other, then even if we knew the parameters of the distributions, it would be impossible to classify the points correctly with high confidence. Therefore, Dasgupta [Das99] introduced the notion of a separation condition, which is a promise that each pair of distributions is sufficiently different according to some measure. Given points from a mixture of distributions and a separation condition, the goal is to find the parameters of the mixture \( D \), and cluster all but a small fraction of the points correctly. A commonly used separation measure is the distance between the centers of the distributions parameterized by the maximum directional variance, \( \sigma \), of any distribution in the mixture.

A common approach to learning the mixtures and therefore, clustering the high-dimensional cloud of points is to find a few interesting vectors, such that points from different distributions are sharply separated upon projection onto these vectors. Various distance-based methods [AK01, LO82, DLR77] are then applied to cluster in the resulting low-dimensional subspace. The state-of-the-art, in practice, is to use the vectors corresponding to the top few directions of highest variance of the mixture and to hope that it contains most of the separation between the centers. This is computed by a Singular Value Decomposition (SVD) of the matrix of samples. This approach has been theoretically analyzed by [VW02] for spherical distributions, and for more general distributions in [KSV05, AM05]. The latter show that the maximum variance directions are indeed the interesting directions when the separation is \( \Theta(\frac{1}{\sqrt{w_{\min}}} \) \), where \( w_{\min} \) is the smallest mixing weight of any distribution.

This is the best possible result for SVD-based approaches; the directions of maximum variance may well not be the directions in which the centers are separated, but instead may be the directions of very high noise, as illustrated in Figure 1(b). This problem is exacerbated when the mixing weights \( w_i \) are skewed – because a distribution with low mixing weight diminishes the contribution to the variance along a direction that separates the centers.

This bound is suboptimal for two reasons. Although mixtures with skewed mixing weights arise naturally in practice (see [PS00] for an example), given enough samples, mixing weights have no bearing on the separability of distributions. Consider two mixtures \( D' \) and \( D'' \) of distributions \( D_1 \) and \( D_2 \); in \( D' \), \( w_1 = w_2 = 1/2 \), and in \( D'' \), \( w_1 = 1/4 \) and \( w_2 = 3/4 \). Given enough computational resources, if we can learn \( D' \) from 50 samples, we should be able to learn \( D'' \) from 100 samples. This does not necessarily hold for SVD-based methods. Secondly, regardless of \( \sigma \), an algorithm, which has prior knowledge of the subspace containing the centers of the distributions, should be able to learn the mixture when the separation is proportional to \( \sigma^2 \), the maximum directional standard deviation of any distribution in the subspace containing the centers. An example in which \( \sigma \) and \( \sigma_{\ast} \) are significantly different is shown in Figure 1(b).

In this paper, we study the problem of learning mixtures of product distributions. A product distribution over \( \mathbb{R}^n \) is one in which each coordinate is distributed independently of any others. In practice, mixtures of product distributions have been used as mathematical models for data and learning mixtures of product distributions specifically has been studied [FM99, FOS05, FOS06, DHKS05] – see the Related Work section for examples and details. However, even under this seemingly restrictive assumption, providing an efficient algorithm that does better than the bounds of [AM05, KSV05] turns out to be quite challenging. The main challenge is to find a low-dimensional subspace that contains most of the separation between the centers; although the independence assumption can (sometimes) help us identify which coordinates contribute to the distance between some pair of centers, the problem of actually finding the low-dimensional space still requires more involved techniques.

In this paper, we present an algorithm for learning mixtures of product distributions, which is stable in the presence of skewed mixing weights, and, under certain conditions, in the presence of high variance outside the subspace containing the centers. In particular, the dependence of the separation required by our algorithm on skewed mixing weights is only logarithmic. Additionally, with arbitrarily small separation, (i.e., even when the separation is not enough for classification), with enough samples, we can approximate the subspace containing the centers. Previous techniques failed to do so for non-spherical distributions regardless of the number of samples, unless the separation was sufficiently large. Our algorithm works for binary product distributions and axis-aligned Gaussians. We require that the distance between the centers be spread across \( \Theta(T \log \Lambda) \) coordinates, where \( \Lambda \) depends polynomially on the maximum dis-
tance between centers and $\omega_{\text{min}}$. For our algorithm to classify the samples correctly, we further need the separation between centers to be $\Theta(\sigma \sqrt{T \log \Lambda})$.

In addition, if a stronger version of the spreading condition is satisfied, then our algorithm requires a separation of only $\Theta(\sigma \sqrt{T \log \Lambda})$ to ensure correct classification of the samples. The stronger spreading condition, discussed in more detail later, ensures that when we split the coordinates randomly into two sets, the maximum directional variance of any distribution in the mixture along the projection of the subspace containing the centers into the subspaces spanned by the coordinate vectors in each set, is comparable to $\sigma^2$.

In summary, compared to [AM05, KSV05], our algorithm is much (exponentially) less susceptible to the imbalance in mixture weights and, when the stronger spreading condition holds, to high variance noise outside the subspace containing the centers. However, our algorithm requires a spreading condition and coordinate-independence, while [AM05, KSV05] are more general. We note that for perfectly spherical distributions, the results of [VW02] are better than our results – however, these results do not apply even for distributions with bounded eccentricity. Finally unlike the results of [Das99, AK01, DS00], which require the separation to grow polynomially with dimension, our separation only grows logarithmically with the dimension.

Our algorithm is based upon two key insights. The first insight is that if the centers are separated along several coordinates, then many of these coordinates are correlated with each other. To exploit this observation, we choose the coordinates randomly, and search the space of this half for directions of high variance. We use the remaining half of coordinates to filter the found directions. If a found direction separates the centers, it is likely to have some correlation with coordinates in the remaining half, and therefore is preserved by the filter. If, on the other hand, the direction found is due to noise, coordinate independence ensures that there will be no correlation with the second half of coordinates, and therefore such directions get filtered away.

The second insight is that the tasks of searching for and filtering the directions can be simultaneously accomplished via a singular value decomposition of the matrix of covariances between the two halves of coordinates. In particular, we show that the top few directions of maximum variance of the covariance matrix approximately capture the subspace containing the centers. Moreover, we show that the covariance matrix has low singular value along any noise direction. By combining these ideas, we obtain an algorithm that is almost insensitive to mixing weights, a property essential for applications like population stratification [CHRZ07], and which can be implemented using the heavily optimized and thus, efficient, SVD procedure, and which works with a separation condition closer to the information theoretic bound.

### Related Work

The first provable results for learning mixtures of Gaussians are due to Dasgupta [Das99] who shows how to learn mixtures of spherical Gaussians with a separation of $\Theta(\sigma / \sqrt{n})$ in an $n$-dimensional space. An EM based algorithm by Dasgupta and Schulman [DS00] was shown to apply to more situations, and with a separation of $\Theta(\sigma n^{1/4})$. Arora and Kannan [AK01] show how to learn mixtures of distributions of arbitrary Gaussians whose centers are separated by $\Theta(n^{1/4} \sigma)$. Their results apply to many other situations, for example, concentric Gaussians with sufficiently different variance.

The first result that removed the dependence on $n$ in the separation requirement was that of Vempala and Wang [VW02] who use SVD to learn mixtures of spherical Gaussians with $O(\sigma T^{1/4})$ separation. They project to a subspace of dimension $T$ using an SVD and use a distance based method in the low dimensional space. If the separation is not enough for classification, [VW02] can also find, given enough samples, a subspace approximating the subspace containing the centers. While the results of [VW02] are independent of the imbalance on mixing weights, they apply only to perfectly spherical Gaussians, and cannot be extended to Gaussians with bounded eccentricity. In further work Kannan, Salmasian, and Vempala[KSV05] and Achlioptas and McSherry [AM05] show how to cluster general Gaussians using SVD. While these results are weaker than ours, they apply to a mixture of general Gaussians, axis-aligned or not. We note that their analysis also applies to binary product distributions again with polynomial dependence on the imbalance in mixing weights. In contrast, our separation requirement is $\Omega(\sigma \sqrt{T \log \Lambda})$, i.e., is logarithmically dependent on the mixing weights and dimension and the maximum variance in noise directions.

There is also ample literature on specifically learning mixtures of product distributions. Freund and Mansour [FM99] show an algorithm which generates distributions that are $\epsilon$-close to a mixture of two product distributions over $\{0,1\}^n$ in time polynomial in $n$ and $1/\epsilon$. Feldman, O’Donnell, and Servedio show how to generate distributions that are $\epsilon$-close to a mixture of $T$ product distributions [FOS05] and axis-aligned Gaussians [FOS06]. Like [FM99], they have no separation requirements, but their algorithm while polynomial in $1/\epsilon$ takes $n^{O(T^3)}$ time. Dasgupta et. al [DHKS05] provide an algorithm for learning mixtures of heavy-tailed product distributions which works with a separation of $\Theta(R/T)$, where $R$ is the maximum half-radius of any

---

1 They do not directly address binary product distributions in their paper, but their techniques apply.
distribution in the mixture. While their separation requirement does not depend polynomially on \( \frac{w_{\min}}{w_{\max}} \), their algorithm runs in time exponential in \( \Theta\left(\frac{w_{\min}}{w_{\max}}\right) \). They also require a slope, which is comparable to our spreading condition. Chaudhuri et al. [CHRZ07] show an iterative algorithm for learning mixtures of two product distributions that implicitly uses the notion of co-ordinate independence to filter out noise directions. However, the algorithm heavily uses the two distribution restriction to find the appropriate directions, and does not work when \( T > 2 \).

More broadly, the problem of analyzing mixture models data has received a great deal of attention in statistics, see for example, [MB88, TSM85], and has numerous applications. We present three applications where data is modelled as a mixture of product distributions. First, the problem of population stratification in population genetics has been posed as learning mixtures of binary product distributions in [SRH07]. In their work, the authors develop an MCMC method for addressing the problem and their software embodiment is widely used. A second application is in speech recognition [Rey95, PFK02], which models acoustic features at a specific location model [BNJ03]. Here, documents are modelled as distributions over topics which, in turn, are distributions over words. Subsequent choices of topics and words are assumed to be independent. (For words, this is referred to as the “bag of words” assumption.) [BNJ03] develops variational techniques that provide interesting results for various corpora. Interestingly, the same model was used by Kleinberg and Sandler [KS04] to model user preferences for purchasing goods (users correspond to documents, topics to categories, and words to goods). Their algorithm, which provides provably good performance in this model, also uses SVD-like clustering algorithms as a subroutine.

**Discussion**

**The Spreading Condition.** The spreading condition loosely states that the distance between each pair of centers is spread along \( \Theta(T \log \Lambda) \) coordinates. We demonstrate by an example, that a spread of \( \Omega(T) \) is a natural limit for all methods that use linear correlations between coordinates, such as our methods and SVD based methods [VW02, KSV05, AM05]. We present, as an example, two distributions: a mixture \( D_1 \) of \( T \) binary product distributions, and a single binary product distribution \( D_2 \), which have exactly the same covariance matrix. Our example is based on the Hadamard code, in which a codeword for a \( k \)-bit message is \( 2^k \) bits long, and includes a parity bit for each subset of the bits of the message. The distributions comprising \( D_1 \) are defined as follows. Each of the \( T = 2^k \) centers is a code-word for a \( k \)-bit string appended by a string of length \( n - k \) in which each coordinate has value \( 1/2 \). Notice that the last \( n - k \) bits are noise. Thus, the centers are separated by \( T/2 \) coordinates. \( D_2 \) is the uniform distribution over the \( n \)-dimensional hypercube. As there are no linear correlations between any two bits in the Hadamard code, the covariance of \( D_1 \) along any two directions is 0, and each direction has the same variance. As this is also the case for \( D_2 \), any SVD-based or correlation-based algorithm will fail to distinguish between the two mixtures. We also note that learning binary product distributions with minimum separation 2 and average separation \( 1 + \frac{1}{4} \log T \) would allow one to learn parities of \( \log T \) variables with noise. Finally, we note that when the spreading condition fails, one has only a few coordinates that contain most of the distance between centers. One could enumerate the set of possible coordinates to deal with this case, and is exponential in \( T \log n \log \Lambda \), [FOS05] on the other hand takes time exponential in \( T^3 \log n \), and works with no separation requirement.

**2 A Summary of Our Results**

We begin with some preliminary definitions about distributions drawn over \( n \) dimensional spaces. We use \( f, g, \ldots \) to range over coordinates, and \( i, j, \ldots \) to range over distributions. For any \( x \in \mathbb{R}^n \), we write \( x^f \) for the \( f \)-th coordinate of \( x \). For any subspace \( H \) (resp. vector \( v \)), we use \( H \) (resp. \( v \)) to denote the orthogonal complement of \( H \) (resp. \( v \)). For a subspace \( H \) and a vector \( v \), we write \( P_H(v) \) for the projection of \( v \) onto the subspace \( H \). For any vector \( x \), we use \( ||x|| \) for the Euclidean norm of \( x \). For any two vectors \( x \) and \( y \), we use \( (x, y) \) for the dot-product of \( x \) and \( y \).

**Mixtures of Distributions.** A mixture of distributions \( \mathcal{D} \), is a collection of distributions, \( \{D_1, \ldots, D_T\} \), over points in \( \mathbb{R}^n \), and a set of mixing weights \( w_1, \ldots, w_T \) such that \( \sum w_i = 1 \). In the sequel, \( n \) is assumed to be much larger than \( T \). In a product distribution over \( \mathbb{R}^n \), each coordinate is distributed independently of the others. When working with a mixture of binary product distributions, we assume that the \( f \)-th coordinate of a point drawn from distribution \( D_i \) is 1 with probability \( \mu^f_i \), and 0 with probability \( 1 - \mu^f_i \). When working with a mixture of axis-aligned Gaussian distributions, we assume that the \( f \)-th coordinate of a point drawn from distribution \( D_i \) is distributed as a Gaussian with mean \( \mu^f_i \) and standard deviation \( \sigma^f_i \).

**Centers.** We define the center of a distribution \( i \) as the vector \( \mu_i \), and the center of mass of the mixture as the vector \( \bar{\mu} \) where \( \bar{\mu} \) is the mean of the mixture for the coordinate \( f \). We write \( \mathcal{C} \) for the subspace containing \( \mu_1, \ldots, \mu_T \).

**Directional Variance.** We define \( \sigma^2 \) as the maximum
variance of any distribution in the mixture along any direction. We define $\sigma^2$ as the maximum variance of any distribution in the mixture along any direction in the subspace containing the centers of the distributions. We write $\sigma_{\text{max}}^2$ as the maximum variance of the entire mixture in any direction. This may be more than $\sigma^2$ due to contribution from the separation between the centers.

**Spread.** We say that a unit vector $v$ in $\mathbb{R}^n$ has spread $S$ if $\sum f(v^f)^2 \geq S \cdot \max_f (v^f)^2$.

**Distance.** Given a subspace $\mathcal{K}$ of $\mathbb{R}^n$ and two points $x, y$ in $\mathbb{R}^n$, we write $d_{\mathcal{K}}(x, y)$ for the square of the Euclidean distance between $x$ and $y$ projected along the subspace $\mathcal{K}$.

The spreading condition and effective distance. The spreading condition tells us that the distance between each $\mu_i$ and $\mu_j$ should not be concentrated along a few coordinates. One way to ensure this is to demand that for all $i, j$, the vector $\mu_i - \mu_j$ has high spread. This is comparable to the slope condition used in [DHKS05]. However, we do not need such a strong condition for dealing with mixtures with imbalanced mixing weights. Our spreading condition therefore demands that for each pair of centers $\mu_i, \mu_j$, the norm of the vector $\mu_i - \mu_j$ high, even if we ignore the contribution of the top few (about $T \log T$) coordinates. Due to technicalities in our proofs, the number of coordinates we can ignore needs to depend (logarithmically) on this distance.

We therefore define the spreading condition as follows. We define parameters $c_{ij}$ and a parameter $\Lambda$ as:

$$\Lambda > \frac{\min_{\text{max}} T \log^2 n}{w_{\min} (\min_{ij} c_{ij})}$$

and $c_{ij}$ is the maximum value such that there are $49T \log \Lambda$ coordinates $f$ with $|\mu_i^f - \mu_j^f| > c_{ij}$. We note that $\Lambda$ is bounded by a polynomial in $T, \sigma, 1/w_{\min}, 1/c_{ij}$ and logarithmic in $n$.

We define $c_{\min}$ to be the minimum over all pairs $i, j$ of $c_{ij}$. Given a pair of centers $i$ and $j$, let $\Delta_{ij}$ be the set of coordinates $f$ such that $|\mu_i^f - \mu_j^f| > c_{ij}$, and let $\nu_{ij}$ be defined as: $\nu_{ij}^f = \mu_i^f - \mu_j^f$, if $f \notin \Delta_{ij}$, and $\nu_{ij}^f = c_{ij}$ otherwise. We define $d(\mu_i, \mu_j)$, the effective distance between $\mu_i$ and $\mu_j$, to be the square of the $L_2$ norm of $\nu_{ij}$. In contrast, the square of the norm of the vector $\mu_i - \mu_j$ is the actual distance between centers $\mu_i$ and $\mu_j$, and is always greater than or equal to the effective distance between $\mu_i$ and $\mu_j$. Moreover, given $i$ and $j$ and the subspace $\mathcal{K}$, we define $d_{\mathcal{K}}(\mu_i, \mu_j)$ as the square of the norm of the vector $\nu_{ij}$ projected onto the subspace $\mathcal{K}$.

Under these definitions, our spreading condition now requires that $d(\mu_i, \mu_j) \geq 49c_{ij}^2T \log \Lambda$ and our stronger spreading condition requires that every vector in $\mathcal{C}$ has spread $32T^2 \log \frac{\sigma}{\Delta}$.

A Formal Statement of Our Results. Our main contribution is Algorithm CORR-CLUSTER, a correlation based algorithm for learning mixtures of binary product distributions and axis-aligned Gaussians. The input to the algorithm is a set of samples from a mixture of distributions, and the output is a clustering of the samples.

The main component of Algorithm CORR-CLUSTER is Algorithm CORR-SUBSPACE, which, given samples from a mixture of distributions, computes an approximation to the subspace containing the centers of the distributions. The motivation for approximating the latter space is as follows. In the $T$-dimensional subspace containing the centers of the distributions, the distance between each pair of centers $\mu_i$ and $\mu_j$ is the same as their distance in $\mathbb{R}^n$; however, because of the low dimensionality, the magnitude of the noise is small. Therefore, provided the centers of the distributions are sufficiently separated, projection onto this subspace will sharply separate samples from different distributions. SVD-based algorithms [VW02, AM05, KSV05] attempt to approximate this subspace by the top $T$ singular vectors of the matrix of samples. However, for product distributions, our Algorithm CORR-SUBSPACE can approximate this subspace correctly under more restrictive separation conditions.

The properties of Algorithms CORR-SUBSPACE and CORR-CLUSTER are formally summarized in Theorem 1 and Theorem 2 respectively.

**Theorem 1 (Spanning centers)** Suppose we are given a mixture of distributions $\mathcal{D} = \{D_1, \ldots, D_T\}$, with mixing weights $w_1, \ldots, w_T$. Then with at least constant probability, the subspace $\mathcal{K}$ of dimension at most $2T$ output by Algorithm CORR-SUBSPACE has the following properties.

1. If, for all $i$ and $j$, $d(\mu_i, \mu_j) \geq 49c_{ij}^2T \log \Lambda$, then, for all pairs $i, j$,

   $$d_{\mathcal{K}}(\mu_i, \mu_j) \geq \frac{99}{100} (d(\mu_i, \mu_j) - 49Tc_{ij}^2 \log \Lambda)$$

2. If, in addition, every vector in $\mathcal{C}$ has spread $32T^2 \log \frac{\sigma}{\Delta}$, then with at least constant probability, the maximum directional variance in $\mathcal{K}$ of any distribution $D_i$ in the mixture is at most $11\sigma^2_d$.

The number of samples required by Algorithm CORR-SUBSPACE is polynomial in $\frac{1}{\sigma^2_d}$, $T$, $\eta, \sigma$ and $\frac{1}{\text{max}_i}$, and the algorithm runs in time polynomial in $n, T$, and the number of samples.

The subspace $\mathcal{K}$ computed by Algorithm CORR-SUBSPACE approximates the subspace containing the centers of the distributions in the sense that the distance between each pair of centers $\mu_i$ and $\mu_j$ is high along $\mathcal{K}$. Theorem 1 states that Algorithm CORR-SUBSPACE computes an approximation to the subspace containing the centers of
Theorem 2 (Clustering) Suppose we are given a mixture of distributions \( D = \{D_1, \ldots, D_T\} \), with mixing weights \( w_1, \ldots, w_T \). Then, Algorithm CORR-CLUSTER has the following properties.

1. If for all \( i \) and \( j \), \( d(\mu_i, \mu_j) \geq 49T\sigma^2 \log \Lambda \), and for all \( i, j \) we have:
   \[
   d(\mu_i, \mu_j) > 59\sigma^2 T(\log \Lambda + \log n)
   \]
   (for axis-aligned Gaussians)
   \[
   d(\mu_i, \mu_j) > 59T(\log \Lambda + \log n)
   \]
   (for binary product distributions)
   then with probability \( 1 - \frac{1}{n} \) over the samples and with constant probability over the random choices made by the algorithm, Algorithm CORR-CLUSTER computes a correct clustering of the sample points.

2. For axis-aligned Gaussians, if every vector in \( C \) has spread at least \( 32T \log \frac{\Lambda}{\sigma^2} \), and for all \( i, j \):
   \[
   d(\mu_i, \mu_j) \geq 150\sigma^2 T(\log \Lambda + \log n)
   \]
   then, with constant probability over the randomness in the algorithm, and with probability \( 1 - \frac{1}{n} \) over the samples, Algorithm CORR-CLUSTER computes a correct clustering of the sample points.

Algorithm CORR-CLUSTER runs in time polynomial in \( n \) and the number of samples required by Algorithm CORR-CLUSTER is polynomial in \( \sigma^2, T, n, \sigma \) and \( \frac{1}{\alpha_{\min}} \).

We note that because we are required to do classification here, we do require an absolute lower bound on the distance between each pair of centers in Theorem 2.

The second theorem follows from the first and the distance concentration Lemmas of [AM05] as described in Section 5.3 of the Appendix. The Lemmas show that once the points are projected onto the subspace computed in Theorem 1, a distance-based clustering method suffices to correctly cluster the points.

A Note on the Stronger Spreading Condition. The motivation for requiring the stronger spreading condition is as follows. Our algorithm splits the coordinates randomly into two sets \( F \) and \( G \). If \( C_F \) and \( C_G \) denote the restriction of \( C \) to the coordinates in \( F \) and \( G \) respectively, then our algorithm requires that the maximum directional variance of any distribution in the mixture is close to \( \sigma \) in \( C_F \) and \( C_G \) respectively. Notice that this does not follow from the fact that the maximum directional variance along \( C \) is \( \sigma^2 \). Suppose \( C \) is spanned by \( (0,1,0,1,1) \) and \( (0,1,1,1,1) \), variances of \( D_1 \) along the axes are \((10,1,1,1,1)\), and \( F = \{1,2\} \). Then, \( \sigma^2 \) is about 2.8, while the variance of \( D_1 \) along \( C_F \) is 10. However, as Lemma 7 shows, the required condition is ensured by the strong spreading condition.

However, in general, the maximum directional variance of any \( D_i \) in the mixture along \( C_F \) and \( C_G \) may still be close to \( \sigma^2 \), even though strong spreading condition is far from being met. For example: if \( C \) is the space spanned by the first \( T \) coordinate vectors \( e_1, \ldots, e_T \), then with probability \( 1 - \frac{1}{n^T} \), the maximum variance along \( C_F \) and \( C_G \) is also \( \sigma^2 \).

3 Algorithm CORR-CLUSTER

Our clustering algorithm follows the same basic framework as the SVD-based algorithms of [VW02, KSV05, AM05]. The input to the algorithm is a set \( S \) of samples, and the output is a pair of clusterings of the samples according to source distribution.

<table>
<thead>
<tr>
<th>CORR-CLUSTER(S)</th>
</tr>
</thead>
<tbody>
<tr>
<td>1. Partition ( S ) into ( S_A ) and ( S_B ) uniformly at random.</td>
</tr>
<tr>
<td>2. Compute: ( K_A = Corr - Subspace(S_A) ), ( K_B = Corr - Subspace(S_B) )</td>
</tr>
<tr>
<td>3. Project each point in ( S_B ) (resp. ( S_A )) on the subspace ( K_A ) (resp. ( K_B )).</td>
</tr>
<tr>
<td>4. Use a distance-based clustering algorithm [AK01] to partition the points in ( S_A ) and ( S_B ) after projection.</td>
</tr>
</tbody>
</table>

The first step in the algorithm is to use Algorithm CORR-SUBSPACE to find a \( O(T) \)-dimensional subspace \( K \) which is an approximation to the subspace containing the centers of the distributions. Next, the samples are projected onto \( K \) and a distance-based clustering algorithm is used to find the clusters.

We note that in order to preserve independence the samples we project onto \( K \) should be distinct from the ones we use to compute \( K \). A clustering of the complete set of points can then be computed by partitioning the samples into two sets \( A \) and \( B \). We use \( A \) to compute \( K_A \), which is used to cluster \( B \) and vice-versa.

We now present our algorithm which computes a basis for the subspace \( K \). With slight abuse of notation we use \( K \) to denote the set of vectors that form the basis for
the subspace \( \mathcal{K} \). The input to \textsc{Corr-Subspace} is a set \( S \) of samples, and the output is a subspace \( \mathcal{K} \) of dimension at most \( 2T \).

\textbf{Algorithm \textsc{Corr-Subspace}:}

\begin{itemize}
  \item \textbf{Step 1: Initialize and Split} Initialize the basis \( \mathcal{K} \) with the empty set of vectors. Randomly partition the coordinates into two sets, \( \mathcal{F} \) and \( \mathcal{G} \), each of size \( n/2 \). Order the coordinates as those in \( \mathcal{F} \) first, followed by those in \( \mathcal{G} \).
  \item \textbf{Step 2: Sample} Translate each sample point so that the center of mass of the set of sample points is at the origin. Let \( \tilde{F} \) (respectively \( \tilde{G} \)) be the matrix which contains a row for each sample point, and a column for each coordinate in \( \mathcal{F} \) (respectively \( \mathcal{G} \)). For each matrix, the entry at row \( x \), column \( f \) is the value of the \( f \)-th coordinate of the sample point \( x \) divided by \( \sqrt{|S|} \).
  \item \textbf{Step 3: Compute Singular Space} For the matrix \( \tilde{F}^T \tilde{G} \), compute \( \{v_1, \ldots, v_T\} \), the top \( T \) left singular vectors, \( \{y_1, \ldots, y_T\} \), the top \( T \) right singular vectors, and \( \{\lambda_1, \ldots, \lambda_T\} \), the top \( T \) singular values.
  \item \textbf{Step 4: Expand Basis} For each \( i \), we abuse notation and use \( v_i \) (resp. \( y_i \)) to denote the vector obtained by concatenating \( v_i \) with the \( 0 \) vector in \( n/2 \) dimensions (\( 0 \) vector in \( n/2 \) dimensions concatenated with \( y_i \) respectively). For each \( i \), if the singular value \( \lambda_i \) is more than a threshold \( \tau = O \left( \frac{\min \hat{\sigma}_i^2}{\log N} \cdot \sqrt{\log N} \right) \), we add \( v_i \) and \( y_i \) to \( \mathcal{K} \).
  \item \textbf{Step 5: Output} Output the set of vectors \( \mathcal{K} \).
\end{itemize}

The main idea behind our algorithm is to use half the coordinates to compute a subspace which approximates the subspace containing the centers, and the remaining half to validate that the subspace computed is indeed a good approximation. We critically use the coordinate independence property of product distributions to make this validation possible.

\section{Analysis of Algorithm \textsc{Corr-Cluster}}

This section is devoted to proving Theorems 1, and 2. We use the following notation.

\textbf{Notation.} We write \( \mathcal{F} \)-space (resp. \( \mathcal{G} \)-space) for the \( n/2 \) dimensional subspace of \( \mathbb{R}^n \) spanned by the coordinate vectors \( \{e_f \mid f \in \mathcal{F}\} \) (resp. \( \{e_g \mid g \in \mathcal{G}\} \)). We write \( \hat{\mathcal{C}} \) for the subspace spanned by the set of vectors \( \mu_i \). We write \( \mathcal{C}_\mathcal{F} \) for the space spanned by the set of vectors \( \mathbf{P}_\mathcal{F}(\mu_i) \). We write \( \mathbf{P}_\mathcal{F}(\hat{\mathcal{C}}_\mathcal{F}) \) for the orthogonal complement of \( \mathcal{C}_\mathcal{F} \) in the \( \mathcal{F} \)-space. Moreover, we write \( \mathcal{C}_\mathcal{F} \cap \mathcal{G} \) for the subspace of dimension \( 2T \) spanned by the union of a basis of \( \mathcal{C}_\mathcal{F} \) and a basis of \( \mathcal{C}_\mathcal{G} \). Next, we define a key ingredient of the analysis.

\textbf{Covariance Matrix.} Let \( N \) be a large number. We define \( \hat{F} \) (resp. \( \hat{G} \)), the perfect sample matrix with respect to \( F \) (resp. \( G \)) as the \( N \times n/2 \) matrix whose rows from \( (w_1 + \ldots + w_{i-1})N + 1 \) through \( (w_1 + \ldots + w_i)N \) are equal to the vector \( \mathbf{P}_\mathcal{F}(\mu_i) / \sqrt{N} \) (resp. \( \mathbf{P}_\mathcal{G}(\mu_i) / \sqrt{N} \)). For a coordinate \( f \), let \( X_f \) be a random variable which is distributed as the \( f \)-th coordinate of the mixture \( D \). As the entry in row \( f \) and column \( g \) in the matrix \( \hat{F}^T \hat{G} \) is equal to \( \text{Cov}(X_f, X_g) \), the covariance of \( X_f \) and \( X_g \), we call the matrix \( \hat{F}^T \hat{G} \) the covariance matrix of \( F \) and \( G \).

\textbf{Proof Structure.} The overall structure of our proof is as follows. First, we show that the centers of the distributions in the mixture have a high projection on the subspace of highest correlation between the coordinates. To do this, we first assume in Section 4.1 that the input to the algorithm in Step 2 are the perfect sample matrices \( F \) and \( G \). Of course, we cannot directly feed in the matrices \( F, G \) as the values of the centers are not known in advance. Next, we show in Section 4.2 that this holds even when the matrices \( F \) and \( G \) in Step 2 of Algorithm \textsc{Corr-Subspace} are obtained by sampling. In Section 4.3, we combine these two results and prove Theorem 1. Finally, in Section 5.3 of the Appendix, we show that distance concentration algorithms work in the low-dimensional subspace produced by Algorithm \textsc{Corr-Cluster}, and complete the analysis by proving Theorem 2.

\subsection{The Perfect Sample Matrix}

The goal of this section is to prove Lemmas 3 and 5, which establish a relationship between directions of high correlation of the covariance matrix constructed from the perfect sample matrix, and directions which contain a lot of separation between centers. Lemma 3 shows that a direction which contains a lot of effective distance between some pair of centers, is also a direction of high correlation.

Lemma 5 shows that a direction \( v \in \mathbf{P}_\mathcal{F}(\hat{\mathcal{C}}_{\mathcal{F}}) \), which is perpendicular to the space containing the centers, is a direction with 0 correlation. In addition, we show in Lemma 6, another property of the perfect sample matrix – the covariance matrix constructed from the perfect sample matrix has rank at most \( T \). We conclude this section by showing in Lemma 7 that when every vector in \( \mathcal{C} \) has high spread, the directional variance of any distribution in the mixture along \( \mathcal{F} \)-space or \( \mathcal{G} \)-space is of the order of \( \sigma_z^2 \).

We begin by showing that if a direction \( v \) contains a lot of the distance between the centers, then, for most ways of splitting the coordinates, the magnitude of the covariance of the mixture along the projection of \( v \) on \( \mathcal{F} \)-space and the projection of \( v \) \( \mathcal{G} \)-space is high. In other words, the projections of \( v \) along \( \mathcal{F} \)-space and \( \mathcal{G} \)-space
are directions of high correlation.

**Lemma 3** Let \( v \) be any vector in \( C_{F, G} \) such that for some \( i \) and \( j \), \( d_v(\mu_i, \mu_j) \geq 49T \epsilon_n^2 \log \Lambda \). If \( v_F \) and \( v_G \) are the normalized projections of \( v \) to \( F \)-space and \( G \)-space respectively, then, with probability at least \( 1 - \frac{1}{T} \) over the splitting step, for all such \( v \), \( v_F^T \hat{F} v_G \geq \tau \) where \( \tau = O\left( \frac{w_{\min}^2 \epsilon_n^2}{T \log^2 n} \cdot \sqrt{\log \Lambda} \right) \).

The main ingredient of the proof, which is in the Appendix, is Lemma 4.

**Lemma 4** Let \( v \) be a fixed vector in \( C \) such that for some \( i \) and \( j \), \( d_v(\mu_i, \mu_j) \geq 49T c^2 \log \Lambda \). If \( v_F \) and \( v_G \) are the projections of \( v \) to \( F \)-space and \( G \)-space respectively, then, with probability at least \( 1 - \Lambda^{-2} \) over the splitting step, \( v_F^T \hat{F} \hat{G} v_G \geq 2\tau \) where \( \tau = O\left( \frac{w_{\min}^2 \epsilon_n^2}{T \log^2 n} \cdot \sqrt{\log \Lambda} \right) \).

Let \( \hat{F}_v \) (\( \hat{G}_v \) respectively) be the \( s \times n/2 \) matrix obtained by projecting each row of \( \hat{F} \) (respectively \( \hat{G} \)) on \( v_F \) (respectively \( v_G \)). Then,

\[
v_F^T \hat{F}_v \hat{G}_v v_G = \sum_i w_i(v_F, P_{v_F}(\mu_i - \hat{\mu})) \langle v_G, P_{v_G}(\mu_i - \hat{\mu}) \rangle
\]

\[
= v_F^T \hat{F} \hat{G} v_G
\]

Moreover, for any pair of vectors \( x \) in \( F \)-space and \( y \) in \( G \)-space such that \( \langle x, v_F \rangle = 0 \) and \( \langle y, v_G \rangle = 0 \),

\[
x^T \hat{F} \hat{G} y = \sum_i w_i(x, P_{v_F}(\mu_i - \hat{\mu})) \langle y, P_{v_G}(\mu_i - \hat{\mu}) \rangle = 0
\]

Therefore, \( \hat{F}_v \hat{G}_v \) has rank at most 1.

The proof strategy for Lemma 4 is to show that if \( d_v(\mu_i, \mu_j) \) is large then the matrix \( \hat{F}_v \hat{G}_v \) has high norm. We require the following notation. For each coordinate \( i \) define a \( T \)-dimensional vector \( z \) as

\[
z_f = [\sqrt{w_i} P_v(\mu_i' - \hat{\mu}'), \ldots, \sqrt{w_T} P_v(\mu_T' - \hat{\mu}')]^T
\]

Notice that for any two coordinates \( f, g \):

\[
\langle z_f, z_g \rangle = \text{Cov}(P_v(X_f), P_v(X_g))
\]

, computed over the entire mixture. We also observe that

\[
\sum_f ||z||^2 = \sum_i w_i \cdot d_v(\mu_i, \hat{\mu})
\]

The RHS of this equality is the weighted sum of the squares of the Euclidean distances between the centers of the distributions and the center of mass. By the triangle inequality, this quantity is at least \( 49w_{\min}^2 \epsilon_n^2 T \log \Lambda \).

We also require two technical lemmas - Lemmas 17 and 18, which are stated and proved in the Appendix.

**Proof:** (Of Lemma 4) From the definition of effective distance, if the condition: \( d_v(\mu_i, \mu_j) > 49c^2 \epsilon_n^2 T \log \Lambda \) holds then there are at least \( 49T \log \Lambda \) vectors \( z \) with total squared norm at least \( 98w_{\min}^2 \epsilon_n^2 T^2 \log \Lambda \). In the sequel we will scale down each vector \( z \) with norm greater than \( c_\epsilon \sqrt{w_{\min}} \) so that its norm is exactly \( c_\epsilon \sqrt{w_{\min}} \).

We divide the vectors into \( \log n \) groups as follows: group \( B_k \) contains vectors which have norm between \( c_\epsilon \sqrt{w_{\min}} \) and \( 4c_\epsilon \sqrt{w_{\min}} \).

We will call a vector small if its norm is less than \( \frac{w_{\min}^2 \epsilon_n^2 \log \Lambda}{2 \log n} \), and otherwise, we call the vector big. We observe that there exists a set of vector \( B \) with the following properties: (1) the cardinality of \( B \) is more than \( 49T \log \Lambda \), (2) the total sum of squares of the norm of the vectors in \( B \) is greater than \( 49T \log \Lambda \), and, (3) the ratio of the norms of any two vectors in \( B \) is at most \( 2 \sqrt{\log n} \).

**Case 1:** Suppose there exists a group \( B_k \) of small vectors the squares of whose norms sum to a value greater than \( \frac{w_{\min}^2 \epsilon_n^2 \log \Lambda}{2 \log n} \). By definition, such a group has more than \( 49T \log \Lambda \) vectors, and the ratio is at most 2.

**Case 2:** Otherwise, there are at least \( 49T \log \Lambda \) big vectors. By definition, the sum of the squares of their norms exceeds \( 49T \log \Lambda \). Due to the scaling, the ratio is at most \( 2 \sqrt{\log n} \).

We scale down the vectors in \( B \) so that each vector has squared norm \( \frac{w_{\min}^2 \epsilon_n^2 \log \Lambda}{4 \log n} \) in case 1, and, squared norm \( \frac{w_{\min}^2 \epsilon_n^2 \log \Lambda}{4 \log n} \) in case 2. Due to (2) and (3), the total squared norm of the scaled vectors is at least \( 49T \log \Lambda \).

Due to (1), we can now apply Lemmas 17 and 18 on the vectors to conclude that for some constant \( a_1 \), with probability \( 1 - \Lambda^{-2T \log n} \cdot \sum_{f \in F, g \in G} \epsilon_n^2 \geq a_1 \cdot \left( \frac{w_{\min}^2 \epsilon_n^2 \log \Lambda}{4 \log^2 n} \right) \).

The above sum is the square of the Frobenius norm \( \| \hat{F}_v \hat{G}_v \|^2 \) of the matrix \( \hat{F}_v \hat{G}_v \). Since \( \hat{F}_v \hat{G}_v \) has rank at most 1, an application of Lemma 15 completes the proof, for \( \tau = O\left( \frac{w_{\min}^2 \epsilon_n^2 \log \Lambda}{4 \log^2 n} \cdot \sqrt{\log \Lambda} \right) \). \( \square \)

Next we show that a vector \( x \in \mathbb{P}_F(\hat{C}_F) \) is a direction of 0 correlation. A similar statement holds for a vector \( y \in \mathbb{P}_G(\hat{C}_G) \).

**Lemma 5** If at Step 2 of Algorithm CORR-SUBSPACE, the values of \( F \) and \( G \) are respectively \( \hat{F} \) and \( \hat{G} \), and for some \( k \), the top \( k \)-th left singular vector is \( v_k \) and the corresponding singular value \( \lambda_k \) is more than \( \tau \), then for any vector \( x \in \mathbb{P}_F(\hat{C}_F), \langle v_k, x \rangle = 0 \).

**Proof:** We first show that for any \( x \in \mathbb{P}_F(\hat{C}_F) \), and any
\[ y, x^T \hat{F}^T \hat{G} y = 0. \]
\[ x^T \hat{F}^T \hat{G} y = \sum_{i=1}^{T} w_i (P_F(\mu_i), x) \cdot (P_G(\mu_i), y) \]

Since \( x \) is in \( P_F(\tilde{C}_F) \), \( P_F(\mu_i), x \) = 0, for all \( i \), and hence \( x^T \hat{F}^T \hat{G} y = 0 \) for all \( x \) in \( P_F(\tilde{C}_F) \). We now prove the Lemma by induction on \( k \).

**Base case** \((k = 1)\). Let \( v_1 = u_1 + x_1 \), where \( u_1 \in \tilde{C}_F \) and \( x_1 \in P_F(\tilde{C}_F) \). Let \( y_{k1} \) be the top right singular vector of \( \hat{F}^T \hat{G} \), and let \(|x_1| > 0\). Then, \( \hat{F}^T \hat{G} y_{k1} = u_1^T \hat{F}^T \hat{G} y_1 + u_1/|u_1| \) is a vector of norm 1 such that \( \frac{1}{|u_1|} u_1^T \hat{F}^T \hat{G} y_1 > v_1^T \hat{F}^T \hat{G} y \), which contradicts the fact that \( v_1 \) is the top left singular vector of \( \hat{F}^T \hat{G} \).

**Inductive case.** Let \( v_k = u_k + x_k \), where \( u_k \in \tilde{C}_F \) and \( x_k \in P_F(\tilde{C}_F) \). Let \( y_{k} \) be the top \( k \)-th right singular vector of \( \hat{F}^T \hat{G} \), and let \(|x_k| > 0\). We first show that \( u_k \) is orthogonal to each of the vectors \( v_1, \ldots, v_{k-1} \). Otherwise, suppose there is some \( j \), \( 1 \leq j \leq k-1 \), such that \( \langle u_k, v_j \rangle \neq 0 \). Then, \( \langle u_k, v_j \rangle = \langle x_k, v_j \rangle + \langle u_k, v_j \rangle \neq 0 \). This contradicts the fact that \( v_k \) is a left singular vector of \( \hat{F}^T \hat{G} \). Therefore, \( v_k^T \hat{F}^T \hat{G} y_{k} = u_k^T \hat{F}^T \hat{G} y_{k} + u_k/|u_k| \) is a vector of norm 1, orthogonal to \( v_1, \ldots, v_{k-1} \) such that \( \frac{1}{|u_k|} v_k^T \hat{F}^T \hat{G} y_{k} > v_k^T \hat{F}^T \hat{G} y_k \). This contradicts the fact that \( v_k \) is the top \( k \)-th left singular vector of \( \hat{F}^T \hat{G} \). The Lemma follows. \( \square \)

**Lemma 6** The covariance matrix \( \hat{F}^T \hat{G} \) has rank at most \( T \).

**Proof:** For each distribution \( i \), define \( y_i^T \) as an \( n/2 \) dimensional vector, whose \( f \)-th element is \( \sqrt{n} \mu_f - \bar{\mu}_f \), i.e., the \( i \)-th element of \( z_f \). Similarly, for each distribution \( i \), we define \( y_i^G \) as an \( n/2 \) dimensional vector, whose \( g \)-th element is \( \sqrt{n} \mu_g - \bar{\mu}_g \), i.e., the \( i \)-th element of \( z_g \). We observe that the value of \( \hat{F}^T \hat{G} \) equals \( \sum_{i=1}^{T} y_i^T \cdot (y_i^G)^T \). As each outer product of the sum is a rank 1 matrix, the sum, i.e., the covariance matrix, has rank at most \( T \). \( \square \)

Finally, we show that if the spread of every vector in \( C \) is high, then with high probability over the splitting of coordinates in Step 1 of Algorithm CORR-SUBSPACE, the maximum directional variances of any distribution \( D_i \) in \( \tilde{C}_F \) and \( \tilde{C}_G \) are high. This means that there is enough information in both \( F \)-space and \( G \)-space for correctly clustering the distributions through distance concentration.

**Lemma 7** If every vector \( v \) in \( C \) has spread at least \( 32T \log \frac{\sigma}{\sigma^2} \), then, with constant probability over the splitting of coordinates in Step 1 of Algorithm CORR-SUBSPACE, the maximum variance along any direction in \( \tilde{C}_F \) or \( \tilde{C}_G \) is at most \( 5\sigma^2 \).

**Proof:** Of Lemma 7 Let \( v \) and \( v' \) be two unit vectors in \( C \), and let \( v_F \) (resp. \( v'_F \)) and \( v_G \) (resp. \( v'_G \)) denote the normalized projections of \( v \) (resp. \( v' \)) on \( F \)-space and \( G \)-space respectively. If \(|v_F - v'_F|^2 < \frac{\sigma^2}{2} \), then, the directional variance of any \( D_i \) in the mixture along \( v'_F \) can be written as:
\[ \mathbb{E}[(v'_F, x - \mathbb{E}[x])^2] = \mathbb{E}[(v_F, x - \mathbb{E}[x])^2] + 2\mathbb{E}[(v_F, x - \mathbb{E}[x])\mathbb{E}[(v'_F - v_F, x - \mathbb{E}[x])]] \leq \mathbb{E}[(v_F, x - \mathbb{E}[x])^2] + ||v_F - v'_F||^2\sigma^2 \]
Thus, the directional variance of any distribution in the mixture along \( v' \) is at most the directional variance along \( v \) plus an additional \( \sigma^2 \). Therefore, to show this lemma, we need to show that if \( v \) is any vector on a \( \frac{\sigma}{2\sqrt{T}} \)-cover of \( C \), then with high probability over the splitting of coordinates in Step 1 of Algorithm CORR-SUBSPACE, the directional variances of any \( D_i \) in the mixture along \( v_F \) and \( v_G \) are at most \( 4\sigma^2 \).

We show this in two steps. First we show that for any \( v \) in a \( \frac{\sigma^2}{4T} \)-cover of \( C \), \( \frac{1}{T} \leq \sum_{f \in F} (v_f)^2 \leq \frac{1}{4} \). Then, we show that this condition means that for this vector \( v \), the maximum directional variances along \( v_F \) and \( v_G \) are at most \( 4\sigma^2 \).

Let \( v \) be any fixed unit vector in \( C \). We first show that with probability \( 1 - \left(\frac{\sigma}{2\sqrt{T}}\right)^{2T} \) over the splitting of coordinates in Step 1 of Algorithm CORR-SUBSPACE, \( \frac{1}{T} \leq \sum_{f \in F} (v_f)^2 \leq \frac{1}{4} \). To show this bound, we apply the Method of Bounded Differences (Theorem 13 in the Appendix). Since we split the coordinates into \( F \) and \( G \) uniformly at random, \( \mathbb{E}[(\sum_{f \in F} (v_f)^2)] = \frac{1}{2} \). Let \( \gamma_F \) be the change in \( \sum_{f \in F} (v_f)^2 \) when the inclusion or exclusion of coordinates \( f \) in the set \( F \) changes. Then, \( \gamma_F = \langle v_f \rangle^2 \) and \( \gamma = \sum_f \gamma_F \). Since the spread of vector \( v \) is at least \( 32T \log \frac{\sigma}{\sigma^2} \), \( \gamma = \sum_f \gamma_f \leq \frac{1}{4T} \log \frac{1}{\sigma^2} \), and from the Method of Bounded Differences,
\[ \Pr[||\sum_{f \in F} (v_f)^2 - \mathbb{E}[(\sum_{f \in F} (v_f)^2)] > \frac{1}{4}] \leq e^{-1/32\gamma} \]
By taking an union bound over all \( v \) on a \( \frac{\sigma}{2\sqrt{T}} \)-cover of \( C \), we deduce that for any such \( v, \frac{1}{4} \leq \sum_{f \in F} (v_f)^2 \leq \frac{1}{4} \).

Since the maximum directional variance of any distribution \( D_i \) in the mixture in \( C \) is at most \( \sigma^2 \),
\[ \sum_f (v_f)^2 \leq \frac{\sigma^2}{4T} \]
Therefore the maximum variance along \( v_F \) as well as \( v_G \) can be computed as:
\[ ||v_F||^2 \sum_{f \in F} (v_f)^2 \leq \frac{1}{||v_F||^2} \sum_f (v_f)^2 \leq 4\sigma^2 \]

The lemma follows. \( \square \)
4.2 Working with Real Samples

In this section, we show that given sufficient samples, the properties of the matrix $F^T G$, where $F$ and $G$ are generated by sampling in Step 2 of Algorithm CORR-CLUSTER are very close to the properties of the matrix $F^T G$. The lemmas are stated below, and most of the proofs are in the Appendix. The proofs use the Method of Bounded Differences (when the input is a mixture of binary product distributions) and the Gaussian Concentration of Measure Inequality (for axis-aligned Gaussians).

The central lemma of this section is Lemma 8, which shows that, if there are sufficiently many samples, for any set of $2m$ vectors, $\{v_1, \ldots, v_m\}$ and $\{y_1, \ldots, y_m\}$, $\sum_k v_k^T F^T G y_k$ and $\sum_k v_k^T F^T G y_k$ are very close. This lemma is then used to prove Lemmas 9 and 10. Lemma 9 shows that the top few singular vectors of $F^T G$ output by Algorithm CORR-SUBSPACE have very low projection on $P_\mathcal{C}(\mathcal{C})$ or $P_\mathcal{G}(\mathcal{G})$. Lemma 10 shows that the rank of the matrix $F^T G$ is almost $T$, in the sense that the $T+1$th singular value of this matrix is very low.

**Lemma 8** Let $U = \{u_1, \ldots, u_m\}$ and $Y = \{y_1, \ldots, y_m\}$ be any two sets of orthonormal vectors, and let $F$ and $G$ be the matrices generated by sampling in Step 2 of the algorithm. If the number of samples $|S|$ is greater than $\Omega\left(\frac{m^2 n^2 \log n \log(\sigma_{\max}/\delta)}{\delta^2}\right)$ (for Binary Product Distributions), and $\Omega\left(\max\left(\frac{\sigma_{\max}^2 m^2 \log^2 n \log(\sigma_{\max}/\delta)}{\delta^2}\right)\right)$ (for axis-aligned Gaussians), then, with probability at least $1 - 1/n$, 
\[
|\sum_k u_k^T (F^T G - E[F^T G]) y_k| \leq \delta
\]

**Lemma 9** Let $F$ and $G$ be the matrices generated by sampling in Step 2 of the algorithm, and let $v_1, \ldots, v_m$ be the vectors output by the algorithm in Step 4. If the number of samples $|S|$ is greater than $\Omega\left(\frac{m^2 n^2 \log n \log(\Lambda + \log 1/\delta)}{\delta^2}\right)$ (for Binary Product Distributions), and $\Omega\left(\max\left(\frac{\sigma_{\max}^2 m^2 \log^2 n \log(\Lambda/\delta)}{\delta^2}\right)\right)$ (for axis-aligned Gaussians), then, for each $k$, and any $x \in P_\mathcal{F}(\mathcal{C})$, $\langle v_k, x \rangle \leq \epsilon$.

**Lemma 10** Let $F$ and $G$ be the matrices generated by sampling in Step 2 of Algorithm CORR-SUBSPACE. If the number of samples $|S|$ is greater than $\Omega\left(\frac{T^2 \sigma_{\max}^2 m^2 \log^2 n \log(\Lambda/\delta)}{\delta^2}\right)$ (for binary product distributions) and $\Omega\left(\max\left(\frac{\sigma_{\max}^2 m^2 \log^2 n \log(\Lambda/\delta)}{\delta^2}\right)\right)$ for axis-aligned Gaussians, then, $\lambda_{T+1}$, the $T+1$th singular value of the matrix $F^T G$ is at most $\tau/8$.

4.3 The Combined Analysis

In this section, we combine the lemmas proved in Sections 4.1 and 4.2 to prove Theorem 1.

We begin with a lemma which shows that if every vector in $\mathcal{C}$ has spread $32 T \log \frac{\tau}{\delta}$, then the maximum directional variance in $\mathcal{C}$, the space output by Algorithm CORR-SUBSPACE, is at most $11 \sigma^2$.

**Lemma 11** Let $\mathcal{C}$ be the subspace output by the algorithm, and let $v$ be any vector in $\mathcal{C}$. If every vector in $\mathcal{C}$ has spread $32 T \log \frac{\tau}{\delta}$, then the number of samples $|S|$ is greater than $\Omega\left(\max\left(\frac{\sigma_{\max}^2 m^2 \log^2 n \log(\Lambda/\delta)}{\delta^2}\right)\right)$ then for any $i$ the maximum variance of $D_i$ along $v$ is at most $11 \sigma^2$.

**Proof** (Of Lemma 11) We first show that for any $v_k$ (or $y_k$) in the set $\mathcal{K}$ output by Algorithm CORR-SUBSPACE, and for any distribution $D_i$ in the mixture, the maximum variance of $D_i$ along $v_k$ (or $y_k$) is at most $11 \sigma^2$.

Let $v_k = u_k + x_k$ where $u_k$ is in $\mathcal{C}$ and $x_k$ is in $P_{\mathcal{F}}(\mathcal{C})$. From Lemma 9, we deduce that $||x_k|| \leq \frac{25}{4\sigma} \Delta$. Let $M_i$ be the matrix in which each row is a sample from the distribution $D_i$. Then the variance of distribution $D_i$ along the direction $v$ is the square of the norm of the vector $(M_i - E[M_i])v$. This norm can be written as:

\[
||M_i - E[M_i]v_k||^2 = ||M_i - E[M_i]v_k||^2 + ||(M_i - E[M_i])x_k||^2 + 2\langle (M_i - E[M_i])x_k, (M_i - E[M_i])v_k \rangle \leq 5\sigma^2 ||x_k||^2 + \sigma^2 ||x_k||^2 + 10 \sigma \lambda_1 ||x_k||^2 \leq 2(5\sigma^2 + ||x_k||^2) \leq 11 \sigma^2
\]

The third line follows from Lemma 7, and the last step follows from the bound on $||x_k||$ from Lemma 9. Now, let $v = \sum \alpha_i v_k + \alpha_i y_k$ be any unit vector in the space $\mathcal{K}$. Then, $\sum \alpha_i^2 + \alpha_i^2 = 1$, and for any $i$, the variance of $D_i$ along $v$ is

\[
\sum \alpha_i^2 ||(M_i - E[M_i])v_k||^2 + \alpha_i^2 ||(M_i - E[M_i])y_k||^2 \leq 11 \sigma^2
\]

The lemma follows. $\square$

The above Lemmas are now combined to prove Theorem 1.

**Proof** (Of Theorem 1) Suppose $\mathcal{K} = \mathcal{K}_L \cup \mathcal{K}_R$, where $\mathcal{K}_L = \{v_1, \ldots, v_m\}$, the top $m$ left singular vectors of $F^T G$ and $\mathcal{K}_R = \{y_1, \ldots, y_m\}$ are the corresponding right singular vectors. We abuse notation and use $v_k$ to denote the vector $v_k$ concatenated with a vector consisting of $n/2$ zeros, and use $y_k$ to denote the vector consisting of $n/2$ zeros concatenated with $y_k$. Moreover, we use $\mathcal{K}$, $\mathcal{K}_L$, and $\mathcal{K}_R$ interchangeably to denote sets of vectors and the subspace spanned by those sets of vectors.
We show that with probability at least $1 - \frac{1}{T}$ over the splitting step, there exists no vector $v \in C_{F \cup G}$ such that $(1)$ $v$ is orthogonal to the space spanned by the vectors $K$ and $(2)$ there exists some pair of centers $i$ and $j$ such that $d(v, \mu_i, \mu_j) > 49T c_{ij}^2 \log \Lambda$. For contradiction, suppose there exists such a vector $v$.

Then, if $v_F$ and $v_G$ denote the normalized projections of $v$ onto $F$-space and $G$-space respectively, from Lemma 3, $v_F T F^T G v_G \geq \tau$ with probability at least $1 - \frac{1}{T}$ over the splitting step. From Step 8, if the number of samples $|S|$ is greater than $\Omega\left(T^2 \frac{n^2 \log n \log \Lambda}{\tau^2}\right)$ for binary product distributions, and if $|S|$ is greater than $\Omega\left(\max\left(\frac{\alpha^2 n^2 \log n \log \Lambda}{\tau^2}, \frac{\sigma^2 n^2 \log n \log \Lambda}{\tau^2}\right)\right)$ for axis-aligned Gaussians, $v_F T F^T G_{ij} x \geq \frac{\tau}{2}$ with at least constant probability. Thus, $v$ is orthogonal to the space spanned by $K$, $v_F$ is orthogonal to $K_L$, and $v_G$ is orthogonal to $K_R$. As $\lambda_{\alpha+1}$ is the maximum value of $x F^T G y$ over all vectors $x$ orthogonal to $K_L$ and $y$ orthogonal to $K_R$, $\lambda_{\alpha+1} \geq \frac{\tau}{2}$, which is a contradiction. Moreover, from Lemma 10, $\lambda_{\alpha+1} < \frac{\tau}{2}$, and hence $m \leq T$.

Let us construct an orthonormal series of vectors $v_1, \ldots, v_m$, where $m$ is a constant. To extend the vectors output by Algorithm CORMUTSUBSPACE. We inductively define $v_l$ as follows. Suppose for each $k$, $v_k = u_k + x_k$, where $u_k \in C_F$ and $x_k \in P_F(C_F)$. Let $u_l$ be a unit vector in $C_F$ which is perpendicular to $u_1, \ldots, u_{l-1}$. Then, $v_l = u_l$. By definition, this vector is orthogonal to $u_1, \ldots, u_{l-1}$. In addition, for any $k \neq l$, $\langle v_k, v_l \rangle = \langle u_k, u_l \rangle + \langle u_k, x_l \rangle = 0$, and $v_l$ is also orthogonal to $v_1, \ldots, v_{l-1}$. Moreover, if $\epsilon < \frac{1}{T}$ and $v_1, \ldots, v_m$ are linearly independent, we can always find $\dim(C_F)$ such vectors. Similarly, we construct a set of vectors $y_1, y_2, \ldots$. Let us call the combined set of vectors $C^*$.

We now show that if there are sufficient samples, $d(C^*, \mu_i, \mu_j) \leq c_{ij}^2$. Note that for any unit vector $v^*$ in $C^*$, and any unit $x \in C_{F \cup G}$, $\langle v, x \rangle \leq me$. Also, note that for any $u_k$ and $u_i$, $k \neq i$, $\langle u_k, u_i \rangle \leq 2$, and $\|u_k\|^2 \geq l - \epsilon^2$. Let $v = \sum_k \alpha_k x_k$ be any unit vector in $C_{F \cup G}$, and $u_l = \|v\|^2 = \sum_{k,l} \alpha_k \alpha_l \langle u_k, u_l \rangle = \sum_k \alpha_k^2 \|u_k\|^2 - \Omega(T^2 \epsilon^2)$. The projection of $v$ on $C^*$ can be written as:

$$\sum_k \langle v, v_k \rangle^2 = \sum_k \langle v, x_k \rangle^2 = \sum_k \alpha_k^2 \|u_k\|^2 - \Omega(T^2 \epsilon^2)$$

The last step follows because for each $k$, $\|u_k\|^2 \geq 1 - \epsilon^2$. If the number of samples $|S|$ is greater than $\Omega\left(\frac{m^2 n^2 \log n (\log \Lambda + \log 100T)}{\epsilon^2}\right)$ (for Binary Product Distributions), and $\max\left(\frac{\alpha^2 m^2 \log n (\log 100T \Lambda)}{\epsilon^2}, \frac{\sigma^2 m^2 \log n (\log 100T \Lambda)}{\epsilon^2}\right)$ (for axis-aligned Gaussians), then $\epsilon < 1/100T$. Therefore,

$$d(C^*, \mu_i, \mu_j) \leq \frac{1}{100} d(\mu_i, \mu_j)$$

For any $i, j$.

Since vectors $v_{n+1}, \ldots, n_0, \ldots, v_{y_{n+1}}, \ldots, v_{y_{n+1}}$ all belong to $C_{F \cup G}$ (as well as $C^* \setminus K$), there exists no $v \in C^* \setminus K$ with the Conditions (1) and (2) in the previous paragraph, and $d(C^*, \mu_i, \mu_j) \leq 49T c_{ij}^2 \log \Lambda$. That is, the actual distance between $\mu_i$ and $\mu_j$ in $C_{F \cup G} \setminus K$ (as well as $C^* \setminus K$) is at most the contribution to $d(\mu_i, \mu_j)$ from the top $49T c_{ij}^2 \log \Lambda$ coordinates, and the contribution to $d(\mu_i, \mu_j)$ from $K$ and $C^*$ is at least the contribution from the rest of the coordinates. Since $d(C^*, \mu_i, \mu_j) \leq \frac{1}{100} d(\mu_i, \mu_j)$, the distance between $\mu_i$ and $\mu_j$ in $K$ is at least $\frac{99}{100} d(\mu_i, \mu_j) - 49T \log \Lambda c_{ij}^2$. The first part of the theorem follows.

The second part of the theorem follows directly from Lemma 11.


5 Appendix

The Appendix is organised as follows. First, in Section 5.1, we state some inequalities from probability and linear algebra that we use in our proofs. In Section 5.2, we state some of the proofs that did not fit into the main body of the paper due to space constraints. Finally, in Section 5.3, we show how distance concentration lemmas from [AM05] can be combined with Theorem 1 to prove Theorem 2.

5.1 Some Inequalities

We use the following Lemma, which follows easily from Markov’s Inequality.

Lemma 12 Let $X$ be a random variable such that $0 \leq X \leq \Gamma$. Then, $\Pr[X \leq E[X]/2] \leq \frac{2\Gamma - 2E[X]}{2\Gamma - E[X]}$.

Proof: (Of Lemma 12) Let $Z = \Gamma - X$. Then, $Z > 0$, $E[Z] = \Gamma - E[X]$, and when $X < E[X]/2$, $Z > \Gamma - E[X]/2$. We can apply Markov’s Inequality on $Z$ to
conclude that for any $\alpha$,
\[
\Pr[Z \geq \alpha E[Z]] \leq \frac{1}{\alpha}
\]

The lemma follows by plugging in $\alpha = \frac{\Gamma - E[X]}{T - E[X]}$. \hfill $\square$

In Section 4.2, we use the Method of Bounded Differences [PD05], and the Gaussian Concentration of Measure Theorem, stated below.

**Theorem 13 (Method of Bounded Differences)** Let $X_1, \ldots, X_n$ be arbitrary independent random variables and let $f$ be any function of $X_1, \ldots, X_n$. For each $i$, if for any $a$ and $a'$,
\[
|f(X_1, \ldots, X_i = a, \ldots, X_n) - f(X_1, \ldots, X_i = a', \ldots, X_n)| \leq \gamma_i
\]
and $\gamma = \sum_i \gamma_i^2$, then,
\[
\Pr[|f(X_1, \ldots, X_n) - E[f(X_1, \ldots, X_n)]| > t] \leq 2e^{-t^2/2\gamma}
\]

**Theorem 14** [Led00] Let $F(X_1, \ldots, X_n)$ be any function of independent Gaussian variables $X_1, \ldots, X_n$, such that $\nabla F \leq \gamma$, for any value of $X_1, \ldots, X_n$. Then,
\[
\Pr[|F(t) - E[F(t)]| > t] \leq 2e^{-t^2/2\gamma}
\]

Finally, we use the following theorem, which relates the Frobenius norm and the top singular value of any matrix.

**Theorem 15** ([GL96]) For a matrix $M$ of rank $r$, the top singular value is greater than or equal to $|M|_F / \sqrt{r}$.

5.2 Proofs

Proof (Of Lemma 3) From Lemma 4, for a fixed vector $v \in C_{\mathcal{F} \cup \mathcal{G}}$, $v^T F^T \hat{G} v \geq 2\tau$, with probability $1 - \Lambda^{-2t}$. Let $u$ be a vector, and $u_{\mathcal{F}}$ and $u_{\mathcal{G}}$ be its normalized projection on $\mathcal{F}$-space and $\mathcal{G}$-space respectively, such that $\|u_{\mathcal{F}} - v_{\mathcal{F}}\| < \frac{1}{2}$ and $\|u_{\mathcal{G}} - v_{\mathcal{G}}\| < \frac{1}{2}$. Then,
\[
u^T F^T \hat{G} v_{\mathcal{G}} = v^T F^T \hat{G} v_{\mathcal{G}} + (u^T - v^T) F^T \hat{G} u_{\mathcal{G}}
\geq v^T F^T \hat{G} u_{\mathcal{G}} - \delta \sigma_{\max}
\]
\]
as $\sigma_{\max}$ is the maximum value of $x^T F^T \hat{G} y$, over all unit vectors $x$ and $y$. Since $v^T F^T \hat{G} u_{\mathcal{G}} > 2\tau$, if $\delta < \frac{\tau}{\sigma_{\max}}$, $v^T F^T \hat{G} u_{\mathcal{G}} > \tau$. To show that the lemma holds for all $v \in C_{\mathcal{F} \cup \mathcal{G}}$, it is therefore sufficient to show that the statement holds for all unit vectors in a $(\frac{\tau}{\sigma_{\max}})$-cover of $C_{\mathcal{F} \cup \mathcal{G}}$. Since $C_{\mathcal{F} \cup \mathcal{G}}$ has dimension at most $2T$, there are at most $(\frac{\tau}{\sigma_{\max}})^{2T}$ vectors in such a cover. As $\Lambda > \frac{\tau}{\sigma_{\max}}$, the lemma follows from Lemma 4 and a union bound over the vectors in a $(\frac{\tau}{\sigma_{\max}})$-cover of $C_{\mathcal{F} \cup \mathcal{G}}$. \hfill $\square$

Now we formally prove Lemma 4. The proof requires Lemmas 17, 18 and 16, which we prove first.

**Lemma 16** If $A$ is a set of coordinates whose cardinality $|A|$ is greater than $T$, such that for each $f \in A$, the norms $|z_f|$ are equal and $\sum_{f \in A} |z_f|^2 = D$, then
\[
\sum_{f,g \in A, f \neq g} \langle z_f, z_g \rangle^2 \geq \frac{D^2}{2T^2}(1 - \frac{T}{|A|})
\]

Proof (Of Lemma 16) The proof uses the following strategy. First, we group the coordinates into sets $B_1, B_2, \ldots$ such that for each $B_k$, the set $\{z_f \mid f \in B_k\}$ approximately form a basis of the $T$ dimensional space spanned
by the set of vectors \( \{ z_f \mid f \in A \} \). Next, for each basis \( B_k \) we estimate the value of \( \sum_{f \in B_k} \sum_{g \in B_k} \langle z_f, z_g \rangle^2 \) and finally, we prove the Lemma, by using the estimates, and the fact that: \( \sum_{f \in A, f \neq g} \langle z_f, z_g \rangle^2 \) can be written as \( \sum_{B_1} \sum_{f \in B_1} \sum_{g \in B_1} \langle z_f, z_g \rangle^2 \). Consider the following procedure which returns a series of sets of vectors, \( B_1, B_2, \ldots \).

\[
S = A; \ k = 0 \\
\text{while } S \neq \emptyset \\
\quad k = k + 1; \ B_k = \emptyset; \\
\quad \text{for each } z \in S \\
\quad \quad \text{if } \|z\|^2 \leq \text{projection of } z \text{ onto } B_k \text{ then } \\
\quad \quad \quad \text{remove } z \text{ from } S, \text{ add } z \text{ to } B_k \\
\]

We observe that by construction, for each \( k \),

\[
\sum_{f \in B_k, g \in B_k, k' > k} \langle z_f, z_g \rangle^2 \geq \frac{D}{|A|} \sum_{g \in B_k, k' > k} \frac{\|z_g\|^2}{2}
\]

This inequality follows since the basis is formed of vectors of length \( |D|/|A| \) and the fact that the projection of each \( z_g \) on \( B_k \) preserves at least half of its norm.

Since the vectors in \( T \) dimensions there are at least \( |A|/T \) sets \( B_k \). Notice that for all \( f \in A \), we have \( \|z_f\|^2 = |D|/|A| \), as all the \( \|z_f\| \) are equal. Hence, the average contribution to the sum from each \( B_k \) is at least \( \frac{|A|-T}{|A|} \cdot \frac{D^2}{2} \), from which the Lemma follows. \( \square \)

**Lemma 17** Let \( A \) be a set of coordinates with cardinality more than \( 144T^2 \log \Lambda \) such that for each \( f \in A \), \( \|z_f\| \) is equal and \( \sum_{f \in A} \|z_f\|^2 = D \). Then, \( \sum_{f \in A, f \neq g} \langle z_f, z_g \rangle^2 \geq \frac{288T^2 \log \Lambda}{D^2} \) and \( \|z_f\|^2 \geq \frac{\log \Lambda}{1152T^2} \log \Lambda \).

**Proof:** (Of Lemma 17) We can partition \( A \) into \( 72T \log \Lambda \) groups \( A_1, A_2, \ldots \) each with at least \( 2T \) coordinates, such that for each group \( A_k \), \( \sum_{f \in A_k} \|z_f\|^2 \geq \frac{D}{2} \). From Lemma 16, for each such group \( A_k \),

\[
\sum_{f, f' \in A_k, f \neq f'} \langle z_f, z_{f'} \rangle^2 \geq \left( \frac{D}{2(72T \log \Lambda)} \right)^2 \cdot \frac{1}{4T} \\
\geq \frac{D^2}{20736T^3 \log \Lambda}
\]

Summing over all the groups,

\[
\sum_{f, f' \in A, f \neq f'} \langle z_f, z_{f'} \rangle^2 \geq \sum_{k} \sum_{f, f' \in A_k, f \neq f'} \langle z_f, z_{f'} \rangle^2 \\
\geq \frac{D^2}{288T^2 \log \Lambda}
\]

from which the first part of the lemma follows.

For each group \( A_k \), \( \sum_{f \in \mathcal{F} \cap A_k} \sum_{g \in \mathcal{G} \cap A_k} \langle z_f, z_g \rangle^2 \) is a random variable whose value depends on the outcome of the splitting in Step 1 of the algorithm. The maximum value of this random variable is \( \sum_{f \in \mathcal{F} \cap A_k, f \neq g} \langle z_f, z_g \rangle^2 \), and the expected value is \( \frac{1}{2} \sum_{f \in \mathcal{F} \cap A_k, f \neq g} \langle z_f, z_g \rangle^2 \). Therefore, by Lemma 12, with probability \( \frac{1}{2} \),

\[
\sum_{f \in \mathcal{F} \cap A_k, g \in \mathcal{G} \cap A_k} \langle z_f, z_g \rangle^2 \geq \frac{1}{4} \sum_{f \in \mathcal{F} \cap A_k, f \neq g} \langle z_f, z_g \rangle^2 \\
\geq \frac{D^2}{82944T^3 \log \Lambda}
\]

Moreover, as each group \( A_k \) is disjoint, the splitting process in each group is independent. Since there are \( 72T \log \Lambda \) groups, using the Chernoff bounds, we can conclude that with probability \( 1 - \Lambda^{-2T} \), for at least \( \frac{6}{7} \) fraction of the groups \( A_k \),

\[
\sum_{f \in \mathcal{F} \cap A_k, g \in \mathcal{G} \cap A_k} \langle z_f, z_g \rangle^2 \geq \frac{D^2}{82944T^3 \log \Lambda}
\]

, from which the second part of the lemma follows. \( \square \)

**Lemma 18** Let \( A \) be a set of coordinates such that for each \( f \in A \), \( \|z_f\| \) is equal and \( \sum_{f \in A} \|z_f\|^2 = D \). If \( 48T \log \Lambda + T < |A| \leq 144T^2 \log \Lambda \), then \( (1) \),

\[
\sum_{f \in \mathcal{F} \cap A, g \in \mathcal{G} \cap A} \langle z_f, z_g \rangle^2 \geq \frac{D^2}{1152T^4 \log \Lambda}
\]

and \( (2) \) with probability \( 1 - \Lambda^{-2T} \) over the splitting in Step 1,

\[
\sum_{f \in \mathcal{F} \cap A} \langle z_f, z_g \rangle^2 \geq \frac{D^2}{4608T^4 \log \Lambda}
\]

**Proof:** (Of Lemma 18) The proof strategy here is to group the coordinates in \( A \) into tuples \( A_k = \{ f_k, f_k' \} \) such that for all but \( T \) such tuples, \( \langle z_f, z_{f'} \rangle^2 \) is high. For each \( A_k \) we estimate the quantity \( \langle z_{f_k}, z_{f_k'} \rangle^2 \). The proof of the first part of the lemma follows by using the estimates and the fact that

\[
\sum_{f, f' \in A, f \neq f'} \langle z_f, z_{f'} \rangle^2 \geq \sum_{k} \langle z_{f_k}, z_{f_k'} \rangle^2
\]

Consider the following procedure, which returns a series of sets of vectors \( A_1, A_2, \ldots \).

\[
S = A; \ k = 0 \\
\text{while } S \neq \emptyset \\
\quad k = k + 1; \ A_k = \emptyset; \\
\quad \text{if there exists } z_{f_k}, z_{f_k'} \in S \text{ with } \langle z_{f_k}, z_{f_k'} \rangle^2 \geq \frac{D}{2T|A|}, \text{ then } \\
\quad \quad A_k = \{ f_k, f_k' \} \\
\quad \quad \text{remove } z_{f_k}, z_{f_k'} \text{ from } S
\]

\[14\]
follows.
there are the splitting process in each group is independent. As
B vectors D unit vectors such that
| \langle f_k, \mathbf{P}_G(y) \rangle | \leq \frac{\delta}{2} occurs with high probability for all U and Y where U and Y contain vectors from a \( \frac{\delta}{8\sigma_{\max}} \)-cover of \( \mathbb{R}^n \).

Let us consider a fixed set of vectors U and Y. We can write
\[
\sum_k u_k^T F^T G y_k = \frac{1}{|S|} \sum_k \sum_{y \in S} \langle u_k, \mathbf{P}_F(x) \rangle \cdot \langle y_k, \mathbf{P}_G(y) \rangle
\]

We now apply concentration inequalities to bound the deviation of \( \sum_k u_k^T F^T G y_k \) from its mean. For

Since the vectors in S have dimension at most T, if |S| > T, there is at least one vector \( z \notin B \). By construction, the projection of any such \( z \) on B preserves at least half its norm. Therefore, there exists some vector \( z_f \in B \) such that \( \langle z, z_f \rangle \geq \frac{D^2}{2|A|} \), which is a contradiction to the failure of Step 4 of our procedure. Thus,
\[
\sum_k (z_{f_k}, z_{f_k}^T) \geq \frac{D^2}{4T^2|A|} \geq \frac{D^2}{8T^2|A|} > \frac{1}{152T^2 \log \Lambda}
\]
from which the first part of the lemma follows.

With probability \( \frac{1}{2} \) over the splitting in Step 1 of the algorithm, for any tuple \( A, f_k \) and \( f_k \) will belong to different sets, and thus contribute \( \frac{D^2}{4T^2|A|} \) to \( \sum_{f \in F \cap A, y \in G \cap A} (z_f, z_f^T) \). Since the groups are disjoint, the splitting process in each group is independent. As there are \( |A|^2 \) groups, and \( |A| \geq 48T^2d \), we can use the Chernoff Bounds to conclude that with probability at least \( 1 - \Lambda^{-2T} \), at least \( 1/4 \) fraction of the tuples contribute to the sum, from which the second part of the lemma follows.

5.2.1 Working with Real Samples

**Proof:** (Of Lemma 8) Let \( u, v \) and \( u', v' \) be two pairs of unit vectors such that \( |\langle u - u', v \rangle| < \frac{\delta}{2\sigma_{\max}} \), and \( |\langle y - y' \rangle| < \frac{\delta}{2\sigma_{\max}} \). Then, if \( \delta < \frac{\sigma_{\max}}{2\sigma_{\max}} \),
\[
|u^T F^T G y - u'^T F^T G y'| \\
\leq |(u - u')^T F^T G y + u'^T F^T G (y - y')| \\
\leq |(u - u')| \sigma_{\max} + |y - y'| \sigma_{\max} < \frac{\delta}{2}
\]
The second line follows as \( \sigma_{\max} = \max_{x, y} u^T F^T G y \). It is therefore sufficient to show that the event \( \sum_k u_k^T (F^T G - E[F^T G]) y_k \) holds for all sets of \( m \) unit vectors U and Y on a \( \frac{\delta}{2\sigma_{\max}} \)-cover of \( \mathbb{R}^n \). We show this by showing that if \(|S|\) is large enough, the

For axis-aligned Gaussians, we use a similar proof involving the Gaussian Concentration of Measure Theorem (Theorem 14 in the Appendix).

If a sample x is generated by distribution \( D_i \),
\[
\langle y_k, \mathbf{P}_G(x) \rangle = \langle y_k, \mathbf{P}_G(\mu) \rangle + \langle y_k, \mathbf{P}_G(x - \mu) \rangle
\]
Since \( y_k \) is a unit vector, and the maximum directional variance of any \( D_i \) is at most \( \sigma^2 \), \( \langle y_k, \mathbf{P}_G(x - \mu) \rangle \) is distributed as a Gaussian with mean 0 and standard deviation at most \( \sigma \). Therefore for each \( x \in S \),
\[
\Pr[|\langle y_k, \mathbf{P}_G(x - \mu) \rangle| > \sigma \sqrt{mn \log(\sigma_{\max}/\delta) \log n}] \\
\leq \left( \frac{\delta}{8\sigma_{\max}} \right)^{2mn \log n}
\]

and the probability that this happens for all $x \in S$ is at most $\frac{1}{m} \cdot \left(\frac{\delta}{\sigma_{\max}}\right)^{2mn}$, provided $\delta < \sigma_{\max}$ and $|S|$ is polynomial in $n$. Thus, for any $x \in S$ generated from distribution $D_x$, and any $k$, $(\langle y_k, P_G(x) \rangle)^2 \leq 2(\langle y_k, \mu \rangle)^2 + 8\sigma^2 mn \log(\sigma_{\max}/\delta) \log n$.

Therefore, for a specific $x \in S$, as $P_F(x)$ and $P_G(x)$ are independently distributed, except with very low probability, the distribution of $(\langle u_k, P_F(x) \rangle (\langle y_k, P_G(x) \rangle)$ is dominated by the distribution of $\sqrt{2(\langle y_k, \mu \rangle)^2 + 8\sigma^2 mn \log(\sigma_{\max}/\delta) \log n} (\langle u_k, P_F(x) \rangle)$. Note that $(\langle u_k, P_F(x) \rangle)$ is distributed as a Gaussian with mean $(\langle u_k, \mu \rangle)$ and variance at most $\sigma^2$.

Let $\gamma_x$ be the derivative of $\sum_{x \in S} \sum_k \sqrt{2(\langle y_k, \mu \rangle)^2 + 8\sigma^2 mn \log(\sigma_{\max}/\delta) \log n} \cdot (\langle u_k, P_F(x) \rangle)$ with respect to the value of sample $x \in S$. Then, the gradient $\gamma$ of this expression is at most:

$$
\gamma^2 = \sum_x \gamma_x^2 \\
\leq \frac{1}{|S|^2} \sum_{x \in S} 2\sigma^2 mn (\langle y_k, \mu \rangle)^2 + 4\sigma^2 mn \log(\sigma_{\max}/\delta) \log n \\
+ \frac{1}{|S|^2} (8\sigma^4 m^3 n \log(8\sigma_{\max}/\delta) \log n \\
+ m^2 \sigma^2 \sigma_{\max}^2) |S| \\
\leq \frac{m^2 \sigma^2}{|S|} (8\sigma^4 m^3 n \log(8\sigma_{\max}/\delta) \log n + 2\sigma_{\max}^2).
$$

Applying the Gaussian Concentration of Measure Theorem,

$$
\Pr \left[ \left| \sum_k u_k^T (F^T G - E[F^T G]) y_k \right| \right] \\
\leq \frac{1}{|S|} \left( 2\sigma^2 mn \log(8\sigma_{\max}/\delta) \log n + \sigma_{\max}^2 \right) \\
\leq e^{-t^2/2}
$$

Plugging in $t = \sqrt{2mn \log(\sigma_{\max}/\delta) \log n}$, we see that the quantity on the left-hand side is at most $\delta$ when $|S| \geq \Omega \left( \max \left( \frac{\sigma^4 m^2 n^2 \log^2 n \log^2(\sigma_{\max}/\delta)}{\delta^2}, \frac{\sigma^2 m^2 \sigma_{\max}^4 m^3 n \log(\sigma_{\max}/\delta)}{\delta^2} \right) \right)$

In both cases, the lemma now follows by applying a Union Bound over $\left(8\sigma_{\max}\right)^{2mn}$ choices of $U$ and $Y$. □

Proof (Of Lemma 9) Let $\{v_1, \ldots, v_m\}$ and $\{y_1, \ldots, y_m\}$ be the top $m$ left and right singular vectors, and let $\{\lambda_1, \ldots, \lambda_m\}$ be the top $m$ singular values of $E[F^T G]$. Let $\{y_1, \ldots, y_m\}$ be the top right singular vectors of $F^T G$. If, for any $k$ and any $x \in P_F(C_\tau)$, $(\langle v_k, x \rangle \leq \epsilon$, then, for any set of orthonormal vectors $\{y_1, \ldots, y_m\}$,

$$
\sum_{k=1}^m v_k^T E[F^T G] y_k = \sum_{k \neq l} v_k^T E[F^T G] y_k + (\langle v_k, (\langle v_l, x \rangle) \rangle E[F^T G] y_k \\
\leq \sum_{k=1}^m \lambda_k - (1 - \sqrt{1 - \epsilon^2}) \lambda_m
$$

If $\epsilon$ is smaller than $1/4$, the left-hand side can be bounded as

$$
\sum_{k=1}^m v_k^T E[F^T G] y_k \leq \sum_{k=1}^m \lambda_k - \frac{\epsilon^2}{4} \lambda_m
$$

Next we apply Lemma 8, which states that if the number of samples $|S|$ is large enough,

$$
\left| \sum_{k=1}^m v_k^T (F^T G - E[F^T G]) y_k \right| \leq \delta
$$

Therefore, if $\delta \leq \epsilon^2 \lambda_m/16$,

$$
\sum_{k=1}^m v_k^T F^T G y_k \leq \sum_{k=1}^m \lambda_k - \frac{\epsilon^2}{16} \lambda_m
$$

On the other hand, using Lemma 8,

$$
\sum_{k=1}^m v_k^T F^T G y_k \geq \sum_{k=1}^m \lambda_k - \frac{\epsilon^2}{16} \lambda_m
$$

This contradicts the fact that $\{v_1, \ldots, v_m\}$ and $\{y_1, \ldots, y_m\}$ are the top singular vectors of $F^T G$, and hence the lemma follows. □

Proof (Of Lemma 10) For any $k$, let $\lambda_k$ denote the $k$-th top singular value of the matrix $F^T G$ and let $\lambda_k$ denote the $k$-th top singular value of the matrix $E[F^T G]$. From Lemma 6, $E[F^T G] = F^T \tilde{G}$ has rank at most $T$. Thus $\hat{\lambda}_{T+1} = 0$.

Let $\{v_1, \ldots, v_m\}$ (resp. $\{\hat{v}_1, \ldots, \hat{v}_m\}$) and $\{y_1, \ldots, y_m\}$ (resp. $\{\hat{y}_1, \ldots, \hat{y}_m\}$) denote the top $m$ left and right singular vectors of $F^T G$ (resp. $E[F^T G]$). From Lemma 8, if $|S|$ is greater than $\Omega \left( T^{\frac{1}{2} + \frac{n^2 \log(\log N) \log N}{\sigma^2} \log(n \log N) \right)$ (for binary product distributions) and $|S|$ is greater than $\Omega \left( \max \left( \frac{\sigma^2 m^2 n^2 \log^2 \log \log N}{\delta^2}, \frac{\sigma^2 m^2 \sigma_{\max}^4 m^3 n \log(\sigma_{\max}/\delta)}{\delta^2} \right) \right)$ (for axis-aligned Gaussians), then,

$$
\hat{\lambda}_1 + \ldots + \hat{\lambda}_T \geq \sum_{k=1}^T v_k^T E[F^T G] y_k \geq \hat{\lambda}_1 + \ldots \hat{\lambda}_T - \frac{\tau}{16}
$$

Moreover, from Lemma 8,

$$
\hat{\lambda}_1 + \ldots + \hat{\lambda}_{T+1} \geq \sum_{k=1}^T v_k^T E[F^T G] y_k \geq \hat{\lambda}_1 + \ldots + \hat{\lambda}_{T+1} - \frac{\tau}{16}
$$
Combining the above two equations, and the fact that \( \lambda_{r+1} = 0, \lambda_{r+1} \leq 2 \frac{\lambda_r}{n} \leq \frac{\lambda}{n} \).

5.3 Distance Concentration

In this section, we show how to prove Theorem 2 by combining Theorem 1 and distance-concentration methods. We begin with the following distance-concentration lemmas of [AM05] and [McS01], which we prove for the sake of completeness.

**Lemma 19** Let \( K \) be a \( d \)-dimensional subspace of \( \mathbb{R}^n \), and \( x \) be a point drawn from axis-aligned Gaussian. Let \( \sigma_K^2 \) be the maximum variance of \( x \) along any direction in the subspace \( K \). Then,

\[
\Pr[|\|P_K(x - E[x])\| > \sigma_K \sqrt{2d \log(d/\delta)}] \leq \delta
\]

**Proof:** Let \( v_1, \ldots, v_d \) be an orthonormal basis of \( K \). Since the projection of a Gaussian is a Gaussian, the projection of the distribution of \( x \) along any \( v_k \) is a Gaussian with variance at most \( \sigma_K^2 \). By the properties of the normal distribution,

\[
\Pr[|\langle x - E[x], v_k \rangle | > \sigma_K \sqrt{2 \log(d/\delta)}] \leq \frac{\delta}{d}
\]

Since \( |\|P_K(x - E[x])\||^2 = \sum_k (x - E[x], v_k)^2 \), the lemma follows by a Union Bound over \( v_1, \ldots, v_d \). \qed

**Lemma 20** Let \( K \) be a \( d \)-dimensional subspace of \( \mathbb{R}^n \), and \( x \) be a point drawn from a binary product distribution. Then, \( \Pr[|\|P_K(x - E[x])\| > \sqrt{2d \log(d/\delta)}] \leq \delta \).

**Proof:** Let \( v_1, \ldots, v_d \) be an orthonormal basis of \( K \). For a fixed \( v_k \), we bound \( |\langle v_k, x - E[x] \rangle| \), where \( x \) is generated by a binary product distribution by applying the Method of Bounded Differences. Let \( \gamma_f \) be the change in \( \langle v_k, x - E[x] \rangle \), when the value of coordinate \( f \) in \( x \) changes. Then, \( \gamma_f = v_k^f \), and \( \gamma = \sum_f \gamma_f^2 = \|v_k\|^2 = 1 \). From the Method of Bounded Differences,

\[
\Pr[|\langle x - E[x], v_k \rangle | > \sqrt{2 \log(d/\delta)}] \leq \frac{\delta}{d}
\]

Since \( |\|P_K(x - E[x])\||^2 = \sum_k (v_k, x - E[x])^2 \), the lemma follows by a union bound over \( v_1, \ldots, v_d \). \qed

Now we are ready to prove Theorem 2.

**Proof:** (Of Theorem 2) From Theorem 1, if for all \( i \) and \( j \), \( d(\mu_i, \mu_j) \geq 49Tc^2_j \log \Lambda \), then, for all \( i \) and \( j \), with constant probability,

\[
d_K(\mu_i, \mu_j) \geq \frac{99}{100} (d(\mu_i, \mu_j) - 49Tc^2_j \log \Lambda)
\]

From the separation conditions in Theorem 2, this means that for all \( i \) and \( j \), we have that \( d_K(\mu_i, \mu_j) \geq 9T(\log \Lambda + \log n) \) for binary product distributions, and \( d_K(\mu_i, \mu_j) \geq 9\sigma^2T(\log \Lambda + \log n) \) for axis-aligned Gaussians.

Applying Lemma 20, for binary product distributions, any two samples from a fixed distribution \( D_i \) in the mixture are at a distance of at most \( 4 \sqrt{T(\log T + \log n)} \) in \( K \), with probability \( 1 - \frac{1}{n} \). On the other hand, two points from different distributions are at distance \( 5\sqrt{T(\log T + \log n)} \). Therefore, with probability \( 1 - \frac{1}{n} \), the distance concentration algorithm succeeds.

Similarly, for axis-aligned Gaussians, from Lemma 19, any two samples from a fixed distribution \( D_i \) in the mixture are at a distance of at most \( 4\sigma^2\sqrt{T(\log T + \log n)} \) in \( K \), with probability \( 1 - \frac{1}{n} \). On the other hand, two points from different distributions are at distance \( 5\sigma^2\sqrt{T(\log T + \log n)} \). Distance concentration therefore works, and the first part of the theorem follows.

If every vector in \( C \) has spread at least \( 49T \log \Lambda \), from Theorem 1, the maximum variance of any \( D_i \) in \( K \) is at most \( 11\sigma^2 \). The Theorem now follows by the same arguments as above. \qed