Linear regression without correspondence

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Abstract

This article considers algorithmic and statistical aspects of linear regression when the correspondence between the covariates and the responses is unknown. First, a fully polynomial-time approximation scheme is given for the natural least squares optimization problem in any constant dimension. Next, in an average-case and noise-free setting where the responses exactly correspond to a linear function of i.i.d. draws from a standard multivariate normal distribution, an efficient algorithm based on lattice basis reduction is shown to exactly recover the unknown linear function in arbitrary dimension. Finally, lower bounds on the signal-to-noise ratio are established for approximate recovery of the unknown linear function by any estimator.

1 Introduction

Consider the problem of recovering an unknown vector \(\hat{w} \in \mathbb{R}^d\) from noisy linear measurements when the correspondence between the measurement vectors and the measurements themselves is unknown. The measurement vectors (i.e., covariates) from \(\mathbb{R}^d\) are denoted by \(x_1, x_2, \ldots, x_n\); for each \(i \in [n] := \{1, 2, \ldots, n\}\), the \(i\)-th measurement (i.e., response) \(y_i\) is obtained using \(x_{\bar{\pi}(i)}\):

\[
y_i = \hat{w}^T x_{\bar{\pi}(i)} + \varepsilon_i, \quad i \in [n].
\]

Above, \(\bar{\pi}\) is an unknown permutation on \([n]\), and the \(\varepsilon_1, \varepsilon_2, \ldots, \varepsilon_n\) are unknown measurement errors.

This problem, which has been called unlabeled sensing \([20]\) and linear regression with an unknown permutation \([16]\), arises in many settings; see the aforementioned references for more details. In short, sensing limitations may create ambiguity in or even completely lose the ordering of measurements. The problem is also interesting because the missing correspondence makes an otherwise well-understood problem into one with very different computational and statistical properties.

Prior works. Unnikrishnan et al. \([20]\) study conditions on the measurement vectors that permit recovery of any target vector \(\hat{w}\) under noiseless measurements. They show that when the entries of the \(x_i\) are drawn i.i.d. from a continuous distribution, and \(n \geq 2d\), then almost surely, every vector \(\hat{w} \in \mathbb{R}^d\) is uniquely determined by noiseless correspondence-free measurements as in (1). (Under noisy measurements, it is shown that \(\hat{w}\) can be recovered when an appropriate signal-to-noise ratio tends to infinity.) It is also shown that \(n \geq 2d\) is necessary for such a guarantee that holds for all vectors \(\hat{w} \in \mathbb{R}^d\).

Pananjady et al. \([16]\) study statistical and computational limits on recovering the unknown permutation \(\bar{\pi}\). On the statistical front, they consider necessary and sufficient conditions on the signal-to-noise ratio \(\text{snr} := \|\hat{w}\|^2 / \sigma^2\) when the measurement errors \((\varepsilon_i)_{i=1}^n\) are i.i.d. draws from the normal distribution \(\mathcal{N}(0, \sigma^2)\) and the measurement vectors \((x_i)_{i=1}^n\) are i.i.d. draws from the standard multivariate normal distribution \(\mathcal{N}(0, I_d)\). Roughly speaking, exact recovery of \(\bar{\pi}\) is possible via maximum likelihood when \(\text{snr} \geq n^c\) for some absolute constant \(c > 0\), and approximate recovery is impossible for any method when \(\text{snr} \leq n^{c'}\) for some
other absolute constant $c' > 0$. On the computational front, they show that the least squares problem (which is equivalent to maximum likelihood problem)

$$\min_{w, \gamma} \sum_{i=1}^{n} (w^\top x_{x(i)} - y_i)^2$$

given arbitrary $x_1, x_2, \ldots, x_n \in \mathbb{R}^d$ and $y_1, y_2, \ldots, y_n \in \mathbb{R}$ is NP-hard when $d = \Omega(n)$, but admits a polynomial-time algorithm (in fact, an $O(n \log n)$-time algorithm based on sorting) when $d = 1$.

Elhami et al. [7] give a $O(dn^{d+1})$-time algorithm that, in dimension $d = 2$, is guaranteed to approximately recover $w$ when the measurement vectors are chosen in a very particular way from the unit circle and the measurement errors are uniformly bounded.

Contributions. We make progress on both computational and statistical aspects of the problem.

1. We give an approximation algorithm for the least squares problem from (2) that, any given $\left( x_i \right)_{i=1}^{n}$, $\left( y_i \right)_{i=1}^{n}$, and $\epsilon \in (0, 1)$, returns a solution with objective value at most $1 + \epsilon$ times that of the minimum in time $(n/\epsilon)^{O(d)}$. This a fully polynomial-time approximation scheme for any constant dimension.

2. We give an algorithm that exactly recovers $\bar{w}$ in the measurement model from (1), under the assumption that there are no measurement errors and the covariates $\left( x_i \right)_{i=1}^{n}$ are i.i.d. draws from $N(0, I_d)$. The algorithm, which is based on a reduction to a lattice problem and employs the lattice basis reduction algorithm of Lenstra et al. [14], runs in $\text{poly}(n, d)$ time when the covariate vectors $\left( x_i \right)_{i=1}^{n}$ and target vector $w$ are appropriately quantized. This result may also be regarded as for each-type guarantee for exactly recovering a fixed vector $\bar{w}$, which complements the for all-type results of Umnikrishnan et al. [20] concerning the number of measurement vectors needed for recovering all possible vectors.

3. We show that in the measurement model from (1) where the measurement errors are i.i.d. draws from $N(0, \sigma^2)$ and the covariate vectors are i.i.d. draws from $\text{N}(0, I_d)$, then no algorithm can approximately recover $w$ unless $\text{snr} \geq C \min\{1, d/\log\log(n)\}$ for some absolute constant $C > 0$. We also show that when the covariate vectors are i.i.d. draws from the uniform distribution on $[-1/2, 1/2]^d$, then approximate recovery is impossible unless $\text{snr} \geq C'$ for some other absolute constant $C' > 0$.

Our algorithms are not meant for practical deployment, but instead are intended to shed light on the computational difficulty of the least squares problem and the average-case recovery problem. Indeed, note that a naïve brute-force search over permutations requires time $\Omega(n!) = n^{\Omega(n)}$, and the only other previous algorithms (already discussed above) were restricted to $d = 1$ [16] or only had some form of approximation guarantee when $d = 2$ [7]. We are not aware of previous algorithms for the average-case problem in general dimension $d$.

Our lower bounds on $\text{snr}$ stand in contrast to what is achievable in the classical linear regression model (where the covariate/response correspondence is known): in that model, the $\text{snr}$ requirement for approximately recovering $\bar{w}$ scales as $d/n$, and hence the problem becomes easier with $n$. The lack of correspondence thus drastically changes the difficulty of the problem.

2 Approximation algorithm for the least squares problem

In this section, we consider the least squares problem from Equation (2). The inputs are an arbitrary matrix $X = [x_1 | x_2 | \cdots | x_n] \in \mathbb{R}^{n \times d}$ and an arbitrary vector $y = (y_1, y_2, \ldots, y_n) \in \mathbb{R}^n$, and the goal is to find a vector $w \in \mathbb{R}^d$ and permutation matrix $\Pi \in \mathcal{P}_n$ (where $\mathcal{P}_n$ denotes the space of $n \times n$ permutation matrices) to minimize $\|Xw - \Pi y\|_2^2$. This problem is NP-hard in the case where $d = \Omega(n)$ [16]. We give an approximation scheme that, for any $\epsilon \in (0, 1)$, returns a $(1 + \epsilon)$-approximation in time $(n/\epsilon)^{O(k)} + \text{poly}(n, d)$, where $k := \text{rank}(X) \leq \min\{n, d\}$.  

\footnote{A recent algorithm of Pananjady et al. [17] exploits a similar average-case setting but only for a somewhat easier variant of the problem where more information about the unknown correspondence is provided.}

\footnote{Each permutation matrix $\Pi \in \mathcal{P}_n$ corresponds to a permutation $\pi$ on $[n]$; the $(i, j)$-th entry of $\Pi$ is one if $\pi(i) = j$ and is zero otherwise.}
Algorithm 1 Approximation algorithm for least squares problem

**Input.** Covariate matrix $X = [x_1 | x_2 | \cdots | x_n] \in \mathbb{R}^{n \times k}$; response vector $y = (y_1, y_2, \ldots, y_n)^T \in \mathbb{R}^n$; approximation parameter $\epsilon \in (0, 1)$.

**Assume.** $X^T X = I_k$.

**Output.** Weight vector $\hat{w} \in \mathbb{R}^k$ and permutation matrix $\hat{\Pi} \in \mathcal{P}_n$.

1. Run "Row Sampling" algorithm with input matrix $X$ to obtain a matrix $S \in \mathbb{R}^{r \times n}$ with $r = 4k$.
2. Let $B$ be the set of vectors $b = (b_1, b_2, \ldots, b_n)^T \in \mathbb{R}^n$ satisfying the following: for each $i \in [n],$
   - if the $i$-th column of $S$ is all zeros, then $b_i = 0$;
   - otherwise, $b_i \in \{y_1, y_2, \ldots, y_n\}$.
3. Let $c := 1 + 4(1 + \sqrt{n/(4k)})^2$.
4. For each $b \in B$ do
   5. Compute $\hat{w}_b = \arg\min_{w \in \mathbb{R}^k} \|S(Xw - b)\|_2^2$, and let $r_b := \min_{\Pi \in \mathcal{P}_n} \|X\hat{w}_b - \Pi^T y\|_2^2$.
   6. Construct a $\sqrt{cr_b}/c$-net $N_b$ for the Euclidean ball of radius $\sqrt{cr_b}$ around $\hat{w}_b$, so that for each $v \in \mathbb{R}^d$ with $\|v - \hat{w}_b\|_2 \leq \sqrt{cr_b}$, there exists $v' \in N_b$ such that $\|v - v'\|_2 \leq \sqrt{cr_b}/c$.
7. end for
8. return $\hat{w} \in \arg\min_{w \in \bigcup_{b \in B} N_b} \|Xw - \Pi^T y\|_2^2$ and $\hat{\Pi} \in \arg\min_{\Pi \in \mathcal{P}_n} \|Xw - \Pi^T y\|_2^2$.

We assume without loss of generality that $X \in \mathbb{R}^{n \times k}$ and $X^T X = I_k$. This is because we can always replace $X$ with its matrix of left singular vectors $U \in \mathbb{R}^{n \times k}$, obtained via singular value decomposition $X = U \Sigma V^T$, where $U^T U = V^T V = I_k$ and $\Sigma \succ 0$ is diagonal. A solution $(w, \Pi)$ for $(U, y)$ has the same cost as the solution $(V \Sigma^{-1} w, \Pi)$ for $(X, y)$, and a solution $(w, \Pi)$ for $(X, y)$ has the same cost as the solution $(\Sigma V^T w, \Pi)$ for $(U, y)$.

2.1 Algorithm

Our approximation algorithm, shown as Algorithm 1, uses a careful enumeration to beat the naive brute-force running time of $\Omega(|\mathcal{P}_n|) = \Omega(n!)$. It uses as a subroutine a “Row Sampling” algorithm of Boutsidis et al. [4] (described in Appendix A), which has the following property.

**Theorem 1** (Specialization of Theorem 12 in [4]). There is an algorithm (“Row Sampling”) that, given any matrix $A \in \mathbb{R}^{n \times k}$ with $n \geq k$, returns in $\text{poly}(n, k)$ time a matrix $S \in \mathbb{R}^{r \times n}$ with $r = 4k$ such that the following hold.

1. Every row of $S$ has at most one non-zero entry.
2. For every $b \in \mathbb{R}^n$, every $w' \in \arg\min_{w \in \mathbb{R}^k} \|S(Aw - b)\|_2^2$ satisfies $\|Aw - b\|_2^2 \leq c \cdot \min_{w \in \mathbb{R}^k} \|Aw - b\|_2^2$ for $c = 1 + 4(1 + \sqrt{n/(4k)})^2 = O(n/k)$.

The matrix $S$ returned by Row Sampling determines a (weighted) subset of $O(k)$ rows of $A$ such that solving a (ordinary) least squares problem (with any right-hand side $b$) on this subset of rows and corresponding right-hand side entries yields an $O(n/k)$-approximation to the least squares problem over all rows and right-hand side entries. Row Sampling does not directly apply to our problem because (1) it does not minimize over permutations of the right-hand side, and (2) the approximation factor is too large. However, we are able to use it to narrow the search space in our problem.

An alternative to Row Sampling is to simply enumerate all subsets of $k$ rows of $X$. This is justified by a recent result of Dereziński and Warmuth [6], which shows that for any right-hand side $b \in \mathbb{R}^n$, using “volume sampling” [2] to choose a matrix $S \in \{0, 1\}^{k \times k}$ (where each row has one non-zero entry) gives a similar guarantee as that of Row Sampling, except with the $O(n/k)$ factor replaced by $k + 1$ in expectation.

2.2 Analysis

The approximation guarantee of Algorithm 1 is given in the following theorem.
Theorem 2. Algorithm 1 returns \( \hat{w} \in \mathbb{R}^k \) and \( \hat{\Pi} \in \mathcal{P}_n \) satisfying
\[
\left\| X \hat{w} - \hat{\Pi}^\top y \right\|_2^2 \leq (1 + \epsilon) \min_{w \in \mathbb{R}^k, \Pi \in \mathcal{P}_n} \left\| X w - \Pi^\top y \right\|_2^2 .
\]

Proof. Let \( \text{opt} := \min_{w, \Pi} \left\| X w - \Pi^\top y \right\|_2^2 \) be the optimal cost, and let \((w_\star, \Pi_\star)\) denote a solution achieving this cost. The optimality implies that \( w_\star \) satisfies the normal equations \( X^\top X w_\star = X^\top \Pi_\star^\top y \). Observe that there exists a vector \( b_\star \in \mathcal{B} \) satisfying \( S b_\star \approx S \Pi_\star^\top y \). By Theorem 1 and the normal equations, the vector \( \hat{w}_{b_\star} \), and cost value \( r_{b_\star} \), satisfy
\[
\text{opt} \leq r_{b_\star} \leq \left\| X \hat{w}_{b_\star} - \Pi_\star^\top y \right\|_2^2 = \left\| X (\hat{w}_{b_\star} - w_\star) \right\|_2^2 + \text{opt} \leq c \cdot \text{opt} .
\]
Moreover, since \( X^\top X = I_k \), we have that \( \left\| X \hat{w}_{b_\star} - w_\star \right\|_2 \leq \sqrt{(c-1) \text{opt}} \leq \sqrt{c r_{b_\star}} \). By construction of \( N_{b_\star} \), there exists \( w \in N_{b_\star} \) satisfying \( \left\| w - w_\star \right\|_2 \leq \| X (w - w_\star) \|_2 \leq \epsilon r_{b_\star} / c \leq \epsilon \text{opt} \). For this \( w \), the normal equations imply
\[
\min_{\Pi \in \mathcal{P}_n} \left\| X w - \Pi^\top y \right\|_2^2 \leq \left\| X w - \Pi_\star^\top y \right\|_2^2 = \left\| X (w - w_\star) \right\|_2^2 + \text{opt} \leq (1 + \epsilon) \text{opt} .
\]
Therefore, the solution returned by Algorithm 1 has cost no more than \((1 + \epsilon) \text{opt}\).

Running time. We now consider the running time of Algorithm 1. There is the initial cost for singular value decomposition (as discussed at the beginning of the section), and also for “Row Sampling”; both of these take \( \text{poly}(n, d) \) time. For the rest of the algorithm, we need to consider the size of \( \mathcal{B} \) and the size of the net \( N_{b_\star} \) for each \( b \in \mathcal{B} \). First, we have \( |\mathcal{B}| \leq n' = n^{O(k)} \), since \( S \) has only \( 4k \) rows and each row has at most a single non-zero entry. Next, for each \( b \in \mathcal{B} \), we construct the \( \delta \)-net \( N_{b_\star} \) (for \( \delta := \sqrt{\epsilon r_{b_\star} / c} \)) by constructing a \( \delta / \sqrt{c} \)-net for the \( \ell_\infty \)-ball of radius \( \sqrt{c r_{b_\star}} \) centered at \( \hat{w}_{b_\star} \) (using an appropriate axis-aligned grid). This has size \( |N_{b_\star}| \leq (4c^2 k / \epsilon)^{k / 2} = (n/c)^{O(k)} \). Finally, each \( \arg \min_{w \in \mathbb{R}^k} \) computation takes \( O(n k^2) \) time, and each \( \arg \min_{\Pi \in \mathcal{P}_n} \) takes \( O(nk + n \log n) \) time [16] (also see Appendix A). So, the overall running time is \((n/c)^{O(k)} + \text{poly}(n, d)\).

3 Exact recovery algorithm in noiseless Gaussian setting

To counter the intractability of the least squares problem in (2) confronted in Section 2, it is natural to explore distributional assumptions that may lead to faster algorithms. In this section, we consider the noiseless measurement model where the \( (x_i)_{i=1}^n \) are i.i.d. draws from \( N(0, I_d) \) (as in [16]). We give an algorithm that exactly recovers \( \bar{w} \) with high probability when \( n \geq d + 1 \). The algorithm runs in \( \text{poly}(n, d) \)-time when \( (x_i)_{i=1}^n \) and \( \bar{w} \) are appropriately quantized.

It will be notationally simpler to consider \( n + 1 \) covariate vectors and responses
\[
y_i = \bar{w}^\top x_{\pi(i)} , \quad i = 0, 1, \ldots, n .
\]
Here, \( (x_i)_{i=0}^n \) are \( n + 1 \) i.i.d. draws from \( N(0, I_d) \), the unknown permutation \( \pi \) is over \( \{0, 1, \ldots, n\} \), and the requirement of at least \( d + 1 \) measurements is expressed as \( n \geq d \).

In fact, we shall consider a variant of the problem in which we are given one of the values of the unknown permutation \( \pi \). Without loss of generality, assume we are given that \( \pi(0) = 0 \). Solving this variant of the problem suffices because there are only \( n + 1 \) possible values of \( \pi(0) \): we can try them all, incurring just a factor \( n + 1 \) in the computation time. So henceforth, we just consider \( \pi \) as an unknown permutation on \([n]\).
Algorithm 2 Find permutation

**input** Covariate vectors \( \mathbf{x}_0, \mathbf{x}_1, \mathbf{x}_2, \ldots, \mathbf{x}_n \) in \( \mathbb{R}^d \); response values \( y_0, y_1, y_2, \ldots, y_n \) in \( \mathbb{R} \); confidence parameter \( \delta \in (0, 1) \); lattice parameter \( \beta > 0 \).

**assume** there exists \( \mathbf{w} \in \mathbb{R}^d \) and permutation \( \hat{\pi} \) on \( [n] \) such that \( y_i = \mathbf{w}^\top \mathbf{x}_{\pi(i)} \) for each \( i \in [n] \), and that \( y_0 = \mathbf{w}^\top \mathbf{x}_0 \).

**output** Permutation \( \hat{\pi} \) on \( [n] \) or failure.
1. Let \( \mathbf{X} = [\mathbf{x}_1 | \mathbf{x}_2 | \cdots | \mathbf{x}_n]^\top \in \mathbb{R}^{n \times d} \), and its pseudoinverse be \( \mathbf{X}^\dagger = [\bar{\mathbf{x}}_1 | \bar{\mathbf{x}}_2 | \cdots | \bar{\mathbf{x}}_n] \).
2. Create Subset Sum instance with \( n^2 \) source numbers \( c_{i,j} := y_i \bar{\mathbf{x}}_j^\top \mathbf{x}_0 \) for \( (i, j) \in [n] \times [n] \) and target sum \( y_0 \).
3. Run Algorithm 3 with Subset Sum instance and lattice parameter \( \beta \).
4. if Algorithm 3 returns a solution \( \mathcal{S} \subseteq [n] \times [n] \) then
   5. return any permutation \( \hat{\pi} \) on \([n]\) such that \( \hat{\pi}(i) = j \) implies \((i, j) \in \mathcal{S} \).
6. else
7. return failure.
8. end if

Algorithm 3 Lagarias and Odlyzko [10] subset sum algorithm

**input** Source numbers \( \{c_i\}_{i \in \mathcal{I}} \subseteq \mathbb{R} \); target sum \( t \in \mathbb{R} \); lattice parameter \( \beta > 0 \).

**output** Subset \( \hat{\mathcal{S}} \subseteq \mathcal{I} \) or failure.
1. Construct lattice basis \( \mathbf{B} \in \mathbb{R}^{(|\mathcal{I}|+2) \times (|\mathcal{I}|+1)} \) where
   \[
   \mathbf{B} := \begin{bmatrix} \mathbf{I}_{|\mathcal{I}|+1} \\ \beta t \mathbf{I}_{|\mathcal{I}|+1} \end{bmatrix} \in \mathbb{R}^{(|\mathcal{I}|+2) \times (|\mathcal{I}|+1)}. \]
2. Run basis reduction [e.g., 14] to find non-zero lattice vector \( \mathbf{v} \) of length at most \( 2^{|\mathcal{I}|/2} \cdot \lambda_1(\mathbf{B}) \).
3. if \( \mathbf{v} = z(1, \chi_{\hat{\mathcal{S}}}^\top, 0)^\top \), with \( z \in \mathbb{Z} \) and \( \chi_{\hat{\mathcal{S}}} \in \{0, 1\}^\mathcal{I} \) is characteristic vector for some \( \hat{\mathcal{S}} \subseteq \mathcal{I} \) then
   4. return \( \hat{\mathcal{S}} \).
5. else
   6. return failure.
7. end if

3.1 Algorithm

Our algorithm, shown as Algorithm 2, is based on a reduction to the Subset Sum problem. An instance of Subset Sum is specified by an unordered collection of source numbers \( \{c_i\}_{i \in \mathcal{I}} \subseteq \mathbb{R} \), and a target sum \( t \in \mathbb{R} \). The goal is to find a subset \( \mathcal{S} \subseteq \mathcal{I} \) such that \( \sum_{i \in \mathcal{S}} c_i = t \). Although Subset Sum is NP-hard in the worst case, it is tractable for certain structured instances [10, 8]. We prove that Algorithm 2 constructs such an instance with high probability. A similar algorithm based on such a reduction was recently used by Andoni et al. [1] for a different but related problem.

Algorithm 2 proceeds by (i) solving a Subset Sum instance based on the covariate vectors and response values (using Algorithm 3), and (ii) constructing a permutation \( \hat{\pi} \) on \([n]\) based on the solution to the Subset Sum instance. With the permutation \( \hat{\pi} \) in hand, we (try to) find a solution \( \mathbf{w} \in \mathbb{R}^d \) to the system of linear equations \( y_i = \mathbf{w}^\top \mathbf{x}_{\pi(i)} \) for \( i \in [n] \). If \( \hat{\pi} = \pi \), then there is a unique such solution almost surely.

3.2 Analysis

The following theorem is the main recovery guarantee for Algorithm 2.

**Theorem 3.** Pick any \( \delta \in (0, 1) \). Suppose \( (\mathbf{x}_i)_{i=0}^n \) are i.i.d. draws from \( \mathcal{N}(\mathbf{0}, \mathbf{I}_d) \), and \( (y_0)_{i=1}^n \) follow the noiseless measurement model from (3) for some \( \mathbf{w} \in \mathbb{R}^d \) and permutation \( \hat{\pi} \) on \([n]\) and \( \hat{\pi}(0) = 0 \), and that \( n \geq d \). Furthermore, suppose Algorithm 2 is run with inputs \( (\mathbf{x}_i)_{i=0}^n \), \( (y_i)_{i=0}^n \), \( \delta \), and \( \beta \), and also that \( \beta \geq 2n^2/\varepsilon \) where \( \varepsilon \) is defined in Equation (8). With probability at least \( 1 - \delta \), Algorithm 2 returns \( \hat{\pi} = \pi \).
Remark 1. The value of \( \varepsilon \) from Equation (8) is directly proportional to \( \|\bar{w}\|_2 \), and Algorithm 2 requires a lower bound on \( \varepsilon \) (in the setting of the lattice parameter \( \beta \)). Hence, it suffices to determine a lower bound on \( \|\bar{w}\|_2 \). Such a bound can be obtained from the measurement values: a standard tail bound (Lemma 6 in Appendix B) shows that with high probability, \( \sqrt{\sum_{i=1}^{n} y_i^2 / (2n)} \) is a lower bound on \( \|\bar{w}\|_2 \), and is within a constant factor of it as well.

Remark 2. Algorithm 2 strongly exploits the assumption of noiseless measurements, which is expected given the \textit{snr} lower bounds of Pananjady et al. [16] for recovering \( \hat{\pi} \). The algorithm, however, is also very brittle and very likely fails in the presence of noise.

Remark 3. The recovery result does not contradict the results of Unnikrishnan et al. [20], which show that a collection of \( 2d \) measurement vectors are necessary for recovering all \( \bar{w} \), even in the noiseless measurement model of (3). Indeed, our result shows that for a fixed \( \bar{w} \in \mathbb{R}^d \), with high probability \( d + 1 \) measurements in the model of (3) suffice to permit exactly recovery of \( \bar{w} \), but this same set of measurement vectors (when \( d + 1 < 2d \)) will fail for some other \( \bar{w}' \).

The proof of Theorem 3 is based on the following theorem—essentially due to Lagarias and Odlyzko [10] and Frieze [8]—concerning certain structured instances of Subset Sum that can be solved using the lattice basis reduction algorithm of Lenstra et al. [14]. Given a basis \( B = \{b_1, b_2, \ldots, b_k\} \subset \mathbb{R}^n \times \mathbb{k} \) for a lattice

\[
\mathcal{L}(B) := \left\{ \sum_{i=1}^{k} z_i b_i : z_1, z_2, \ldots, z_k \in \mathbb{Z} \right\} \subset \mathbb{R}^n,
\]

this algorithm can be used to find a non-zero vector \( v \in \mathcal{L}(B) \setminus \{0\} \) whose length is at most \( 2^{(k-1)/2} \) times that of the shortest non-zero vector in the lattice

\[
\lambda_1(B) := \min_{v \in \mathcal{L}(B) \setminus \{0\}} \|v\|_2.
\]

Theorem 4 ([10, 8]). Suppose the Subset Sum instance specified by source numbers \( \{c_i\}_{i \in I} \subset \mathbb{R} \) and target sum \( t \in \mathbb{R} \) satisfy the following properties.

1. There is a subset \( S^* \subset I \) such that \( \sum_{i \in S^*} c_i = t \).

2. Define \( R := 2^{n/2} \sqrt{|S^*|} + 1 \) and \( Z_R := \{(z_0, z) \in \mathbb{Z} \times \mathbb{Z}^n : 0 < z_0^2 + \sum_{i \in I} z_i^2 \leq R^2 \} \). There exists \( \varepsilon > 0 \) such that \( z_0^2 + \sum_{i \in I} z_i^2 \leq R^2 \) holds for each \( (z_0, z) \in Z_R \) that is not an integer multiple of \( (1, \chi^*) \), where \( \chi^* \in \{0, 1\}^n \) is the characteristic vector for \( S^* \).

Let \( B \) be the lattice basis \( B \) constructed by Algorithm 3, and assume \( \beta \geq 2^{2n/2} / \varepsilon \). Then every non-zero vector in the lattice \( \Lambda(B) \) with length at most \( 2^{2n/2} \) times the length of the shortest non-zero vector in \( \Lambda(B) \) is an integer multiple of the vector \( (1, \chi_{S^*}, 0) \), and the basis reduction algorithm of Lenstra et al. [14] returns such a non-zero vector.

The Subset Sum instance constructed in Algorithm 2 has \( n^2 \) source numbers \( \{c_{i,j} : (i, j) \in [n] \times [n]\} \) and target sum \( y_0 \). We need to show that it satisfies the two conditions of Theorem 4.

Let \( S_{\pi} := \{(i, j) : \pi(i) = j\} \subset [n] \times [n] \), and let \( \Pi = (\Pi_{i,j})_{(i, j) \in [n] \times [n]} \in \mathcal{P}_n \) be the permutation matrix with \( \Pi_{i,j} := 1 \{\pi(i) = j\} \) for all \( (i, j) \in [n] \times [n] \). Note that \( \Pi \) is the “characteristic vector” for \( S_{\pi} \). Define

\[
R := 2^{n^2/2} \sqrt{n} + 1
\]

and

\[
Z_R := \left\{ (z_0, Z) \in \mathbb{Z} \times \mathbb{Z}^{n \times n} : 0 < z_0^2 + \sum_{1 \leq i, j \leq n} Z_{i,j}^2 \leq R^2 \right\}.
\]

A crude bound shows that \( |Z_R| \leq 2^{O(n^4)} \).

The following lemma establishes the first required property in Theorem 4.

Lemma 1. The random matrix \( X \) has rank \( d \) almost surely, and the subset \( S_{\pi} \) satisfies \( y_0 = \sum_{(i, j) \in S_{\pi}} c_{i,j} \).
Proof. That $X$ has rank $d$ almost surely follows from the fact that the probability density of $X$ is supported on all of $\mathbb{R}^{n \times d}$. This implies that $X^\dagger X = \sum_{j=1}^{n} \bar{x}_j x_j^\top = I_d$, and

$$y_0 = \sum_{j=1}^{n} x_0^\top \bar{x}_j x_j^\top \bar{w} = \sum_{1 \leq i, j \leq n} x_0^\top \bar{x}_j \cdot y_i \cdot 1\{\hat{\pi}(i) = j\} = \sum_{1 \leq i, j \leq n} c_{i,j} \cdot 1\{\hat{\pi}(i) = j\}. \quad \square$$

The next lemma establishes the second required property in Theorem 4. Here, we use the fact that the Frobenius norm $\|z_0 \bar{\Pi} - Z\|_F$ is at least one whenever $(z_0, Z) \in \mathbb{Z} \times \mathbb{Z}^n$. We also use the fact that $\hat{\Pi}$ is supported on $\{\pi(i) : (i, j) \in [n] \times [n]\}$.

Lemma 2. Pick any $\eta, \eta' > 0$ such that $3|Z_R| \eta + \eta' < 1$. With probability at least $1 - 3|Z_R| \eta - \eta'$, every $(z_0, Z) \in \mathbb{Z}_R$ with $Z = (Z_{i,j})_{i,j} \in [n] \times [n]$ satisfies

$$\left|z_0 \cdot y_0 - \sum_{i,j} Z_{i,j} \cdot c_{i,j}\right| \geq \frac{(\pi/4) \cdot \sqrt{(d-1)/n \cdot \eta^{2+\epsilon}}} {\sqrt{n + \sqrt{d} + \sqrt{2 \ln(1/\eta')}}} \cdot \|z_0 \bar{\Pi} - Z\|_F \cdot \|\bar{w}\|_2.$$

Proof. By Lemma 1, the matrix $\bar{\Pi}$ satisfies $y_0 = \sum_{i,j} \bar{\Pi}_{i,j} \cdot c_{i,j}$. Fix any $(z_0, Z) \in \mathbb{Z}_R$ with $Z = (Z_{i,j})_{i,j} \in [n] \times [n]$. Then

$$z_0 \cdot y_0 - \sum_{i,j} Z_{i,j} \cdot c_{i,j} = \sum_{i,j} (z_0 \cdot \bar{\Pi}_{i,j} - Z_{i,j}) \cdot x_0^\top \bar{x}_j \cdot \bar{w}^\top x_{\hat{\pi}(i)}.$$

Using matrix and vector notations, this can be written compactly as the inner product $x_0^\top (X^\dagger (z_0 \bar{\Pi} - Z)^\top \bar{\Pi} X \bar{w})$. Since $x_0 \sim \mathcal{N}(0, I_d)$ and is independent of $X$, the distribution of the inner product is normal with mean zero and standard deviation equal to $\|X^\dagger (z_0 \bar{\Pi} - Z)^\top \bar{\Pi} X \bar{w}\|_2$. By Lemma 7 (in Appendix B), with probability at least $1 - \eta$,

$$\|X^\dagger (z_0 \bar{\Pi} - Z)^\top \bar{\Pi} X \bar{w}\|_2 \geq \sqrt{n \cdot \eta}.$$

Observe that $X^\dagger = (X^\top X)^{-1} X^\top$ since $X$ has rank $d$ by Lemma 1, so

$$\|X^\dagger (z_0 \bar{\Pi} - Z)^\top \bar{\Pi} X \bar{w}\|_2 \geq \sqrt{n \cdot \eta} \cdot \|X\|_2^2.$$

By Lemma 4 (in Appendix B), with probability at least $1 - \eta'$,

$$\|X\|_2^2 \leq \left(\sqrt{n + \sqrt{d} + \sqrt{2 \ln(1/\eta')}}\right)^2.$$

And by Lemma 9 (in Appendix B), with probability at least $1 - 2\eta$,

$$\|X^\top (z_0 \bar{\Pi} - Z)^\top \bar{\Pi} X \bar{w}\|_2 \geq \|X^\top (z_0 \bar{\Pi} - Z)^\top \bar{\Pi} \bar{w}\|_2 \geq \sqrt{(d-1)/n \cdot \eta^{1+1/(d-1)}}.$$

Since $\bar{\Pi}$ is orthogonal, we have that $\|((z_0 \bar{\Pi} - Z)^\top \bar{\Pi}\|_F = \|z_0 \bar{\Pi} - Z\|_F$. Combining this with (4), (5), (6), and (7), and union bounds over all $(z_0, Z) \in \mathbb{Z}_R$ proves the claim. \(\square\)

Proof of Theorem 3. Lemma 1 and Lemma 2 (with $\eta' := \delta/2$ and $\eta := \delta/(6|Z_R|)$) together imply that with probability at least $1 - \delta$, the source numbers $\{c_{i,j} : (i, j) \in [n] \times [n]\}$ and target sum $y_0$ satisfy the conditions of Theorem 4 with

$$S^* := \{(i, j) \in [n] \times [n] : \hat{\pi}(i) = j\},$$

$$\varepsilon := \frac{(\pi/4) \cdot \sqrt{(d-1)/n \cdot \delta/(6|Z_R|)^{2+\epsilon}}} {\sqrt{n + \sqrt{d} + \sqrt{2 \ln(2/\delta)}}} \cdot \|\bar{w}\|_2 \geq 2 \cdot \text{poly}(n, \log(1/\delta)) \cdot \|\bar{w}\|_2.$$

Thus, in this event, Algorithm 3 (with $\beta$ satisfying $\beta \geq 2n^{2/2}/\varepsilon$) returns $\hat{S} = S^*$, which uniquely determines the permutation $\hat{\pi} = \hat{\pi}$ returned by Algorithm 2. \(\square\)
Running time. The basis reduction algorithm of Lenstra et al. [14] is iterative, with each iteration primarily consisting of Gram-Schmidt orthogonalization and another efficient linear algebraic process called “size reduction”. The total number of iterations required is

\[ O \left( \frac{k(k+1)}{2} \log \left( \sqrt{k} \cdot \max_{i \in [k]} \| b_i \|_2 \right) \right). \]

In our case, \( k = n^2 \) and \( \lambda_1(B) = \sqrt{n+1} \), and by Lemma 10 (in Appendix B), each of the basis vectors constructed has squared length at most \( 1 + \beta^2 \cdot \text{poly}(d, \log(n), 1/\delta) \cdot \| \hat{w} \|_2^2 \). Using the tight setting of \( \beta \) required in Theorem 3, this gives a poly\((n, d, \log(1/\delta))\) bound on the total number of iterations as well as on the total running time.

However, the basis reduction algorithm requires both arithmetic and rounding operations, which are typically only available for finite precision rational inputs. Therefore, a formal running time analysis would require the idealized real-valued covariate vectors \((x_i)_{i=0}^n\) and unknown target vector \(\hat{w}\) to be quantized to finite precision values. This is doable, and is similar to using a discretized Gaussian distribution for the distribution of the covariate vectors (and assuming \(\hat{w}\) is a vector of finite precision values), but leads to a messier analysis incomparable to the setup of previous works. Nevertheless, it would be desirable to find a different algorithm that avoids lattice basis reduction that still works with just \(d+1\) measurements.

4 Lower bounds on signal-to-noise for approximate recovery

In this section, we consider the measurement model from (1) where \((x_i)_{i=1}^n\) are i.i.d. draws from either \(N(0, I_d)\) or the uniform distribution on \([-1/2, 1/2]^d\), and \((\varepsilon_i)_{i=1}^n\) are i.i.d. draws from \(N(0, \sigma^2)\). We establish lower bounds on the signal-to-noise ratio (snr),

\[ \text{snr} = \frac{\| \hat{w} \|_2^2}{\sigma^2}, \]

required by any estimator \(\hat{w} = \hat{w}((x_i)_{i=1}^n, (y_i)_{i=1}^n)\) for \(\hat{w}\) to approximately recover \(\bar{w}\) in expectation. The estimators may have a priori knowledge of the values of \(\| \hat{w} \|_2\) and \(\sigma^2\).

Theorem 5. Assume \((\varepsilon_i)_{i=1}^n\) are i.i.d. draws from \(N(0, \sigma^2)\).

1. There is an absolute constant \(C > 0\) such that the following holds. If \(n \geq 3\), \(d \geq 22\), \((x_i)_{i=1}^n\) are i.i.d. draws from \(N(0, I_d)\), \((y_i)_{i=1}^n\) follow the measurement model from (1), and

\[ \text{snr} \leq C \cdot \min \left\{ \frac{d}{\log \log(n)}, 1 \right\}, \]

then for any estimator \(\hat{w}\), there exists some \(\bar{w} \in \mathbb{R}^d\) such that

\[ \mathbb{E} \| \hat{w} - \bar{w} \|_2 \geq \frac{1}{24} \| \bar{w} \|_2. \]

2. If \((x_i)_{i=1}^n\) are i.i.d. draws from the uniform distribution on \([-1/2, 1/2]^d\), and \((y_i)_{i=1}^n\) follow the measurement model from (1), and

\[ \text{snr} \leq 2, \]

then for any estimator \(\hat{w}\), there exists some \(\bar{w} \in \mathbb{R}^d\) such that

\[ \mathbb{E} \| \hat{w} - \bar{w} \|_2 \geq \frac{1}{2} \left( 1 - \frac{1}{\sqrt{2}} \right) \| \bar{w} \|_2. \]

Note that in the classical linear regression model where \(y_i = \hat{w}^T x_i + \varepsilon_i\) for \(i \in [n]\), the maximum likelihood estimator \(\hat{w}_{mle}\) satisfies \(\mathbb{E} \| \hat{w}_{mle} - \bar{w} \|_2 \leq C \sigma \sqrt{d/n}\), where \(C > 0\) is an absolute constant. Therefore, the snr requirement to approximately recover \(\bar{w}\) up to (say) Euclidean distance \(\| \hat{w} \|_2 / 24\) is \(\text{snr} \geq 24^2 C d / n\).
Compared to this setting, Theorem 5 implies that with the measurement model of (1), the $\text{snr}$ requirement (as a function of $n$) is at substantially higher $(d/\log \log(n))$ in the normal covariate case, or a constant not even decreasing with $n$ in the uniform covariate case.

For the normal covariate case, Pananjady et al. [16] show that if $n > d$, $\epsilon < \sqrt{n}$, and

$$\text{snr} \geq n^{c - \frac{1}{\sqrt{n}}} + \epsilon,$$

then the maximum likelihood estimator $(\hat{w}_{\text{mle}}, \hat{\pi}_{\text{mle}})$ (i.e., any minimizer of (2)) satisfies $\hat{\pi}_{\text{mle}} = \pi$ with probability at least $1 - c' n^{-2\epsilon}$. (Here, $c > 0$ and $c'$ are absolute constants.) It is straightforward to see that, on the same event, we have $\|\hat{w}_{\text{mle}} - \hat{w}\|_2 \leq C \sigma \sqrt{d/n}$ for some absolute constant $C > 0$. Therefore, the necessary and sufficient conditions on $\text{snr}$ for approximate recovery of $\hat{w}$ lie between $C'd/\log \log(n)$ and $n^{C''}$ (for absolute constants $C', C'' > 0$). Narrowing this range remains an interesting open problem.

A sketch of the proof in the normal covariate case is as follows. Without loss of generality, we restrict attention to the case where $\hat{w}$ is a unit vector. We construct a $1/\sqrt{d}$-packing of the unit sphere in $\mathbb{R}^d$; the target $\hat{w}$ will be chosen from from this set. Observe that for any distinct $u, u' \in U$, each of $(x_i^\top u)^n_{i=1}$ and $(x_i^\top u')^n_{i=1}$ is an i.i.d. sample from $\mathcal{N}(0,1)$ of size $n$; we prove that they therefore determine empirical distributions that are close to each other in Wasserstein-2 distance with high probability. We then prove that conditional on this event, the resulting distributions of $(y_i)^n_{i=1}$ under $\hat{x} = u$ and $\hat{x} = u'$ (for any pair $u, u' \in U$) are close in Kullback-Leibler divergence. Hence, by (a generalization of) Fano’s inequality [see, e.g., 9], no estimator can determine the correct $u \in U$ with high probability.

The proof for the uniform case is similar, using $U = \{e_1, -e_1\}$ where $e_1 = (1, 0, \ldots, 0)^\top$. The full proof of Theorem 5 is given in Appendix C.

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A Additional details for approximation algorithm

This section provides some additional details on subroutines used in Algorithm 1.

Row sampling. First, we give the details of the “Row Sampling” algorithm of Boutsidis et al. [4] used in Section 2. The pseudocode is presented as Algorithm 4, and uses the following notations:

- For each $i \in [n]$, $e_i$ is the $i$-th coordinate basis vector in $\mathbb{R}^n$.
- $L(x, \delta_L, A, \ell) := \frac{x^\top (A - (\ell + \delta_L)I_k)^{-2} x}{\phi(\ell + \delta_L, A) - \Phi(\ell, A)} - (\ell + \delta_L)I_k^{-1} x$, where $\phi(\ell, A) := \sum_{i=1}^k \frac{1}{\lambda_i(A)}$ and $(\lambda_i(A))_{i=1}^k$ are the eigenvalues of $A$.
- $\hat{U}(x, \delta, B, u) := \frac{x^\top (B - u' I_r)^{-2} x}{\phi'(u, B) - \phi'(u', B)} - x^\top (B - u' I_r)^{-1} x$, where $u' = u + \delta$ and $\phi'(u, B) := \sum_{i=1}^r \frac{1}{u - \lambda_i(B)}$ and $(\lambda_i(B))_{i=1}^k$ are the eigenvalues of $B$.

Algorithm 4 “Row Sampling” algorithm of Boutsidis et al. [4]

| input | Matrix $X = [x_1 | x_2 | \cdots | x_n]^\top \in \mathbb{R}^{n \times k}$ such that $X^\top X = I_k$; integer $r \geq k$. |
| output | Matrix $S = (S_{i,j})_{i,j} \in [r] \times [n] \in \mathbb{R}^{r \times n}$. |
| 1: | Set $A_0 = 0_{k \times k}$, $B_0 = 0_{r \times n}$, $S = 0_{r \times n}$. |
| 2: | for $\tau = 0$ to $r - 1$ do |
| 3: | Let $\ell_\tau = \tau - \sqrt{r/k}$ and $u_\tau = \delta(\tau + \sqrt{\alpha r})$. |
| 4: | Select $i_\tau \in [n]$ and number $t_\tau > 0$ such that $\hat{U}(e_{i_\tau}, \delta, B_{t_\tau}, u_\tau) \leq \frac{1}{\tau} \leq L(x_{i_\tau}, \delta_L, A_{t_\tau}, \ell_\tau)$. |
| 5: | Set $A_{t_\tau} = A_{t_\tau} + t_\tau x_{i_\tau} x_{i_\tau}^\top$, $B_{t_\tau} = B_{t_\tau} + t_\tau e_{i_\tau} e_{i_\tau}^\top$, and $S_{t_\tau} = \sqrt{r^{-1}(1 - \sqrt{k/r})/\sqrt{t_\tau}}$. |
| 6: | end for |
| 7: | return $S$. |

One may also consider using leverage score sampling (i.e., sample a row of $X$ proportional to its squared length) instead of this Row Sampling algorithm. This would work, but would require selecting $O(k \log k)$ rows as opposed to just $O(k)$ [21], this leads to an overall running time of $(n/e)^{O(k \log k)} + \text{poly}(n, d)$.

One-dimensional permutation problem. Next, we explain how to solve the optimization problem

$$\min_{\Pi \in \mathcal{P}_n} \|a - \Pi^\top b\|_2^2$$

for any given $a, b \in \mathbb{R}^n$. Let $(a(i))_{i=1}^n$ denote the non-decreasing ordering $a(1) \leq a(2) \leq \cdots \leq a(n)$ of the entries of $a$, and let $(b(i))_{i=1}^n$ be analogously defined. By Lemma 11, we have

$$\min_{\Pi \in \mathcal{P}_n} \|a - \Pi^\top b\|_2^2 = \sum_{i=1}^n (a(i) - b(i))^2.$$

Hence, if $\Pi_a$ (respectively, $\Pi_b$) is the permutation matrix that rearranges the entries of $a$ (respectively, $b$) in non-decreasing order, then

$$\sum_{i=1}^n (a(i) - b(i))^2 = \|\Pi_a a - \Pi_b b\|_2^2 = \|\Pi_a^\top (\Pi_a a - \Pi_b b)\|_2^2 = \|a - \Pi_a^\top \Pi_b b\|_2^2,$$

where the second and third equalities use the fact that permutation matrices are orthogonal. Thus, the minimizing permutation matrix is $\Pi = \Pi_b^\top \Pi_a$. This can be found by sorting the entries of $a$ and of $b$ in $O(n \log n)$ time.
B Probability inequalities

This section collects several probability inequalities used in the analysis of Algorithm 2. Let $\sigma_i(M)$ denote the $i$-th largest singular value of the matrix $M$.

Extreme singular values of Gaussian random matrices.

**Lemma 3** (Eq. 3.2 in [19]). Let $A$ be an $n \times d$ matrix whose entries are i.i.d. $N(0, 1)$ random variables and $n \geq d$. For any $\eta \in (0, 1)$,

$$\Pr\left(\sigma_d(A) \leq \frac{\eta}{\sqrt{d}}\right) \leq \eta.$$  

**Lemma 4** (Theorem II.13 in [5]). Let $A$ be an $n \times d$ matrix whose entries are i.i.d. $N(0, 1)$ random variables. For any $\eta \in (0, 1)$,

$$\Pr\left(\sigma_1(A) \geq \sqrt{n} + \sqrt{d + \sqrt{2\ln(1/\eta)}}\right) \leq \eta.$$  

Tail bounds for Gaussian and $\chi^2$ random variables.

**Lemma 5.** Let $Z \sim N(0, 1)$. For any $\eta \in (0, 1)$, $\Pr(Z^2 \geq 2\ln(2/\delta)) \leq \eta$.

*Proof.* This follows from the standard Chernoff bounding method. □

**Lemma 6** (Lemma 1 in [11]). Let $W \sim \chi^2_k$. For any $\eta \in (0, 1)$, $\Pr(W \geq k + 2\sqrt{k\ln(1/\eta) + 2\ln(1/\eta)}) \leq \eta$.

Anti-concentration bounds for Gaussian and $\chi^2$ random variables.

**Lemma 7.** Let $Z \sim N(0, 1)$. For any $\eta \in (0, 1)$, $\Pr(Z^2 \leq \pi\eta^2/2) \leq \eta$.

*Proof.* This follows from direct integration. □

**Lemma 8** (Lemma 9 in [16]). Let $W \sim \chi^2_k$. For any $\eta \in (0, 1)$, $\Pr(W \leq k\eta^2/k/4) \leq \eta$.

**Lemma 9.** Let $x \in \mathbb{R}^d$ be any vector, $M \in \mathbb{R}^{n \times n}$ be any matrix, and $A$ a random $n \times d$ matrix of i.i.d. $N(0, 1)$ random variables. For any $\eta \in (0, 1/2)$,

$$\Pr\left(\|A^\top MAx\|_2 \leq \|M\|_F \cdot \|x\|_2 \cdot \sqrt{\frac{(d-1)\pi}{8n}} \cdot \eta^{1/(d-1)}\right) \leq 2\eta.$$  

*Proof.* Let $u_1 := x/\|x\|_2$, and extend to an orthonormal basis $u_1, u_2, \ldots, u_d$ for $\mathbb{R}^d$. Let $g_i := Au_i$ for each $i \in [d]$, so $g_1, g_2, \ldots, g_d$ are i.i.d. $N(0, I_n)$ random vectors. We first show that

$$\Pr\left(\|Mg_1\|_2 \leq \|M\|_F \cdot \sqrt{\frac{\pi}{2n}} \cdot \eta\right) \leq \eta. \tag{9}$$  

To see this, note that the distribution of $\|Mg_1\|_2^2$ is the same as that of $\sum_{i=1}^n \sigma_i(M)^2 Z_i^2$, where $Z_1, Z_2, \ldots, Z_n$ are i.i.d. $N(0, 1)$ random variables. Therefore, Lemma 7 and the fact $\|M\|_F^2 \geq \|M\|_k^2/n$ proves the claim in (9).

Next, observe that

$$\|A^\top MAx\|_2^2 = \|x\|_2^2 \cdot u_1^\top A^\top MA \left(\sum_{i=1}^d u_i u_i^\top\right) A^\top MA u_1$$

$$= \|x\|_2^2 \cdot g_1^\top M^\top \left(\sum_{i=1}^d g_i g_i^\top\right) M g_1$$

$$\geq \|x\|_2^2 \cdot \sum_{i=2}^d (g_i^\top M g_1)^2.$$  

(10)
Conditional on $g_1$, the final right-hand side in (10) has the same distribution as $\|x\|_2^2 \cdot |Mg_1|_2^2 \cdot W$, where $W \sim \chi^2_{d-1}$ is a chi-squared random variable with $d - 1$ degrees of freedom. Therefore, Lemma 8 implies

$$\Pr \left( \|A^TMAx\|_2 \leq \|x\|_2 \cdot |Mg_1|_2 \cdot \frac{\sqrt{d-1}}{2} \cdot \eta^{1/(d-1)} \right) \leq \eta.$$  

Combining this inequality with the inequality from (9) and a union bound proves the claim. □

**Lattice basis size.** The following lemma is used to bound the size of the lattice basis vectors constructed by Algorithm 2 (via Algorithm 3). Recall that there are $n^2 + 1$ basis vectors; one has length $\sqrt{1 + \beta^2 y_0^2}$, and the remaining $n^2$ have length $\sqrt{1 + \beta^2 c_{i,j}^2}$ for $(i, j) \in [n] \times [n]$.

**Lemma 10.** For any $\eta \in (0, 1/5)$, with probability at least $1 - 5\eta$,

$$|y_0| \leq \|\tilde{w}\|_2 \sqrt{2 \ln (2/\eta)}.$$  

$$|c_{i,j}| \leq \|\tilde{w}\|_2 \cdot \sqrt{2 \ln (2n/\eta)} \cdot \frac{d}{\eta^2} \cdot \sqrt{d + 2 \sqrt{d \ln(n/\eta)} + 2 \ln(n/\eta)} \cdot \sqrt{2 \ln (2n/\eta)}, \quad (i, j) \in [n] \times [n].$$

**Proof.** By Lemma 3, Lemma 5, and Lemma 6, with probability at least $1 - 5\eta$,

$$\|(X^T X)^{-1}\|_2 \leq \frac{d}{\eta^2},$$

$$|x_0^\top \tilde{w}| \leq \|\tilde{w}\|_2 \sqrt{2 \ln (2/\eta)},$$

$$|x_{i(i)}^\top \tilde{w}| \leq \|\tilde{w}\|_2 \sqrt{2 \ln (2n/\eta)}, \quad i \in [n],$$

$$|x_j^\top x_0| \leq \|x_j\|_2 \sqrt{2 \ln (2n/\eta)}, \quad j \in [n],$$

$$|x_j^\top x_0| \leq \sqrt{d + 2 \sqrt{d \ln(n/\eta)} + 2 \ln(n/\eta)} \cdot \sqrt{2 \ln (2n/\eta)}, \quad j \in [n].$$

In this event, we have for each $(i, j) \in [n] \times [n]$,

$$|c_{i,j}| = |x_{i(i)}^\top \tilde{w}| \cdot |x_j^\top x_0|$$

$$\leq \|\tilde{w}\|_2 \cdot \sqrt{2 \ln (2n/\eta)} \cdot \|X^T e_j\|_2 \cdot \sqrt{2 \ln (2n/\eta)}$$

$$= \|\tilde{w}\|_2 \cdot \sqrt{2 \ln (2n/\eta)} \cdot \|(X^T X)^{-1} X^T e_j\|_2 \cdot \sqrt{2 \ln (2n/\eta)}$$

$$\leq \|\tilde{w}\|_2 \cdot \sqrt{2 \ln (2n/\eta)} \cdot \frac{d}{\eta^2} \cdot \sqrt{d + 2 \sqrt{d \ln(n/\eta)} + 2 \ln(n/\eta)} \cdot \sqrt{2 \ln (2n/\eta)},$$

and $|y_0| \leq \|\tilde{w}\|_2 \sqrt{2 \ln (2/\eta)}.$ □

**C Proof of signal-to-noise lower bounds**

This section provides the proof of Theorem 5.

Since $\hat{\pi}$ is unknown in the measurement model from (1), we may assume that $y_1, y_2, \ldots, y_n$ are provided as an unordered multiset, denoted by $\{y_i\}_{i=1}^n$. Below, for any vector $\mathbf{a} = (a_1, a_2, \ldots, a_n)^\top$, we use the notation $(a_{(i)})_{i=1}^n$ to denote the non-decreasing ordering $a_{(1)} \leq a_{(2)} \leq \cdots \leq a_{(n)}$ of its entries, and $(a)^\top := (a_{(1)}, a_{(2)}, \ldots, a_{(n)})^\top$ to denote the vector of the entries in this order.

We use the following representation for the Kantorovich transport distance with respect to Euclidean metric (i.e., Wasserstein-2 distance, denoted by $W_2$).

**Lemma 11** (Lemma 4.1 in [3]). Let $\mu_n$ be the empirical measure on $a_1, a_2, \ldots, a_n \in \mathbb{R}$, and $\nu_n$ be the empirical measure on $b_1, b_2, \ldots, b_n \in \mathbb{R}$. Then

$$W_2(\mu_n, \nu_n)^2 = \min_\pi \frac{1}{n} \sum_{i=1}^n (a_i - b_{\pi(i)})^2 = \frac{1}{n} \sum_{i=1}^n (a_{(i)} - b_{(i)})^2,$$

where $\min_\pi$ denotes minimization over permutations $\pi$ on $[n]$.  

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For probability measures μ and ν, we use \( \text{KL}(\mu, \nu) \) to denote the Kullback-Leibler divergence between μ and ν, and \( \| \mu - \nu \|_{\text{tv}} \) to denote the total variation distance between μ and ν.

### C.1 Normal case

We first consider the case where \( (x_i)_{i=1}^n \) are i.i.d. draws from \( N(0, I_d) \). By homogeneity, we may assume without loss of generality that \( \| \bar{\mathbf{w}} \|_2 = 1 \), so \( \text{snr} = 1/d^2 \).

The proof is based on the Generalized Fano method of Han and Verdú [9] as described by Yu [22].

**Lemma 12** (Lemma 3 in [22]). Let \( (\Theta, \rho) \) be a pseudometric space, and let \( \bar{\Theta} \subseteq \Theta \) index a collection of probability measures \( (P_\theta)_{\theta \in \bar{\Theta}} \) such that \( \rho(\theta, \theta') \geq \alpha \) and \( \text{KL}(P_\theta, P_{\theta'}) \leq \beta \) for all distinct pairs \( \theta, \theta' \in \bar{\Theta} \). Then for any estimator \( \hat{\theta} \) taking values in \( \Theta \),

\[
\max_{\theta \in \bar{\Theta}} E_{P_\theta} [\rho(\hat{\theta}, \theta)] \geq \frac{\alpha}{2} \left( 1 - \frac{\beta + \ln 2}{\ln |\Theta|} \right),
\]

where \( E_{P_\theta} \) denotes expectation with respect to data drawn from \( P_\theta \).

We apply Lemma 12 with \( (\Theta, \rho) = (S^{d-1}, \| \cdot \|_2) \). We construct a packing \( U \) of the unit sphere \( S^{d-1} := \{ \mathbf{u} \in \mathbb{R}^d : \| \mathbf{u} \|_2 = 1 \} \) using the following variant of the Gilbert-Varsamov bound.

**Lemma 13** (Lemma 10 in [15]). For every \( h \in [d] \) such that \( h \leq d/4 \), there exists a subset \( C \) of \( \{0, 1\}^d \) such that (i) the Hamming weight of each \( \mathbf{c} \in C \) is \( h \), (ii) the Hamming distance between every distinct pair \( \mathbf{c}, \mathbf{c}' \in C \) is more than \( h/2 \), and (iii) the cardinality of \( C \) satisfies \( |C| \geq 0.233 h \ln(d/h) \).

We take \( C \subseteq \{0, 1\}^d \) as guaranteed by Lemma 13 with \( h := \lceil d/4 \rceil \), and let

\[
U := \left\{ \mathbf{c}/\sqrt{h} : \mathbf{c} \in C \right\} \subset S^{d-1}.
\]

Observe that \( U \) is a \( (1/\sqrt{2}) \)-packing of \( S^{d-1} \) (i.e., every distinct pair \( \mathbf{u}, \mathbf{u}' \in U \) satisfies \( \| \mathbf{u} - \mathbf{u}' \|_2 > 1/\sqrt{2} \)), and

\[
\ln |U| \geq 0.233 \left( \frac{d}{4} - 1 \right) \ln 4.
\]

For any \( \mathbf{u} \in U \), let \( P_{\mathbf{u}} \) denote the (joint) probability distribution of \( (x_i)_{i=1}^n, (y_i)_{i=1}^n \) when \( \bar{\mathbf{w}} = \mathbf{u} \). Also, let \( Q_{\mathbf{u}} \) denote the conditional probability distribution of \( \{\mathbf{u}^\top x_{\pi_u(i)} + \varepsilon_{\pi_u(i)}\}_{i=1}^n \) given \( (x_i)_{i=1}^n \), where \( \pi_u \) is the permutation on \( [n] \) such that

\[
\mathbf{u}^\top x_{\pi_u(1)} \leq \mathbf{u}^\top x_{\pi_u(2)} \leq \cdots \leq \mathbf{u}^\top x_{\pi_u(n)}.
\]

Note that because \( \varepsilon_1, \varepsilon_2, \ldots, \varepsilon_n \) are i.i.d., \( Q_{\mathbf{u}} \) is also the conditional distribution of \( \{y_i\}_{i=1}^n \) given \( (x_i)_{i=1}^n \) when \( \bar{\mathbf{w}} = \mathbf{u} \).

For any distinct \( \mathbf{u}, \mathbf{u}' \in U \),

\[
\text{KL}(Q_{\mathbf{u}}, Q_{\mathbf{u}'}) \leq \frac{1}{2\sigma^2} \sum_{i=1}^n \left( \mathbf{u}^\top x_{\pi_u(i)} - (\mathbf{u}')^\top x_{\pi_{u'}(i)} \right)^2
\]

by the data processing inequality for KL-divergence. We define \( \mathcal{E} \) to be the event in which

\[
\sum_{i=1}^n \left( \mathbf{u}^\top x_{\pi_u(i)} - (\mathbf{u}')^\top x_{\pi_{u'}(i)} \right)^2 \leq \left( \sqrt{C_0 \log \log(n) + 4\ln(|U|^2)} \right)^2
\]

for all distinct \( \mathbf{u}, \mathbf{u}' \in U \), where \( C_0 > 0 \) is the absolute constant from Lemma 15 (below). By Equation (11), Lemma 15, and a union bound, we have \( \text{Pr}(\mathcal{E}) \geq 1/2 \). Therefore, by Lemma 12, for any estimator \( \hat{\mathbf{w}} \),

\[
\max_{\mathbf{u} \in U} E_{P_{\mathbf{u}}} \left[ \| \hat{\mathbf{w}} - \mathbf{u} \|_2 \right] \geq \max_{\mathbf{u} \in U} E_{P_{\mathbf{u}}} \left[ \| \hat{\mathbf{w}} - \mathbf{u} \|_2 \mid \mathcal{E} \right] \cdot \text{Pr}(\mathcal{E})
\]

\[
\geq \frac{1}{2\sqrt{2}} \left( 1 - \frac{C_0 \log \log(n) + 16 \ln |U|}{\sigma^2 \ln |U|} - \frac{\ln 2}{\ln |U|} \right) \cdot \frac{1}{2}
\]

\[
= \frac{1}{4\sqrt{2}} \left( 1 - \frac{C_0 \log \log(n)}{\sigma^2 \ln |U|} - \frac{16}{\sigma^2} - \frac{\ln 2}{\ln |U|} \right).
\]
Plugging in the lower bound for $\ln|U|$ and the upper bound on $\text{snr} = 1/\sigma^2$ completes the proof.  

\[ \square \]

C.2 Uniform case

We now consider the case where $(x_i)_{i=1}^n$ are drawn i.i.d. from the uniform distribution on $[-1/2, 1/2]^d$.\(^3\) Again, by homogeneity, we assume without loss of generality that $\|\tilde{w}\|_2 = 1$, so $\text{snr} = 1/\sigma^2$.

The proof is based on the two-point method of Le Cam [12] as described by Yu [22].

**Lemma 14** (Lemma 1 in [22]). Let $(\Theta, \rho)$ be a pseudometric space, and let $\theta_1, \theta_2 \in \Theta$ correspond to probability measures $P_{\theta_1}$ and $P_{\theta_2}$ on the same space. Then for any estimator $\hat{\theta}$ taking values in $\Theta$,

$$\max_{\theta \in \{\theta_1, \theta_2\}} \mathbb{E}_{P_{\theta}} \left[ \rho(\hat{\theta}, \theta) \right] \geq \frac{1}{2} \rho(\theta_1, \theta_2) \left( 1 - \|P_{\theta_1} - P_{\theta_2}\|_\text{tv} \right),$$

where $\mathbb{E}_{P_\theta}$ denotes expectation with respect to data drawn from $P_\theta$.

We apply Lemma 14 with $(\Theta, \rho) = (S^{d-1}, \|\cdot\|_2)$. Let $P_{e_1}$ be the (joint) distribution of $((x_i)_{i=1}^n, (y_i)_{i=1}^n)$ when $x = e_1$, and $P_{-e_1}$ be the corresponding distribution when $x = -e_1$. Here, $e_1 = (1, 0, \ldots, 0)$. As before, let $Q_u$ (for $u \in \{e_1, -e_2\}$) denote the conditional probability distribution of $\{(u^\top x_{\pi_u(i)} + e_{\pi_u(i)})_{i=1}^n\}$ given $(x_i)_{i=1}^n$, where $\pi_u$ is the permutation on $[n]$ such that

$$u^\top x_{\pi_u(1)} \leq u^\top x_{\pi_u(2)} \leq \cdots \leq u^\top x_{\pi_u(n)}.$$ 

Again, $Q_u$ is also the conditional distribution of $(y_i)_{i=1}^n$ given $(x_i)_{i=1}^n$ when $\tilde{w} = u$.

Let $E$ be the event in which

$$\sum_{i=1}^n \left( e_1^\top x_{\pi_{e_1}(i)} - (-e_1)^\top x_{\pi_{-e_1}(i)} \right)^2 = \sum_{i=1}^n \left( e_1^\top x_{\pi_{e_1}(i)} + e_1^\top x_{\pi_{e_1}(n+1-i)} \right)^2 \leq 1.$$

By Lemma 19 (below), $\Pr(E) \geq 1/2$. Moreover, since $P_{e_1}(E) = P_{-e_1}(E) = \Pr(E)$,

$$\|P_{e_1} - P_{-e_1}\|_\text{tv} \leq \|P_{e_1}(\cdot | E) - P_{-e_1}(\cdot | E)\|_\text{tv} \Pr(E) + (1 - \Pr(E))$$

$$\leq \sqrt{\frac{1}{2} \text{KL}(P_{e_1}(\cdot | E), P_{-e_1}(\cdot | E)) \Pr(E)} + (1 - \Pr(E))$$

$$\leq \sqrt{\frac{1}{2} \frac{1}{2\sigma^2} \Pr(E) + (1 - \Pr(E))}$$

$$\leq \frac{1}{2} \left( 1 + \frac{1}{\sqrt{2}} \right).$$

Above, the second inequality follows from Pinsker’s inequality; the third inequality uses Equation (11) and the fact that

$$\text{KL}(Q_{e_1}, Q_{-e_1}) \leq \frac{1}{2\sigma^2} \sum_{i=1}^n \left( e_1^\top x_{\pi_{e_1}(i)} - (-e_1)^\top x_{\pi_{-e_1}(i)} \right)^2 \leq \frac{1}{2\sigma^2}$$

on the event $E$; the fourth inequality uses the assumption that $\text{snr} = 1/\sigma^2 \leq 2$ and the fact $\Pr(E) \geq 1/2$. We conclude by Lemma 14 that

$$\max_{u \in \{e_1, -e_2\}} \mathbb{E}_{P_u} \left[ \|\tilde{w} - u\|_2 \right] \geq \frac{1}{2} \cdot \frac{1}{2} \left( 1 - \frac{1}{2} \left( 1 + \frac{1}{\sqrt{2}} \right) \right) = \frac{1}{2} \left( 1 - \frac{1}{\sqrt{2}} \right),$$

completing the proof.  

\[ \square \]

\(^3\)We actually just need that the marginal distribution of the first coordinate of each $x_i$ be uniform on $[-1/2, 1/2]$. 

C.3 Auxiliary results

**Lemma 15.** There is an absolute constant $C_0 > 0$ such that the following holds. Let $n \geq 3$, and let $X$ be a random $n \times d$ matrix of i.i.d. $N(0, 1)$ random variables. For any unit vectors $u, u' \in S^{d-1}$ and $\delta \in (0, 1)$,

$$
\Pr \left( \| (Xu)^\top - (Xu')^\top \|_2 \geq \sqrt{C_0 \log \log(n) + \sqrt{8 \ln(1/\delta)}} \right) \leq \delta.
$$

The proof of Lemma 15 uses the following lemmas.

**Lemma 16** (Corollary 6.14 in [3]). There is an absolute constant $C > 0$ such that the following holds. If $n \geq 3$, $\mu$ is the standard Gaussian measure on $\mathbb{R}$, and $\mu_n$ is the empirical measure for a size-$n$ i.i.d. sample from $\mu$, then

$$
\mathbb{E} \left[ W_2(\mu_n, \mu)^2 \right] \leq \frac{C \log \log(n)}{n}.
$$

**Lemma 17** (Eq. 2.35 in [13]). Let $Z \sim N(0, I_p)$ be a standard normal random vector in $\mathbb{R}^p$, and $f : \mathbb{R}^p \to \mathbb{R}$ be $L$-Lipschitz with respect to the Euclidean metric. Then for any $t > 0$,

$$
\Pr \left( f(Z) \geq \mathbb{E} f(Z) + t \right) \leq e^{-t^2/(2L^2)}.
$$

**Proof of Lemma 15.** Fix unit vectors $u$ and $u'$. Observe that the entries of each of $Xu$ and $Xu'$ comprises an i.i.d. sample from $N(0,1) =: \mu$; let $\mu_n$ and $\nu_n$ denote the respective empirical measures. Define the function $f : \mathbb{R}^{n \times d} \to \mathbb{R}$ by

$$
f(A) := \| (Au)^\top - (Au')^\top \|_2.
$$

Then, by Lemma 11, the triangle inequality, Jensen’s inequality, and Lemma 16,

$$
\frac{\mathbb{E}f(X)}{\sqrt{n}} = \mathbb{E}W_2(\mu_n, \nu_n) \leq \mathbb{E}W_2(\mu_n, \mu) + \mathbb{E}W_2(\nu_n, \mu) \leq 2\sqrt{\mathbb{E}W_2(\mu_n, \mu)^2} \leq \sqrt{\frac{C_0 \log \log(n)}{n}}.
$$

Moreover, for any $A, A' \in \mathbb{R}^{n \times d}$,

$$
f(A) - f(A') \leq \| (Au)^\top - (Au')^\top - (A'u)^\top + (A'u')^\top \|_2
\leq \| (Au)^\top - (A'u)^\top \|_2 + \| (A'u)^\top - (A'u')^\top \|_2
\leq \| Au - A'u \|_2 + \| Au' - A'u' \|_2
\leq 2 \| A - A' \|_F,
$$

where the first two steps follow from the triangle inequality, the third step uses Lemma 11, and $\| \cdot \|_F$ denotes the Frobenius norm. Therefore, $f$ is 2-Lipschitz with respect to the Euclidean metric on $\mathbb{R}^{n \times d}$. By Lemma 17, for any $\delta \in (0, 1)$,

$$
\Pr \left( f(X) \geq \mathbb{E} f(X) + \sqrt{8 \ln(1/\delta)} \right) \leq \delta.
$$

Combining this with the upper bound on $\mathbb{E} f(X)$ completes the proof. \qed

**Lemma 18** (Eqs. 1.7.3 and 1.7.5 in [18]). Let $X_1, X_2, \ldots, X_n$ be i.i.d. draws from the uniform distribution on $[0, 1]$. For any $r \in [n]$,

$$
\mathbb{E}[X_{(r)}] = \frac{r}{n+1},
$$

and for any $r, s \in [n]$ with $r \leq s$,

$$
cov(X_{(r)}, X_{(s)}) = \frac{r}{n+1} \cdot \left(1 - \frac{s}{n+1}\right) \cdot \frac{1}{n+2}.
$$
Lemma 19. Let $U_1, U_2, \ldots, U_n$ be i.i.d. draws from the uniform distribution on $[-1/2, 1/2]$. Then

$$\Pr \left( \sum_{i=1}^{n} \left( U_{(1)} + U_{(n+1-i)} \right)^2 \geq 1 \right) \leq \frac{1}{2}.$$ 

Proof. It suffices to show the expectation bound

$$\mathbb{E} \left[ \sum_{i=1}^{n} \left( U_{(1)} + U_{(n+1-i)} \right)^2 \right] \leq \frac{1}{2},$$

since the claim then follows by Markov’s inequality. Expanding the square and using linearity of expectation gives

$$\mathbb{E} \left[ \sum_{i=1}^{n} \left( U_{(1)} + U_{(n+1-i)} \right)^2 \right] = 2 \sum_{i=1}^{n} \mathbb{E} \left[ U_{(1)}^2 \right] + 2 \sum_{i=1}^{n} \mathbb{E} \left[ U_{(i)} U_{(n+1-i)} \right]$$

$$= \frac{n}{6} + 2 \sum_{i=1}^{n} \mathbb{E} \left[ U_{(i)} U_{(n+1-i)} \right].$$

By Lemma 18, we have for $i \leq (n+1)/2$,

$$\mathbb{E} \left[ U_{(i)} U_{(n+1-i)} \right] = - \left( \frac{i}{n+1} - \frac{1}{2} \right)^2 + \frac{i^2}{(n+1)^2(n+2)};$$

and for $i > (n+1)/2$,

$$\mathbb{E} \left[ U_{(i)} U_{(n+1-i)} \right] = - \left( \frac{i}{n+1} - \frac{1}{2} \right)^2 + \frac{(n+1-i)^2}{(n+1)^2(n+2)}.$$

Plugging-in and simplifying gives

$$\mathbb{E} \left[ \sum_{i=1}^{n} \left( U_{(1)} + U_{(n+1-i)} \right)^2 \right] = \begin{cases} \frac{1}{2} \left( 1 - \frac{1}{n+1} \right) & \text{if } n \text{ is even}, \\ \frac{1}{2} \left( 1 - \frac{1}{n+2} \right) & \text{if } n \text{ is odd.} \end{cases}$$

$\square$