Lower bounds on the performance of polynomial-time algorithms for sparse linear regression

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Abstract

Under a standard assumption in complexity theory (NP \(\not\subseteq\) P/poly), we demonstrate a gap between the minimax prediction risk for sparse linear regression that can be achieved by polynomial-time algorithms, and that achieved by optimal algorithms. In particular, when the design matrix is ill-conditioned, the minimax prediction loss achievable by polynomial-time algorithms can be substantially greater than that of an optimal algorithm. This result is the first known gap between polynomial and optimal algorithms for sparse linear regression, and does not depend on conjectures in average-case complexity.

1. Introduction

The past decade has witnessed a flurry of results on the performance of polynomial-time procedures, many of them based on convex relaxation, that aim at solving challenging optimization problems that arise in statistics. The large majority of these results have been of the positive variety, essentially guaranteeing that a certain polynomial-time procedure produces an estimate with low statistical error; see the overviews (Chandrasekaran et al., 2012; Negahban et al., 2012) for results of this type. Moreover, in many cases, the resulting bounds have been shown to be minimax-optimal, meaning that no estimator can achieve substantially smaller error. More recently, however, this compelling story has begun to develop some wrinkles, in that gaps have been established between the performance of convex relaxations and the performance of optimal methods, notably in the context of sparse PCA and related sparse-low-rank matrix problems (e.g. Amini and Wainwright, 2009; Berthet and Rigollet, 2012, 2013; Krauthgamer et al., 2013; Oymak et al., 2012). The main contribution of this paper is to add an additional twist to this ongoing story, in particular by demonstrating a fundamental gap between the performance of polynomial-time methods and optimal methods for high-dimensional sparse linear regression. Notably, in contrast with the recent work of Berthet and Rigollet (2013) on sparse PCA, and subsequent results on matrix detection (Ma and Wu, 2013), both of which are based on average-case complexity, our result is based only on a standard conjecture in worst-case complexity theory.

Linear regression is a canonical problem in statistics: it is based on observing a response vector \(y \in \mathbb{R}^n\) and a design matrix \(X \in \mathbb{R}^{n \times d}\) that are linked via the linear relationship

\[
y = X\theta^* + w.
\]
Here the vector \( w \in \mathbb{R}^n \) is some form of observation noise, and our goal is to estimate the unknown vector \( \theta^* \in \mathbb{R}^d \), known as the regression vector. Throughout this paper, we focus on the standard Gaussian model, in which the entries of the noise vector \( w \) are i.i.d. \( N(0, \sigma^2) \) variates, and the case of deterministic design, in which the matrix \( X \) is viewed as non-random. In the sparse variant of this model, the regression vector is assumed to have a relatively small number of non-zero coefficients. In particular, for some positive integer \( k < d \), the vector \( \theta^* \) is said to be \( k \)-sparse if it has at most \( k \) non-zero coefficients. Thus, our model is parameterized by the triple \((n, d, k)\) of sample size \( n \), ambient dimension \( d \), and sparsity \( k \).

Given a \( k \)-sparse regression problem, the most direct approach would be to seek a \( k \)-sparse minimizer to the least-squares cost \( \|y - X\theta\|_2^2 \), thereby obtaining the \( \ell_0 \)-based estimator
\[
\hat{\theta}_{\ell_0} := \arg \min_{\theta \in \mathbb{B}_0(k)} \|y - X\theta\|_2^2.
\] (2)

Note that this estimator involves minimization over the \( \ell_0 \)-“ball”
\[
\mathbb{B}_0(k) := \{ \theta \in \mathbb{R}^d \mid \sum_{j=1}^{d} I[\theta_j \neq 0] \leq k \}
\] (3)
of \( k \)-sparse vectors. This estimator is not easy to compute in a brute force manner, since there are \( \binom{d}{k} \) subsets of size \( k \) to consider. More generally, it is known that computing a sparse solution to a set of linear equations is an NP-hard problem (Natarajan, 1995), and this intractability result has motivated the use of various heuristic algorithms and approximations. Recent years have witnessed an especially intensive study of methods based on \( \ell_1 \)-relaxation, including the basis pursuit and Lasso estimators (Tibshirani, 1996; Chen et al., 1998), as well as the Dantzig selector (Candes and Tao, 2007). Essentially, these methods are based on replacing the \( \ell_0 \)-constraint (3) with its \( \ell_1 \)-equivalent, in either a constrained or penalized form. All of these estimators are based on relatively simple convex optimization problems (linear or quadratic programs), and so can be computed in polynomial time. Moreover, in certain cases, the performance of \( \ell_1 \)-based methods have been shown to meet minimax-optimal lower bounds (Raskutti et al., 2011).

Despite this encouraging progress, there remain some intriguing gaps in the performance of \( \ell_1 \)-based procedures, perhaps most notably when assessed in terms of their mean-squared prediction error \( \frac{1}{n} \|X\hat{\theta} - X\theta^*\|_2^2 \). In order to bring this issue into sharper focus, given an estimator \( \hat{\theta} \), suppose that we evaluate its performance in terms of the quantity
\[
\mathcal{M}_{n,k,d}(\hat{\theta}; X) := \sup_{\theta^* \in \mathbb{B}_0(k)} \frac{1}{n} \mathbb{E}[\|X\hat{\theta} - X\theta^*\|_2^2],
\] (4)
where the design matrix \( X \) remains fixed, and expectation is taken over realizations of the noise vector \( w \sim N(0, \sigma^2 I_{n \times n}) \).

The criterion (4) assesses the performance of the estimator \( \hat{\theta} \) uniformly over the set of all \( k \)-sparse regression vectors. In terms of this uniform measure, the \( \ell_0 \)-based estimator (2) is known (Bunea et al., 2007; Raskutti et al., 2011) to satisfy the bound
\[
\mathcal{M}_{n,k,d}(\hat{\theta}_{\ell_0}; X) \lesssim \frac{\sigma^2 k \log d}{n},
\] (5)
where \( \lesssim \) denotes an inequality up to a universal constant, meaning independent of all problem dependent quantities. A noteworthy point is that the upper bound (5) holds for any fixed design matrix \( X \).
By way of contrast, most \(\ell_1\)-based guarantees involve imposing certain conditions on the design matrix \(X\). One of the most widely used conditions is the restricted eigenvalue (RE) condition (Bickel et al., 2009; Van De Geer and B"uhlmann, 2009), which lower bounds the quadratic form defined by \(X\) over a subset of sparse vectors (see equation (10) to follow for a precise definition). Under such an RE condition, it can be shown that the Lasso-based estimator \(\hat{\theta}_{\ell_1}\) satisfies a bound of the form

\[
\mathcal{M}_{n,k,d}(\hat{\theta}_{\ell_1}; X) \lesssim \frac{1}{\gamma^2(X)} \frac{\sigma^2 k \log d}{n},
\]

where \(\gamma(X) \leq 1\) denotes the restricted eigenvalue constant (10). Comparison of this bound to the earlier \(\ell_0\)-based guarantee (5) shows that the only difference is the RE constant, which is a measure of the conditioning of the matrix \(X\). However, from a fundamental point of view, the conditioning of \(X\) has no effect on whether or not a good sparse predictor exists; for instance, a design matrix with two duplicated columns is poorly conditioned, but the duplication would have no effect on sparse prediction performance.

The difference between the bounds (5) and (6) leaves open various questions, both about the performance of the Lasso (and other \(\ell_1\)-based methods), as well as polynomial-time methods more generally. Beginning with \(\ell_1\)-based methods, one possibility is that existing analyses of prediction error are overly conservative, but that the Lasso can actually achieve the bound (5), without the additional RE term. When the regression vector \(\theta^*\) has a bounded \(\ell_1\)-norm, then it is possible to achieve a prediction error bound that does not involve the RE constant (Bunea et al., 2007), but the resulting rate is “slow”, decaying as \(1/\sqrt{n}\) instead of the rate \(1/n\) given in equation (6). Foygel and Srebro (2011) asked whether this slow rate could be improved without an RE condition, and gave a partial negative answer in the case \(k = 2\), constructing a 2-sparse regression vector and a design matrix violating the RE condition for which the Lasso prediction error is lower bounded by \(1/\sqrt{n}\). In this paper, we ask whether the same type of gap persists if we allow for all polynomial-time estimators, instead of just the Lasso. Our main result is to answer this question in the affirmative: we show that there is a family of design matrices \(X_{bad}\) such that, under a standard conjecture in computational complexity (\(NP \not\subset P/poly\)), for any estimator \(\hat{\theta}_{poly}\) that can be computed in polynomial time, its mean-squared error is lower bounded as

\[
\mathcal{M}_{n,k,d}(\hat{\theta}_{poly}; X_{bad}) \gtrsim \frac{1}{\gamma^2(X_{bad})} \frac{\sigma^2 k^{1-\delta} \log d}{n},
\]

where \(\delta > 0\) is an arbitrarily small positive scalar. Consequently, we see that there is a fundamental gap between the performance of polynomial-time methods and that of the optimal \(\ell_0\)-based method.

The remainder of this paper is organized as follows. We begin in Section 2 with background on sparse linear regression and restricted eigenvalue conditions. We then introduce some background on complexity theory, followed by the statement of our main result in Section 3. The proof of the main theorem is given in Section 4, with more technical results deferred to the appendices.

2. Background and problem set-up

We begin with background on sparse linear regression, then introduce restricted eigenvalue conditions. These notions allow us to give a precise characterization of the mean-squared prediction error that can be achieved by the \(\ell_0\)-based algorithm and by a thresholded version of the Lasso algorithm.
2.1. Estimators for sparse linear regression

As previously described, an instance of the sparse linear regression problem is based on observing a pair \( (y, X) \in \mathbb{R}^n \times \mathbb{R}^{n \times d} \) that are linked via the linear equation (1), where the unknown regression vector \( \theta^* \in \mathbb{R}^d \) is assumed to be \( k \)-sparse, and so belongs to the \( \ell_0 \)-ball \( B_0(k) \). An estimator \( \hat{\theta} \) of the regression vector is a (measurable) function \( (y, X) \mapsto \hat{\theta} \in \mathbb{R}^d \), and our goal is to determine an estimator that is both \( k \)-sparse, and has low prediction error \( \frac{1}{n} \|X \hat{\theta} - X \theta^*\|_2^2 \). Accordingly, we let \( \mathcal{A}(k) \) denote the family of all estimators that return vectors in \( B_0(k) \). Note that the \( \ell_0 \)-based estimator \( \hat{\theta}_0 \), as previously defined in equation (2), belongs to the family \( \mathcal{A}(k) \) of estimators. The following result provides a guarantee for this estimator:

**Proposition 1 (Prediction error for \( \hat{\theta}_0 \))** There are universal constants \( c_j, j = 1, 2 \) such for any design matrix \( X \), the \( \ell_0 \)-based estimator \( \hat{\theta}_0 \) satisfies

\[
\frac{1}{n} \|X \hat{\theta}_0 - X \theta^*\|_2^2 \leq c_1 \frac{\sigma^2 k \log d}{n} \quad \text{for any } \theta^* \in B_0(k)
\]

with probability at least \( 1 - 2e^{-c_2 k \log d} \).

We also consider another member of the family \( \mathcal{A}(k) \)—namely, a thresholded version of the Lasso estimator (Tibshirani, 1996; Chen et al., 1998). The ordinary Lasso estimate \( \hat{\theta}_\lambda \) based on regularization parameter \( \lambda > 0 \) is given by

\[
\hat{\theta}_\lambda := \text{arg min}_{\theta \in \mathbb{R}^d} \left\{ \frac{1}{2n} \|y - X \theta\|_2^2 + \lambda \|\theta\|_1 \right\}.
\]

In general, this estimator need not be \( k \)-sparse, but a thresholded version of it can be shown to have similar guarantees. Overall, we define the **thresholded Lasso** estimator \( \hat{\theta}_\mathcal{T}_l \) based on the following two steps:

(a) Compute the ordinary Lasso estimate \( \hat{\theta}_\lambda \) with \( \lambda = 4\sigma \sqrt{\frac{\log d}{n}} \).

(b) Truncate \( \hat{\theta}_\lambda \) to its \( k \) entries that are the largest in absolute value, thereby obtaining the estimate \( \hat{\theta}_\mathcal{T}_l \).

By construction, the estimator \( \hat{\theta}_\mathcal{T}_l \) belongs to the family \( \mathcal{A}(k) \). The choice of regularization parameter given in step (a) is a standard one for the Lasso.

2.2. Restricted eigenvalues and \( \ell_1 \)-guarantees

We now define the notion of a (sparse) restricted eigenvalue (RE), and then discuss guarantees on the Lasso-based estimator \( \hat{\theta}_\mathcal{T}_l \) that hold under such an RE condition. Restricted eigenvalues are defined in terms of subsets \( S \) of the index set \( \{1, 2, \ldots, d\} \), and a cone associated with any such subset. In particular, letting \( S^c \) denote the complement of \( S \), we define the cone

\[
\mathbb{C}(S) := \left\{ \theta \in \mathbb{R}^d \mid \|\theta_{S^c}\|_1 \leq 3\|\theta_S\|_1 \right\}.
\]

Here \( \|\theta_{S^c}\|_1 := \sum_{j \in S^c} |\theta_j| \) corresponds to the \( \ell_1 \)-norm of the coefficients indexed by \( S^c \), with \( \|\theta_S\|_1 \) defined similarly. Note that any vector \( \theta^* \) supported on \( S \) belongs to the cone \( \mathbb{C}(S) \); in addition, it includes vectors whose \( \ell_1 \)-norm on the “bad” set \( S^c \) is small relative to their \( \ell_1 \)-norm on \( S \).
**Definition 1 (Restricted eigenvalue (RE) condition)** Given triplet \((n, d, k)\), the matrix \(X \in \mathbb{R}^{n \times d}\) is said to satisfy a uniform \(\gamma\)-RE condition if
\[
\frac{1}{n} \|X\theta\|_2^2 \geq \gamma \|\theta\|_2^2 \quad \text{for all } \theta \in \bigcup_{|S|=k} C(S).
\] (10)
Moreover, the restricted eigenvalue constant of \(X\), denoted by \(\gamma(X)\), is the greatest \(\gamma\) such that \(X\) satisfies the condition (10).

The RE condition (10) and related quantities have been studied extensively in past work on basis pursuit and the Lasso (e.g. Bickel et al., 2009; Meinshausen and Yu, 2009; Raskutti et al., 2011; see the paper (Van De Geer and Bühlmann, 2009) for an overview of the different types of RE parameters. Note that it characterizes the curvature of the quadratic form specified by \(X\) when restricted to a certain subset of relatively sparse vectors. When the RE constant \(\gamma(X)\) is close to zero, there are relatively sparse vectors \(\tilde{\theta}\) such that \(\|X\tilde{\theta} - X\theta^*\|_2\) is small but \(\|\tilde{\theta} - \theta^*\|_2\) is large. Given that we observe only a noisy version of the product \(X\theta^*\), it is then difficult to distinguish \(\theta^*\) from other sparse vectors, which makes estimating \(\theta^*\) from the data difficult. Thus, it is natural to impose an RE condition if the goal is to produce an estimate \(\hat{\theta}\) such that the \(\ell_2\)-norm error \(\|\hat{\theta} - \theta^*\|_2\) is small. Indeed, for \(\ell_2\)-norm estimation, Raskutti et al. (2011) show that a closely related condition is necessary for any method.

In contrast, it is worth noting that the RE condition is not a necessary condition for minimizing the prediction loss \(\|X\hat{\theta} - X\theta^*\|_2\), since an estimator far apart from \(\theta^*\) may still achieve small prediction error. However, the RE condition turns to be an important criterion for \(\ell_1\)-based methods to guarantee good prediction performance. Under the normalization condition:
\[
\frac{\|X\theta\|_2^2}{n} \leq \|\theta\|_2^2 \quad \text{for all } \theta \in \mathbb{B}_0(2k),
\] (11)
the following result provides such a guarantee for the thresholded Lasso estimator:

**Proposition 2 (Prediction error for thresholded Lasso)** There are universal constants \(c_j, j = 3, 4\) such that for any design matrix \(X\) satisfying the normalization condition (11) and having the RE constant \(\gamma(X) > 0\), the thresholded Lasso estimator \(\hat{\theta}_{TL}\) satisfies
\[
\frac{1}{n} \|X\hat{\theta}_{TL} - X\theta^*\|_2^2 \leq \frac{c_3}{\gamma^2(X)} \frac{\sigma^2 k \log d}{n} \quad \text{for any } \theta^* \in \mathbb{B}_0(k)
\] (12)
with probability at least \(1 - 2e^{-c_4 k \log d}\).

See Appendix B for the proof of this claim. Apart from different numerical constants, the main difference between the guarantee (12) and our earlier bound (8) for the \(\ell_0\)-based estimator is the \(1/\gamma^2(X)\) term. The RE constant \(\gamma(X)\) is a dimension-dependent quantity, since it is a function of the \(n \times d\) design matrix.

### 3. Main result and its consequences

Thus far, we have considered two estimators in the class \(A(k)\)—namely, the \(\ell_0\)-constrained estimator \(\hat{\theta}_{\ell_0}\) and the thresholded Lasso estimator \(\hat{\theta}_{TL}\). We also proved associated guarantees on their prediction error (Propositions 1 and 2 respectively), showing a \(1/\gamma^2(X)\) gap between their respective guarantees. Our main result shows that this gap is not a coincidence: more fundamentally, it
is a characterization of the gap between an optimal algorithm and the class of all polynomial-time algorithms.

In order to state our main result, we need to make precise a particular notion of a polynomial-efficient estimator. Since the observation \((X, y)\) consists of real numbers, any efficient algorithm can only take a finite-length representation of the input. Consequently, we need to introduce an appropriate notion of discretization, as has been done in past work on matrix detection (Ma and Wu, 2013). We begin by defining an appropriate form of input quantization: for any input value \(x\) and integer \(\tau\), the operator

\[
[x]_\tau := 2^{-\tau}[2^\tau x]
\]

represents a \(2^{-\tau}\)-precise quantization of \(x\). (Here \([u]\) denotes the largest integer smaller than or equal to \(u\).) Given a real value \(x\), an efficient estimator is allowed to take \([x]_\tau\) as its input for some finite choice \(\tau\). We denote by \(\text{size}(x; \tau)\) the length of binary representation of \([x]_\tau\), and denote by \(\text{size}(X, y; \tau)\) the total length of the discretized matrix vector pair \((X, y)\).

The following definition of efficiency is parameterized in terms of three quantities: (i) a positive integer \(b\), corresponding to the number of bits required to implement the estimator as a computer program; (ii) a polynomial function \(G\) of the triplet \((n, d, k)\), corresponding to the discretization accuracy of the input, and (iii) a polynomial function \(H\) of input size, corresponding to the runtime of the program.

**Definition 2 (Polynomial-efficient estimators)** Given polynomial functions \(G : (\mathbb{Z}_+)^3 \to \mathbb{R}_+\), \(H : \mathbb{Z}_+ \to \mathbb{R}_+\) and a positive integer \(b \in \mathbb{Z}_+\), an estimator \((y, X) \mapsto \hat{\theta}(y, X)\) is said to be \((b, G, H)\)-efficient if:

- It can be represented by a computer program that is encoded in \(b\) bits.
- For every problem of scale \((n, d, k)\), it accepts inputs quantized to accuracy \([\cdot]_\tau\) where the quantization level is bounded as \(\tau \leq G(n, d, k)\).
- For every input \((X, y)\), it is guaranteed to terminate in time \(H(\text{size}(X, y; \tau))\).

According to this definition, if we choose a sufficiently large code length—say \(b = 10^{16}\)—and polynomial functions \(G\) and \(H\) that grow sufficiently fast—say \(G(n, d, k) = (ndk)^{100}\) and \(H(s) = s^{100}\)—then the class of \((b, G, H)\)-efficient algorithms allow estimators to take sufficiently accurate input, and covers all estimators that are reasonably efficient in terms of storage and running time.

To present the main result, we require a few more notions from complexity theory. See the book (Arora and Barak, 2009) for a more detailed introduction. In complexity theory, the class \(\mathbf{P}\) corresponds to problems that are solvable in polynomial time by a Turing machine. A closely related class denoted by \(\mathbf{P}/\text{poly}\), corresponds to all problems solvable in polynomial time by a Turing machine with a so-called advice string (meaning a side-input to the machine) that is of polynomial length. Although it is known that the class \(\mathbf{P}/\text{poly}\) is strictly bigger than the class \(\mathbf{P}\) (e.g. Arora and Barak, 2009), it is widely believed that \(\mathbf{NP} \not\subset \mathbf{P}/\text{poly}\). Accordingly, throughout this paper, we impose the following assumption:

**Assumption A** The class \(\mathbf{NP}\) is not contained in the class \(\mathbf{P}/\text{poly}\).

Based on Assumption A, we are ready to present the main result. In this statement, we use \(c_j, j = 5, 6\) to denote universal constants independent of \((n, d, k)\), \((F, G, H)\) and \((\gamma, \sigma, \delta)\).
Theorem 3 If \( NP \not\subseteq \mathsf{P/poly} \), then for any positive integer \( b \), any scalar \( \delta \in (0, 1) \), any polynomial functions \( G : (\mathbb{Z}^+)_3 \to \mathbb{R}_+ \) and \( F, H : \mathbb{Z}_+ \to \mathbb{R}_+ \), there is a sparsity level \( k \geq 1 \) such that the following holds:

For any dimension \( d \in [4k, F(k)] \), any sample size \( n \) in the interval \( [c_5 k \log d, F(k)] \), and any polynomial functions \( G : (\mathbb{Z}^+)_3 \to \mathbb{R}_+ \) and \( F, H : \mathbb{Z}_+ \to \mathbb{R}_+ \), there is a sparsity level \( k \geq 1 \) such that the following holds:

(a) It satisfies the normalization condition \( (11) \), and has an RE constant \( \gamma(X) \) that is bounded as \( |\gamma(X) - \gamma| \leq 2^{-G(n,d,k)} \).

(b) For any \((b, G, H)\)-efficient estimator \( \hat{\theta} \in A(k) \), the mean-squared prediction error is lower bounded as

\[
\max_{\theta^* \in \mathbb{B}_0(k)} \mathbb{E}\left[ \frac{\|X(\hat{\theta} - \theta^*)\|_2^2}{n} \right] \geq \frac{c_6 \sigma^2 k^{1-\delta} \log d}{n}.
\] (14)

To understand the consequence of Theorem 3, suppose that we have chosen \( b, F, G, H \) sufficiently large and have chosen \( \delta \) very close to zero. Under this setting, Theorem 3(b) shows that as long as the triplet \((n, d, k)\) is sufficiently large, then there is a explicitly constructed design matrix \( X \) such that any \((b, G, H)\)-polynomial-efficient estimator has prediction risk lower bounded by inequality \( (14) \). Part (a) guarantees that the constructed design matrix \( X \) satisfies the normalization \( (11) \) and RE conditions \( (10) \), so that Proposition 2 can be applied to the thresholded Lasso estimator for this instance of the sparse regression problem. Since the parameter \( \delta \in (0, 1) \) may be chosen arbitrarily close to zero, the lower bound \( (14) \) essentially matches the upper bound of Proposition 2, thereby confirming that Theorem 3 gives a tight lower bound.

Overall, Theorem 3 establishes that the inverse dependence on the RE constant \( \gamma(X) \) is unavoidable for the class of polynomial-time algorithms. In contrast, the \( \ell_0 \)-based estimator—a method that is not polynomially efficient—does not exhibit this dependence, as shown by Proposition 1.

It is worth contrasting Theorem 3 with past work on the computational hardness of sparse linear systems. As mentioned previously, Natarajan (1995) showed that finding sparse solutions to linear systems is an NP-hard problem. To contrast with our result, this earlier result applies to the problem of solving \( X\theta = y \), where \( X \) is the worst-case matrix selected by an adversary who knows the algorithm in advance. In contrast, in our Theorem 3, the design matrix \( X \) is fixed ahead of time, and the associated hardness result applies to all polynomial-time algorithms. Furthermore, Theorem 3 requires \( X \) to satisfy the normalization condition \( (11) \) as well as the RE condition \( (10) \), so that the lower bound is achievable by the Lasso-based estimator in Proposition 2. These two conditions are not satisfied in the earlier construction (Natarajan, 1995). Finally, in Theorem 3 we have proved a lower bound of the order \( \frac{\sigma^2 k \log d}{\gamma(X)n} \), which is stronger than the \( \frac{1}{n} \) lower bound proved in the original paper (Natarajan, 1995).

4. Proof of Theorem 3

We now turn to the proof of our main result. In broad outline, the proof involves three main steps. We begin in Section 4.1 by constructing a particular matrix \( M \) for which the sparse linear regression problem is NP-hard, doing so by comparison to the exact 3-set cover problem. In Section 4.2, we then extend this worst-case hardness result to the probabilistic setting that is of interest in Theorem 3. Section 4.3 is devoted to the actual construction of the design matrix \( X \), using the matrix \( M \) as a
building block, and verifies that it satisfies the conditions in part (a). Sections 4.4 contains the proof of part (b) of the theorem. In all cases, we defer the proofs of more technical lemmas to the appendices.

4.1. An NP-Hard Problem

In this section, we construct a linear system and prove that solving it under a specific sparsity condition is NP-hard. The NP-hardness is establishing by reducing the exact 3-set cover problem (Natarajan, 1995) to this linear regression problem. The exact 3-set cover (X3C) problem is stated as follows: given an integer \( m \geq 3 \) divisible by three, a set \( S = \{1, \ldots, m\} \) and a collection \( C \) of 3-element subsets of \( S \), find \( m/3 \) subsets in \( C \) that exactly cover \( S \), assuming such an exact cover exists.

Throughout our development, we make use of the convenient shorthand notation \( p := \binom{m}{3} \), as well as \([m] := \{1,2,\ldots,m\}\) and \([p] := \{1,2,\ldots,p\}\). The following lemma shows the existence of a sequence of matrices \( \{M\} \), such that solving a sparse vector \( u \) based on observing \( (M,Mu) \) is NP-hard. The proof is deferred to Appendix A.1.

**Lemma 4** Given any integer \( m \in \mathbb{Z}_+ \) divisible by three, we can construct a particular matrix \( M \in \{-1,0,1\}^{(m+3p) \times 4p} \) and define the associated set of vectors

\[
\mathcal{V} := \left\{ v \in \mathbb{R}^{m+3p} \mid v = Mu \text{ for some } u \in \{0,1\}^{4p} \cap B_0(m/3+p) \right\},
\]

such that the following problem is NP-hard: Given any matrix \( M \) constructed above and any vector \( v \in \mathcal{V} \), compute a vector \( u \in \{0,1\}^{4p} \cap B_0(m/3+p) \) such that \( \|Mu - v\|_2 < \frac{1}{2} \).

4.2. Probabilistic Hardness

As in the previous section, we use the shorthand notation \( p = \binom{m}{3} \) throughout. In this section, we extend the result of Section 4.1 to the probabilistic case. In particular, we show that there is always a “bad” distribution \( Q \) over \( u \) such that solving for \( u \) based on observing \( (M,Mu) \) is hard in the average case. The proof is deferred to appendix A.2.

**Lemma 5** Let \( f, g \) and \( h \) be arbitrary polynomial functions defined on \( \mathbb{Z}_+ := \{1,2,\ldots\} \). Then for any \( m_0 \in \mathbb{Z}_+ \), there is an integer \( m > m_0 \) and distribution \( Q \) over \( u \in \{0,1\}^{4p} \cap B_0(m/3+p) \) such that for any solver \( \hat{u} \) that can be encoded in \( f(m) \) bits and terminates in \( g(m) \) time, we have

\[
Q\left[ \|M\hat{u} - Mu\|_2 < \frac{1}{2} \right] \leq \frac{1}{h(m)},
\]

(15)

4.3. Proof of part (a)

We now turn to the construction of the design matrix \( X \) specified in Theorem 3. Lemma 5 implies that given arbitrary polynomial functions \( f, g \) and \( h \), there is a positive integer \( m > m_0 \) (where \( m_0 \) can be arbitrarily large), a matrix \( M \in \mathbb{R}^{(m+3p) \times 4p} \) with \( p = \binom{m}{3} \), and a distribution \( Q \) over the set of binary vectors \( u \in \{0,1\}^{4p} \cap B_0(m/3+p) \) such that solving for \( u \) based on observing \( (M,Mu) \) is hard in the average case. We leave \( f, g, h \) and \( m_0 \) to be specified later, assuming that they are sufficiently large. Based on integer \( m \) and for a second positive integer \( t \) to be chosen, we define the sparsity level

\[
k := t\left(\frac{m}{3} + p\right) = t\left(\frac{m}{3} + \binom{m}{3}\right).
\]
Using the matrix \(M\) as a building block, we construct a design matrix \(X \in \mathbb{R}^{n \times d}\) that satisfies the conditions of Theorem 3. To this end, our first step is to construct a matrix \(A_k\). Recall the matrix \(M \in \mathbb{R}^{(m+3p) \times 4p}\) constructed in Section 4.1. For the given integer \(t > 1\), we consider the rescaled matrix \(\sqrt{t}M\) and replicate this matrix \(t\) times in diagonal blocks to build the following matrix:

\[
A_k := \text{blkdiag} \left\{ \sqrt{t}M, \sqrt{t}M, \ldots, \sqrt{t}M \right\} \in \mathbb{R}^{3k \times 4pt}.
\]

By construction, the matrix \(A_k\) has \(3k = t(m + 3p)\) rows and \(4pt\) columns. Since \(4pt \leq 4k \leq d\), the matrix

\[
B_k := \frac{1}{2} \left[ A_k \ 0_{3k \times (d - 4pt)} \right]
\]

has dimensions \(3k \times d\).

We now use \(B_k\) to construct the design matrix \(X \in \mathbb{R}^{n \times d}\). (To simplify the argument, we assume that \(n\) is divisible by \(6k\); when this condition does not hold, an obvious modification of the argument yields the same result modulo different constants.) Assume that \(R \in \mathbb{R}^{(n/2) \times d}\) is a random Gaussian matrix whose rows are sampled i.i.d. from the Gaussian distribution \(N(0, I_d\times d)\). We define a parameterized family of random matrices \(\{C_x \in \mathbb{R}^{n \times d}, x \geq 0\}\) such that:

- its top \(n/2\) rows consist of \(n/6k\) copies of the matrix \(B_k\).
- for each \(x \geq 0\), the bottom \(n/2\) rows of \(C_x\) are equal to \(xR\).

Given a matrix \(C \in \mathbb{R}^{n \times d}\), let \(C^\uparrow\) and \(C^\downarrow\) represent its top \(n/2\) rows and the bottom \(n/2\) rows, respectively. In addition, we use \([C]_L\) to represent the matrix formed by quantizing every entry of \(C\) by the operator \([\cdot]_L\). In addition, we define \(\varepsilon := 2^{-G(n,d,k)}\) as a shorthand quantity and note that \(\gamma \in [\varepsilon, 1/24\sqrt{2}]\).

The following lemma shows that there exists a particular value of \(x\) such that the matrix \(C_x\) has nice properties.

**Lemma 6** For any integer \(L \geq \max\{\log(12\sqrt{d}), \log(\sqrt{nd}/\varepsilon)\}\), there is a realization of the random matrix \(R\) and a particular value of \(x\) such that the matrix \([C_x]_L\) has RE constant \(|\gamma([C_x]_L) - \gamma| \leq \varepsilon\), and satisfies the upper bounds

\[
\frac{||[C_x]_L\theta||_2}{\sqrt{n}} \leq ||\theta||_2 \quad \text{and} \quad \frac{||[C_x^\downarrow]_L\theta||_2}{\sqrt{n}} \leq 25\gamma ||\theta||_2 \quad \text{for all} \ \theta \in \mathbb{B}_0(2k).
\] (16)

According to Lemma 6, if we choose an integer \(L \geq \max\{\log(12\sqrt{d}), \log(\sqrt{nd}/\varepsilon)\}\) and define the design matrix \(X := [C_x]_L\), then the matrix \(X\) satisfies part (a) of the theorem. We leave the concrete value of \(L\) to be specified later, assuming now that it is sufficiently large.

4.4. Proof of part (b)

We begin with some notation for this section. Let \(X^\uparrow\) and \(X^\downarrow\) represent the top \(n/2\) rows and the bottom \(n/2\) rows of \(X\), respectively. Given a vector \(\theta \in \mathbb{R}^d\), we take its first \(4pt\) coordinates and partition them evenly into \(t\) segments, each of length \(4p\). For \(i \in [t]\), we use \(\theta_i \in \mathbb{R}^{4p}\) to denote the \(i^{th}\) segment, corresponding to the coordinates indexed \(4p(i - 1) + 1\) to \(4pi\). Using this notation, the proof involves the subvectors \(\{\hat{\theta}_i, i \in [t]\}\) and \(\{\hat{\theta}_i^*, i \in [t]\}\) defined by \(\hat{\theta}\) and \(\theta^*\), respectively.
Lemma 5 guarantees the existence of a distribution \( Q \) over \( \{0, 1\}^4 \cap \mathbb{E}_0(\frac{m}{4} + p) \) such that, if \( u^* \) is sampled from \( Q \), then solving \( Mu = Mu^* \) is computationally hard. Now define the quantity

\[
L := \left\lceil \max \left\{ \log(12\sqrt{d}), \log(\sqrt{nd}/\varepsilon), \log(2\sqrt{k}/t/r) \right\} \right\rceil \quad \text{and} \quad \rho := \frac{r}{\sqrt{m/3 + p}},
\]

(17)

where \( r > 0 \) is a constant to be specified later. Suppose that each subvector \( \theta^*_i / \rho \) is independently drawn from the distribution \( Q \)—to be precise, a subvector is drawn from the distribution \( Q \), and is rescaled by a factor of \( \rho \) to obtain \( \theta^*_i \). Note that

\[
\max_{\theta^* \in \mathbb{R}^d} \mathbb{E}_{\theta^*} \left[ \|X\hat{\theta}(X, y) - X\theta^*\|_2^2 \right] \geq \mathbb{E}_{\theta} \left[ \|X\hat{\theta}(X, y) - X\theta^*\|_2^2 \right],
\]

so that it suffices to lower bound the right-hand side.

Denote by \( P_1 \) the problem defined in Lemma 5, and denote by \( P_2 \) an instance of the linear regression problem with fixed design matrix \( X \). Our reduction from \( P_1 \) to \( P_2 \) consists of two main steps:

- We first define an intermediate problem \( P_2' \) that is close to but slightly different from our regression problem \( P_2 \). As we demonstrate, any \( P_2' \) instance can be constructed from a \( P_1 \) instance, and hence by Lemma 5, it is hard to solve \( P_2' \) up to a particular accuracy.

- Our next step is to show that (a) any estimator \( \hat{\theta} \) solving the regression problem \( P_2 \) also solves the intermediate problem \( P_2' \); and (b) the prediction error associated with problems \( P_2 \) and \( P_2' \) are of the same magnitude. Consequently, there is a lower bound on the error of \( \hat{\theta} \) for solving the regression problem \( P_2 \).

In detail, there are two main differences between the linear equation \( Mu = Mu^* \) from Lemma 5 and the linear observation model (see equation (18)) that underlies \( P_2' \). First, in addition to the basic matrix \( M \), the design matrix \( X \) contains additional \( \sqrt{n}/2 \) rows consisting of independent Gaussian entries, Second, in the linear regression problem, the response vector \( y \) is corrupted by Gaussian noise, whereas the original linear system \( Mu = Mu^* \) is based on a noiseless response vector (namely, \( Mu^* \)). Our reduction from \( P_1 \) to \( P_2' \) bridges these gaps.

**Constructing a \( P_2' \) instance from a \( P_1 \) instance:** Suppose that we are given an instance \( (M, Mu^*) \) of problem \( P_1 \), where the vector \( u^* \) is sampled from \( Q \). The input of the \( P_2' \) problem takes the form \((X, y')\), where \( X \in \mathbb{R}^{n \times d} \) is a matrix depending on \( M \), and \( y' \in \mathbb{R}^n \) is a vector depending on \((M, Mu^*)\). The goal of both \( P_1 \) and \( P_2' \) is to recover the unknown vector \( u^* \). We now construct a set of \( P_2' \) instances \((X, y_i)\), indexed by \( i \in \{1, 2, \ldots, t\} \), and we show that solving any of them provides a solution to \( P_1 \). We have already discussed how to construct the design matrix \( X \) in Section 4.3, so that it remains to construct the response vector \( y'_i \). By focusing on a concrete index \( i \), we suppress the dependence on the index \( i \) and write \( y'_i := y'_i \) for the conciseness of writing. Recalling our notation \( \{\theta_j, j \in [t]\} \) for the subvectors of the vector \( \theta \in \mathbb{R}^d \), our construction consists of the following three steps:

(i) Take as input a constant vector \( \tilde{\theta}^i \in \mathbb{R}^d \) provided by an external source.

(ii) Form the vector \( \xi^i \in \mathbb{R}^d \) with subvectors specified as

\[
\xi^i_{\ell} := \begin{cases} 
\rho u^* & \text{if } \ell = i \\
\tilde{\theta}^i & \text{if } \ell \in [t] \setminus i,
\end{cases}
\]
with all other coordinates set to zero.

(iii) Form the response vector

$$y' = \begin{bmatrix} X^\top \xi^i \\ X^\top \theta^i \end{bmatrix} + w,$$

where $w \sim N(0, \sigma^2I_{n \times n})$ is a Gaussian noise vector. Note that this is not a standard linear regression model, since a different regression vector is used for the top and bottom $n/2$ responses (the regression vectors $\xi^i$ and $\theta^i$, respectively).

It is important to note that even though $\xi^i$ contains $pu^*$ as a subvector, constructing the matrix vector product $X^\top \xi^i$ only requires the knowledge of matrix vector product $Mu^*$. The external constant vector $\theta^i$ from step (i) can be seen as an “advice string”, as permitted for solving problems in the $\mathcal{P}/\text{poly}$ class (recall the definition in Section 3). Note that its value is independent of the $P_1$ instance $(M, Mu^*)$.

**Solving $P_2'$ via an estimator for $P_2$:** Having constructed a $P_2'$ instance $(X, y')$ from $(M, Mu^*)$, we now show how an estimator for solving a $P_2$ instance can be used to solve the $P_2'$ instance, and hence to solve the original $P_1$ instance. More precisely, given an estimator $\hat{\theta}$ for problem $P_2$, we output $u^i_\hat{\theta} = \hat{\theta}(X, y')/\rho$ as a solution to $P_2'$, estimating the unknown vector $u^*$. The following lemma provides an upper bound on the probability of the “success event” that the outputted solution $u^i_\hat{\theta}$ is $(k/t)$-sparse, and achieves a small prediction error. Its proof exploits the hardness result from Lemma 5.

**Lemma 7** Assume that $\max\{r, 1/r\}$ is bounded by a polynomial function of $(n, d, k)$, and let $h$ be an arbitrary polynomial function. Then for any vector $\theta \in \mathbb{R}^d$ such that $\|\theta\|_\infty \leq \rho$ and $\theta = [\theta]_L$, we have

$$\mathbb{P}\left[u^i_\hat{\theta} \in \mathbb{B}_0(k/t) \text{ and } \|Mu^i_\hat{\theta} - Mu^*\|_2 < 1/2 \mid \hat{\theta} = \theta\right] \leq \frac{1}{h(m)}.$$

(19)

Note that upper bound (19) holds for any assignment to $\hat{\theta}^i$. Thus, it holds if we randomize it by drawing $\hat{\theta}^i \in \mathbb{R}^d$ randomly from a probability distribution. In particular, suppose that we construct an arbitrary vector $\hat{\theta}^i$ as follows:

- for each $j \in [t] \setminus \{i\}$, we draw the $j^{th}$ sub-vector $\hat{\theta}^j/\rho \in \mathbb{R}^{4p}$ independently from distribution $Q$.
- All other coordinates of $\hat{\theta}^i$ are set to zero.

It is straightforward to verify that the random vector $\hat{\theta}^i$ constructed in this way satisfies $\|\hat{\theta}^i\|_\infty \leq \rho$ and $\hat{\theta}^i = [\hat{\theta}^i]_L$. Thus, integrating over the randomness of $\hat{\theta}^i$, Lemma 7 implies

$$\mathbb{P}\left[u^i_\hat{\theta} \in \mathbb{B}_0(k/t) \text{ and } \|Mu^i_\hat{\theta} - Mu^*\|_2 < 1/2\right] \leq \frac{1}{h(m)}.$$

(20)

We assume for the rest of the proof that $\hat{\theta}^i$ is drawn from the specified distribution. By this construction, it is important to note that $\xi^i$ and $\theta^i$ have the same distribution, because each of their $4p$-length blocks are independently drawn from the same distribution $Q$. To prove the Theorem 3, it remains to convert the hardness result (20) to a lower bound on the mean-squared error $\mathbb{E}\left[\|X\hat{\theta}(X, y) - X\theta^*\|_2^2\right]$. 
Proving a lower bound on P2 via the hardness of P2: For the purposes of analysis, it is convenient to introduce the auxiliary response vector \( y^* := X\hat{\xi} + w \). Since the random vectors \( \hat{\xi} \) and \( \theta^* \) have the same distribution, the response vectors \( y := X\theta^* + w \) and \( y^* \) have the same distribution. We also define an auxiliary estimator \( u_i^\circ := \hat{\theta}_i(X, y^*)/\rho \). Since \( y \) and \( y^* \) share the same probability distribution, we note that the estimator \( u_i^\circ \) shares the same distribution as \( \hat{\theta}_i(X, y)/\rho \). Consequently, we have

\[
\mathbb{E} \left[ \frac{1}{n} \left\| X\hat{\theta}(X, y) - X\theta^* \right\|^2 \right] \geq \mathbb{E} \left[ \frac{1}{n} \left\| X\hat{\theta}(X, y) - X\theta^* \right\|^2 \right]
\]

\[
\quad \overset{(i)}{=} \frac{1}{n} \sum_{i=1}^{t} \mathbb{E} \left[ \left\| M\hat{\theta}(X, y) - M\theta^* \right\|^2 \right] = \frac{t\rho^2}{12k} \sum_{i=1}^{t} \mathbb{E} \left[ \left\| Mu_i^\circ - Mu^* \right\|^2 \right]. \tag{21}
\]

Here, equality (i) holds because the random vectors \( y \) and \( y^* \) are of the same order of magnitude. Note that the estimator \( u_i^\circ \) is equal to \( \hat{\theta}_i(X, y^*)/\rho \), whereas the estimator \( u_i^\circ \) takes \( (X, y^*) \) as input, whereas the estimator \( u_i\) takes \( (X, y) \) as input. Since the top \( n/2 \) coordinates of \( y_i^\circ \) are identical to those of \( y_i \), one might expect that the outputs of \( u_i^\circ \) and \( u_i \) are strongly correlated; it remains to prove this intuition.

Using \( y_{i\text{bot}}^\circ \) and \( y_{i\text{bot}} \) to denote the bottom \( n/2 \) co-ordinates of \( y_i^\circ \) and \( y_i \) respectively, it is easy to check that \( y_{i\text{bot}}^\circ \sim N(X^i\hat{\xi}, \sigma^2 I_{n/2 \times n/2}) \) and \( y_{i\text{bot}} \sim N(X^i\hat{\theta}, \sigma^2 I_{n/2 \times n/2}) \). The means of \( y_{i\text{bot}}^\circ \) and \( y_{i\text{bot}} \) differ by \( X(\hat{\xi} - \hat{\theta}) \). Moreover, we know by construction that \( \|\xi_i - \hat{\theta}^i\|_2 = \|\rho u_i\|_2 \leq r \). These observations are sufficient for lower bounding the right-hand side of inequality (21).

We now state two lemmas that characterize the properties of the estimator \( u_i^\circ \) and the estimator \( u_i \). The first lemma shows that a large subset of the estimators \( \{u_i, i \in [t]\} \) return sparse outputs with positive probability. See Appendix A.5 for the proof of this claim.

**Lemma 8** There is a set \( T \subset \{1, \ldots, t\} \) with cardinality \( |T| \geq \frac{t^2}{2k+2t} \) such that

\[
\mathbb{P} \left[ u_i \in B_0(k/t) \right] \geq e^{-1} \left\{ \frac{t}{2k+2t} - \frac{25\gamma \sqrt{n}}{\sigma} \right\} \quad \text{for each } i \in T.
\]

The second lemma upper bounds the prediction error of \( u_i \) in terms of the mean-squared prediction error of the auxiliary solver \( u_i^\circ \). See Appendix A.6 for the proof of this claim.

**Lemma 9** As long as \( \frac{\gamma \sqrt{n}}{\sigma} \leq 1/50 \), we are guaranteed that

\[
\|Mu_i^\circ - Mu^*\|_2^2 \leq \frac{2e\sigma}{25\gamma r \sqrt{n}} \mathbb{E} \left[ \|Mu_i^\circ - Mu^*\|_2^2 \right].
\]

with probability at least \( 1 - \frac{50\gamma \sqrt{n}}{\sigma} \).

To understand the consequence of Lemma 8, 9, we examine the probability that \( \|Mu_i^\circ - Mu^*\|_2 \) is lower bounded by \( 1/2 \). Performing some algebra and applying the union bound, we have

\[
\mathbb{P} \left[ \|Mu_i - Mu^*\|_2 \geq 1/2 \right] \geq \mathbb{P} \left[ u_i \in B_0(k/t) \right] - \mathbb{P} \left[ u_i \in B_0(k/t) \text{ and } \|Mu_i - Mu^*\|_2 < 1/2 \right]. \tag{22}
\]
By Lemma 8, there is a set $T$ with cardinality $|T| \geq \frac{t^2}{2k+t}$, such that $\mathbb{P}[u_i^e \in \mathcal{B}_0(k/t)]$ is lower bounded by a positive constant for all $i \in T$. Combining such lower bound with inequalities (20), (22), we obtain

$$\mathbb{P}\left[ \|Mu_i^* - Mu_i^e\|_2 \geq 1/2 \right] \geq e^{-1}\left\{ \frac{t}{2k+2t} - \frac{25\gamma r \sqrt{n}}{\sigma} \right\} - \frac{1}{h(m)}$$

for each $i \in T$. (23)

Combining Lemma 9 with inequality (23), we obtain a lower bound on the mean-squared prediction error of the auxiliary solver $u_i^e$. In particular, for each index $i \in T$, we have

$$\mathbb{E}\left[ \|Mu_i^e - Mu_i^e\|_2^2 \right] \geq \frac{1}{4} \cdot \frac{25\gamma r \sqrt{n}}{2e\sigma}$$

with probability at least $e^{-1}\left\{ \frac{t}{2k+2t} - \frac{25\gamma r \sqrt{n}}{\sigma} \right\} - \frac{50\gamma r \sqrt{n}}{\sigma} - \frac{1}{h(m)}$. We proceed by assigning particular values to the parameters $t$ and $r$. In particular, we set

$$r := \frac{\sigma}{25 \times 16 \gamma \sqrt{n}} \frac{t}{k + t}, \quad \text{and} \quad t := \left\lceil \left( \frac{m}{3} + p \right)^{\frac{1}{\alpha}} \right\rceil,$$ (25)

where $\alpha \in (0, 1)$ is to be chosen later. We note that the assignment to $r$ satisfies the assumption of Lemma 7. Recalling that $k = t \left( \frac{m}{3} + p \right)$ by definition, we have the sandwich relation $k^{1-\alpha} \leq t \leq k$, as well as the inequalities

$$e^{-1}\left( \frac{t}{2k+2t} - \frac{25\gamma r \sqrt{n}}{\sigma} \right) - \frac{50\gamma r \sqrt{n}}{\sigma} = \left( \frac{7}{16e} - \frac{1}{8} \right) k + t > 0.03 k^{-\alpha}, \quad \text{and} \quad (26a)$$

$$\frac{25\gamma r \sqrt{n}}{2e\sigma} \geq \frac{t}{32ek + 32et} \geq \frac{k^{-\alpha}}{64e}. \quad (26b)$$

By inequality (26a), we verify that the condition of Lemma 9 is satisfied. Combining with lower bound (24), and by choosing a polynomial function $h$ such that $1/h(m) < 0.03k^{-\alpha}$, there is a subset $T$ with cardinality $|T| \geq \frac{1}{4}k^{1-2\alpha}$ such that for each $i \in T$,

$$\mathbb{P}\left[ \mathbb{E}\left[ \|Mu_i^e - Mu_i^*\|_2^2 \right] \geq \frac{k^{-\alpha}}{256e} \right] > 0.$$ (27)

Since $\mathbb{E}[\|Mu_i^e - Mu_i^*\|_2^2]$ is a numerical value, the above inequality holds if and only if it is greater than or equal to $k^{-\alpha}_{256e}$. Plugging this lower bound into inequality (21), and substituting in

$$\rho = \left\lceil \frac{r}{\sqrt{m/3 + p}} \right\rceil \geq \frac{r}{\sqrt{k/t}} - 2^{-\alpha_L} \geq \frac{r}{2\sqrt{k/t}}$$

as well as our previous choice (25) of the parameter $r$, we find that

$$\mathbb{E}\left[ \frac{1}{n} \|X\hat{\Theta}(X, y) - X\Theta^*\|_2^2 \right] \geq \rho^2 t \sum_{i \in T} \mathbb{E}[\|Mu_i^e - Mu_i^*\|_2^2] \geq \frac{\rho^2 t}{12k} \frac{k^{1-2\alpha}}{3} \frac{k^{-\alpha}}{256e} \geq \frac{c \sigma^2 k^{1-7\alpha}}{4 \times 3 \times 256e} \geq \frac{c \sigma^2 k^{1-7\alpha}}{n},$$

where we have also used the sandwich inequalities $k^{1-\alpha} \leq t \leq k$. By assumption, the problem dimension is upper bounded as $d \leq F(k)$ where $F$ is a polynomial function. Since Lemma 5 allows $m_0$ to be arbitrary, we may choose it large enough so as to ensure that $k^{\alpha} \geq \log d$, or equivalently $k^{1-7\alpha} \geq k^{1-8\alpha} \log d$. Finally, setting $\alpha = \delta/8$ completes the proof of the theorem.
References


Appendix A. Technical lemmas for Theorem 3

In this appendix, we collect together the proofs of various technical lemmas involved in the proof of our main theorem.

A.1. Proof of Lemma 4

**Constructing matrix** $M$: Given an integer $m \in 3 \times \mathbb{Z}_+$, we begin by constructing the matrix $M \in \mathbb{R}^{(m+3p) \times 4p}$ in a blockwise manner, namely

\[
M = \begin{bmatrix} A & 0 \\ B & C \end{bmatrix}
\]  

(28)

where the submatrices have dimensions $A \in \mathbb{R}^{m \times p}$, $B \in \mathbb{R}^{3p \times p}$ and $C \in \mathbb{R}^{3p \times 3p}$.

We define the submatrices by considering all possible subsets of the form $(a, b, c) \in [m]^3$, where the elements are all required to be distinct. Since $p = \binom{m}{3}$, each such subset can be labeled with a unique index $j \equiv j_{abc} \in [p]$. For each $j \in [p]$, this indexing can be used to define the $j$th column of the submatrix $A$ as follows:

\[ A_{aj} = A_{bj} = A_{cj} = 1 \quad \text{and} \quad A_{dj} = 0 \quad \text{for all} \quad d \in [m] \setminus \{a, b, c\}. \]

In other words, the $j$th column of $A$ is the binary indicator vector for membership in the subset $\{a, b, c\}$ indexed by $j$.

The submatrices $B$ and $C$ are defined in terms of their rows. For each $j \in [p]$, we define the following three rows of $B$

\[ B_j = e_j, \quad B_{p+j} = e_j \quad \text{and} \quad B_{2p+j} = 0, \]

where $e_j \in \mathbb{R}^p$ is the $j$th canonical basis vector (i.e., with $j$-th entry equal to 1 and all other entries set to 0). Similarly, we define the rows of $C$ by

\[ C_j = -f_j, \quad C_{p+j} = f_{p+j} \quad \text{and} \quad C_{2p+j} = f_{2p+j}, \]

where $f_j \in \mathbb{R}^{3p}$ is the $j$th canonical basis vector in $\mathbb{R}^{3p}$. Having defined the matrix $M$, it remains to show that solving the linear system $Mu = v$ for all $v \in \mathcal{V}$ is NP-hard.

We reduce the X3C problem to the sparse linear regression problem. Let’s recall the statement of the X3C problem: given an integer $m \geq 3$ divisible by three, a set $\mathcal{S} = \{1, \ldots, m\}$ and a collection $\mathcal{C}$ of 3-element subsets of $\mathcal{S}$, the goal is to find $m/3$ subsets in $\mathcal{C}$ that exactly cover $\mathcal{S}$, assuming such an exact cover exists.

We summarize our strategy for proving the reduction as follow: Given $(\mathcal{S}, \mathcal{C})$, we use it to construct a response vector $v \in \mathbb{R}^{m+3p}$, and then consider the linear system $Mu = v$. We then show that any method that leads to a sparse vector $u$ such that $\|Mu - v\|_2 < \frac{1}{2}$ can also be used to solve the X3C problem.

**Constructing a response vector from X3C**: Given $(\mathcal{S}, \mathcal{C})$, we now construct the response vector $v$. Recalling that $v \in \mathbb{R}^{m+3p}$, we let the first $m$ coordinates of the vector $v$ equal to 1. Since any triplet of distinct elements $(a, b, c) \in [m]^3$ can be associated with a unique index $j_{abc} \in [p]$. We use this
correspondence to define the remaining $3p$ entries of $v$ as follows:

\begin{align*}
\text{If } (a, b, c) &\in C: \quad \text{set } v_{m+jabc} = 0, v_{m+p+jabc} = 1 \text{ and } v_{m+2p+jabc} = 0, \quad (29a) \\
\text{If } (a, b, c) &\notin C: \quad \text{set } v_{m+jabc} = 0, v_{m+p+jabc} = 0, \text{ and } v_{m+2p+jabc} = 1. \quad (29b)
\end{align*}

Assuming the existence of an exact cover, we now show that $v \in V$, in particular by constructing a vector $u \in \{0, 1\}^{4p} \cap B_{0}(\frac{n}{3} + p)$ such that $Mu = v$. Given a fixed triplet $(a, b, c)$, let us introduce the shorthand notation

$$
\alpha := u_{jabc}, \quad \beta := u_{p+iabc}, \quad \delta := u_{2p+jabc} \quad \text{and} \quad \gamma := u_{3p+jabc}.
$$

Observe that the linear equation $Mu = v$ holds if and only if the following conditions hold:

\begin{align*}
\text{For } (a, b, c) &\in C: \quad |\alpha - \beta| = |\alpha + \delta - 1|, \text{ and } |\gamma| = 0. \quad (31a) \\
\text{For } (a, b, c) &\notin C: \quad |\alpha - \beta| = 0, \quad |\alpha + \delta - 1| = 0, \quad \text{and} \quad |\gamma| = 1. \quad (31b) \\
\text{For any } i &\in S: \quad \text{exactly one triplet } (a, b, c) \text{ satisfies } i \in \{a, b, c\} \text{ and } |\alpha| = 1. \quad (31c)
\end{align*}

Thus, it is sufficient to construct a vector $u \in \{0, 1\}^{4p} \cap B_{0}(\frac{n}{3} + p)$ which satisfies the above conditions. If $(a, b, c) \notin C$, then we let $(\alpha, \beta, \delta, \gamma) = (0, 0, 0, 1)$. If $(a, b, c) \in C$ but it does not belong to the exact cover, then we let $(\alpha, \beta, \delta, \gamma) = (0, 0, 1, 0)$. Otherwise, if $(a, b, c) \in C$ and it is selected in the exact cover, then we let $(\alpha, \beta, \delta, \gamma) = (1, 1, 0, 0)$. Given these specifications, it is straightforward to verify that $u \in \{0, 1\}^{4p}$ and exactly $m/3 + p$ entries of $u$ are equal to 1. It is also easy to verify that $u$ satisfies conditions (31a)-(31c). Overall, we conclude that $v \in V$, and hence is a valid response vector for the sparse linear regression problem.

**Solving X3C using a sparse linear system solver:** Suppose that there is a sparse solution $u \in \{0, 1\}^{4p} \cap B_{0}(\frac{n}{3} + p)$ which satisfies $\|Mu - v\|_{2} < 1/2$, where the response vector $v$ was previously defined in equations (29a) and (29b). We now use $u$ to construct an exact 3-set cover. Recalling the notation (30), we observe that the condition $\|Mu - v\|_{2} < 1/2$ leads to the following restrictions:

\begin{align*}
\text{For } (a, b, c) &\in C: \quad |\alpha - \beta| < 1/2, \quad |\alpha + \delta - 1| < 1/2, \text{ and } |\gamma| < 1/2. \\
\text{For } (a, b, c) &\notin C: \quad |\alpha - \beta| < 1/2, \quad |\alpha + \delta| < 1/2, \text{ and } |\gamma - 1| < 1/2.
\end{align*}

Then, we prove the following two claims:

**Claim 1** There are exactly $m/3$ nonzero entries in $\{u_{1}, \ldots, u_{p}\}$. In addition, for any $i \in [p]$ such that $u_{i} > 0$, we must have $u_{i} > 1/2$.

**Proof** For any triplet $(a, b, c)$, we claim that there is at least one nonzero entry in $(\beta, \delta, \gamma)$. Thus, the fact that $u \in B_{0}(m/3 + p)$ implies that there are at most $m/3$ nonzero entries among $\{u_{1}, \ldots, u_{p}\}$. To prove this claim, consider two cases $(a, b, c) \notin C$ and $(a, b, c) \in C$. If $(a, b, c) \notin C$, then we have a restriction $|\gamma - 1| < 1/2$, which makes $\gamma \neq 0$. If $(a, b, c) \in C$, then we have restrictions $|\alpha - \beta| < 1/2$ and $|\alpha + \delta - 1| < 1/2$. Thus,

$$
1 > |\alpha - \beta| + |\alpha + \delta - 1| \geq |\beta + \delta - 1|,
$$

which implies that either $\beta$ or $\delta$ must be nonzero.

On the other hand, we claim that the number of nonzero entries in $\{u_{1}, \ldots, u_{p}\}$ is at least $m/3$. Assume that there are $k$ nonzero entries in $\{u_{1}, \ldots, u_{p}\}$. We notice that each nonzero entry
in \{u_1, \ldots, u_p\} selects a column in matrix \( A \) that has three nonzero coordinates. Thus, there are at most 3\(k \) nonzero entries in \((Mu)_i\) for all \( 1 \leq i \leq m \). Hence, \((Mu)_i\) must be greater than 1/2. Thus, we have 3\(k \geq m \) which establishes the claim.

Hence, there are exactly m/3 nonzero entries in \( \{u_1, \ldots, u_p\} \). Each nonzero entry covers exactly three coordinates of \( v \). Since all elements in \((Mu)_1, \ldots, (Mu)_m\) must be greater than 1/2, we conclude that the nonzero \( u_i \)'s must be greater than 1/2.

**Claim 2**  
*For any triplet \((a, b, c) \notin C\), we have \( u_{abc} = 0 \).*

**Proof** We proceed via proof by contradiction. Suppose that \((a, b, c) \notin C\) and \( u_{abc} \neq 0 \). By Claim 1, we have \( \alpha = u_{abc} > 1/2 \). The condition \( \|Mu - v\|_2 < 1/2 \) imposes the restrictions \( |\alpha - \beta| < 1/2 \) and \( |\gamma - 1| < 1/2 \), which in turn imply that \( \beta \neq 0 \) and \( \gamma \neq 0 \). In the proof of Claim 1, we have shown that for any triplet \((a, b, c)\) there is at least one nonzero entry in \((\beta, \delta, \gamma)\). For this specific triplet \((a, b, c)\), there are at least two nonzero terms among \((\beta, \delta, \gamma)\). Hence, the total number of nonzero entries in \( \{u_1, \ldots, u_p\} \) is at most m/3 - 1, which contradicts Claim 1.

According to Claim 2, all nonzero entries in \( \{u_1, \ldots, u_p\} \) correspond to 3-sets in \( C \). According to Claim 1, exactly m/3 of these 3-sets are selected. The condition \( \|Mu - v\|_2 < 1/2 \) implies that all entries in \((Mu)_1, \ldots, (Mu)_m\) are nonzero, which means that \( S \) is covered by the union of these 3-sets. Thus, these m/3 3-sets form an exact cover of \( S \).

**A.2. Proof of Lemma 5**

To prove Lemma 5, we rely on the following auxiliary result, which is proved in Appendix A.7:

**Lemma 10**  
*Let \( f, g \) be arbitrary polynomial functions, let \( m_0 \in \mathbb{Z}_+ \) be an arbitrary positive integer, and suppose that Assumption A holds. Then for any NP-hard problem \( P \), there is an input length \( m > m_0 \) such that any algorithm that can be encoded in \( f(m) \) bits and that terminates in \( g(m) \) time must fail on at least one input with probability greater than 1/2.*

We are now ready to prove the main result. We proceed via proof by contradiction. In particular, for a given \( m_0 \), suppose that for every \( m > m_0 \) and for any distribution over \( \{0, 1\}^{4p} \cap B_0(\frac{m}{3} + p) \), there is a solver \( \hat{u} \) that can be encoded in \( f(m) \) bits and terminates in \( g(m) \) time, and such that inequality \((15)\) is violated.

For a positive integer \( N \) to be specified, we now construct a sequence of distributions \( \{P_t\}_{t=1}^N \) that lead to a contradiction. Let \( P_1 \) the uniform distribution over \( S_1 := \{0, 1\}^{4p} \cap B_0(\frac{m}{3} + p) \) and let \( \hat{u}_1 \) be the solver which achieves \( P_1[\|M\hat{u} - Mu\|_2 < \frac{1}{2}] > \frac{1}{n(m)} \). For \( t \geq 1 \), we define the set \( S_{t+1} \) recursively as

\[
S_{t+1} = S_t \cap \{u \in S_0 \mid \|M\hat{u} - Mu\|_2 \geq \frac{1}{2}\}.
\]

When \( S_{t+1} \) is non-empty, we let \( P_{t+1} \) be the uniform distribution over \( S_{t+1} \), and let \( \hat{u}_{t+1} \) denote the solver that satisfies the bound \( P_{t+1}[\|M\hat{u} - Mu\|_2 < \frac{1}{2}] > \frac{1}{n(m)} \). If \( S_{t+1} \) is empty, then we simply set \( \hat{u}_{t+1} = \hat{u}_t \).
For any integer \( N \geq 1 \), this construction yields a sequence of solvers \( \{\hat{u}_t\}_{t=1}^N \), and we use them to define a combined solver \( u^*_N \) as follows:

\[
u^*_N = \begin{cases} \hat{u}_t & \text{for first } t \in \{1, \ldots, N\} \text{ such that } \|M\hat{u}_t - Mu\| < \frac{1}{2} \\ 0 & \text{otherwise.} \end{cases}
\]

(32)

If \( S_t \) is empty for some \( t \in [N] \), then \( \|Mu^*_N - Mu\|_2 < 1/2 \) for all \( u \in S_0 \). Otherwise, we may assume that \( S_t \) is non-empty for each \( t \in [N] \). Since \( S_t \) is the set of vectors \( u \) for which all the solvers \( \{\hat{u}_1, \ldots, \hat{u}_{t-1}\} \) fail, and \( P_t \) is the uniform distribution over \( S_t \), the chance of success for every solver \( \hat{u}_t \) in the construction (32) is at least \( \frac{1}{1/2} \). Consequently, if we consider the chance of success for the solver \( u^*_N \) under the uniform distribution over \( S_1 \), we have

\[
P_1\left[ \|Mu^*_N - Mu\|_2 < \frac{1}{2} \right] \geq 1 - \left( 1 - \frac{1}{h(m)} \right)^N \geq 1 - \left( \frac{1}{2} \right)^{4p+1},
\]

(33)

where inequality (ii) follows by choosing \( N := \left\lceil \frac{-4p+1}{\log_2(1/h(m))} \right\rceil \). Since the distribution \( P_1 \) is uniform over a support of at most \( 2^{4p} \) points, every point has probability mass no less than \( \left( \frac{1}{2} \right)^{4p} \). This fact, combined with the lower bound (33), implies that for any \( u \in S_1 \) the solver \( u^*_N \) satisfies the bound \( \|Mu^*_N - Mu\|_2 < 1/2 \) with probability greater than \( 1/2 \).

Note that the solver \( u^*_N \) can be encoded in \( N f(m) \) bits and it terminates in \( N g(m) \) time, both are polynomial functions of \( m \). According to Lemma 10 and Lemma 4, there is an integer \( m > m_0 \) such that on at least one input \( u \), the solver \( u^*_N \) fails to achieve \( \|Mu^*_N - Mu\|_2 < 1/2 \) with probability \( 1/2 \). However, this contradicts the previous conclusion, which implies that our starting assumption was incorrect.

A.3. Proof of Lemma 6

Let \( \tau = 8\sqrt{2}\gamma \), then let \( C^\tau_\triangleright \) and \( C^\tau_\blacktriangleleft \) represent the top \( \frac{n}{2} \) rows and the bottom \( \frac{n}{2} \) rows of \( C_\tau \), respectively. We first prove an auxiliary lemma about the random matrix \( C_\tau \).

**Lemma 11** There are universal positive constants \( c_1, c_2 \) such that, with probability at least \( 1 - c_1 \exp(-c_2n) \), the random matrix \( C_\tau \) has the following properties:

- It satisfies the RE condition (10) with parameter \( \gamma \).
- It satisfies the following upper bounds:

\[
\frac{\|C^\tau_\triangleright \theta\|_2^2}{n} \leq \frac{1}{3} \|\theta\|_2^2 \quad \text{and} \quad \frac{\|C^\tau_\blacktriangleleft \theta\|_2^2}{n} \leq (24\gamma)^2 \|\theta\|_2^2 \leq \frac{1}{2} \|\theta\|_2^2 \quad \text{for all } \theta \in \mathbb{B}_0(2k)
\]

(34)

**Proof** Dealing first with the upper block of the matrix \( C \), we have

\[
\|C^\tau_\blacktriangleleft \theta\|_2^2 = \frac{1}{4} \sum_{k=1}^{3k} \|A_k\theta\|_2^2.
\]

(35)

Recall that the diagonal blocks of \( A_k \) are matrices \( M \). The following claim provides an upper bound on its singular values:

**Claim 3** For any vector \( u \in \mathbb{R}^{4p} \), we have \( \|Mu\|_2 \leq 8p \|u\|_2^2 \).
**Proof** Since the entries of $M$ all belong to $\{-1, 0, 1\}$, if we let $M_i$ denote the $i^{th}$ row of $M$, then \[(M_i u)^2 \leq \left( \sum_{j \mid |M_{ij}| = 1} |u_j|^2 \right) \leq \left| \{ j \mid |M_{ij}| = 1 \} \right| \|u\|_2^2.\]

Summing over all indices $i = 1, 2, \ldots, m + 3p$, we find that \[\|Mu\|_2^2 \leq \left| \{(i, j) \mid |M_{ij}| = 1 \} \right| \|u\|_2^2.\]

Since there are $8p$ nonzero entries in $M$, the claim follows.

Returning to the main thread, Claim 3 implies that $A_k$ has singular values bounded by $\sqrt{8tp}$.

Putting together the pieces, we have \[\|A_k\|_2^2 \leq 8tp \|\theta\|_2^2 \leq 8k \|\theta\|_2^2,\] where the final inequality follows since $tp \leq t(m/3 + p) = k$. In conjunction, inequalities (35) and (36) imply that \[\left| \sum_{i,j} |M_{ij}| = 1 \right| \|u\|_2^2 \leq \left| \{(i,j) \mid |M_{ij}| = 1 \} \right| \|u\|_2^2.\]

Our next step is to prove that $\|C^t\|_2^2 \leq \frac{n}{2} \|\theta\|_2^2$ with high probability. Since $n \gtrsim k \log d$, standard results for Gaussian random matrices (see Lemma 13 in Appendix C) imply that \[\left| \sum_{i} |C_{ij}| = 1 \right| \|\theta\|_2^2 \leq \frac{n}{2} \|\theta\|_2^2\] with probability at least $1 - c_1 \exp(-c_2n)$, where inequality (i) follows since $\tau = 8\sqrt{2} \gamma \leq 1/3$.

Combining the upper bounds for $C^t$ and the upper bound for $C^t$, we find that the normalization conditions (34) hold with high probability.

It remains to show that the RE condition (10) holds with high probability. By Lemma 13, we have \[\|C_\tau\|_2^2 \geq \|C^t\|_2^2 \geq \frac{\left( \sqrt{n/2} \|\theta\|_2 \right)^2}{8} = \frac{n}{2} \|\theta\|_2^2,\] a bound that holds uniformly for all $\theta \in \bigcup_{|S| = k} S \subset \{1, \ldots, d\}$, as required by the condition (10). 

According to Lemma 11, there is a realization of the random matrix $R$ such that the matrix $C_\tau$ satisfies both the condition (34), and the RE condition (10) with parameter $\gamma$. We take this realization of $R$ to define a concrete realization of $C_\tau$. Consequently, the RE constant of this matrix, namely $\gamma(C_\tau)$, satisfies $\gamma(C_\tau) \geq \gamma$. For arbitrary matrices $X, Y \in \mathbb{R}^{n \times d}$, it is straightforward that \[|\gamma(X) - \gamma(Y)| \leq \|X - Y\|_{op},\] showing that the function $X \mapsto \gamma(X)$ is a Lipschitz function with parameter 1, and hence continuous. Consequently, $\gamma(C_x)$ is a continuous function of $x$. We also claim that it satisfies the condition \[\gamma(C_0) = 0,\] (39)
a claim proved at the end of this section. Based on the continuity property (38) and the initial condition (39), there is a constant \( \tau' \in (0, \tau] \) such that \( \gamma(C_{\tau'}) = \gamma \). Since \( C_{\tau} \) satisfies the normalization condition (34) and \( \tau' \leq \tau \), we have

\[
\frac{\|C_{\tau'} \theta\|^2}{n} \leq \frac{\|C_{\tau} \theta\|^2}{n} \leq \frac{5}{6} \|\theta\|^2 \quad \text{and} \quad \frac{\|C^\dagger_{\tau'} \theta\|^2}{n} \leq \frac{\|C^\dagger_{\tau} \theta\|^2}{n} \leq (24\gamma)^2 \|\theta\|^2 \quad \text{for all } \theta \in \mathbb{B}_0(2k).
\] (40)

Next we consider the quantized matrix \( \lfloor C_{\tau'} \rfloor_L \). Since the quantization operator approximates every entry of \( C_{\tau'} \) to precision \( 2^{-L} \), we have

\[
\|\lfloor C_{\tau'} \rfloor_L - C_{\tau'}\|_\text{op} \leq \|\lfloor C_{\tau'} \rfloor_L - C_{\tau'}\|_F \leq 2^{-L} \sqrt{nd}.
\] (41)

Combining this inequality with the Lipschitz condition (38), we find that

\[
|\gamma(\lfloor C_{\tau'} \rfloor_L) - \gamma| \leq 2^{-L} \sqrt{nd},
\]

showing that the quantized version satisfies the stated RE condition as long as \( L \geq \log(\sqrt{nd}/\varepsilon) \).

Turning to the normalization conditions, inequality (40) and inequality (41) imply that

\[
\frac{\|\lfloor C_{\tau'} \rfloor_L \theta\|^2}{\sqrt{n}} \leq \frac{\|C_{\tau'} \theta\|^2 + 2^{-L} \sqrt{nd} \|\theta\|^2}{\sqrt{n}} \leq (\sqrt{5/6} + 2^{-L} \sqrt{d}) \|\theta\|^2 \quad \text{and}
\]

\[
\frac{\|\lfloor C^\dagger_{\tau'} \rfloor_L \theta\|^2}{\sqrt{n}} \leq \frac{\|C^\dagger_{\tau'} \theta\|^2 + 2^{-L} \sqrt{nd} \|\theta\|^2}{\sqrt{n}} \leq (24\gamma + 2^{-L} \sqrt{d}) \|\theta\|^2 \quad \text{for all } \theta \in \mathbb{B}_0(2k).
\]

It is straightforward to verify that if \( L \geq \max\{\log(12\sqrt{d}), \log(\sqrt{nd}/\varepsilon)\} \), then the normalization conditions (16) stated in the lemma are also satisfied.

**Proof of equation (39):** Notice that \( \text{rank}(C_0) = \text{rank}(A_k) \leq 3k \), but the number of nonzero columns of \( C_0 \) is greater than \( 3k \). Hence, there is a \((3k + 1)\)-sparse nonzero vector \( \theta^+ \) such that

\[
C_0 \theta^+ = 0.
\] (42)

Let \( S \) be a set of indices that correspond to the \( k \) largest entries of \( \theta^+ \) in absolute value. Then \( |S| = k \) and \( \|\theta^+_S\|_1 \leq \frac{2k+1}{k} \|\theta^+_S\|_1 \). The later inequality implies that

\[
|S| = k \quad \text{and} \quad \theta^+ \in \mathbb{C}(S).
\] (43)

Combining expressions (42) and (43) yields that \( \gamma(C_0) = 0 \).

**A.4. Proof of Lemma 7**

Recall that the noise vector \( w \) follows the normal distribution \( N(0, \sigma^2 I_{n \times n}) \). Consequently, if we define the event \( \mathcal{E}_1 := \{|w|_\infty \leq \mu\} \), where

\[
\mu := \max \left\{ \frac{2n}{\sqrt{2\pi}}, \sqrt{2\log(2h(m))} \right\},
\]

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then standard tail bounds and the union bound imply that
\[ \mathbb{P}[E_1^c] \leq n \left\{ \frac{2}{\sqrt{2\pi \mu}} e^{-\mu^2/2} \right\} \leq \frac{1}{2h(m)}. \]

Let \( E_2 \) denote the event that \( u'_i \) being \((k/t)\)-sparse and that \( \|Mu'_i - Mu^*\|_2 < 1/2 \). We split the probability \( \mathbb{P}[E_2] \) into two terms. We find that
\[
\mathbb{P}[E_2] = \int_{\|z\| \leq \mu} \mathbb{P}[E_1|w = z] \phi(z; 0, \sigma^2 I_{n \times n}) dz + \mathbb{P}[E_2|E_1^c] \mathbb{P}[E_1^c]
\]
where \( \phi(z; 0, \sigma^2 I_{n \times n}) \) denotes the density of the \( N(0, \sigma^2 I_{n \times n}) \) distribution. Consequently, there must be some specific value \( \|w_0\|_\infty \leq \mu \) such that
\[
\mathbb{P}[E_2] \leq \mathbb{P}[E_2|w = w_0] + \frac{1}{2h(m)}. \tag{44}
\]

Let \( \hat{u}_i^* \) be the solver that uses the deterministic argument \( w_0 \) rather than drawing \( w \) randomly according to the specified procedure. The remaining steps followed by \( \hat{u}_i^* \) are exactly the same as those of \( u'_i \). We now demonstrate that the solver \( \hat{u}_i^* \) has polynomial complexity in terms of code length and running time:

**Code length:** We begin by encoding the design matrix \( X \) and fixed vectors \((\hat{\theta}^i, w_0)\) into the program of the solver \( \hat{u}_i^* \). Recall that both \( X \) and \( \hat{\theta}^i \) are discretized to \( 2^{-L} \) precision. For any discrete value \( x \) of \( 2^{-L} \) precision, encoding it takes at most \( \log(|x| + 1) + L \) bits. Recalling our choice (17) of \( L \), is straightforward to verify that \( L = \mathcal{O}(\log(\text{poly}(n, d, k)) + \log(1/\varepsilon)) \). Moreover, the components of the matrix \( X \) and vector \( \hat{\theta}^i \) are adequately controlled: in particular, by inequality (16), equation (17) and the definition of \( \hat{\theta}^i \), each such entry \( x \) has an encoding bounded by \( \log(|x| + 1) = \mathcal{O}(\text{poly}(n, d, k)) + \log(1/\varepsilon)) \), Since \( \log(1/\varepsilon) = G(n, d, k) \) is a polynomial functions of \( m \), we find that both \( X \) and \( \hat{\theta}^i \) have polynomial code length.

Now the noise vector \( w_0 \) is involved in computing the response vector
\[
y' = \begin{bmatrix} X^\top \hat{\xi}^i \\ X^\top \hat{\theta}^i \end{bmatrix} + w_0.
\]
Since the estimator \( \hat{\theta} \) takes \([y']_{G(n,d,k)} \) as input, it suffices to quantize \( w_0 \) sufficiently finely so as to ensure that the quantized response \([y']_{G(n,d,k)} \) is not altered by the quantization. Note that the components of \( X \), \( \hat{\theta}^i \) and \( \xi^i \) all have \( 2^{-L} \) precision. Consequently, their products have at most \( 2^{-2L} \) precision, and we have
\[
\begin{bmatrix} X^\top \hat{\xi}^i \\ X^\top \hat{\theta}^i \end{bmatrix} + w_0 \right]_G = \begin{bmatrix} X^\top \hat{\xi}^i \\ X^\top \hat{\theta}^i \end{bmatrix} + [w_0]_{\max\{2L,G\}}_G,
\]
showing that it suffices to quantize the noise \( w_0 \) at level \([\cdot]_{\max\{2L,G\}} \). On the other hand, by construction, the absolute values of \( w_0 \) are uniformly bounded by \( \mu \). Hence, the vector \( w_0 \) can also be encoded in polynomial code length.
Finally, we have assumed that code length of $\hat{\theta}$ is bounded by $b$, and observed that the program that constructs $y'$ is of constant code length. Consequently, the total code length is bounded by a polynomial function of $m$, which we denote by $f(m)$.

**Running time:** Since the response vector $y'$ can be constructed from $(M, Mu^*, X, \hat{\theta}^i, w_0)$ in polynomial time and the solver $\hat{\theta}$ terminates in polynomial time (bounded by $H(size(X, y'; G))$), we conclude that $\hat{u}^*_i$ also terminates in polynomial time, say upper bounded by the polynomial $g(m)$.

To complete the argument, we utilize the hardness result of Lemma 5. Indeed, this lemma guarantees that a polynomial-complexity solver cannot achieve probability of success greater than $1/h(m)$, which yields

$$\Pr[\mathcal{E}_2|w = w_0] \leq \frac{1}{2h(m)}.$$

Combining this upper bound with inequality (44) completes the proof.

**A.5. Proof of Lemma 8**

Let $q'$ and $q^\circ$ denote the probability density functions of $y'$ and $y^\circ$. With this notation, we have

$$\Pr[u'_i \in B_0(k/t)] = \int_{\mathbb{R}^n} \mathbb{I}[\hat{\theta}(X, z) \in B_0(k/t)] q'(z)dz$$

$$= \int_{\mathbb{R}^n} \frac{q'(z)}{q^\circ(z)} \mathbb{I}[\hat{\theta}(X, z) \in B_0(k/t)] q^\circ(z)dz$$

$$= \mathbb{E} \left[ \frac{q'(y^\circ)}{q^\circ(y^\circ)} \mathbb{I}[u^\circ_i \in B_0(k/t)] \right]$$

$$\geq e^{-1} \Pr[u^\circ_i \in B_0(k/t) \text{ and } \frac{q'(y^\circ)}{q^\circ(y^\circ)} \geq e^{-1}]$$

$$\geq e^{-1} \left\{ \Pr[u^\circ_i \in B_0(k/t)] - \Pr[\frac{q'(y^\circ)}{q^\circ(y^\circ)} < e^{-1}] \right\}.$$

We claim that

$$\Pr[\frac{q'(y^\circ)}{q^\circ(y^\circ)} < e^{-1}] \leq \frac{25\gamma r \sqrt{n}}{\sigma}.$$  \hspace{1cm} (45)

Deferring the proof of upper bound (45), we now focus on its consequence. If $\Pr[u^\circ_i \in B_0(k/t)] \geq \frac{t}{2k + 2t}$, then we are guaranteed that

$$\Pr[u'_i \in B_0(k/t)] \geq e^{-1} \left( \frac{t}{2k + 2t} - \frac{25\gamma r \sqrt{n}}{\sigma} \right).$$  \hspace{1cm} (46)

Finally, let $T \subseteq \{1, \ldots, t\}$ be the subset of indices for which $u^\circ_i \in B_0(k/t)$ with probability at least $\frac{t}{2k + 2t}$. Using $\hat{\theta}_i$ as a shorthand for $\hat{\theta}_i(X, y)$, note that $u^\circ_i$ and $\hat{\theta}_i/\rho$ share the same distribution, and hence

$$\Pr[u^\circ_i \in B_0(k/t)] = \Pr[\hat{\theta}_i \in B_0(k/t)].$$  \hspace{1cm} (47)
Since \( \hat{\theta} \in \mathbb{B}_0(k) \) by construction, we have

\[
  k \geq \|\hat{\theta}\|_0 \geq \sum_{i=1}^{t} \|\hat{\theta}_i\|_0 \geq \left( \frac{k}{t} + 1 \right) \sum_{i=1}^{t} I[\hat{\theta}_i \notin \mathbb{B}_0(k/t)] = \frac{k + t}{t} \left( t - \sum_{i=1}^{t} I[\hat{\theta}_i \in \mathbb{B}_0(k/t)] \right).
\]

(48)

Following some algebra, we find that \( \sum_{i=1}^{t} I[\hat{\theta}_i \in \mathbb{B}_0(k/t)] \geq \frac{t^2}{k + t} \), and hence

\[
  \sum_{i=1}^{t} \mathbb{P}[\hat{\theta}_i \in \mathbb{B}_0(k/t)] \geq \frac{t^2}{k + t}.
\]

Let \( N \) be the number of indices such that \( \mathbb{P}[\hat{\theta}_i \in \mathbb{B}_0(k/t)] \geq \frac{t}{2k + 2t} \). By definition, there are \( t - N \) indices for which \( \mathbb{P}[\hat{\theta}_i \in \mathbb{B}_0(k/t)] < \frac{t}{2k + 2t} \), and hence

\[
  \frac{t^2}{k + t} \leq \sum_{i=1}^{t} \mathbb{P}[\hat{\theta}_i \in \mathbb{B}_0(k/t)] \leq N + (t - N) \frac{t}{2k + 2t},
\]

which implies that \( N \geq \frac{t^2}{2k + t} \), as claimed. By the equivalence (47), we have \( N = |T| \), so that the proof is complete.

**Proof of inequality (45)** Recall that \( y' \) and \( y^o \) differ in distribution only in their last \( n/2 \) entries, denoted by \( y'_{\text{bot}} \) and \( y^o_{\text{bot}} \) respectively. Since \( y'_{\text{bot}} \sim N(X^i \hat{\theta}^i, \sigma^2 I) \) and \( y^o_{\text{bot}} \sim N(X^i \tilde{\xi}^i, \sigma^2 I) \), we have

\[
  \frac{q'(y^o)}{q'(y')} = \exp \left( - \frac{(X^i \hat{\theta}^i - y'_{\text{bot}})^2}{2\sigma^2} + \frac{(X^i \tilde{\xi}^i - y^o_{\text{bot}})^2}{2\sigma^2} \right).
\]

(49)

Define the scalar \( a := \|X^i(\tilde{\xi}^i - \hat{\theta}^i)\|_2/\sigma \), and define the event \( \mathcal{E} := \{ Z < \frac{a}{2} - \frac{1}{a} \} \) where \( Z \sim N(0, 1) \) is a standard normal variate. From the representation (49), some algebra shows that the probability that the upper bound \( \frac{q'(y^o)}{q'(y')} \leq e^{-1} \) holds is equal to the probability that the event \( \mathcal{E} \) holds. Letting \( \Phi \) denote the CDF of a standard normal variate, we have \( \mathbb{P}[\mathcal{E}] = \Phi \left( \frac{a}{2} - \frac{1}{a} \right) \). It can be verified that \( \Phi \left( \frac{a}{2} - \frac{1}{a} \right) \leq a \) for all \( a \geq 0 \), whence

\[
  \mathbb{P}[\mathcal{E}] \leq a = \frac{\|X^i(\tilde{\xi}^i - \hat{\theta}^i)\|_2}{\sigma} \leq \frac{25\gamma \sqrt{n} \|\xi^i - \hat{\theta}^i\|_2}{\sigma},
\]

where the last step uses inequality (16). Since \( \tilde{\xi}^i \) and \( \hat{\theta}^i \) differ only in one subvector—more precisely, by \( \rho u^* \)—we have

\[
  \mathbb{P}[\mathcal{E}] \leq \frac{25\gamma \sqrt{n} \|\rho u^*\|_2}{\sigma} \leq \frac{25\gamma r \sqrt{n}}{\sigma},
\]

where the last step follows since \( \|\rho u^*\|_2 \leq r \).
A.6. Proof of Lemma 9

Letting $q'$ and $q^\diamond$ denote the probability density functions of $y'$ and $y^\diamond$, respectively, we have

$$
E \left[ \|M u'_i - M u^*\|_2^2 \right] = \int_{\mathbb{R}^n} \|M \hat{\theta}_i(X, z) - M u^*\|_2^2 q^\diamond(z) dz
$$

$$
= \int_{\mathbb{R}^n} \frac{q^\diamond(z)}{q'(z)} \|M \hat{\theta}_i(X, z) - M u^*\|_2^2 q'(z) dz
$$

$$
= E \left[ \frac{q^\diamond(y')}{q'(y')} \|M u'_i - M u^*\|^2 \right]
$$

$$
\geq e^{-1} E \left[ \|M u'_i - M u^*\|_2^2 \mid \frac{q^\diamond(y')}{q'(y')} \geq e^{-1} \right] P \left[ \frac{q^\diamond(y')}{q'(y')} \geq e^{-1} \right]. \quad (50)
$$

Following precisely the proof of inequality (45), we can prove the upper bound

$$
P \left[ \frac{q^\diamond(y')}{q'(y')} < e^{-1} \right] \leq \frac{25 \gamma r \sqrt{n}}{\sigma}.
$$

As a consequence, we have

$$
P \left[ \frac{q^\diamond(y')}{q'(y')} \geq e^{-1} \right] \geq 1 - \frac{25 \gamma r \sqrt{n}}{\sigma} \geq \frac{1}{2},
$$

where the last inequality follows since $\frac{\gamma r \sqrt{n}}{\sigma} \leq \frac{1}{50}$ by assumption. Plugging this lower bound into inequality (50) yields

$$
E \left[ \|M u'_i - M u^*\|_2^2 \mid \frac{q^\diamond(y')}{q'(y')} \geq e^{-1} \right] \leq 2e E \left[ \|M u'_i - M u^*\|_2^2 \right].
$$

Applying Markov’s inequality yields

$$
P \left[ \|M u'_i - M u^*\|_2^2 \leq \frac{\sigma}{25 \gamma r \sqrt{n}} \right] \cdot 2e E \left[ \|M u'_i - M u^*\|_2^2 \mid \frac{q^\diamond(y')}{q'(y')} \geq e^{-1} \right] \geq 1 - \frac{25 \gamma r \sqrt{n}}{\sigma}.
$$

Consequently, with probability $(1 - \frac{25 \gamma r \sqrt{n}}{\sigma})^2 \geq 1 - \frac{50 \gamma r \sqrt{n}}{\sigma}$, we have

$$
\|M u'_i - M u^*\|_2^2 \leq \frac{2e \sigma}{25 \gamma r \sqrt{n}} E \left[ \|M u'_i - M u^*\|_2^2 \right],
$$

as claimed.

A.7. Proof of Lemma 10

Consider a problem $\mathcal{P}$ for which, for any integer $m > m_0$, there is a Turing machine $\mathcal{T}_m$ such that:

(a) the code length of $\mathcal{T}_m$ is at most $f(m)$ bits, and;

(b) it solves every input of length $m$ with probability at least $1/2$, and terminates in time at most $g(m)$.

Under Assumption A, it suffices to show that these two conditions certify that $\mathcal{P} \in \mathbf{P}/\text{poly}$. 

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In order to prove this claim, let \( m > m_0 \) be arbitrary. Given a sufficiently long random sequence, we may execute \( T_m \) a total of \( m+1 \) times, with the randomness in each execution being independent. By property (b), for any binary input of length \( m \), each execution has success probability at least \( 1/2 \). By independence, the probability that at least one of the \( m + 1 \) executions is successful is \( 1 - 2^{-m-1} \). Since the total number of \( m \)-length inputs is at most \( 2^m \), the union bound implies that the \( m + 1 \) executions of \( T_m \) succeed on all inputs with probability at least \( 1 - 2^{-m-1}2^m = 1/2 > 0 \). Consequently, the probabilistic method implies that there exists some realization of the random sequence—call it \( R_m \)—under which \( m + 1 \) executions of \( T_m \) solves the problem for all inputs of length \( m \).

Let \( C_m \) be the code defining machine \( T_m \), and consider a Turing machine \( T'_m \) that takes the pair \((C_m, R_m)\) as an advice string, and then uses the string \( R_m \) to simulate the execution of \( T_m \) a total of \( m+1 \) times. From our previous argument, the Turing machine \( T'_m \) solves the problem on every input of length \( m > m_0 \). Notice that the length of string \((C_m, R_m)\) is of the order \( O(f(m) + mg(m)) \) bits, and the running time of \( T'_m \) is of the order \( O(mg(m)) \). Finally, for all input lengths \( m \leq m_0 \), a constant-size Turing machine \( T'_m \) can solve the problem in constant time. The combination of \( T'_m \) and \( T_m \) provides a Turing machine that certifies \( P \in \text{P/poly} \).

**Appendix B. Proof of Proposition 2**

We begin by defining the error vectors \( \Delta := \hat{\theta}_{\lambda_n} - \theta^* \) of the ordinary Lasso, as well as that \( \Delta := \hat{\theta}_n - \theta^* \) of the thresholded Lasso. By construction, the error vector \( \Delta \) is at most \( 2k \)-sparse, so that the normalization condition (11) guarantees that

\[
\frac{\|X\Delta\|^2}{n} \leq \|\Delta\|^2 \tag{51}
\]

The following lemma shows that the truncated Lasso is essentially as good as the non-truncated Lasso:

**Lemma 12** The error of the truncated Lasso is bounded as \( \|\Delta\|^2 \leq 5\|\Delta\|^2 \).

We return to prove this intermediate lemma at the end of this section.

Taking Lemma 12 as given for the moment, we complete the proof by bounding the error of the ordinary Lasso using Corollary 2 in (Negahban et al., 2012). With the specified choice of \( \lambda_n \), this corollary implies that

\[
\|\Delta\|^2 \leq \frac{64}{\gamma^2(X)} \frac{\sigma^2 k \log d}{n} \tag{52}
\]

with probability at least \( 1 - 2e^{-c_k \log d} \). Combining the bound (52), Lemma 12 and inequality (51) yields the claim of the proposition.

Finally, we return to prove Lemma 12. Throughout this proof, we use \( \hat{\theta} \) as a shorthand for the thresholded estimator \( \hat{\theta}_n \). Since \( \hat{\theta} \) is the \( k \)-truncated version of \( \hat{\theta}_{\lambda_n} \), if \( |\text{supp}(\hat{\theta}_{\lambda_n})| \leq k \), we must have \( \hat{\theta} = \hat{\theta}_{\lambda_n} \), in which case the proof is complete. Otherwise, we may assume that \( |\text{supp}(\hat{\theta}_{\lambda_n})| > k \), and then define the set \( G := \{ \text{supp}(\theta^*) \cap \text{supp}(\hat{\theta}_{\lambda_n}) \} \setminus \text{supp}(\hat{\theta}) \), corresponding to “good” entries in the estimate \( \hat{\theta}_{\lambda_n} \) that were lost in the truncation. Introduce the notation \( G = \{i_1, \ldots, i_t \} \) where
In this section, we provide some background on the singular value of Gaussian random matrices, Appendix C. Singular values of random matrices

Combining the two cases completes the proof of inequality (55b).

Now consider the decomposition
\[
\|\hat{\theta} - \theta^*\|_2^2 = \sum_{i=1}^{d} (\hat{\theta}_i - \theta_i^*)^2 = \sum_{s=1}^{t} \left( (\hat{\theta}_{is} - \theta_i^*)^2 + (\hat{\theta}_{js} - \theta_j^*)^2 \right) + \sum_{i \notin B \cup G} (\hat{\theta}_i - \theta_i^*)^2. \tag{54}
\]

The proof will be complete if we can establish the two inequalities

\[
(\hat{\theta}_i - \theta_i^*)^2 \leq ((\hat{\theta}_{\lambda_n})_i - \theta_i^*)^2 \quad \text{if } i \notin B \cup G, \quad \text{and} \quad (\hat{\theta}_{is} - \theta_i^*)^2 + (\hat{\theta}_{js} - \theta_j^*)^2 \leq 5((\hat{\theta}_{\lambda_n})_i - \theta_i^*)^2 + 5((\hat{\theta}_{\lambda_n})_j - \theta_j^*)^2 \quad \text{for } s = 1, \ldots, t. \tag{55a}
\]

Inequality (55a) is straightforward, since since \(i \notin B \cup G\) implies either \(\hat{\theta}_i = (\hat{\theta}_{\lambda_n})_i\) or \(\hat{\theta}_i = \theta_i^* = 0\). Turning to the inequality (55b), note that \(\hat{\theta}_{is} = \theta_j^* = 0\), so that it is equivalent to upper bound \((\theta_{is}^*)^2 + (\hat{\theta}_{\lambda_n})_j^2\). We divide the proof of inequality (55b) into two cases:

**Case 1:** First, suppose that \(|(\hat{\theta}_{\lambda_n})_i - \theta_i^*| \leq \frac{1}{2} |\theta_i^*|\). In this case, we have \(|(\hat{\theta}_{\lambda_n})_{is}| \geq \frac{1}{2} |\theta_i^*|\), and hence by inequality (53), we obtain \(|(\hat{\theta}_{\lambda_n})_{js}| \geq \frac{1}{2} |\theta_j^*|\). Combining the pieces yields
\[
(\theta_{is}^*)^2 + ((\hat{\theta}_{\lambda_n})_j)^2 \leq 5((\hat{\theta}_{\lambda_n})_i)^2 \leq 5((\hat{\theta}_{\lambda_n})_i - \theta_i^*)^2 + 5((\hat{\theta}_{\lambda_n})_j - \theta_j^*)^2.
\]

**Case 2:** On the other hand, if \(|(\hat{\theta}_{\lambda_n})_{is} - \theta_i^*| > \frac{1}{2} |\theta_i^*|\), then we have
\[
(\hat{\theta}_{is} - \theta_i^*)^2 + (\hat{\theta}_{js} - \theta_j^*)^2 = (\theta_i^*)^2 + ((\hat{\theta}_{\lambda_n})_j)^2
\leq 4((\hat{\theta}_{\lambda_n})_i - \theta_i^*)^2 + ((\hat{\theta}_{\lambda_n})_j)^2
\leq 5((\hat{\theta}_{\lambda_n})_i - \theta_i^*)^2 + 5((\hat{\theta}_{\lambda_n})_j - \theta_j^*)^2.
\]

Combining the two cases completes the proof of inequality (55b).

**Appendix C. Singular values of random matrices**

In this section, we provide some background on the singular value of Gaussian random matrices, required in the proof of Theorem 3(a). Our results apply to a random matrix \(A \in \mathbb{R}^{n \times d}\) formed of i.i.d. \(N(0, 1)\) entries. Recall the set \(C(S) = \{ \theta \in \mathbb{R}^d \mid \|\theta_S\|_1 \leq 3\|\theta_S\|_1 \}\) that is involved in the RE condition (10).
Lemma 13 Suppose that \( n > c_0 k \log d \) for a sufficiently large constant \( c_0 \). Then there are universal constants \( c_j, j = 1, 2 \) such that
\[
\frac{\|A\theta\|_2}{\sqrt{n}} \leq 3\|\theta\|_2 \quad \text{for all } \theta \in \mathbb{B}_0(2k), \quad \text{and} \tag{56a}
\]
\[
\frac{\|A\theta\|_2}{\sqrt{n}} \geq \frac{\|\theta\|_2}{8} \quad \text{for all } \theta \in \bigcup_{S \subset \{1, \ldots, d\} \atop |S| = k} \mathbb{C}(S), \quad \text{and} \tag{56b}
\]
where both bounds hold with probability at least \( 1 - c_1 \exp(-c_2 n) \).

Proof To prove the upper bound (56a), it suffices to show that
\[
\max_{|S| = 2k} \frac{\|A_T\|_{\text{op}}}{\sqrt{n}} \leq 3,
\]
where \( A_T \) denotes the columns of \( A \) whose indices are in set \( T \), and \( \| \cdot \|_{\text{op}} \) denotes the maximum singular value of a matrix. By known bounds on singular values of Gaussian random matrices (Rudelson and Vershynin, 2010), for any subset \( T \) of cardinality \( 2k \), we have
\[
P\left[ \frac{\|A_T\|_{\text{op}}}{\sqrt{n}} \geq \sqrt{n} + \sqrt{2k} + \delta \right] \leq 2 \exp(-\delta^2/2) \quad \text{for all } \delta > 0.
\]
If we let \( \delta = \sqrt{n} \) and use the assumption that \( 2k \leq n \), the above bound becomes
\[
P\left[ \|A_T\|_{\text{op}} \geq 3\sqrt{n} \right] \leq 2 \exp(-n/2).
\]
Noting that there are \( \binom{d}{2k} \leq e^{2k \log(\frac{ed}{2k})} \) subsets of cardinality \( 2k \), the claim thus follows by union bound.

On the other hand, the lower bound (56b) is implied by the main result of Raskutti et al. (2010). \( \blacksquare \)