

Efficient Robust Constrained Signal Detection via Kolmogorov Width Approximations

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Abstract

Robust statistical inference often faces a severe computational-statistical gap when dealing with complex parameter spaces. We investigate minimax signal detection in the Gaussian sequence model under strong ϵ -contamination, where the signal belongs to a general prior constraint K . Existing optimal tests require computing the exact Kolmogorov k -width of K , a computationally intractable task for general non-trivial sets. We bridge this gap by proposing a polynomial-time testing framework that universally applies to balanced, type-2, and exactly 2-convex constraints. By leveraging a semidefinite programming relaxation and a modified ellipsoid method equipped with an approximate subgradient oracle, we efficiently approximate the Kolmogorov widths. Remarkably, our unconditional efficient algorithm achieves a robust detection boundary that matches existing upper bounds up to a mere polylogarithmic factor. This establishes a computationally tractable testing solution for a broad class of structured signals without requiring prior knowledge of their exact geometric complexity.

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1 Introduction

Robust Statistics

A central challenge in the application of modern statistical estimation and hypothesis testing lies in specifying the true underlying data distribution. Traditionally, statistical frameworks operate under the idealized assumption that observations are sampled from a well-behaved parametric family, such as the Gaussian distribution. However, this assumption is frequently violated in real-world applications. These violations typically arise from experimental measurement inaccuracies (e.g., Leek et al. [2010]), heavy-tailed phenomena inherent in quantitative finance (e.g., Cont [2001]), or deterministic adversarial outliers, such as data poisoning in machine learning (e.g., Steinhardt et al. [2017], Yin et al. [2018]). In the presence of such anomalies, even a small fraction of corrupted data can catastrophically compromise standard inference procedures. To illustrate, consider the fundamental task of estimating the population mean from a sample x_1, x_2, \dots, x_n drawn from $\mathcal{N}(\mu, 1)$. While the sufficient statistic $\bar{x} = \frac{1}{n} \sum_{i=1}^n x_i$ is minimax optimal in both the estimation problem (under the ℓ_2 norm loss) and the hypothesis testing problem (for prescribed Type I and Type II error), an adversarial perturbation to just a single observation in $\{x_i\}_{i=1}^n$ can arbitrarily skew \bar{x} (e.g., forcing it to zero), effectively nullifying its reliability. Conversely, robust alternatives such as the sample median of $\{x_i\}_{i=1}^n$ exhibit significantly greater resilience against such data modifications. Recognizing such vulnerabilities, statisticians have long emphasized the critical role of robustness in methodological development. One of the foundational works formalizing this field was presented by Huber [1964], who introduced the *Huber contamination model* (also known as the ϵ -contamination model). In this model, rather than observing samples exclusively from a presumed nominal distribution \mathbb{P} , one assumes the data generation process is corrupted by an unknown, adversarially chosen distribution \mathbb{Q} . Consequently, the observations are drawn from the mixture distribution $(1 - \epsilon)\mathbb{P} + \epsilon\mathbb{Q}$ for some contamination proportion $\epsilon \in (0, 1)$. Crucially, the adversary’s choice of \mathbb{Q} is permitted to depend on \mathbb{P} . To achieve robust estimation under this framework, Huber also introduced the Huber loss—a piecewise function combining quadratic and linear penalties—as a resilient estimation criterion. An even more formidable adversary operates within the *strong ϵ -contamination model*. Under this framework, for a given contamination fraction $\epsilon \in (0, 1)$, the adversary is permitted to inspect the original uncorrupted samples and replace up to an ϵ fraction of them with arbitrary malicious values. This replacement process is highly adaptive, as the adversary possesses full knowledge of the original sample realization, the underlying data-generating mechanism, and even the downstream inference procedures to be deployed. This full adaptivity

renders robust inference significantly more challenging than under the standard Huber contamination model. The strong ϵ -contamination model serves as the primary focus of this work (see Definition 1.1 for a formal formulation). Also see [Diakonikolas and Kane \[2019\]](#) for a comprehensive overview of various contamination models. Alternatively, even in the absence of an explicit adversary, the true underlying distribution may inherently deviate from standard parametric assumptions. A prominent example is the heavy-tailed distribution, where the tail probability decays polynomially rather than exponentially (as is characteristic of Gaussian or exponential distributions). In such scenarios, one typically relies on milder moment assumptions. A substantial body of literature has been developed within this framework. For instance, [Catoni \[2012\]](#) demonstrates that it is possible to recover sub-Gaussian concentration bounds for heavy-tailed distributions with only finite variance. Furthermore, [Lugosi and Mendelson \[2019\]](#) extend this approach to high-dimensional heavy-tailed settings. For related work on the estimation of covariance matrices under similar conditions, see [Minsker \[2018\]](#). Below we introduce the concepts of Kolmogorov k -width and the ellipsoid method algorithm, which play central role in the derivation of the approximation framework in this work.

Kolmogorov k -Width

The Kolmogorov k -width (often referred to as the Kolmogorov n -width) quantifies the optimal degree to which a subset of a Banach space can be approximated by a k -dimensional subspace under a specified norm (see Definition 2.1). Originally introduced by [Kolmogoroff \[1936\]](#), it has emerged as a foundational concept in approximation theory. The analytical strength of this measure lies in its generality: rather than evaluating a specific projection, it takes an infimum over all possible k -dimensional subspaces, thereby establishing a fundamental information-theoretic limit. As a result, it serves as an indispensable tool for deriving minimax error bounds for linear estimators, particularly when approximating functions subject to specified smoothness conditions ([Pinkus \[2012\]](#), [Tsybakov \[2009\]](#)). This paradigm later inspired the development of the Gelfand k -width ([Tikhomirov \[1960\]](#)), a metric that now plays a pivotal role in the theory of compressed sensing ([Donoho \[2006\]](#)). Beyond classical statistics, the Kolmogorov k -width is extensively utilized in the numerical analysis of partial differential equations (PDEs), where a rapid decay rate of the width theoretically guarantees the existence of highly accurate approximate solutions ([Cohen et al. \[2010\]](#)). For modern extensions integrating this concept with neural networks and optimal transport in PDE settings, see the recent works by [Papapicco et al. \[2022\]](#) and [Florian Arbes and Urban \[2025\]](#).

Ellipsoid Method

Independently pioneered by [Shor \[1977\]](#) and [Yudin and Nemirovski \[1976\]](#), the celebrated ellipsoid method represents a cornerstone in convex optimization. Algorithmically, it operates by querying a separation oracle at each iteration to prune a region mathematically guaranteed to exclude the global optimizer, thereby producing a sequence of enclosing ellipsoids with exponentially decaying volumes. The method achieved historic prominence when [Khachiyan \[1979\]](#) leveraged it to establish the first polynomial-time complexity bound for linear programming. A profound theoretical consequence of this framework is the polynomial-time equivalence of the separation and optimization problems ([Grötschel et al. \[1981\]](#)). From the view of information-based complexity, the ellipsoid algorithm tightly characterizes the fundamental minimax lower bounds on oracle complexity for general non-smooth convex optimization ([Nemirovski and Yudin \[1983\]](#)). While it has largely been supplanted by interior-point methods in practical implementations due to their superior empirical speed, the

ellipsoid method remains important in theoretical computer science for resolving the computational complexity and theoretical tractability of algorithmic problems.

1.1 Problem Formulation and Our Contribution

This paper investigates the minimax ℓ_2 -norm signal detection problem in the Gaussian sequence model. Departing from classical unconstrained paradigms, we impose a prior structural constraint on the mean vector, restricting it to a known set K , i.e., $\mu \in K$. More broadly, we formulate this problem within a robust framework, wherein the pristine observations are subject to adversarial corruption. Namely, we assume the samples can be arbitrarily modified by an adversary \mathcal{C} operating under the strong ϵ -contamination model, with a known upper bound ϵ on the contamination fraction. Specifically, assume $\tilde{Y}_i = \mu + \sigma \cdot \xi_i \in \mathbb{R}^d, i = 1, 2, \dots, N$, where $\sigma > 0$ is the noise scale and $\xi_i \sim \mathcal{N}(0, \mathbf{I}_d)$ are standard multivariate Gaussian random variables. Denote the observations (contaminated samples) as $\mathbf{Y} := \mathcal{C}(\tilde{\mathbf{Y}}) = \left\{ \mathcal{C}(\tilde{Y}_1), \dots, \mathcal{C}(\tilde{Y}_N) \right\}$. The goal of the problem is to test the following hypotheses with uniform small Type I and Type II errors over $\mu \in K$ and \mathcal{C} based on \mathbf{Y} :

$$\begin{aligned} H_0 &: \mu = 0, \\ H_1 &: \|\mu\|_2 \geq \rho, \mu \in K. \end{aligned} \tag{1.1}$$

From an information-theoretic perspective, if the signal μ is too close to the origin (i.e., the separation ρ is insufficiently large), H_0 and H_1 are statistically indistinguishable, even in the pristine, uncontaminated setting. Consequently, our goal is not to solve (1.1) for all possible $\rho > 0$, but to characterize the fundamental minimal separation ρ that renders the testing problem theoretically possible and computationally tractable. We investigate this fundamental limit through the minimax framework. Given a desired Type I error bound α (e.g., $\alpha = 0.05$), we introduce the family of valid level- α tests, which uniformly control the false positive rate over all arbitrary adversaries \mathcal{C} :

$$A_s(N, d, K, \epsilon, \alpha, \sigma) := \left\{ \phi : \sup_{\mathcal{C}} \mathbb{P}_0 \left(\phi(\mathcal{C}(\tilde{\mathbf{Y}})) = 1 \right) \leq \alpha \right\}, \tag{1.2}$$

where \mathbb{P}_0 represents the distribution of the observations $\tilde{\mathbf{Y}}$ under the null hypothesis H_0 . When there is no ambiguity, we suppress the explicit parameter dependence and simply write A_s . With the robust level- α tests established, we define the critical minimax separation rate $\rho(N, d, K, \epsilon, \alpha, \sigma)$ by

$$\rho(N, d, K, \epsilon, \alpha, \sigma) := \inf_{\rho} \left\{ \rho : \rho > 0, \sup_{\|\mu\|_2 \geq \rho, \mu \in K} \sup_{\mathcal{C}} \mathbb{P}_{\mu}(\phi(\mathcal{C}(\mathbf{Y})) = 0) \leq \alpha, \phi \in A_s \right\}. \tag{1.3}$$

This rate characterizes the fundamental limit of testability: it is the smallest radius $\rho > 0$ ensuring that there exists a robust test controlling the Type II error by α , regardless of the true signal $\mu \in K$ and the adversarial strategy \mathcal{C} . When no ambiguity arises, we simply denote this minimax rate by ρ_{critical} . This paradigm, which involves fixing the Type I error and subsequently minimizing the Type II error, is known as the Neyman-Pearson framework (Neyman and Pearson [1933]).

In this work, we exclusively consider the [strong \$\epsilon\$ -contamination model](#), which is formally introduced below.

Definition 1.1 (Strong ϵ -contamination model). Consider a set of N uncorrupted observations $\tilde{\mathbf{Y}} = \{\tilde{Y}_1, \dots, \tilde{Y}_N\}$ generated independently from a true underlying distribution \mathcal{P} under either H_0 or H_1 . For a known contamination fraction $\epsilon \in [0, \frac{1}{2})$, an omniscient adversary \mathcal{C} with full knowledge of \mathcal{P} is permitted to inspect the exact realization of $\tilde{\mathbf{Y}}$. The adversary may then maliciously substitute up to $\lfloor \epsilon N \rfloor$ pristine samples with arbitrary values. The resulting corrupted dataset $\mathbf{Y} = \mathcal{C}(\tilde{\mathbf{Y}})$ is provided to the statistician while the corruption index set C is unknown.

While the process of contamination is very common, as mentioned, \mathcal{C} from the strong ϵ -contamination model represents the most powerful agent for the reason that it has the full knowledge of not only the original data $\tilde{\mathbf{Y}}$, but also the data generation mechanism $\mathbb{P}_0, \mathbb{P}_\mu$ and even the downstream developed algorithms. Consequently, it leads to the strongest minimax lower bounds for the testing problem (1.1).

We note here that for the testing problem formulated above, when the constraint K is a quadratically convex orthosymmetric (QCO) set (see Definition 2.9 for a formal definition), the optimal minimax rate has been established by Li and Neykov [2026], where the information-theoretic lower bound is matched by a theoretical algorithm and an efficient algorithm except for poly-logarithmic factors in $N, d, \frac{1}{\alpha}, \frac{1}{\epsilon}$ (see Section 4 and Appendix A for a comprehensive review of their results). While the efficient algorithm enjoys a polynomial-time complexity, the implementation is conditional — the efficiency depends on the computability of the Kolmogorov k -width. Though the computation is feasible and efficient for some common constraint sets such as the whole space \mathbb{R}^d , axes-aligned hyperrectangles $H_d(a_1, \dots, a_d)$, ellipsoids $E_d(a_1, \dots, a_d)$ and some other QCO sets with explicit expressions, the general computability is not guaranteed efficient or even not permissible (see the discussions in Section 2.2). Therefore, this work aims to extend their results from the computational aspect. Specifically, we show that it is possible to relax the definition of the Kolmogorov k -width, and consequently reduce the computation of the Kolmogorov widths to another optimization problem where approximate and efficient solution exists. The constraints that qualify the approximation above is even beyond the QCO sets. We also show that such relaxation and approximation do not harm the original upper bounds significantly in the sense that it only introduces additional poly-logarithmic factors. In this view, the unconditional efficient algorithms for the testing problem (1.1) exist for a wide range of constraints.

1.2 Related Work

We review some related works in this section. The minimax testing problem in Gaussian models without adversary has been intensively studied since last century. Some pioneer works include Ingster [1982], where Ingster studies and establishes one of the first results in the nonparametric signal detection problem for the unconstrained case; for the constrained case, Ingster [1993a] studies the testing problem with general ellipsoidal constraints; Suslina [1996] studies the problem under the constraint of a ℓ_q ellipsoid with a ℓ_p ball removed; Baraud [2002] studies such testing problem for $\mu \in \mathbb{R}^\infty$ with sparse constraint, and the minimax rate was also extended to the ellipsoidal constraint. For the counterpart estimation problem, Donoho et al. [1990] studies the problem with the hyperrectangle constraints in the minimax sense and proved that hyperrectangles show certain least favorable property. For recent advances in the estimation problem, Prasad and Neykov [2025] establish the optimal minimax rate of estimation for star-shape constraint. Neykov [2026] establishes a computationally tractable framework for the robust estimation. While their approach similarly capitalizes on Kolmogorov width approximations, it relies on a distinctly different optimization

machinery — one that would entail an intractable exponential-time complexity if naively adapted to our hypothesis testing framework. Finally, see [Ingster \[1993a,b,c\]](#) for the significant difference between the testing and estimation problems.

For the robust statistics, since the fundamental work by [Huber \[1964\]](#), various works under such framework emerged. Some recent advances include [Chen et al. \[2015\]](#), where the authors establish a general decision theory under the Huber contamination model and applied such theory to the high-dimensional covariance matrix using the Scheffé estimate, and [Du et al. \[2018\]](#) where the authors construct a local binning median for the robust nonparametric regression problem and proved its minimax optimality over the Hölder function class with smoothness parameters smaller or equal to 1. The works related to the strong ϵ -contamination model ([Diakonikolas and Kane \[2019\]](#)) thrive after 2010s. Some important works include [Narayanan \[2022\]](#), [Canonne et al. \[2020, 2023\]](#), where the optimal minimax rates of the mean testing problem under oblivious adversaries (including the Huber contamination model) and adaptive adversaries (including the [strong contamination model](#)) are established in the unconstrained settings. A recent work by [Li and Neykov \[2026\]](#), as mentioned, establishes the optimal minimax rate for the robust testing problem (1.1) under the strong ϵ -contamination model when the constraint K is a QCO set, which includes hyperrectangles and general ellipsoids as special cases.

1.3 Notation

In this work, we use the following notations for consistency. We assume the constraint $K \subset \mathcal{X} = \mathbb{R}^d$ without special notice. Let $\tilde{\mathbf{Y}} := \{\tilde{Y}_1, \tilde{Y}_2, \dots, \tilde{Y}_N\}$ denote the original authentic samples, where N is the sample size. \mathcal{C} is the adversary from the [strong \$\epsilon\$ -contamination model](#) and $\mathcal{C}(\tilde{\mathbf{Y}})$ is the observations contaminated by \mathcal{C} . Let C be the index set on which the samples are contaminated. $[n]$ is used to represent the set $\{1, 2, \dots, n\}$. When $I \subset [N]$ is a index set, let \mathbf{Y}_I (or $\tilde{\mathbf{Y}}_I$) denote the submatrix of \mathbf{Y} (or $\tilde{\mathbf{Y}}$) where the rows are selected according to I . By our assumptions on the contamination process, we know $|C| \leq \epsilon N$ and we have $\tilde{\mathbf{Y}}_{[N] \setminus C} = \mathbf{Y}_{[N] \setminus C}$, where $|S|$ is the cardinality of S when S is a set. ρ denotes the ℓ_2 norm of the mean vector μ under the alternative. We use $\langle \cdot, \cdot \rangle$ and $(\cdot)^\top (\cdot)$ alternatively to denote the inner product in a Hilbert space, probably with a subscript. The common notations $\mathcal{O}, \Theta, \Omega$ are used to represent the relative asymptotic orders with a possible subscript indicating the related quantity. The lowercases c, c_1, c_2, \dots represent the constants that might vary from line to line.

1.4 Organization

This paper is structured as follows. Section 2 formally introduces the concepts of the Kolmogorov k -width, optimal dimensions and optimal projections. This section also introduces the relaxed optimization problem related to the Kolmogorov k -width for the sake of the approximation later. A significant result related to the approximate solution to the constrained quadratic maximization by [Bhattiprolu et al. \[2021\]](#), which is a crucial tool for this work, is also reviewed. Section 3 introduces the modified ellipsoid method leveraged in the approximation to the Kolmogorov k -width mentioned above. It also contains the convergence analysis for the theoretical guarantee of the efficiency. Section 4 is the core of the main text, where it introduces the unconditional efficient testing procedures for the problem (1.1) based on the results from Section 2 and 3 and the work by [Li and Neykov \[2026\]](#). Finally, Section 5 contains the relevant discussions and the possible directions for future works. We place the theoretical algorithm and proof of the corresponding upper bound in

Appendix A for readers’ reference. Some probabilistic preliminaries and the proofs of the lemmas, theorems, and corollaries that are not suitable for the main text are deferred to Appendix B.

2 Linear Approximations

2.1 Kolmogorov k -Width and Optimal Dimensions

Definition 2.1 below formally introduces the Kolmogorov k -width leveragd in this work.

Definition 2.1 (Kolmogorov k -width). Let \mathcal{X} be a Banach space equipped with the norm $\|\cdot\|$, and $K \subset \mathcal{X}$ is any subset. The k -dimensional Kolmogorov width is defined as

$$D_k(K) = \inf_{P \in \mathcal{P}_k} \sup_{\theta \in K} \|\theta - P\theta\|, \quad (2.1)$$

where \mathcal{P}_k is the set of all projection operators that project a vector onto a subspace in \mathcal{X} with intrinsic dimension k .

Remark. For a fixed $\theta \in \mathcal{X}$, the “projection operator” in the definition above should be understood as finding the point $P\theta \in K$ such that the the distance between θ and $P\theta$ is (nearly) minimal. When the norm in (2.1) is the ℓ_2 norm, it can be characterized by some projection matrix.

Throughout this work, we operate in the Euclidean space $\mathcal{X} = \mathbb{R}^d$ and define the Kolmogorov width with respect to the ℓ_2 -norm, in line with standard practice. We also restrict the candidate set \mathcal{P}_k to orthogonal projections. We further note that in the ℓ_2 norm separation, the Kolmogorov width is defined originally as in Definition 2.1. However, in the case of ℓ_p norm separation, such definition should be discretized, i.e., only consider the projections that are aligned with the axes. This topic is beyond the scope of this work, and we kindly refer the reader to Section 4 of Li and Neykov [2026] for more details.

By definition, the Kolmogorov width $D_k(K)$ is monotonically non-increasing with respect to the subspace dimension k . A crucial determinant for characterizing the geometric and statistical properties of the subset K within the ambient space \mathcal{X} is the exact rate of this decay. Intuitively, a rapid decay rate implies that K can be efficiently approximated by low-dimensional subspaces, reflecting an intrinsically lower structural complexity. For instance, the class of analytic functions on a compact interval exhibits an exponential decay rate, bounded by $D_k(K) \asymp e^{-ck}$ for some absolute constant $c > 0$. Conversely, a slow decay rate indicates that the set resists efficient linear approximation. A canonical example of this is the Sobolev space W_p^r , whose Kolmogorov width decays polynomially as $D_k(K) \asymp k^{-\alpha}$, where the exponent α is governed by the smoothness parameter of the space. For a comprehensive exposition on the exact calculation of Kolmogorov widths for various reproducing kernel Hilbert spaces and Sobolev spaces, we refer the reader to Pinkus [2012].

In the context of statistical inference and hypothesis testing, however, the intrinsic geometric properties of the subset K and the whole space \mathcal{X} are of secondary interest compared to their interaction with the true data-generating mechanism. This distinction, coupled with the concept of approximation decay rates, motivates us to define new complexity measures. These quantities are specifically designed to bridge the geometric structure of the parameter space—characterized by its decay rate—with the statistical nature of the underlying distributions.

Definition 2.2 (Optimal dimension). Assume $\tilde{\mathbf{Y}}$ are drawn from the distribution $\mathcal{P} \in \mathcal{P}$. Let the function $f(\cdot, \cdot) : \mathbb{N} \times \mathcal{P} \mapsto \mathbb{R}^+$, then a optimal dimension regarding K and f is defined as

$$k_f^* := \max \{j \mid j \geq 0, D_{j-1}(K) > f(j, \mathcal{P})\},$$

where we additionally define $D_k(K) = \infty$ for $k < 0$ and $D_k(K) = 0$ when $k \geq d = \dim(\mathcal{X})$. The corresponding optimal projection with dimension k is denoted as P_k^* .

Remark. Such definition can be extended to ℓ_2 with minimal effort.

The index $j - 1$ in the definition is set for the convenience of the computation and is not essentially important. Since $D_k(K) = 0$ for $k \geq d$ and $D_0(K) = \sup_{\theta \in K} \|\theta\|_2 > 0$ when $K \neq \{\mathbf{0}\}$, such k_f^* always

exists given $K \neq \{\mathbf{0}\}$. It is notable that $D_k(K)$ is irrelevant with $\tilde{\mathbf{Y}}$, \mathbf{Y} , and \mathcal{P} , while $f(j, \mathcal{P})$ is irrelevant with $\tilde{\mathbf{Y}}$, \mathbf{Y} , \mathcal{X} and K . Therefore, k_f^* is a quantity characterizing the interaction between K , \mathcal{X} , and \mathcal{P} , and does not depend on the realization $\tilde{\mathbf{Y}}$ and \mathbf{Y} once K , \mathcal{X} and \mathcal{P} are assumed.

The optimal dimension is closely related to the relevant minimax problem. In [Li and Neykov \[2026\]](#), the authors set $f(j, \mathcal{P})$ as $\frac{j^{1/4}}{\sqrt{N}}\sigma$ and $\frac{j^{1/4}\sqrt{\epsilon}}{N^{1/4}}\sigma$ in [Definition 2.2](#), respectively, and the resulting optimal dimensions (denoted as k_1^* and k_2^*) are crucial for the derived optimal minimax rate for QCO constraints. See [Section 4.2](#) and [4.3](#) for more details; In [Prasadan and Neykov \[2025\]](#), the authors obtain the optimal minimax rate for the ℓ_2 robust estimation problem with star-shape constraints. Though they do not explicitly specify the optimal dimension, there results can be summarized via such concept. Let $f(j, \mathcal{P}) = \sqrt{\frac{j}{N}}\sigma$ in [Definition 2.2](#) and denote the resulting optimal dimension by k_e^* . [Prasadan and Neykov \[2025\]](#) prove that

$$\inf_{\hat{\nu}} \sup_{\mu \in K} \mathbb{E}_\mu \|\hat{\nu}(\mathbf{Y}) - \mu\|_2^2 \asymp \sigma^2 \max \left\{ \frac{k_e^*}{N}, \epsilon^2 \right\}, \quad (2.2)$$

where inf is taken with respect to all measurable functions $\hat{\nu}$ of the observations \mathbf{Y} .

2.2 Relaxed Optimization Problem

When K is an ellipsoid defined by a positive semidefinite matrix and the distance is measured under the Euclidean norm, calculating the optimal k -dimensional approximation is equivalent to principal component analysis (PCA). However, computing the Kolmogorov k -width for general sets and arbitrary norms is a formidable task, typically incurring a super-polynomial time complexity. As an example in \mathbb{R}^d , [Brieden \[2002\]](#) demonstrate that computing $D_0(K)$ for convex polytopes (even symmetric ones) cannot be approximated in polynomial time within a factor of 1.090 unless $P = NP$. Even when restricting the metric to the ℓ_2 -norm, the geometric representation of K can be highly complex; in certain scenarios, K might only be accessible via a membership oracle rather than an explicit analytical expression, further exacerbating the computational intractability of the Kolmogorov widths.

To circumvent the computational intractability of exactly evaluating the Kolmogorov widths and identifying their associated optimal projections, we propose a computationally efficient alternative. In particular, we formulate the following semidefinite programming (SDP) relaxation, which allows us to tightly approximate the Kolmogorov widths in polynomial time.

Definition 2.3 (Approximate Kolmogorov k -Width). Let $\mathcal{X} = \mathbb{R}^d$ equipped with the Euclidean norm. For any $K \subset \mathcal{X}$ and $0 \leq k \leq d$, we define the approximate Kolmogorov k -width $\tilde{D}_k^2(K)$ as the optimal value of the following SDP.

$$\begin{aligned} & \min_{X \in \mathbb{R}^{d \times d}} \max_{\theta \in K} \theta^\top X \theta, \\ & \text{subject to } \text{tr}(X) = d - k, \mathbf{0} \preceq X \preceq \mathbf{I}_d, X^\top = X. \end{aligned} \quad (2.3)$$

We further denote the solution to (2.3) as X_k^* .

Indeed, (2.3) is a relaxed optimization problem compared with the original definition by the following fact.

Lemma 2.1. For any $0 \leq k \leq d$ and $\theta \in K$, we have $\theta^\top X_k^* \theta \leq \tilde{D}_k^2(K) \leq D_k^2(K)$.

Proof of Lemma 2.1. For any fixed k , let $P_{d-k} := \mathbf{I}_d - P_k^*$, where P_k^* satisfies $\sup_{\theta \in K} \|(\mathbf{I}_d - P_k)\theta\|_2^2 = D_k^2(K)$. This is equivalent to $\theta^\top P_{d-k} \theta \leq D_k^2(K)$. Since P_{d-k} is a feasible point of (2.3), we know $\tilde{D}_k^2(K) \leq D_k^2(K)$. The other inequality follows directly from the definition of (2.3). \square

We should be aware that $\tilde{D}_k(K)$ is possibly significantly smaller than $D_k(K)$. Hence, (2.3) is a non-trivial relaxation.

Example 1. Consider the case when $K = B_d(0, 1)$ i.e., the unit ball in \mathbb{R}^d . It is not hard to see that

$$D_k(B_d(0, 1)) = \begin{cases} 1 & (0 \leq k \leq d - 1), \\ 0 & (k = d). \end{cases}$$

However, we have $\tilde{D}_k(K) = \min_X \sqrt{\lambda_{\max}(X)}$ subject to $\text{tr}(X) = d - k, \mathbf{0} \preceq X \preceq \mathbf{I}_d, X = X^\top$.

Minimal calculation shows that $\tilde{D}_k(K) = \sqrt{1 - \frac{k}{d}}$. The ratio between $D_k(K)$ and $\tilde{D}_k(K)$ is $\sqrt{d} \gtrsim 1$ when $k = d - 1$.

The SDP (2.3) formulated above exhibits a bilevel optimization structure with respect to X and θ . Focusing on the inner problem, for any fixed feasible matrix $X \in \mathbb{R}^{d \times d}$, we aim to solve the following classical constrained quadratic maximization:

$$\max_{\theta} \theta^\top X \theta, \quad \text{subject to } \theta \in K. \quad (2.4)$$

While (2.4) is conventionally formulated over an origin-symmetric convex set K , identifying its exact global optimum is analytically intractable. Even in the seemingly benign case where X is positive semidefinite and K is a polytope in \mathbb{R}^d , maximizing a convex function over a polytope remains notoriously NP-hard due to the potentially exponential growth of the number of vertices with respect to d . To circumvent this, [Bhattiprolu et al. \[2021\]](#) recently proposed a polynomial-time approximation framework based on a modified ellipsoid method (which is distinct from the algorithm discussed later in Section 3). This framework ensures that a κ -approximate optimal value for (2.4) can be efficiently obtained. Crucially, [Bhattiprolu et al. \[2021\]](#) establish that satisfying the type-2 condition on the set K is a strict prerequisite for the tractability of this approximation.

Definition 2.4 (Minkowski gauge). Let \mathcal{X} be a vector space \mathcal{X} and $K \subset \mathcal{X}$ is any fixed subset. The Minkowski gauge of K is defined as a function $\rho_K : \mathcal{X} \rightarrow [0, \infty]$, where

$$\rho_K(\theta) = \inf_r \{r \mid r \in \mathbb{R}^+, \theta \in rK, \theta \in \mathcal{X}\}. \quad (2.5)$$

If such r does not exist for a given θ , we exclusively define $\rho_K(\theta) = \infty$.

Definition 2.5 (Type-2 condition). Let $\{\epsilon_i\}_{i=1}^m$ be i.i.d. Rademacher random variables for some $m \in \mathbb{N}^+$. For $K \subset \mathcal{X}$ above, K is type-2 with constant $T_2(K)$ if for any $m \in \mathbb{N}^+$ and any $\theta_i \in \mathbb{R}^d$, we have

$$\mathbb{E}_\epsilon \rho_K^2 \left(\sum_{i=1}^m \epsilon_i \theta_i \right) \leq T_2^2(K) \sum_{i=1}^m \rho_K^2(\theta_i). \quad (2.6)$$

$T_2^2(K)$ is referred as the type-2 constant of K .

Example 2 (ℓ_p unit ball). [Milman and Schechtman \[1986\]](#) establish the following results about the ℓ_p unit ball (i.e., $\rho_K = \|\cdot\|_p$) in \mathbb{R}^d .

$$T_2(K) \asymp \begin{cases} d^{\frac{1}{p}-\frac{1}{2}}, & (1 \leq p \leq 2), \\ \sqrt{\min\{p, \ln d\}}, & (2 \leq p \leq \infty) \end{cases}$$

Crucially, [Bhattachiprolu et al. \[2021\]](#) establish that the originally intractable problem (2.4) can be rigorously reduced to the following equivalent semidefinite program:

$$\begin{aligned} & \max_{\Theta} \langle X, \Theta \rangle, \\ & \text{subject to } \mathbb{E} \left[\rho_K^2(\Theta^{\frac{1}{2}}g) \right] \leq 1, \Theta \succeq \mathbf{0}, \end{aligned} \quad (2.7)$$

where $g \in \mathbb{R}^d$ denotes a standard Gaussian random vector.

Proposition 2.2 (Observation 4.1 of [Bhattachiprolu et al. \[2021\]](#)). *The optimal value of the original SDP (2.4) is equivalent to that of (2.7).*

Proof. Let $f_{\text{ori}}(A)$ denote the optimal value of (2.4), $f_{\text{equ}}(A)$ denote the optimal value of (2.7), and $\theta_1 \in K$ be any vector such that $\theta_1^\top A \theta_1 = f_{\text{raw}}(A)$. Consider $\Theta = \theta_1 \theta_1^\top$. Since $\mathbb{E} \rho_K^2(\Theta^{\frac{1}{2}}g) = \mathbb{E} \rho_K^2 \left(\frac{\theta_1 X_1^\top g}{\|\theta_1\|_2} \right) = \mathbb{E} (\theta_1^\top g)^2 \cdot \rho_K^2 \left(\frac{\theta_1}{\|\theta_1\|_2} \right) = \|\theta_1\|_2^2 \cdot \rho_K^2 \left(\frac{\theta_1}{\|\theta_1\|_2} \right) = \rho_K^2(\theta) \leq 1$, we know that $f_{\text{ori}}(A) \leq f_{\text{equ}}(A)$.

On the other hand, for any \mathbf{X} that is feasible for (2.7), we have

$$\langle A, \Theta \rangle = \mathbb{E}_g \left\langle \Theta^{\frac{1}{2}}g, A \Theta^{\frac{1}{2}}g \right\rangle = \rho_K^2(\Theta^{\frac{1}{2}}g) \cdot \mathbb{E}_g \left(\frac{g^\top \Theta^{\frac{1}{2}}}{\rho_K(\Theta^{\frac{1}{2}}g)} A \frac{\Theta^{\frac{1}{2}}g}{\rho_K(\Theta^{\frac{1}{2}}g)} \right) \stackrel{(i)}{\leq} \rho_K^2(\Theta^{\frac{1}{2}}g) f_{\text{ori}}(A) \leq f_{\text{ori}}(A),$$

where (i) is from the fact that $\frac{\Theta^{1/2}g}{\rho_K(\Theta^{1/2}g)} \in K$ by the definition of the [Minkowski gauge](#). Taking the supreme over Θ on the LHS completes the proof. \square

As a final preliminary, we introduce the notions of sign-invariant norms, exact 2-convexity, and balanced set .

Definition 2.6 (Sign-invariant norm). A norm $\|\cdot\|$ is *sign-invariant* if for any $\theta \in \mathbb{R}^d$ and set of signs $\gamma = (\gamma_1, \dots, \gamma_d), \gamma_i \in \{-1, 1\}, i = 1, \dots, d$, we have $\|\gamma \odot \theta\| = \|\theta\|$, where \odot denotes the entrywise multiplication.

Definition 2.7 (Exact 2-convex norm). The Minkowski gauge $\rho_K(\cdot)$ of a set K is 2-convex if for any $m \in \mathbb{N}^+$ and $\theta_1, \dots, \theta_m \in \mathbb{R}^d$, we have

$$\rho_K \left(\left(\sum_{i=1}^m \theta_i^2 \right)^{\frac{1}{2}} \right) \leq c(K) \left(\sum_{i=1}^m \rho_K^2(\theta_i) \right)^{\frac{1}{2}},$$

where $\theta^2 = (\theta_1^2, \dots, \theta_d^2)^\top$ is the entrywise square of θ and $c(K)$ is a quantity that only depends on K . When $c(K) = 1$, $\rho_K(\cdot)$ is said to be *exact 2-convex*.

Definition 2.8 (Balanced set). For a subset $K \subset \mathcal{X} = \mathbb{R}^d$, K is balanced if there exist $0 < r \leq R \leq \infty$ such that $B_d(0, r) \subset K \subset B_d(0, R)$.

The following theorem from [Bhattiprolu et al. \[2021\]](#) formally establishes the existence of an approximate efficient solution to (2.4).

Theorem 2.3 (Proposition 4.2 and Theorem 7.6 of [Bhattiprolu et al. \[2021\]](#)). *For the equivalent semidefinite program (2.7), if K is type-2 and balanced, and the Minkowski gauge ρ_K is a sign-invariant and exact 2-convex norm, then there is an algorithm $\mathcal{A}(X, R, r)$, on any input $X \in \mathbb{R}^{d \times d}$, \mathcal{A} runs in polynomial time of $d, |\ln R|, |\ln r|, \text{bit}(X)$ and outputs a vector $\theta_1 \in K$ such that $\theta_1^\top X \theta_1 \gtrsim \frac{1}{\kappa} \max_{\theta \in K} \theta^\top X \theta$ with $\kappa \lesssim \ln d$.*

For a comprehensive treatment of this framework—encompassing Theorem 2.3 and the underlying proof techniques—we refer the reader to [Bhattiprolu et al. \[2021\]](#). Moving forward, we operate under the assumption that K is type-2 and balanced, and its associated Minkowski gauge is sign-invariant and exactly 2-convex. As we establish in Lemma 2.4, the QCO sets inherently satisfy these structural prerequisites and thus emerge as a natural special case.

Back to our problem (2.3). For the inner problem, we leverage Theorem 2.3 and know that for any feasible X in (2.3), there exists an oracle $\mathcal{O}(X) \in K$ guaranteed by Theorem 2.3 that approximates the optimal value of (2.3) by a constant κ :

$$\mathcal{O}(X)^\top X \mathcal{O}(X) \geq \frac{1}{\kappa} \max_{\theta \in K} \theta^\top X \theta. \quad (2.8)$$

Importantly, the approximation factor κ remains moderately small, being bounded strictly by a polylogarithmic function of the dimension d , i.e., $\kappa = \text{polylog}(d)$.

Concluding this preparatory section, we introduce the formal definition of QCO sets and establish that they naturally fulfill the prerequisites of Theorem 2.3.

Definition 2.9 (Quadratically convex orthosymmetric (QCO) set). Given a set $K \subset \mathbb{R}^d$, K is a QCO set if it satisfies the following conditions:

- (1), K is convex;
- (2), K is quadratically convex, i.e., K^2 is also convex, where K^2 is defined as

$$K^2 := \left\{ (\theta_1^2, \dots, \theta_d^2)^\top \mid (\theta_1, \dots, \theta_d)^\top \in K \right\};$$

- (3), K is orthosymmetric, i.e., if $\theta = (\theta_1, \dots, \theta_d)^\top \in K$, then for any $\eta = (\eta_1, \dots, \eta_d)^\top, \eta_i \in \{-1, 1\}, i = 1, \dots, d$, we have $\theta_\eta := \eta \odot \theta \in K$.

Lemma 2.4 (Neykov [2026]). *For any QCO set K , K is type-2 and exactly 2-convex. Furthermore, the type-2 constant $T_2(K) = c \ln d$ for some universal constant c .*

The proof of Lemma 2.4 is essentially from Neykov [2026]. We place the proof of Lemma 2.4 in Appendix B.1 for the reader's convenience.

3 Approximate Subgradient and Ellipsoid Method

3.1 Approximate Subgradient and Related Property

In the context of the SDP (2.3), we define the optimal value function $h_S(X) := \max_{\theta \in K} \theta^\top X \theta$ for any feasible X , and denote its corresponding maximizer by θ_X^* (provided it exists). It is straightforward to verify that the outer product $\theta_X^* \theta_X^{*\top}$ serves as a valid subgradient of h_S at X . This directly follows from the fact that for any feasible X_1, X_2 , the following inequality holds: $h_S(X_2) = \theta_{X_2}^{*\top} X_2 \theta_{X_2}^* \geq \theta_{X_1}^{*\top} X_2 \theta_{X_1}^* = h_S(X_1) + \langle \theta_{X_1}^* \theta_{X_1}^{*\top}, X_2 - X_1 \rangle$. Consequently, since the oracle $\mathcal{O}(X)$ only provides an approximation of the optimal value, we anticipate it to return an approximate subgradient, thereby satisfying a relaxed subgradient condition.

Lemma 3.1. *For any feasible X and the corresponding $\mathcal{O}(X)$, $\mathcal{O}(X)\mathcal{O}(X)^\top$ is an ω -subgradient of h_S at X in the sense that for any feasible X' , we have*

$$h_S(X') \geq \frac{1}{\kappa} h_S(X) + \mathcal{O}(X)^\top (X' - X) \mathcal{O}(X). \quad (3.1)$$

Proof. Lemma 3.1 can be directly verified via (2.8). We have

$$\begin{aligned} h_S(X) + \mathcal{O}(X)^\top (X' - X) \mathcal{O}(X) &= h_S(X) - \mathcal{O}(X)^\top X \mathcal{O}(X) + \mathcal{O}(X)^\top X' \mathcal{O}(X) \\ &\leq \left(1 - \frac{1}{\kappa}\right) h_S(X) + h_S(X'). \end{aligned}$$

Cancelling $h_S(X)$ on the both sides yields the results. \square

Lemma 3.1 tells us that while $\mathcal{O}(X)\mathcal{O}(X)^\top$ is no longer a rigorous subgradient, it preserves the property somehow by adding a $\frac{1}{\kappa}$ factor, which is the key to the following ellipsoid method.

3.2 Ellipsoid Method

3.2.1 Problem without Equality Constraints

The ellipsoid method (Shor [1977], Yudin and Nemirovski [1976]) stands as a fundamental iterative tool for constrained convex optimization. Diverging from standard gradient-descent techniques that take local steps towards the optimal solution, the ellipsoid method systematically constructs a sequence of bounding ellipsoids strictly guaranteed to enclose the true minimizer. The algorithm derives its convergence from the continuous, geometric reduction in the volume of the ellipsoids. In this theoretical light, the approach can be understood as a generalization of the one-dimensional bisection method to the Euclidean space \mathbb{R}^d . Crucially, the exponential decay rate of these ellipsoidal volumes is precisely what guarantees polynomial-time computational complexity across a broad spectrum of optimization paradigms.

Consider the unconstrained convex optimization problem:

$$\text{minimize } f(x), x \in \mathbb{R}^d, \quad (3.2)$$

where $f : \mathbb{R}^d \rightarrow \mathbb{R}$ ($d \geq 2$) is convex. We do not require f to be differentiable, assuming only the availability of a subgradient oracle. Let the algorithm be initialized with an ellipsoid $E^{(0)} = \{x \mid (x - x^{(0)})^\top (P^{(0)})^{-1} (x - x^{(0)}) \leq 1\}$ centered at $x^{(0)} \in \mathbb{R}^d$, where $P^{(0)} \in \mathbb{R}^{d \times d}$ is a positive definite matrix. At the k -th iteration, let $x^{(k)}$ and $P^{(k)}$ define the current search region, and compute a subgradient $g^{(k)} \in \partial f(x^{(k)})$. If $\mathbf{0} \in \partial f(x^{(k)})$, the exact global minimizer is successfully identified at $x^{(k)}$. Otherwise, for any non-zero $g^{(k)}$, the subgradient inequality implies that any point x satisfying $\langle x - x^{(k)}, g^{(k)} \rangle > 0$ strictly degrades the objective value relative to $f(x^{(k)})$, fundamentally excluding such points from optimality. Consequently, our objective is to construct an updated ellipsoid $E^{(k+1)}$, parameterized by $x^{(k+1)}$ and $P^{(k+1)}$, that tightly circumscribes the remaining feasible half-ellipsoid:

$$\left\{x \mid x \in E^{(k)}, \langle x - x^{(k)}, g^{(k)} \rangle \leq 0\right\} \subset E^{(k+1)} := \left\{x \mid (x - x^{(k+1)})^\top [P^{(k+1)}]^{-1} (x - x^{(k+1)}) \leq 1\right\}.$$

We select the ellipsoid with smallest volume while satisfying the condition above. To determine such ellipsoid, we can first assume $E^{(k)} = \mathbf{I}_d$, and transform back with the affine transformation $(E^{(k)})^{\frac{1}{2}}$. It is proved (see the original work by [Yudin and Nemirovskii \[1976\]](#)) that the ellipsoid in the $(k+1)$ -th iteration is determined by:

$$\begin{aligned} x^{(k+1)} &= x - \frac{P^{(k)} g^{(k)}}{(d+1) \sqrt{g^{(k)\top} P^{(k)} g^{(k)}}}, \\ P^{(k+1)} &= \frac{d^2}{d^2 - 1} \left(P^{(k)} - \frac{2}{d+1} \frac{P^{(k)} g^{(k)} g^{(k)\top} P^{(k)}}{g^{(k)\top} P^{(k)} g^{(k)}} \right). \end{aligned} \quad (3.3)$$

The updates above can be generalized to convex optimization programs with inequality constraints. Consider the following constrained problem:

$$\text{minimize } f_0(x), \quad \text{subject to } f_i(x) \leq 0, i = 1, \dots, m, \quad (3.4)$$

where the functions f_i for $i = 0, 1, \dots, m$ are convex and equipped with subgradient oracles. Guided by the same geometric intuition, the update procedure at any iteration hinges strictly on the feasibility of the current iterate $x^{(k)}$. In particular: (1) if $x^{(k)}$ violates any constraint, we identify a violated index (e.g., the smallest i satisfying $f_i(x^{(k)}) > 0$) and execute the update using a *feasibility cut* derived from the subgradient of f_i ; (2) if $x^{(k)}$ resides within the feasible region, we proceed with an *objective cut* based on the subgradient of f_0 , mirroring the unconstrained scenario. Thus, from the perspective of the ellipsoid method, constrained optimization is mechanistically equivalent to its unconstrained counterpart; in both paradigms, the algorithm consistently curtails the search space by constructing a ellipsoid with smaller volume.

3.2.2 Convergence Analysis

As noted earlier, the theoretical convergence of the ellipsoid algorithm relies on the monotonic decrease in the volume of the constructed ellipsoids, namely, $\text{Vol}(E^{(k+1)}) < \text{Vol}(E^{(k)})$. In fact, we have the following famous results

Theorem 3.2. For the convex optimization problem (3.2) and (3.4), we have

$$\frac{\text{Vol}(E^{(k+1)})}{\text{Vol}(E^{(k)})} = \frac{d^d}{(d+1)^{\frac{d+1}{2}} \cdot (d-1)^{\frac{d-1}{2}}} < e^{-\frac{1}{2d}}.$$

Furthermore, for the feasible ϵ -suboptimal set

$$\left\{ x \mid x \in \mathbb{R}^d, f_0(x) - f_0(x^*) \leq \epsilon, f_i(x) \leq 0, i = 1, \dots, m \right\}$$

where x^* is the true minimizer, if it has positive volume, then we have

$$\left| f(x^\dagger) - f(x^*) \right| \leq c \cdot \exp \left\{ -\frac{M}{2d^2} \right\},$$

where $x_M^\dagger := \arg \min_{x^{(i)}, i=1,2,\dots,M} f(x^{(i)})$ is the best candidate up to the M -th iteration, and c is a constant that only depends on f_0 .

From Theorem 3.2, we know that the ellipsoid method converges exponentially under very mild conditions. Equivalently, for any given precision ϵ , the ellipsoid method finds at least one feasible point that is ϵ -suboptimal in at most $c'd^2 \ln(\frac{1}{\epsilon})$ iterations for some constant c' . The proof of Theorem 3.2 derives from the updates 3.3 and the computation of their volumes, which can be found in many references (Shor [1970], Yudin and Nemirovskii [1976], Khachiyan [1980]). Note that when computing the volume ratio between $E^{(k+1)}$ and $E^{(k)}$, we can again assume $x^{(k)} = \mathbf{0}, P^{(k)} = \mathbf{I}_d$ since the volume ratio is invariant under the affine transformations.

The pseudo code of the ellipsoid method for the problems (3.2) and (3.4) is provided as follows.

3.2.3 Problem with Equality Constraints

The approximate Kolmogorov width SDP (2.3) can be equivalently formulated as a convex program over symmetric matrices, subject to the trace constraint $\text{tr}(X) = d - k$ and the positive semidefinite bounds $\mathbf{0} \preceq X \preceq \mathbf{I}_d$:

$$\begin{aligned} & \min_X h_S(X), \\ & \text{subject to } \text{tr}(X) = d - k, \mathbf{0} \preceq X \preceq \mathbf{I}_d, X^\top = X. \end{aligned} \tag{3.5}$$

To execute the ellipsoid method in the presence of affine equality constraints, we adopt the framework established by Shah et al. [2001]. Consider a general convex optimization problem with both inequality and equality constraints:

$$\begin{aligned} & \text{minimize } f(x), \\ & \text{subject to } f_i(x) < 0, Ax = b, \end{aligned} \tag{3.6}$$

where $A \in \mathbb{R}^{\text{rank}(A) \times d}$ is a matrix with full row rank. Assuming the initial iterate $x^{(0)}$ lies within the affine subspace defined by $Ax = b$, Shah et al. [2001] derive the following projected update rules

Algorithm 1: ellipsoid method with (inequality) constraints

Input: $d, f_0, \dots, f_m, x^{(0)}, P^{(0)}, \epsilon, c'$
 1 $k = 0, x_{\text{best}} = \mathbf{0}, v_{\text{best}} = \infty$
 2 $M = c' d_A^2 \ln\left(\frac{1}{\epsilon}\right)$
 3 **while** $k < M$ **do**
 4 **if** $x^{(k)}$ *is feasible* **then**
 5 **if** $f(x^{(k)}) < v_{\text{best}}$ **then**
 6 $x_{\text{best}} = x^{(k)}$
 7 $v_{\text{best}} = f(x^{(k)})$
 8 Select any $g^{(k)} \in \partial f_0(x^{(k)})$
 9 **else**
 10 Find the smallest i such that $f_i(x^{(k)}) > 0$
 11 Select any $g^{(k)} \in \partial f_i(x^{(k)})$
 12 Update $x^{(k+1)} = x^{(k)} - \frac{P^{(k)} g^{(k)}}{(d+1)\sqrt{g^{(k)\top} P^{(k)} g^{(k)}}}$
 13 Update $P^{(k+1)} = \frac{d^2}{d^2-1} \left(P^{(k)} - \frac{2}{d+1} \frac{P^{(k)} g^{(k)} g^{(k)\top} P^{(k)}}{g^{(k)\top} P^{(k)} g^{(k)}} \right)$
 14 $k = k + 1$
Output: $x_{\text{best}}, v_{\text{best}}$

for the center $x^{(k)}$ and the shape matrix $P^{(k)}$:

$$\begin{aligned}
 r^{(k)} &:= \frac{\left(P^{(k)} - P^{(k)} A^\top (A P^{(k)} A^\top)^{-1} A P^{(k)} \right) g^{(k)}}{\sqrt{g^{(k)\top} \left(P^{(k)} - P^{(k)} A^\top (A P^{(k)} A^\top)^{-1} A P^{(k)} \right) g^{(k)}}}, \\
 x^{(k+1)} &= x^{(k)} - \frac{1}{d+1} r^{(k)}, \\
 P^{(k+1)} &= \frac{d^2}{d^2-1} \left(P^{(k)} - \frac{2}{d+1} r^{(k)} r^{(k)\top} \right).
 \end{aligned} \tag{3.7}$$

Following standard cutting-plane logic, $g^{(k)} \in \partial f_0(x^{(k)})$ is an objective cut if $x^{(k)}$ is fully feasible. If $x^{(k)}$ violates any constraint, $g^{(k)} \in \partial f_i(x^{(k)})$ functions as a feasibility cut, where i identifies the first constraint satisfying $f_i(x^{(k)}) \geq 0$.

While the update equations in (3.7) appear algebraically intricate, they stem from a straightforward geometric reduction. Specifically, the affine constraint $Ax = b$ allows any feasible solution x to be explicitly parameterized as $x = A^\top (A A^\top)^{-1} b + B y$. Here, $y \in \mathbb{R}^{d-\text{rank}(A)}$ serves as the coordinate vector in the reduced space, and the columns of the matrix $B \in \mathbb{R}^{d \times (d-\text{rank}(A))}$ constitute a basis for the null space $\mathcal{N}(A)$. This transformation seamlessly projects the original problem onto an affine subspace of dimension $d_A := d - \text{rank}(A)$. Because convexity is preserved under affine composition, the functions $f_i(\cdot), i = 0, 1, \dots, m$ remain convex over this restricted domain. By absorbing the equality constraints directly into the definition of the search space, the problem reduces to a standard inequality-constrained framework. As a result, the theoretical guarantees of Theorem 3.2 carry over directly to this projected setting.

Corollary 3.3. Consider the convex optimization problem (3.6), and define the projected ellipsoid $\tilde{E}^{(k)} := E^{(k)} \cap \{x \mid Ax = b\}$. The volumetric decay rate is strictly bounded by

$$\frac{\text{Vol}(\tilde{E}^{(k+1)})}{\text{Vol}(\tilde{E}^{(k)})} = \frac{d^{d_A}}{(d+1)^{\frac{d_A+1}{2}} \cdot (d-1)^{\frac{d_A-1}{2}}} < 1,$$

where $\text{Vol}(\cdot)$ denotes the Lebesgue measure in the d_A -dimensional affine subspace. Furthermore, suppose that for some $\epsilon > 0$, the ϵ -suboptimal feasible region $\{x \in \mathbb{R}^d \mid f(x) - f(x^*) \leq \epsilon, f_i(x) < 0, Ax = b\}$ possesses a strictly positive d_A -dimensional volume, where x^* represents the constrained global minimizer. Letting $x_{(M)}^\dagger := \arg \min_{1 \leq i \leq M} f(x^{(i)})$ denote the best feasible iterate observed up to step M , we obtain the exponential convergence guarantee

$$f(x_{(M)}^\dagger) - f(x^*) \leq c \cdot \exp\left\{-\frac{M}{2d_A^2}\right\},$$

where $c > 0$ is an absolute constant determined solely by f_0 .

The pseudo code of the ellipsoid method when equality constraints present is shown in Algorithm 2.

Algorithm 2: ellipsoid method with (equality) constraints

Input: $d, f_0, \dots, f_m, x^{(0)}, P^{(0)}, \epsilon, c'$

- 1 $k = 0, x_{\text{best}} = \mathbf{0}, v_{\text{best}} = \infty$
- 2 $M = c' d^2 \ln\left(\frac{1}{\epsilon}\right)$
- 3 **while** $k < M$ **do**
- 4 **if** $x^{(k)}$ **is feasible** **then**
- 5 **if** $f_0(x^{(k)}) < v_{\text{best}}$ **then**
- 6 $x_{\text{best}} = x^{(k)}$
- 7 $v_{\text{best}} = f_0(x^{(k)})$
- 8 Select any $g^{(k)} \in \partial f_0(x^{(k)})$
- 9 **else**
- 10 Find the smallest i such that $f_i(x^{(k)}) \geq 0$
- 11 Select any $g^{(k)} \in \partial f_i(x^{(k)})$
- 12 Set $r^{(k)} = \frac{(P^{(k)} - P^{(k)} A^\top (A P^{(k)} A^\top)^{-1} A P^{(k)}) g^{(k)}}{\sqrt{g^{(k)\top} (P^{(k)} - P^{(k)} A^\top (A P^{(k)} A^\top)^{-1} A P^{(k)}) g^{(k)}}$
- 13 Update $x^{(k+1)} = x^{(k)} - \frac{1}{d+1} r^{(k)}$
- 14 Update $P^{(k+1)} = \frac{d^2}{d^2-1} \left(P^{(k)} - \frac{2}{d+1} r^{(k)} r^{(k)\top} \right)$
- 15 $k = k + 1$

Output: $x_{\text{best}}, v_{\text{best}}$

3.3 Ellipsoid Method with Approximate Subgradient

A prerequisite of the classical ellipsoid method is the availability of exact subgradients for the functions f_i ($i = 0, 1, \dots, m$) to accurately determine the cutting planes. In our setting, however,

Lemma 3.1 demonstrates that the subgradient information for $h_S(X)$ is limited to the rank-one matrix $\mathcal{O}(X)\mathcal{O}(X)^\top$, which merely acts as a κ -approximate subgradient for some constant $\kappa > 1$. Geometrically, the ellipsoid method can be viewed as an elimination process: rather than directly isolating the true minimizer, each iteration systematically prunes the region $E^{(k)} \cap (E^{(k+1)})^c$ mathematically precluded from containing it.

Specifically, at iteration k , any point x satisfying the strict halfspace inequality $\langle x - x^{(k)}, g^{(k)} \rangle > 0$ yields an objective value strictly worse than $f(x^{(k)})$, fundamentally disqualifying it as the minimizer. By substituting the exact $g^{(k)}$ with the approximate oracle $\mathcal{O}(X^{(k)})\mathcal{O}(X^{(k)})^\top$ in the update step, we essentially compromise the exactness of this elimination. Consequently, the guarantee of global exact minimization is lost; instead, the algorithm converges to an approximate optimum $X_{M\text{-best}} :=$

$\arg \min_{X \in \{X^{(1)}, \dots, X^{(M)}\}} h_S(X)$ such that $f(X^*) \leq f(X_{M\text{-best}}) \lesssim \kappa f(X^*)$. We rigorously establish this convergence behavior in the following theorem.

Theorem 3.4. *Consider the semidefinite program (2.3) for approximating the Kolmogorov widths. Substituting the exact subgradient $g(\cdot) \in \partial h_S(\cdot)$ with the approximate oracle $\mathcal{O}(\cdot)\mathcal{O}(\cdot)^\top$ in Algorithm 2 yields the following strict volumetric decay:*

$$\frac{\text{Vol}(\tilde{E}^{(k+1)})}{\text{Vol}(\tilde{E}^{(k)})} = \frac{d^{2\tilde{d}}}{(d^2 + 1)^{\frac{\tilde{d}+1}{2}} \cdot (d^2 - 1)^{\frac{\tilde{d}-1}{2}}} < 1.$$

Let $X_{M\text{-best}} := \arg \min_{X \in \{X^{(1)}, \dots, X^{(M)}\}} h_S(X)$ denote the best feasible iterate generated after M steps of

Algorithm 2. The suboptimality gap with respect to the true global minimum X^* is bounded by

$$h_S(X_{M\text{-best}}) \leq \kappa \left(h_S(X^*) + c \cdot \exp \left\{ -\frac{M}{2\tilde{d}^2} \right\} \right).$$

Consequently, to achieve a κ -approximate solution up to an additive precision $\epsilon > 0$, such that

$$h_S(X^*) \leq h_S(X_{M\text{-best}}) \leq \kappa (h_S(X^*) + \epsilon),$$

the algorithm requires at most $M \leq c\tilde{d}^2 \cdot \ln\left(\frac{1}{\epsilon}\right)$ iterations. Here, $\tilde{d} = \frac{(d-1)(d+2)}{2}$ represents the intrinsic dimension of the search space under the trace (consumes 1 degree of freedom) and symmetry constraints (consume $\frac{d(d-1)}{2}$ degrees of freedom).

Remark. One might initially find the coexistence of the ambient dimension d^2 and the intrinsic dimension \tilde{d} in the decay rate counterintuitive. This distinction naturally arises because our algorithm executes the matrix updates in the full ambient space $\mathbb{R}^{d \times d}$ (which introduces the d^2 terms in the update rules), whereas the effective search space — and consequently the Lebesgue measure of the feasible region — is confined to a \tilde{d} -dimensional affine subspace.

The rigorous proof of Theorem 3.4 is deferred to Appendix B.2. While Theorem 3.4 theoretically guarantees that the best historical iterate $X_{M\text{-best}}$ achieves the desired suboptimality bound, this specific matrix is practically unidentifiable since the exact objective function $h_S(X)$ is unknown. To circumvent this, we leverage the approximation oracle $\mathcal{O}(X)$ developed in Section 2 to evaluate the iterates. Specifically, replacing the intractable objective $h_S(X^{(i)})$ with its computable surrogate $\mathcal{O}(X^{(i)})^\top X^{(i)} \mathcal{O}(X^{(i)})$ yields the fully empirical Approximate Ellipsoid Method, detailed in Algorithm

Algorithm 3: ellipsoid method for SDP (2.3)

Input: d , constraints $\text{tr}(X) = d - k$, $\mathbf{0} \preceq X \preceq \mathbf{I}_d$, $X = X^\top$, $X^{(0)}$, $P^{(0)}$, ϵ , c'

- 1 Deduce $f_1, \partial f_1, \dots$ from the inequality constraints
- 2 $k = 0$, $X_k^\dagger = \mathbf{0}$, $v_{\text{best}} = \infty$
- 3 $M = c' d^2 \ln\left(\frac{1}{\epsilon}\right)$
- 4 **while** $k < M$ **do**
- 5 **if** $X^{(k)}$ *is feasible* **then**
- 6 **if** $\mathcal{O}(X^{(k)})^\top X^{(k)} \mathcal{O}(X^{(k)}) < v_{\text{best}}$ **then**
- 7 $X_k^\dagger = X^{(k)}$
- 8 $v_{\text{best}} = \mathcal{O}(X^{(k)})^\top X^{(k)} \mathcal{O}(X^{(k)})$
- 9 $g^{(k)} = \mathcal{O}(X^{(k)}) \mathcal{O}(X^{(k)})^\top$
- 10 **else**
- 11 Find the smallest i such that $f_i(X^{(k)}) \geq 0$
- 12 Select any $g^{(k)} \in \partial f_i(X^{(k)})$
- 13 Set $r^{(k)} = \frac{(P^{(k)} - P^{(k)})A^\top (AP^{(k)}A^\top)^{-1} AP^{(k)} g^{(k)}}{\sqrt{g^{(k)\top} (P^{(k)} - P^{(k)})A^\top (AP^{(k)}A^\top)^{-1} AP^{(k)} g^{(k)}}}$
- 14 Update $X^{(k+1)} = X^{(k)} - \frac{1}{d+1} r^{(k)}$
- 15 Update $P^{(k+1)} = \frac{d^2}{d^2-1} \left(P^{(k)} - \frac{2}{d+1} r^{(k)} r^{(k)\top} \right)$
- 16 $k = k + 1$

Output: $X_k^\dagger, v_{\text{best}}$

3. Noting that the scaling constant c' is chosen to be sufficiently large and depends solely on ϵ and K , we define the practically selected optimum as $X_k^\dagger := \arg \min_{X \in \{X^{(1)}, \dots, X^{(M)}\}} \mathcal{O}(X)^\top X \mathcal{O}(X)$, which

immediately yields the following corollary.

Corollary 3.5. *For the output X_k^\dagger of Algorithm 3, we have*

$$\frac{1}{\kappa} h_S(X^\star) \leq \mathcal{O}(X_k^\dagger)^\top X_k^\dagger \mathcal{O}(X_k^\dagger) \leq \kappa h_S(X^\star). \quad (3.8)$$

Proof. The first inequality follows from (2.8) and the definition of X^\star :

$$\mathcal{O}(X_k^\dagger)^\top X_k^\dagger \mathcal{O}(X_k^\dagger) \geq \frac{1}{\kappa} h(X_k^\dagger) \geq \frac{1}{\kappa} h_S(X^\star).$$

The second inequality follows from the definition of X_k^\dagger , (unknown) $X_{M\text{-best}}$, h_S and Theorem 3.4:

$$\mathcal{O}(X_k^\dagger)^\top X_k^\dagger \mathcal{O}(X_k^\dagger) \leq \mathcal{O}(X_{M\text{-best}})^\top X_{M\text{-best}} \mathcal{O}(X_{M\text{-best}}) \leq h_S(X_{M\text{-best}}) \leq \kappa h_S(X^\star).$$

The proof is completed. \square

Moving forward, we denote $\tilde{h}_S(X_k^\dagger) := \mathcal{O}(X_k^\dagger)^\top X_k^\dagger \mathcal{O}(X_k^\dagger)$ as the resulting approximate optimal value of X_k^\dagger .

4 Testing Procedures

4.1 Approximate Optimal Dimension

Armed with the approximation frameworks developed for the inner maximization (2.4) in Section 2 and the outer minimization in Section 3, we are now prepared to bypass the computational intractability inherent in evaluating the exact Kolmogorov widths. Recall that the SDP (2.3) acts as a rigorous convex relaxation of the original width computation, with X_k^\dagger serving as its computationally efficient approximate solution. Consequently, for general balanced, [type-2](#) and [2-convex](#) sets where $D_k(K)$ remains analytically or computationally elusive, it is a natural progression to substitute the true width with the surrogate quantity $\sqrt{\tilde{h}_S(X_k^\dagger)}$. Building upon this practical substitution, we formally define the *approximate optimal dimension* in terms of $\tilde{h}_S(X_k^\dagger)$ as follows.

Definition 4.1 (Approximate optimal Dimension). Define $\tilde{D}_k(K) := \sqrt{\tilde{h}_S(X_k^\dagger)}$, $k = 1, 2, \dots$ as the approximate Kolmogorov k -width. Assume that the data $\tilde{\mathbf{Y}}$ are drawn from the distribution $\mathcal{P} \in \mathcal{P}$. Let the function $f(\cdot, \cdot) : \mathbb{N} \times \mathcal{P} \mapsto \mathbb{R}^+$, then the approximate optimal dimension regarding K and f is defined as

$$\tilde{k}_f^* := \max \left\{ j \mid j \geq 0, \tilde{D}_{j-1}(K) > f(j, \mathcal{P}) \right\},$$

where similarly, we additionally define $\tilde{D}_k(K) = \infty$ for $k < 0$ and $\tilde{D}_k(K) = 0$ when $k > \dim(\mathcal{X})$.

As previously discussed, the exact first and second optimal dimensions are fundamental for characterizing the minimax rate in the robust ℓ_2 -norm testing framework (4.2). When transitioning to their approximate counterparts—the approximate first and second optimal dimensions, which are obtained by selecting the identical criterion $f(j, \mathcal{P})$ —we establish the following property.

Lemma 4.1. Let $f_1(j, \mathcal{P}) = \frac{j^{1/4}}{\sqrt{N}}\sigma$, $f_2(j, \mathcal{P}) = \frac{j^{1/4}\sqrt{\epsilon}}{N^{1/4}}\sigma$. For the exact first and second optimal dimensions k_1^*, k_2^* and the approximate first and second optimal dimensions $\tilde{k}_1^* := \tilde{k}_{f_1}^*$, $\tilde{k}_2^* := \tilde{k}_{f_2}^*$, we have the following inequalities.

$$\begin{aligned} \tilde{k}_1^* &\leq \min \left\{ \kappa^2(k_1^* + 1) - 1, d \right\}, \\ \tilde{k}_2^* &\leq \min \left\{ \kappa^2(k_2^* + 1) - 1, d \right\}. \end{aligned} \tag{4.1}$$

Remark. Lemma 4.1 states that the ratios between $\tilde{k}_i^*, k_i^*, i \in \{1, 2\}$ are bounded by at most κ^2 . Recall that $\kappa = \text{polylog}(d)$.

Proof. The proof of Theorem 4.1 is a direct application of the non-increasing property of the Kolmogorov widths. We only prove for k_1^* and \tilde{k}_1^* , the same logic holds for k_2^* and \tilde{k}_2^* as well. By definition of k_1^* , we have

$$D_{k_1^*}(K) \leq \frac{(k_1^* + 1)^{\frac{1}{4}}}{\sqrt{N}}\sigma.$$

Recall from Corollary 3.5 and Lemma 2.1, we know (omit the arbitrarily small constant ϵ)

$$\tilde{D}_k(K) = \sqrt{\tilde{h}_S(X_k^\dagger)} \leq \sqrt{\kappa h_S(X_k^*)} \leq \sqrt{\kappa} D_k(K).$$

Consider $j = \kappa^2(k_1^* + 1) \geq k_1^* + 1$, by the non-increasing property of the Kolmogorov widths, we know

$$\tilde{D}_j(K) \leq \sqrt{\kappa} D_j(K) \leq \sqrt{\kappa} D_{k_1^*}(K) \leq \sqrt{\kappa} \frac{(k_1^* + 1)^{\frac{1}{4}}}{\sqrt{N}} \sigma = \frac{j^{\frac{1}{4}}}{\sqrt{N}} \sigma.$$

By definition, we know that $\tilde{k}_1^* \leq j - 1 = \kappa^2(k_1^* + 1) - 1$. $\tilde{k}_1^* \leq d$ is trivial. The proof is completed. \square

4.2 A Simple Case

To build intuition, we first consider a simplified, unconstrained version of (1.1) where $N = 1$, $\epsilon = 0$, and $K = \mathbb{R}^d$. Upon observing a single random vector $Y \in \mathbb{R}^d$ from $\mathcal{N}(\mu, \sigma^2 \mathbf{I}_d)$, we wish to test the following hypotheses:

$$\begin{aligned} H_0 : \mu &= 0, \\ H_1 : \|\mu\|_2 &\geq \rho. \end{aligned} \tag{4.2}$$

This formulation represents the canonical mean testing problem in the Gaussian sequence model, a subject exhaustively investigated in nonparametric statistics. Notably, for any fixed testing error level α , [Ingster and Suslina \[2003\]](#) established that the fundamental minimax separation rate scales as $\rho_{\text{critical}} \asymp d^{\frac{1}{4}} \sigma$. Formally, the following theorem holds:

Theorem 4.2 ([Ingster and Suslina \[2003\]](#)). *When $K = \mathbb{R}^d$, for the testing problem (4.2), denote $R(\rho) := \min_{T: \mathcal{X} \rightarrow \{0,1\}} \max_{\mu: \|\mu\|_2 \geq \rho} R(\mu, T)$ as the minimax risk function under the testing function T , where*

$$R(\mu, T) := \text{Type I error} + \text{Type II error} = \mathbb{P}_0(T(Y) = 1) + \mathbb{P}_\mu(T(Y) = 0).$$

Define $\rho_{\text{critical}} := d^{\frac{1}{4}} \sigma$, then

- (1), when $\rho \gtrsim \rho_{\text{critical}}$, we have $R(\rho) = o(1)$;
- (2), when $\rho \lesssim \rho_{\text{critical}}$, we have $R(\rho) = \Theta(1)$.

When we operate under the constrained setting, i.e., $K \neq \mathbb{R}^d$, intuitively, the non-trivial additional information on the constraint should yield a refined rate of ρ_{critical} compared to Theorem 4.2. [Li and Neykov \[2026\]](#) proved that for general QCO constraints, we have the following results

Theorem 4.3 ([Li and Neykov \[2026\]](#)). *When K is a QCO set, for the testing problem (1.1) with no corruption, define $\rho_{\text{critical}} := (k_1^*)^{\frac{1}{4}} \sigma$, then*

- (1), when $\rho \gtrsim \rho_{\text{critical}}$, we have $R(\rho) = o(1)$;
- (2), when $\rho \lesssim \rho_{\text{critical}}$, we have $R(\rho) = \Theta(1)$,

where k_1^* is the first optimal dimension, i.e., selecting $f(j, \mathcal{P}) = j^{\frac{1}{4}} \sigma$ in Definition 2.2.

Theorem 4.3 attains a sharper rate than Theorem 4.2 since $k_1^* \leq d$. The core intuition is to project the observation Y onto the lower-dimensional subspace $H_{P_{k_1^*}^*}$ spanned by the optimal projection $P_{k_1^*}^*$, thereby reformulating the testing problem within $H_{P_{k_1^*}^*}$. The resulting minimax separation rate in this reduced subspace is then linked to the original problem through the Kolmogorov widths $D_{k_1^*}(K)$. Unfortunately, this oracle-type approach is largely unimplementable, as it strictly requires the knowledge of k_1^* and $P_{k_1^*}^*$, and consequently $D_k(K)$, which is intractable in general as highlighted in Section 2.2. We will include more discussions in Section 4.3.1. Before that, let us first introduce the following scenario as a preliminary.

Consider a random vector $Y \in \mathbb{R}^d$ where $Y = \mu + \zeta$, $\mu \in \mathbb{R}^d$, $\zeta \sim \mathcal{N}(\mathbf{0}, \Sigma)$ with $\Sigma \in \mathbb{R}^{d \times d}$ is the covariance matrix with the constraint $\text{tr}(\Sigma) = k\sigma^2$ and $\mathbf{0} \preceq \Sigma \preceq \sigma^2 \cdot \mathbf{I}_d$, where k and σ are known. We attempt to obtain the minimax rate for ρ in the testing problem (4.2). Define the following testing function of χ^2 -type:

$$T_\chi(Y) = \mathbf{1}_E, E := \left\{ Y \mid \|Y\|_2^2 - k\sigma^2 \geq t \right\}, \quad (4.3)$$

where t is some appropriate threshold to be determined. χ^2 test is known to be minimax optimal for the testing problem (4.2) when $\Sigma = \mathbf{I}_d$. Specifically, we have the following results.

Theorem 4.4. *Given the assumptions above, for the problem (4.2), if $\rho \gtrsim k^{\frac{1}{4}}\sigma$, the test (4.3) achieves uniform small error less than α . Consequently, we have $\rho_{\text{critical}} \lesssim k^{\frac{1}{4}}\sigma$ in this scenario.*

Theorem 4.4 broadens the classical results by relaxing the requirement of an identity covariance matrix; in fact, it operates effectively even without full knowledge of Σ . Denote the d eigenvalues of Σ by $\{\lambda_1\sigma^2, \dots, \lambda_d\sigma^2\}$, sorted as $\lambda_1 \geq \dots \geq \lambda_d$, under the trace constraint $\sum_{i=1}^d \lambda_i = k$. For any fixed covariance matrix Σ , applying the standard Le Cam's method yields a minimax lower bound for the testing problem (4.2) scaling as $\rho_{\text{critical}} \gtrsim \left(\sum_{i=1}^d \lambda_i^2 \right)^{\frac{1}{4}} \sigma$. Imposing the constraints on λ_i for $i = 1, \dots, d$, it follows algebraically that $\frac{k}{d^{1/4}} \leq \left(\sum_{i=1}^d \lambda_i^2 \right)^{\frac{1}{4}} \leq k^{\frac{1}{4}}$. Therefore, when Σ is unknown and the minimax framework takes the supremum over all feasible Σ , the rate established in Theorem 4.4 is minimax optimal. Another consequence of this result is that the separation rate is purely characterized by the energy parameters of the covariance matrix (i.e., k and σ), entirely bypassing the original dimension d .

Theorem 4.4 can be viewed as a natural consequence of the celebrated Hanson–Wright inequality. The proof is deferred to Appendix B.3.

4.3 Robust Testing

In this section, we consider a more realistic scenario. For now, let us assume that we have N observations $\mathbf{Y} = \{Y_1, \dots, Y_N\}$ contaminated by the adversary \mathcal{C} from the **strong ϵ -contamination model**: $\mathbf{Y} = \mathcal{C}(\tilde{\mathbf{Y}})$. For such setting, Li and Neykov [2026] prove the following minimax lower bound when the constraint K is a QCO set.

Theorem 4.5 (Li and Neykov [2026]). *For the mean testing problem (1.1) with potential contamination from a strong adversary \mathcal{C} and a prior constraint on the mean μ induced by a QCO set K , the following condition is necessary to ensure the existence of a valid test whose Type I and Type II errors are both uniformly below a prescribed constant α*

$$\rho^2 \gtrsim \sigma^2 \max \left\{ \frac{\sqrt{k_1^*}}{N}, \epsilon \sqrt{\frac{k_2^*}{N}}, \epsilon^2 \right\}, \quad (4.4)$$

where N is the sample size, σ^2 is the variance of the noise, and ϵ is the fraction of contamination. k_1^* and k_2^* are the previously defined first and second optimal dimensions.

Theorem 4.5 demonstrates that the minimax lower bound is jointly determined by three components: the geometric properties of the set K , which are captured by the optimal dimensions k_1^* and k_2^* ; the characteristics of the original data, which are parameterized by N and σ ; and the corruption process, which is dictated by k_2^* and ϵ . The contamination fraction ϵ intricately connects all three terms in the lower bound. As ϵ increases from 0 to c_0 (assuming $c_0 = \Theta(1)$), each of the three terms in the max operator dominates the other two in turn, thereby forming three distinct phases of model behavior.

4.3.1 Polynomial-Time Algorithm

In this section, we construct a polynomial-time testing procedure centered around the surrogate matrix A_k^\dagger . Specifically, the core intuition is to project the observations \mathbf{Y} into some subspace $H_{P_k^*}$ of \mathbb{R}^d with intrinsic dimension $k \leq d$ induced by k -dimensional optimal projection P_k^* . We then conduct the deduction based on the subspace $H_{P_k^*}$ as if we transfer the testing problem to \mathbb{R}^k where there is no constraint. The resulting upper bound on $\|\mu\|_2$ depends on two parts. (1), the price we have to pay for the distance between the original μ and the projected $P_k^*\mu$, which is $\|\mu - P_k^*\mu\|_2$; (2), the fundamental requirement on $\|P_k^*\mu\|_2$ when we attempt to conduct the testing in $H_{P_k^*}$. The best upper bound possible is achieved by balancing the two parts above, and this is one of the motivations of the optimal dimension and projection. This approach depending on the balance of the two parts already appears in the recent work by Li and Neykov [2026]. However, in this work, as mentioned previously, the fundamental difference is that we do not assume the knowledge of the Kolmogorov widths, and consequently the exact optimal dimensions and projections are unknown, neither. We therefore replace the exact Kolmogorov k -width with the solution to its relaxed version in (2.3), and also replace the exact optimal dimensions k_1^*, k_2^* and projections P_1^*, P_2^* with $\tilde{k}_1^*, \tilde{k}_2^*$ defined in Definition 4.1 and $A_1^\dagger := (\mathbf{I}_d - X_{\tilde{k}_1^*}^\dagger)^{\frac{1}{2}}$, $A_2^\dagger := (\mathbf{I}_d - X_{\tilde{k}_2^*}^\dagger)^{\frac{1}{2}}$ where $X_{\tilde{k}_1^*}^\dagger, X_{\tilde{k}_2^*}^\dagger$ are the outputs of Algorithm 3 when setting $k = \tilde{k}_1^*$ and $k = \tilde{k}_2^*$, respectively. Another issue derived from such approximate projections is that the covariance matrix of the projected samples is no longer identity. (When we use the exact projections, the covariance matrix can be thought as the identity as we might apply some known appropriate rotation to the projected samples.) We are already well prepared for this scenario, as is what Theorem 4.4 serves.

We also note here that while prior literature has explored robust testing and estimation for the unconstrained version of (1.1), these classical procedures often suffer from restrictive assumptions — requiring either a vanishing contamination fraction ($\epsilon \rightarrow 0$) or suboptimal conditions scaling with N and d . Truly general and computationally tractable algorithms have only emerged recently. Notably, Diakonikolas et al. [2019] pioneered the first efficient robust estimation framework via algorithmic “filtering” in the Gaussian sequence model. For the robust testing counterpart, analogous filtering techniques were first successfully adapted by Narayanan [2022], Canonne et al. [2023], and later extended by Li and Neykov [2026]. This algorithmic framework to be leveraged is the basis of robust testing in the unconstrained setting (in $H_{A_1^\dagger}$ and $H_{A_2^\dagger}$). Below we review such framework. We start with the concept of ω -regularity, where a suit of weights is defined to adaptively “filter” the observations \mathbf{Y} .

Definition 4.2 (ω -regularity). Given a weight vector $\omega = (\omega_1, \dots, \omega_N)^\top$ and an integer $1 \leq k \leq d$, \mathbf{Y} is said to be $(\epsilon, \beta_1, \beta_2)$ -regular if for all subsets $S \subset [N]$ with $|S| \leq \epsilon N$, we have the following properties:

$$\begin{aligned}
\text{(i),} & \quad \left| \sum_{i \in S} \|Y_i\|_2^2 - |S|k \right| \leq c\beta_1, \\
\text{(ii),} & \quad \left| \left\| \sum_{i \in S} \sqrt{\omega_i} Y_i \right\|_2^2 - \|\omega_S\|_1 k \right| \leq c\beta_2, \text{ and} \\
\text{(iii),} & \quad \left| \left\langle \sum_{i \in S} \sqrt{\omega_i} Y_i, \sum_{j \in [N]} \sqrt{\omega_j} Y_j \right\rangle - \|\omega_S\|_1 k \right| \leq c\sqrt{N}\beta_1,
\end{aligned}$$

where $\omega_S \in \mathbb{R}^N$ is the restriction of ω on the set S .

Before proceeding, we note a slight abuse of notation: hereafter, Y_i and μ should be implicitly understood as their approximately projected counterparts, $A_k^\dagger Y_i$ and $A_k^\dagger \mu$. Crucially, because the surrogate matrices X_k^\dagger and A_k^\dagger are constructed independently of the observation sets $\tilde{\mathbf{Y}}$ and \mathbf{Y} , the mutual independence among the projected observations $A_k^\dagger \tilde{Y}_i$ and $A_k^\dagger \tilde{Y}_j$ ($i \neq j$) is preserved. To avoid notational clutter, we drop the explicit projection operator A_k^\dagger and simply write $\tilde{\mathbf{Y}}, \mathbf{Y}, Y_i$, and μ in the remaining text. We also assume $\sigma = 1$ in this section for the conciseness of the notation without loss of generality. When $\sigma \neq 1$, we only need to scale all the results by σ .

As a baseline, consider the uncontaminated setting where $Y_i = \tilde{Y}_i$ for all $i \in [N]$. Under this clean-data regime, we naturally adopt the uniform weighting scheme $\omega = \mathbf{1}_{[N]}$ across the entire dataset. By standard concentration inequalities, it can be shown that this uncorrupted empirical distribution inherently satisfies the $(\epsilon, \beta_1, \beta_2)$ -regularity condition, with β_1 and β_2 explicitly given by:

$$\begin{aligned}
\beta_1 &= \epsilon N \left(\sqrt{k} + \sqrt{N} \|\mu\|_2 \right) \sqrt{\ln \left(\frac{N}{\alpha} \right)} + \epsilon N \ln \left(\frac{N}{\alpha} \right) + \epsilon N \sqrt{N} \|\mu\|_2^2, \\
\beta_2 &= \epsilon N \sqrt{\epsilon N k \ln \left(\frac{1}{\epsilon} \right)} + (\epsilon N)^2 \ln \left(\frac{1}{\epsilon} \right) + \|\mu\|_2 (\epsilon N)^2 \sqrt{\ln \left(\frac{1}{\epsilon} \right)} + \|\mu\|_2^2 (\epsilon N)^2.
\end{aligned} \tag{4.5}$$

Consequently, the parameters β_1 and β_2 naturally motivate the formulation of a χ^2 -type test statistic. Nevertheless, such deterministic regularity is invariably violated in the contaminated setting by the malicious adversary \mathcal{C} . To actively enforce the required $(\epsilon, \beta_1, \beta_2)$ -regularity condition on the empirical distribution, we propose the following three-step algorithmic framework.

Algorithm 4: Pre-filtering.

```

1 Set  $\gamma_1 = c \left[ \sqrt{k \ln \left( \frac{N}{\alpha} \right)} + \ln \left( \frac{N}{\alpha} \right) \right]$ , count = 0, i = 0
2 while  $i < N$  do
3   if  $\left| \|Y_i\|_2^2 - k \right| > \gamma_1$  then
4     count = count + 1
5     if  $\text{count} > \epsilon N$  then
6       return None
7     Delete  $Y_i$  from  $\mathbf{Y}$ 
8   i = i + 1
9 return  $\mathbf{Y}$ 

```

Algorithm 4 is designed to explicitly enforce the first condition of ω -regularity. The threshold parameter γ_1 is carefully chosen based on the (sub)-Gaussian concentration of the uncorrupted data. Specifically, for clean samples, the concentration inequality $\left| \left\| \tilde{Y}_i \right\|_2^2 - k \right| \leq \gamma_1$ holds across all $i = 1, 2, \dots, N$ with probability at least $1 - \alpha$. Because the adversarial contamination budget is strictly limited to ϵ , observing strictly more than ϵN violations of this bound provides sufficient statistical evidence to confidently reject the null hypothesis and conclude H_1 . We formalize this pre-filtering logic in the following lemma.

Lemma 4.6. *After the execution of Algorithm 4, we have for any $i \in [N]$,*

$$\left| \left\| Y_i \right\|_2^2 - k^* \right| \leq \gamma_1.$$

Moreover, under H_0 , the algorithm will terminate without rejection with probability at least $1 - \alpha$.

For the second condition, we separate the classical scenario when $k < N$ and the high-dimensional scenario when $k \geq N$. In the following text, let $D(\omega) = \text{diag}\{\omega\} \in \mathbb{R}^{N \times N}$.

Classical scenario ($k < N$): select γ_2 as:

$$\gamma_2 := c \left[\sqrt{Nk} + \sqrt{N \ln \left(\frac{1}{\alpha} \right) + \ln \left(\frac{1}{\alpha} \right) + \epsilon N \ln \left(\frac{1}{\epsilon} \right)} \right], \quad (4.6)$$

where c is a sufficiently large but universal constant. Denote the value of ω in the t -th iteration as $\omega^{(t)}$. Define $\tau_i = \langle v, Y_i \rangle^2 \mathbf{1}_{\{\omega_i > 0\}}$, where v is the unit singular vector associated with $\lambda = \left\| \mathbf{Y}^\top D(\omega) \mathbf{Y} - N \mathbf{I}_k \right\|_2$. The update policy for ω is as

$$\omega_i^{(t+1)} = \begin{cases} \left(1 - \frac{\tau_i^{(t)}}{\tau_1^{(t)}} \right) \omega_i^{(t)} & \text{if } i \leq I, \\ \omega_i^{(t)} & \text{if } i > I. \end{cases} \quad (4.7)$$

See the detailed filtering process for the second condition of ω -regularity under classical scenario in Algorithm 5.

Algorithm 5 guarantees the second condition of ω -regularity under the classical scenario (i.e., $N > k$) within polynomial-time complexity. In particular, we have the following lemma.

Lemma 4.7. *The following facts hold for Algorithm 5:*

- (i), it terminates within finite rounds,
- (ii), under H_0 , it does not output rejection (*None*) with probability greater than $1 - \alpha$;
- (iii), once finished without rejection, for any $S \subset [N]$ with $|S| \leq \epsilon N$, we have $\left\| \sum_{i \in S} \sqrt{\omega_i} Y_i \right\|_2^2 \leq 2\gamma_2 \epsilon N$ with probability higher than $1 - \alpha$.

High-dimensional scenario ($k \geq N$): select γ_3 as:

$$\gamma_3 = c \left(\sqrt{Nk} + \sqrt{k \ln \left(\frac{1}{\alpha} \right) + \ln \left(\frac{1}{\alpha} \right) + \epsilon N \ln \left(\frac{1}{\epsilon} \right)} \right), \quad (4.8)$$

Algorithm 5: Sample filtering when $N > k$.

```

1 Set  $\gamma_2$  as (4.6), and  $\lambda = \|\mathbf{Y}^\top D(\omega)\mathbf{Y} - N\mathbf{I}_k\|_2$ . ( $\omega$  is initialized as  $\mathbf{1}$ .)
2 while  $\lambda \geq \gamma_2$  do
3   Set  $v$  to be the unit singular vector associated with  $\lambda$ 
4   Compute  $\tau_i = \langle v, Y_i \rangle^2 \mathbf{1}_{\{\omega_i > 0\}}$  for  $1 \leq i \leq N$ 
5   Sort  $\{\tau_i\}_{i=1}^N$  according to the decreasing order
6   Set  $I$  be the smallest index such that  $\sum_{i=1}^I \omega_i \geq 2\epsilon N$ 
7   Update  $w$  according to (4.7)
8   if  $\|\omega\|_1 < N(1 - 2\epsilon)$  then
9     return None
10  Set  $\lambda = \|\mathbf{Y}^\top D(\omega)\mathbf{Y} - N\mathbf{I}_k\|_2$ 
11 return  $\omega$ 

```

where again c is some sufficiently large but universal constant. Define $\tau_i = \frac{v_i^2}{\omega_i} \mathbf{1}_{\{\omega_i > 0\}}$ where v is the unit singular vector associated with $\lambda = \left\| \sqrt{D(\omega)} \mathbf{Y} \mathbf{Y}^\top \sqrt{D(\omega)} - kD(\omega) \right\|_2$. The update policy for ω in the t -th iteration is as:

$$\omega_i^{(t+1)} = \left(1 - \frac{\tau_i}{\max_i \tau_i} \right) \omega_i^{(t)}. \quad (4.9)$$

See the detailed filtering process for the second condition of ω -regularity under high-dimensional scenario in Algorithm 6.

Algorithm 6: Sample filtering when $N \leq k$.

```

1 Set  $\gamma_3$  as (4.8), and  $\lambda = \left\| \sqrt{D(\omega)} \mathbf{Y} \mathbf{Y}^\top \sqrt{D(\omega)} - kD(\omega) \right\|_2$ . ( $\omega$  is initialized as  $\mathbf{1}$ .)
2 while  $\lambda \geq \gamma_3$  do
3   Set  $v$  to be the unit singular vector associated with  $\lambda$ 
4   Compute  $\tau_i = \frac{v_i^2}{\omega_i} \mathbf{1}_{\{\omega_i > 0\}}$ 
5   Update  $w$  according to (4.9)
6   if  $\|\omega\|_1 < N(1 - 6\epsilon)$  then
7     return None
8   Set  $\lambda = \left\| \sqrt{D(\omega)} \mathbf{Y} \mathbf{Y}^\top \sqrt{D(\omega)} - kD(\omega) \right\|_2$ 
9 return  $\omega$ 

```

similar to Algorithm 5, the following lemma establishes that Algorithm 6 indeed enforces the second condition of ω -regularity under the high-dimensional scenario (i.e., $N \leq k$).

Lemma 4.8. *The following facts hold for Algorithm 6:*

- (i), it terminates within finite rounds;
- (ii), under H_0 , it does not output rejection (*None*) with probability greater than $1 - \alpha$;

(iii), once finished without rejection, for any $S \subset [N]$ with $|S| \leq \epsilon N$, we have

$$\left| \left\| \sum_{i \in S} \sqrt{\omega_i} Y_i \right\|_2^2 - k \|\omega_S\|_1 \right| \leq c\epsilon N \gamma_3$$

with probability greater than $1 - \alpha$.

We note here that the techniques involving the interaction with the maximal eigenvalue and the corresponding eigenvector of either $\mathbf{Y}^\top D(\omega) \mathbf{Y} - N \mathbf{I}_k$ or $\sqrt{D(\omega)} \mathbf{Y} \mathbf{Y}^\top \sqrt{D(\omega)} - k D(\omega)$ in Algorithm 5 and 6 first appeared in Klivans et al. [2009] and was later adapted to the robust statistics by Diakonikolas et al. [2019, 2017].

Finally, the third condition of ω -regularity is guaranteed by the weight filtering algorithm, presented below.

Algorithm 7: Weight filtering.

- 1 Compute $\tau_i = \left| \left\langle \sqrt{\omega_i} Y_i, \sum_{j=1}^N \sqrt{\omega_j} Y_j \right\rangle - k \omega_i \right| \mathbf{1}_{\{\omega_i > 0\}}, 1 \leq i \leq N$
 - 2 Sort τ_i by decreasing order and find the indices $\{i_1, i_2, \dots, i_{\epsilon N}\}$ corresponding to the first ϵN maximal τ_i
 - 3 Set $\omega_{i_1}, \omega_{i_2}, \dots, \omega_{i_{\epsilon N}}$ to zero
 - 4 **return** ω
-

The input weights ω in Algorithm 7 is the output of either Algorithm 5 or Algorithm 6 according to the relationship between N and k . Algorithm 7 enforces the third condition of ω -regularity, as established in the following lemma.

Lemma 4.9. *For the weights ω outputted by Algorithm 7 and any $S \subset [N]$ with $|S| \leq \epsilon N$, we have*

$$\left| \sum_{i \in S} \sum_{j=1}^N \sqrt{\omega_i \omega_j} \langle Y_i, Y_j \rangle - k \|\omega_S\|_1 \right| \leq c(\sqrt{N} \beta_1 + \beta_2 + \epsilon N \gamma),$$

where β_1 and β_2 are defined in (4.5), γ is equal to either γ_2 or γ_3 accordingly.

As outlined above, the sequential execution of Algorithms 4, 5 (or 6), and 7 allows us to adaptively refine the data weights ω in response to the contaminated observations \mathbf{Y} . This adaptive filtering strictly enforces the ω -regularity condition with respect to the predefined parameters β_1 and β_2 . Crucially, endowed with this enforced regularity, the adversarially corrupted dataset \mathbf{Y} is compelled to statistically mimic the concentration behavior of the uncorrupted samples $\check{\mathbf{Y}}$. We formalize this fundamental property in the following lemma.

Lemma 4.10. *Let ω be the output weights after executing Algorithm 4, 5 (or 6), and 7 without rejection (otherwise, we simply reject H_0), and β_1, β_2 selected as in (4.5). Then, with high probability we have*

$$\left| \left\| \sum_{i=1}^N \sqrt{\omega_i} Y_i \right\|_2^2 - k \|\omega\|_1 - N^2 \|\mu\|_2^2 \right| \leq c \left(\sqrt{N} \beta_1 + \beta_2 + \epsilon N \gamma + \left(N^{\frac{3}{2}} \|\mu\|_2 + \sqrt{Nk} \right) \sqrt{\ln \left(\frac{1}{\alpha} \right) + \ln \left(\frac{1}{\alpha} \right)} \right). \quad (4.10)$$

On the LHS of the preceding inequality, the term $k \|\omega\|_1$ is a known quantity upon completion of the filtering algorithms. To guarantee a strictly detectable separation between H_0 and H_1 , we require the signal term $N^2 \|\mu\|_2^2$ to asymptotically dominate all residual terms on the RHS. This bounding strategy closely mirrors the approach developed for the [theoretical algorithm](#). Through straightforward algebraic manipulations, we formally establish the following equivalence.

Theorem 4.11. *Assume the same conditions as in Theorem A.1, for any $A_k^\dagger, 1 \leq k \leq d$, if*

$$\|A_k^\dagger \mu\|_2^2 \gtrsim \sigma^2 \max \left\{ \frac{\epsilon \ln \left(\frac{N}{\alpha} \right)}{\sqrt{N}}, \epsilon^2 \ln \left(\frac{N}{\alpha} \right), \sqrt{\frac{\epsilon^2 k \ln \left(\frac{N}{\alpha} \right)}{N}}, \frac{\sqrt{k} \ln \left(\frac{1}{\alpha} \right)}{N} \right\} := P_{\text{raw}}^2(\epsilon, K, N, \sigma^2, k), \quad (4.11)$$

we are able to test (1.1) with Type I and Type II errors uniformly less than α .

Since

$$\|A_k^\dagger \mu\|_2 = \sqrt{\mu^\top (A_k^\dagger)^2 \mu} = \sqrt{\mu^\top (\mathbf{I}_d - X_k^\dagger) \mu} = \sqrt{\|\mu\|_2^2 - \mu^\top X_k^\dagger \mu} \geq \sqrt{\|\mu\|_2^2 - \tilde{h}_S(X_k^\dagger)},$$

we come to the following sufficient condition.

Corollary 4.12. *Assume the same assumptions as in Theorem A.2. For any $1 \leq k \leq d$, if*

$$\|\mu\|_2^2 \gtrsim \sigma^2 \left[\tilde{h}_S(X_k^\dagger) + P_{\text{raw}}^2(\epsilon, K, N, \sigma^2, k) \right], \quad (4.12)$$

where $\tilde{k}_1^*, \tilde{k}_2^*$ are defined as in Definition 4.1, then we are able to test the problem (1.1) with uniformly small Type I and Type II errors. Moreover, the testing procedures can be completed within polynomial-logarithmic time of the parameters $(N, d, \frac{1}{\epsilon}, \frac{1}{\alpha})$ for any type-2, exactly 2-convex constraint K .

Since (4.12) holds for any $1 \leq k \leq d$, we are free to optimize over k to obtain the sharpest bound:

$$\begin{aligned} \|\mu\|_2^2 &\gtrsim \min_{1 \leq k \leq d} \sigma^2 \left[\tilde{h}_S(X_k^\dagger) + P_{\text{raw}}^2(\epsilon, K, N, \sigma^2, k) \right] \\ &= \max \left\{ \frac{\epsilon \ln \left(\frac{N}{\alpha} \right)}{\sqrt{N}}, \epsilon^2 \ln \left(\frac{N}{\alpha} \right), \sqrt{\frac{\epsilon^2 \min \{ \tilde{k}_1^*, \tilde{k}_2^* \} \ln \left(\frac{N}{\alpha} \right)}{N}}, \frac{\sqrt{\min \{ \tilde{k}_1^*, \tilde{k}_2^* \} \ln \left(\frac{1}{\alpha} \right)}}{N} \right\}. \end{aligned} \quad (4.13)$$

Proof. Since we are allowed to optimize over k , we select \tilde{k}_1^* . From the definition of \tilde{k}_1^* , we know that $\tilde{h}_S(X_{\tilde{k}_1^*}^\dagger) = \tilde{D}_{\tilde{k}_1^*}^2(K) \leq \frac{\sqrt{\tilde{k}_1^*}}{N} \sigma^2$, which is dominated by the term $\frac{\sqrt{\tilde{k}_1^*} \ln(1/\alpha)}{N} \sigma^2$ in the maximal operator. Therefore, we can safely remove $\tilde{h}_S(X_{\tilde{k}_1^*}^\dagger)$ in this case. The same logic also holds for \tilde{k}_2^* . The obtained two upper bounds are both valid for the problem (1.1), hence we select the smaller one, which leads to (4.13). The proof is completed. \square

When (4.13) is satisfied, let $\tilde{k}^* = \min \{ \tilde{k}_1^*, \tilde{k}_2^* \}$ and $A_{\tilde{k}^*}^\dagger$ be the corresponding approximate optimal projection. To distinguish between H_0 and H_1 , we focus on the following weighted χ^2 -type testing event

$$\left\{ \left| \left\| \sum_{i=1}^N \sqrt{\omega_i} A_{\tilde{k}^*}^\dagger Y_i \right\|_2^2 - \tilde{k}^* \|\omega\|_1 \right| \geq c_2 N^2 P_{\text{raw}}(\epsilon, K, N, \sigma^2, \tilde{k}^*) \right\}.$$

We reject H_0 when the event is true and vice versa.

We benchmark our upper bound (4.12) against the theoretical guarantees established in Li and Neykov [2026]. In an idealized oracle setting where the exact Kolmogorov widths and their corresponding optimal projections are known a priori, the approximate dimensions \tilde{k}_1^* and \tilde{k}_2^* can be perfectly substituted with the exact optimal dimensions k_1^* and k_2^* . This substitution yields a sharper upper bound that matches the minimax lower bound (4.4) up to logarithmic factors (specifically, $\ln N$, $\ln(\frac{1}{\epsilon})$, and $\ln(\frac{1}{\alpha})$) under QCO constraints. Remarkably, even in our practical setting where these Kolmogorov widths are inherently unknown, Lemma 4.1 ensures that the fully empirical upper bound (4.12) incurs merely an additional $\text{polylog}(d)$ multiplicative factor.

The proofs of Lemmas 4.6, 4.7, 4.8, 4.9, and 4.10, which the rationale of the filtering method, share substantial overlap with the analysis in Li and Neykov [2026]. We therefore refer the interested reader to Appendix C.6 of Li and Neykov [2026] for the full derivations. Nevertheless, our framework requires a non-trivial adaptation to accommodate the non-identity covariance structure of the approximately projected observations $A_k^\dagger \tilde{Y}_i$. (Had A_k^\dagger been an exact orthogonal projection, this covariance would trivially reduce to the identity within a rotated subspace). To rigorously circumvent this difficulty, we establish Lemma B.4 and Lemma B.3 to govern the ℓ_2 -operator norms of the sample covariance and Gram matrices, respectively. The proofs of these two lemmas are deferred to Appendix B.5. The pseudo code of the polynomial-time algorithm is shown in the following.

5 Discussion

In this paper, we establish a polynomial-time algorithmic framework for robust signal detection, resolving the testing problem (1.1) with rigorous control over Type I and Type II errors. Our procedure accommodates any balanced, type-2, and 2-convex constraint set K , thereby generalizing the QCO constraints analyzed in recent literature. The computational tractability of our framework hinges crucially on the SDP relaxation of the Kolmogorov widths (2.3), complemented by the approximation schemes detailed in Bhattiprolu et al. [2021] and Section 3. Remarkably, our empirical minimax upper bound derived in Section 4 is near-optimal, incurring merely a $\text{polylog}(d)$ multiplicative penalty compared to the theoretical minimax rates achieved by Li and Neykov [2026] in the oracle QCO setting.

Looking forward, several compelling directions warrant further investigation. A natural extension is the robust ℓ_p -norm testing problem, obtained by replacing the ℓ_2 -norm objective in (1.1) with an ℓ_p -norm. For the regime $1 \leq p < 2$, Li and Neykov [2026] established that the minimax lower bound and corresponding upper bound successfully transfer under p -convex orthosymmetric (PCO) constraints. Nevertheless, given the inherent computational intractability of evaluating ℓ_p -Kolmogorov widths, the existence of an efficient algorithmic counterpart remains unknown. The scenario where $p > 2$ is fundamentally more challenging. Although the χ^p -test achieves minimax optimality in the unconstrained, corruption-free model Ingster and Suslina [2003], the exact behavior of the minimax separation rates under the interplay of structural constraints and adversarial corruption is entirely open.

Another critical avenue is testing against a shifted null hypothesis, $H_0 : \mu = \mu_0$ ($\mu_0 \neq 0$). For clean data, Wei and Wainwright [2020] proved that the minimax lower bound under ellipsoidal constraints is governed by localized Kolmogorov widths. A broader formulation extends this to composite null hypotheses, where $H_0 : \|\mu\|_2 \leq \rho_0$ is tested against $H_1 : \|\mu\|_2 \geq \rho_1$. Relevant unconstrained and uncorrupted baselines can be found in Kania et al. [2025], yet the robust constrained version stands

Algorithm 8: unconditional polynomial-time algorithm.

```

1 import  $\mathbf{Y}, N, K, \sigma = 1$ , Subalgorithm 4, 5, 6, 7

2  $\mathbf{Y} = \mathbf{Y}/\sigma$ 
3 Compute  $\tilde{k}_1^*, A_1^\dagger, \tilde{k}_2^*, A_2^\dagger$  from  $N, K, \sigma = 1, \epsilon$  as in Definition 4.1 and Algorithm 3
4 if  $\tilde{k}_1^* \leq \tilde{k}_2^*$  then
5    $\tilde{k}^* = \tilde{k}_1^*, A^\dagger = A_1^\dagger$ 
6 else
7    $\tilde{k}^* = \tilde{k}_2^*, A^\dagger = A_2^\dagger$ 
8  $\mathbf{Y} = \mathbf{Y}A^\dagger$ 
9 Execute Algorithm 4, and record the returning value  $R_1$ 
10 if  $R_1$  is None then
11   Reject  $H_0$ 
12   return
13 else
14    $\mathbf{Y} = R_1$ 
15 Set  $\omega = \mathbf{1}_{[\tilde{k}^]}$ 
16 if  $N > \tilde{k}^*$  then
17   Execute Algorithm 5, and record the returning value  $R_2$ 
18 else
19   Execute Algorithm 6, and record the returning value  $R_2$ 
20 if  $R_2$  is None then
21   Reject  $H_0$ 
22   return
23 else
24    $\omega = R_2$ 
25 Execute Algorithm 7, and record the returning value  $R_3$ 
26  $\omega = R_3$ 
27 if  $\left| \left\| \sum_{i=1}^N \sqrt{\omega_i} Y_i \right\|_2^2 - \tilde{k}^* \|\omega\|_1 \right| \geq c_2 N^2 P(\epsilon, \tilde{k}^*, N, 1)$  then
28   Reject  $H_0$ 
29 else
30   Accept  $H_0$ 
31 return

```

as an open problem.

Finally, characterizing robust testing over more exotic geometries — such as general symmetric, arbitrarily convex, or non-convex sets—remains a highly non-trivial frontier. Fully resolving this demands bridging two fundamental questions: (i) identifying the absolute information-theoretic limits (i.e., establishing minimax lower bounds and the existence of statistical upper bounds without computational constraints), and (ii) understanding the potential statistical-computational gap by designing efficient algorithms, recognizing that polynomial-time tractability may inherently require sub-optimal statistical conditions.

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A Theoretical Algorithm

To provide a complete theoretical picture of the upper bounds presented in Section 4, this appendix details the fully information-theoretic testing procedure. Although this theoretical algorithm operates in exponential time—necessitating an exhaustive search over all possible subsets—it achieves a sharper statistical upper bound than the polynomial-time filtering developed in Section 4. At its core, this procedure is anchored by a classical χ^2 -type test combined with the high-probability existence of a *consistent subset* among the observations \mathbf{Y} . We formally introduce the definition of this consistency property below.

Definition A.1 (Consistent subset). A subset $S \subset \mathbf{Y} = \{Y_1, \dots, Y_N\}$ is called a consistent subset regarding the contamination fraction ϵ and a test $\phi : 2^S \rightarrow \{0, 1\}$ if

- (i), $|S| \geq (1 - \epsilon)N$,
- (ii), $\phi(S') = \phi(S)$ for any $S' \subset S$ with $|S'| \geq (1 - 2\epsilon)N$.

The intuition driving the [consistent subset](#) formulation is conceptually straightforward. Given the observations \mathbf{Y} from adversarial contamination up to ϵ fraction, it is guaranteed to contain a subset of purely uncorrupted samples (specifically, $\mathbf{Y}_{[N] \setminus C}$, even though the corruption indices C is unknown). Suppose \mathbf{Y}_{S_0} forms a valid consistent subset for some index set $S_0 \subset [N]$. It then becomes possible to link the test statistic evaluated on \mathbf{Y}_{S_0} back to the unobservable original samples $\tilde{\mathbf{Y}}$. This exact correspondence is validated by demonstrating that, under appropriate conditions on ρ , $\tilde{\mathbf{Y}}$ inherently constitutes a consistent subset with high probability. As a result, evaluating the proposed χ^2 -type test ϕ_e yields the crucial equivalence $\phi_e(\mathbf{Y}_{S_0}) = \phi_e(\mathbf{Y}_{S_0 \cap [N] \setminus C}) = \phi_e(\mathbf{Y}_{[N] \setminus C}) = \phi_e(\tilde{\mathbf{Y}})$. In essence, extracting and testing on S_0 allows us to completely bypass the adversarial corruption and perfectly replicate the oracle test on the authentic data $\tilde{\mathbf{Y}}$.

Recall the surrogate projection matrix $A_k^\dagger := \left(\mathbf{I}_d - X_k^\dagger\right)^{\frac{1}{2}}$, and consider the following χ^2 -type rejection region:

$$E_e := \left\{ \left\| A_k^\dagger \sum_{i \in S} Y_i \right\|_2^2 - k |S| \sigma^2 \geq c |S|^2 E_{\text{raw}}^2(\epsilon, K, N, \sigma^2, k) \right\}. \quad (\text{A.1})$$

Here, S denotes a subset of the data (with a slight abuse of notation, S interchangeably refers to the index set within $[N]$ and the corresponding observations in \mathbf{Y}), $c > 0$ is a carefully chosen absolute constant, and the threshold $E_{\text{raw}}^2(\epsilon, K, N, \sigma^2, k)$ is defined as:

$$E_{\text{raw}}^2(\epsilon, K, N, \sigma^2, k) := \sigma^2 \max \left\{ \frac{\sqrt{k}}{N}, \epsilon^2 \ln \left(\frac{1}{\epsilon} \right), \sqrt{\frac{\epsilon^2 \ln \left(\frac{1}{\epsilon} \right) k}{N}} \right\}. \quad (\text{A.2})$$

Evaluated on any subset S , the theoretical test ϕ_e rejects H_0 if and only if the event E_e holds. Given this explicit formulation of the test ϕ_e , the subsequent lemma establishes that the underlying clean dataset $\tilde{\mathbf{Y}}$ naturally forms a consistent subset with high probability.

Theorem A.1. *For any $1 \leq k \leq d$, if $\left\| A_k^\dagger \mu \right\|_2 \gtrsim E_{\text{raw}}(\epsilon, K, N, \sigma^2, k)$, then with probability $1 - o(1)$, $\tilde{\mathbf{Y}} A_k^\dagger$ is consistent with respect to ϕ_e defined above. Moreover, for the original samples $\tilde{\mathbf{Y}}$, ϕ_e achieves $o(1)$ Type I and Type II errors uniformly over $\mu \in K$.*

The framework of the proof of Theorem A.1 is essentially based on the proof of Theorem 3.8 of Li and Neykov [2026], which is deferred to Appendix B.4.

Following the discussions above, we know that when the selected μ and A_k^\dagger satisfy the condition $\|A_k^\dagger \mu\|_2 \gtrsim E_{\text{raw}}(\epsilon, K, N, \sigma^2, k)$, ϕ_e provides us an approach to distinguish between H_0 or H_1 with uniformly small Type I and Type II errors despite the corruptions. similar to the argument for the polynomial-time algorithm, we have $\|A_k^\dagger \mu\|_2 \geq \sqrt{\|\mu\|_2^2 - \tilde{h}_S(X_k^\dagger)}$. Therefore, it suffices that

$$\|\mu\|_2^2 \gtrsim \tilde{h}_S(X_k^\dagger) + E_{\text{raw}}^2(\epsilon, K, N, \sigma^2, k).$$

Notice that the LHS is independent of k . Consequently, we can optimize over k on the RHS for the best bound possible. This leads us to the main upper bound for the theoretical algorithm.

Corollary A.2 (Upper bound of the theoretical algorithm). *For the testing problem (1.1) with potential corruptions from a strong ϵ -contamination adversary \mathcal{C} and the prior constraint K on the mean μ , if*

$$\|\mu\|_2 \gtrsim E^2(\epsilon, K, N, \sigma^2) := \sigma^2 \max \left\{ \frac{\sqrt{\min \{ \tilde{k}_1^*, \tilde{k}_2^* \}}}{N}, \epsilon^2 \ln \left(\frac{1}{\epsilon} \right), \sqrt{\frac{\epsilon^2 \ln \left(\frac{1}{\epsilon} \right) \min \{ \tilde{k}_1^*, \tilde{k}_2^* \}}{N}} \right\}, \quad (\text{A.3})$$

then we are able to distinguish between H_0 and H_1 with error probability $o(1)$. Note that we do not require the computability of the Kolmogorov widths for the constraint K .

Proof. Since we are allowed to optimize over k , we select \tilde{k}_1^* . From the definition of \tilde{k}_1^* , we know that $\tilde{h}_S(X_{\tilde{k}_1^*}^\dagger) = \tilde{D}_{\tilde{k}_1^*}^2(K) \leq \frac{\sqrt{\tilde{k}_1^*}}{N} \sigma^2$, which is one of the terms in the maximal operator. Therefore, we can safely remove $\tilde{h}_S(X_k^\dagger)$ when $k = \tilde{k}_1^*$. The same logic also holds for \tilde{k}_2^* . The obtained two upper bounds are both valid for the problem (4.2), hence we select the smaller one, which leads to (A.3). The proof is completed. \square

Similar to the comparison in the efficient algorithm, $\tilde{k}_1^*, \tilde{k}_2^*$ can be substituted with k_1^* and k_2^* if the Kolmogorov widths are prior knowledge. From Lemma 4.1, we know such replacement is $\mathcal{O}(\kappa^2)$ sub-optimal.

However, as mentioned, one drawback for such testing procedure is that it requires checking the results of ϕ_e on all subsets \mathbf{Y}_S with $|\mathbf{Y}_S| \geq (1 - 2\epsilon)N$, whose amount can be as large as $\left(\frac{e}{2\epsilon}\right)^{2\epsilon N}$ (see Lemma B.2) — a number grows exponentially with N . Hence, the theoretical algorithm is computationally intractable when N is even moderately large and only remains theoretically valuable for the existence of the upper bound.

B Deferred Proofs

B.1 Proof of Lemma 2.4

Proof. Define $\tilde{K}^2 = \{\tilde{x} \mid \tilde{x} \in \mathbb{R}^d, \tilde{x} = \eta \odot x, x \in K^2\}$ as the symmetrized squared set of K , where \odot means entrywise multiplication. Since K is orthosymmetric and quadratically convex, we know that

\tilde{K}^2 is orthosymmetric and convex. Note that under such condition, the Minkowski gauges $\rho_K(\cdot)$ and $\rho_{\tilde{K}^2}(\cdot)$ is well known to be a norm, and satisfy the connection that $\rho_K^2(x) = \rho_{\tilde{K}^2}(x^2), \forall x \in K$. Here again x^2 is the entrywise square.

For any $m \in \mathbb{N}^+$, consider $\sum_{i=1}^m g_i x_i$. Its j -th coordinate is $\left(\sum_{i=1}^m g_i x_i\right)_j = \sum_{i=1}^m g_i x_{ij} \sim \mathcal{N}\left(0, \sum_{i=1}^m x_{ij}^2\right)$.

Therefore, we have $\left(\sum_{i=1}^m g_i x_i\right)^2 \sim \left(\sum_{i=1}^m x_{ij}^2\right) \tilde{g}_j^2$, where $\tilde{g}_j^2 \sim \chi_1^2$. Since \tilde{K}^2 is convex and contains the origin, we know that it is star-shape, and consequently for any $r, s \in \mathbb{R}^d$, if $|r_i| \leq |s_i|, 1 \leq i \leq d$, we have $\rho_{\tilde{K}^2}(r) \leq \rho_{\tilde{K}^2}(s)$. We have

$$\begin{aligned} \mathbb{E} \left[\rho_K^2 \left(\sum_{i=1}^m \epsilon_i x_i \right) \right] &= \mathbb{E} \left[\rho_{\tilde{K}^2} \left(\sum_{i=1}^m \epsilon_i x_i \right)^2 \right] \stackrel{(i)}{\leq} \frac{\pi}{2} \mathbb{E} \left[\rho_{\tilde{K}^2} \left(\sum_{i=1}^m g_i x_i \right)^2 \right] \\ &\leq \frac{\pi}{2} \rho_{\tilde{K}^2} \left(\sum_{i=1}^m x_i^2 \right) \mathbb{E} \max_j \tilde{g}_j^2 \stackrel{(ii)}{\leq} c \frac{\pi}{2} \ln d \sum_{i=1}^m \rho_{\tilde{K}^2}(x_i^2) \\ &= c \frac{\pi}{2} \ln d \sum_{i=1}^m \rho_K^2(x_i^2). \end{aligned}$$

Here (i) is from Lemma 4.5 of [Ledoux and Talagrand \[2013\]](#), and (ii) is from the standard maximum inequality for χ_1^2 random variables. Therefore, the type-2 constant for any QCO set is $T_2(K) \asymp \sqrt{\ln d}$. For the exact 2-convexity, we have:

$$\rho_K^2 \left(\sum_{i=1}^m x_i^2 \right)^{\frac{1}{2}} = \rho_{\tilde{K}^2} \left(\sum_{i=1}^m x_i^2 \right) \leq \sum_{i=1}^m \rho_{\tilde{K}^2}(x_i^2) = \sum_{i=1}^m \rho_K^2(x_i).$$

Taking the square root on both sides proves the exact 2-convexity. The proof is completed. \square

B.2 Proof of Theorem 3.4

Proof. The proof is derived from Lemma 3.1 and Corollary 3.3. First, the optimization problem (2.3) is indeed a convex problem, where the object function h_S is convex, and all the constraints are linear. Second, the feasible region for the problem is bounded. This can be verified by calculating the distance between any two feasible matrices X_1, X_2 : $\|X_1 - X_2\|_F^2 \leq 2\|X_1\|_F^2 + 2\|X_2\|_F^2$. Assuming the singular value decomposition of X_i is $X_i = U_i \Lambda_i U_i^\top$, where $i \in \{1, 2\}$, $\Lambda_i = \text{diag}\{\lambda_{i,1}, \dots, \lambda_{i,d}\}$ and $\lambda_{i,j}$ is the j -th singular value of X_i , then we have $\|X_i\|_F^2 = \text{tr}(X_i^2) = \text{tr}(U_i \Lambda_i^2 U_i^\top) = \sum_{j=1}^d \lambda_{i,j}^2 \stackrel{(i)}{\leq} \sum_{j=1}^d \lambda_{i,j} = d - k$,

where (i) is from the condition that $\mathbf{0} \preceq X \preceq \mathbf{I}_d$. Therefore, for any initial feasible point $X^{(0)}$ (for example, $X^{(0)} = \text{diag}\{\underbrace{1, \dots, 1}_{d-k}, 0, \dots, 0\}$), we can set $E^{(0)} = \{X \mid \|X - X^{(0)}\|_F \leq 2\sqrt{d-k}\}$, and

$E^{(0)}$ is guaranteed to contain the minimizer X^* .

Consider the ϵ -tolerance set

$$E_\epsilon := \left\{ X \mid X \in \mathbb{R}^{d \times d}, h_S(X) - h_S(X^*) \leq \epsilon \right\}.$$

The constraints of X can be expressed as $\{X \mid AX = B\}$, where A, B can be determined from the trace constraint $\sum_{i=1}^d X_{ii} = k$ and the symmetry constraint $X_{ij} = X_{ji}, \forall i \neq j$. The rank of the equality constraints is $1 + \frac{d(d-1)}{2}$, and therefore the dimension of $\mathcal{N}(A) = d^2 - \left(1 + \frac{d(d-1)}{2}\right) = \frac{(d-1)(d+2)}{2}$. It is also not hard to check that for any $\epsilon > 0$, the intersection set $E_\epsilon \cap \{X \mid AX = B\}$ has positive volume when considering in the $\mathbb{R}^{\frac{(d-1)(d+2)}{2}}$ -dimensional manifold determined by A, B in $\mathbb{R}^{d \times d}$. Finally, in the k -th iteration:

- If $X^{(k)}$ is not feasible, then any point in $E^{(k+1)\mathbb{C}} \cap E^{(k)}$ is not feasible as well and will be excluded in the next iteration.
- If $X^{(k)}$ is feasible, then for any point X in $E^{(k+1)\mathbb{C}} \cap E^{(k)} \cap \{X \mid AX = B\}$, from Lemma 3.1 we know that $h_S(X) \geq \frac{1}{\kappa} h_S(X^{(k)}) + \mathcal{O}(X^{(k)})^\top (X - X^{(k)}) \mathcal{O}(X^{(k)}) \geq \frac{1}{\kappa} h_S(X^{(k)})$.

From Corollary 3.3, we know that the volume of $E^{(k)} \cap \{X \mid AX = B\}$ decreases exponentially with the factor $\frac{d^{2\bar{d}}}{(d^2+1)^{(d+1)/2} \cdot (d^2-1)^{(d-1)/2}} < 1$. Therefore, within $M := c' \bar{d}^2 \cdot \ln\left(\frac{1}{\epsilon}\right)$ iterations, we have $\text{Vol}(E^{(M)} \cap \{X \mid AX = B\}) \leq \text{Vol}(E_\epsilon \cap \{X \mid AX = B\})$, which means that there exists a feasible point, denoted by $X_e \in E_\epsilon$, is excluded along the iterations. Assume that it is excluded in i -th iteration. Then from the analysis above, we know that

$$h_S(X^*) \leq h_S(X_{M\text{-best}}) \leq h_S(X^{(i)}) \leq \kappa h_S(X_e) \leq \kappa (h_S(X^*) + \epsilon).$$

The proof is completed. □

B.3 Proof of Theorem 4.4

We start with the following simple fact.

Fact. For a hypothesis testing problem $H_0 : \theta \in \Theta_0, H_1 : \theta \in \Theta_1$ and a test statistic $M(\mathbf{Y})$. If there exists a threshold t independent with θ such that

$$\begin{aligned} \sup_{\theta \in \Theta_0} \mathbb{P}_\theta(M(\mathbf{Y}) \leq t) &\geq 1 - \alpha, \\ \sup_{\theta \in \Theta_1} \mathbb{P}_\theta(M(\mathbf{Y}) > t) &\geq 1 - \alpha, \end{aligned}$$

then the test $\mathbf{1}_{\{M(\mathbf{Y}) \leq t\}}$ is valid with uniform Type I and Type II errors less than α . As an application, if for a quantity v_θ depending on θ such that $\mathbb{P}_\theta(|M(\mathbf{Y}) - \mathbb{E}_\theta(M(\mathbf{Y}))| \leq v_\theta) \geq 1 - \alpha$ holds uniformly for $\theta \in \Theta_0 \cup \Theta_1$ and $\sup_{\theta \in \Theta_0} \mathbb{E}_\theta(M(\mathbf{Y})) + v_\theta \leq \inf_{\theta \in \Theta_1} \mathbb{E}_\theta(M(\mathbf{Y})) - v_\theta$ (or $\inf_{\theta \in \Theta_0} \mathbb{E}_\theta(M(\mathbf{Y})) - v_\theta \geq \sup_{\theta \in \Theta_1} \mathbb{E}_\theta(M(\mathbf{Y})) + v_\theta$), then there exists a valid testing achieving Type I and Type II errors less than α uniformly for $\theta \in \Theta_0 \cup \Theta_1$.

Before proceeding, we also briefly recall the celebrated Hanson–Wright inequality. This fundamental result provides sharp tail bounds for quadratic forms of sub-Gaussian random vectors, making it an indispensable tool for analyzing Gaussian observations with non-identity covariance structures. The original inequality was established by Hanson and Wright [1971], while a modern and streamlined proof can be found in Rudelson and Vershynin [2013].

Lemma B.1 (The Hanson-Wright inequality). *Assume that $\mathbf{X} = (X_1, \dots, X_n) \in \mathbb{R}^n$ is a random vector with independent, zero-mean and sub-Gaussian coordinates. Let A be an $n \times n$ matrix and $K = \max_i \|X_i\|_{\psi_2}$. Then for any $t \geq 0$, we have*

$$\mathbb{P}\left(|X^\top AX - \mathbb{E}X^\top AX| > t\right) \leq 2 \exp\left\{-c \min\left\{\frac{t^2}{K^4 \|A\|_F^2}, \frac{t}{K^2 \|A\|_2}\right\}\right\}, \quad (\text{B.1})$$

where $\|\cdot\|_F$ is the Frobenius norm of a matrix, and $\|\cdot\|_2$ is the ℓ_2 operator norm of a matrix.

We are now ready for the proof of Theorem 4.4.

Proof. We analyze the term $\|Y\|_2^2$ directly. Since $Y = \mu + \zeta$, we have

$$\begin{aligned} \|Y\|_2^2 &= \|\mu\|_2^2 + 2 \underbrace{\mu^\top \zeta}_\text{I} + \underbrace{\|\zeta\|_2^2}_\text{II}, \\ \mathbb{E}\|Y\|_2^2 &= \|\mu\|_2^2 + k\sigma^2. \end{aligned}$$

We bound the variation of the both terms I and II. Note that since $\zeta \sim \mathcal{N}(0, \Sigma)$, we can represent it as $\zeta = \Sigma^{\frac{1}{2}}Z$ for $Z \sim \mathcal{N}(0, \mathbf{I}_d)$.

Bound on I. We have $\mu^\top \zeta = \mu^\top \Sigma^{\frac{1}{2}}Z$. By standard properties of multivariate normal distribution and sub-Gaussian random variables, we know $\mu^\top \zeta \sim \mathcal{N}(0, \mu^\top \Sigma \mu)$ and therefore

$$\mathbb{P}\left(|\mu^\top \zeta| \geq t\right) \leq 2 \exp\left\{-c \frac{t^2}{\mu^\top \Sigma \mu}\right\} \stackrel{(1)}{\leq} 2 \exp\left\{-c \frac{t^2}{\sigma^2 \|\mu\|_2^2}\right\},$$

where (1) is from the condition that $\mathbf{0} \preceq \Sigma \preceq \sigma^2 \mathbf{I}_d$ and c is a universal constant. Therefore, for $t \gtrsim \sigma \|\mu\|_2$, we have

$$\mathbb{P}\left(|\mu^\top \zeta| \geq t\right) = \mathcal{O}(1).$$

Bound on II. Since $\|\zeta\|_2^2 = Z^\top \Sigma Z$, by the Hanson-Wright's inequality (B.1), we know that

$$\mathbb{P}\left(|Z^\top \Sigma Z - k\sigma^2| \geq t\right) \leq 2 \exp\left\{-c \min\left\{\frac{t^2}{\|\Sigma\|_F^2}, \frac{t}{\|\Sigma\|_2}\right\}\right\}.$$

From the condition that $\mathbf{0} \preceq \Sigma \preceq \sigma^2 \mathbf{I}_d$ and $\text{tr}(\Sigma) = k\sigma^2$, we have

$$\begin{aligned} \|\Sigma\|_F^2 &= \sigma^4 \sum_{i=1}^d \lambda_i^2 \leq \sigma^4 \sum_{i=1}^d \lambda_i = k\sigma^4, \\ \|\Sigma\|_2 &= \max\{|\lambda_i|\} = \lambda_{\max} \leq \sigma^2. \end{aligned}$$

Plug in the inequalities above, we have

$$\mathbb{P}\left(|Z^\top \Sigma Z - k\sigma^2| \geq t\right) \leq 2 \exp\left\{-c \min\left\{\frac{t^2}{k\sigma^4}, \frac{t}{\sigma^2}\right\}\right\},$$

where c is again a universal constant. Consequently, for $t \gtrsim \sqrt{k}\sigma^2$, we have

$$\mathbb{P}\left(|Z^\top \Sigma Z - k\sigma^2| \geq t\right) = \mathcal{O}(1).$$

Finally, to ensure the effectiveness of the test (4.3), it suffices that the variations are asymptotically dominated by the mean:

$$\begin{aligned}\sigma \|\mu\|_2 &\lesssim \|\mu\|_2^2, \\ \sqrt{k}\sigma^2 &\lesssim \|\mu\|_2^2,\end{aligned}$$

which is equivalent to requiring that $\|\mu\|_2 \gtrsim k^{\frac{1}{4}}\sigma$. The proof is completed. \square

B.4 Proof of Theorem A.1

Proof. In this section, we provide the underlying rationale for Theorem A.1, which is also the basis of Corollary A.2. The central strategy mirrors the one employed in Appendix B.3: ensuring that the unknown variations are asymptotically dominated by the known expectations. The crucial difference here is that this technique must be integrated with the definition of the [consistent subset](#).

As a preliminary step, we recall a classical estimate for binomial coefficients. This bound will be essential for controlling the number of subsets of $[N]$ with cardinality exceeding $(1 - 2\epsilon)N$ in our subsequent analysis.

Lemma B.2. *For any $0 \leq k \leq n$, we have $\sum_{i=0}^k \binom{n}{i} \leq \left(\frac{en}{k}\right)^k$.*

Proof. By the Taylor's expansion of e^k , we know $e^k = \sum_{i=0}^{\infty} \frac{k^i}{i!} \geq \sum_{i=0}^k \frac{k^i}{i!}$. Therefore, it suffices to

show that $\frac{k^k}{n^k} \binom{n}{i} \leq \frac{k^i}{i!}$ for any $0 \leq i \leq k$, and this is equivalent to showing that $\frac{k^k}{n^k} \prod_{j=0}^{i-1} (n-j) \leq k^i$.

However, the last statement can be directly obtained from the fact that $\frac{k^{k-i}}{n^{k-i}} \leq 1 \leq \frac{n^i}{\prod_{j=0}^{i-1} (n-j)}$. \square

Consider E_e as in (A.1) for $S = S'$, where S' can be any subset of $[N]$ with $|S'| \geq (1 - 2\epsilon)N$. We first focus on the authentic samples (i.e., let $\mathbf{Y} = \tilde{\mathbf{Y}}$), and analyze the square term. By centering the samples, we have

$$\left\| A_k^\dagger \sum_{i \in S'} \tilde{Y}_i \right\|_2^2 = \underbrace{\left\| A_k^\dagger \sum_{i \in S'} (\tilde{Y}_i - \mu) \right\|_2^2}_{:=\text{I}} + \underbrace{|S'|^2 \left\| A_k^\dagger \mu \right\|_2^2}_{:=\text{II}} + 2 \underbrace{|S'| \sum_{i \in S'} \left[A_k^\dagger (\tilde{Y}_i - \mu) \right]^\top \mu}_{:=\text{III}}. \quad (\text{B.2})$$

Bound on I. By representing $S' = [N] \setminus ([N] \setminus S')$, we have

$$\text{I} = \underbrace{\left\| A_k^\dagger \sum_{i \in [N]} (\tilde{Y}_i - \mu) \right\|_2^2}_{:=\text{I}_1} + \underbrace{\left\| A_k^\dagger \sum_{i \in [N] \setminus S'} (\tilde{Y}_i - \mu) \right\|_2^2}_{:=\text{I}_{2,S'}} - 2 \underbrace{\sum_{i \in [N]} (\tilde{Y}_i - \mu)^\top (\mathbf{I}_d - X_k^*) \sum_{i \in [N] \setminus S'} (\tilde{Y}_i - \mu)}_{:=\text{I}_{3,S'}}.$$

For I_1 , the arguments of concentration is from the Hanson-Wright inequality B.1. Since $\mathbb{E}\text{I}_1 = kN\sigma^2 \left\| \mathbf{I}_d - X_k^* \right\|_F^2 \leq k$, $\left\| \mathbf{I}_d - X_k^* \right\|_2 \leq 1$. We know that

$$\mathbb{P} \left(\left| \text{I}_1 - kN\sigma^2 \right| \geq t \right) \leq 2 \exp \left\{ -c \min \left\{ \frac{t^2}{kN^2\sigma^4}, \frac{t}{N\sigma^2} \right\} \right\}.$$

By adjusting t to control the probability less than α , we obtain

$$t \geq c\sqrt{k}\sigma^2N$$

for some constant c that only depends on α .

For $\mathbf{I}_{2,S'}$, we combine the Hanson–Wright inequality with the union bound on S' . For any fixed S' with $|S'| \geq (1 - 2\epsilon)N$, by the Hanson–Wright inequality B.1 and the same conditions of $\mathbf{I}_d - A_k^\dagger$ as in the bound of \mathbf{I}_1 , we have

$$\mathbb{P}(|\mathbf{I}_2 - k|[N] \setminus S'| \sigma^2| \geq t) \leq 2 \exp \left\{ -c \min \left\{ \frac{t^2}{k|[N] \setminus S'|^2 \sigma^4}, \frac{t}{|[N] \setminus S'| \sigma^2} \right\} \right\}.$$

Now consider all such possible S' , whose amount is at most $\sum_{i=0}^{2\epsilon N} \binom{N}{i} = \left(\frac{e}{2\epsilon}\right)^{2\epsilon N}$ by Lemma B.2. Using the union bound, we have $\mathbf{I}_{2,S'} - k|[N] \setminus S'| \sigma^2 \geq t$ holds for some S' with probability less than

$$\sum_{S'} \exp \left\{ -c \min \left\{ \frac{t^2}{k|[N] \setminus S'|^2 \sigma^4}, \frac{t}{|[N] \setminus S'| \sigma^2} \right\} \right\} \leq \left(\frac{e}{2\epsilon}\right)^{2\epsilon N} \cdot \exp \left\{ -c \min \left\{ \frac{t^2}{k\epsilon^2 N^2 \sigma^4}, \frac{t}{\epsilon N \sigma^2} \right\} \right\}.$$

If we select $t_{\mathbf{I}_2}$ to be

$$t_{\mathbf{I}_2} = c \max \left\{ \sqrt{\left[\epsilon N \ln \left(\frac{e}{\epsilon} \right) + \ln \left(\frac{1}{\alpha} \right) \right] k \epsilon^2 N^2 \sigma^4}, \left[\epsilon N \ln \left(\frac{1}{\epsilon} \right) + \ln \left(\frac{1}{\alpha} \right) \right] \epsilon N \sigma^2 \right\}$$

then it can be verified that $|\mathbf{I}_{2,S'} - k|[N] \setminus S'| \sigma^2| \leq t_{\mathbf{I}_2}$ holds for all S' with probability greater than $1 - \alpha$.

For $\mathbf{I}_{3,S'}$, first, observe that the joint distribution of $\sum_{i \in [N] \setminus S'} (\tilde{Y}_i - \mu)$ and $\sum_{i \in [N]} (\tilde{Y}_i - \mu)$ is a multivariate Gaussian

$$\left(\sum_{i \in [N] \setminus S'} (\tilde{Y}_i - \mu)^\top, \sum_{i \in [N]} (\tilde{Y}_i - \mu)^\top \right) \sim \mathcal{N} \left(\mathbf{0}, \begin{pmatrix} |[N] \setminus S'| \mathbf{I}_d & |[N] \setminus S'| \mathbf{I}_d \\ |[N] \setminus S'| \mathbf{I}_d & N \mathbf{I}_d \end{pmatrix} \sigma^2 \right).$$

By the property of multivariate Gaussian distribution, conditional on the event $\left\{ \sum_{i \in [N]} (\tilde{Y}_i - \mu) = y \right\}$,

the conditional distribution of $\sum_{i \in [N] \setminus S'} (\tilde{Y}_i - \mu)$ is

$$\left[\sum_{i \in [N] \setminus S'} (\tilde{Y}_i - \mu) \mid \sum_{i \in [N]} (\tilde{Y}_i - \mu) = y \right] \sim \mathcal{N} \left(\frac{|[N] \setminus S'|}{N} y, \sigma^2 \left(|[N] \setminus S'| - \frac{|[N] \setminus S'|^2}{N} \right) \mathbf{I}_d \right).$$

From the maximal inequality that $\mathbb{P} \left(\max_{1 \leq i \leq N} \tilde{Y}_i \geq \sqrt{2\sigma^2(\ln N + t)} \right) \leq e^{-t}$, and conditional on the

event $\left\{ \sum_{i \in [N]} (\tilde{Y}_i - \mu) = y \right\}$, we have the probability of the following event

$$\left\{ \mathbf{I}_{3,S'} \geq \frac{|[N] \setminus S'|}{N} y^\top (\mathbf{I}_d - X_k^\dagger) y + \sqrt{2y^\top (\mathbf{I}_d - X_k^\dagger) y \left(|[N] \setminus S'| - \frac{|[N] \setminus S'|^2}{N} \right) \sigma^2 \cdot \left(2\epsilon N \ln \left(\frac{e}{2\epsilon} \right) + t \right)} \right\}$$

holding for some $S' \subset [N]$ less than e^{-t} . Also, for any general random variables X, Y , and $a \in \mathbb{R}$, we have

$$\begin{aligned} \mathbb{P}(Y \geq a) &= \mathbb{E} \mathbf{1}_{\{Y \geq a\}} \\ &= \mathbb{E} [\mathbf{1}_{\{Y \geq a\}} \cdot (\mathbf{1}_{\{f(X) \geq a\}} + \mathbf{1}_{\{f(X) < a\}})] \\ &= \mathbb{E} [\mathbb{E} [\mathbf{1}_{\{Y \geq a\}} \cdot \mathbf{1}_{\{f(X) \geq a\}} | X] + \mathbb{E} [\mathbf{1}_{\{Y \geq a\}} \cdot \mathbf{1}_{\{f(X) < a\}} | X]] \\ &\leq \mathbb{E} [\mathbb{E} [\mathbf{1}_{\{f(X) \geq a\}} | X] + \mathbb{E} [\mathbf{1}_{\{Y \geq f(X)\}} \cdot \mathbf{1}_{\{f(X) < a\}} | X]] \\ &\leq \mathbb{P}(f(X) \geq a) + \mathbb{E} [\mathbb{P}(Y \geq f(X) | X)]. \end{aligned}$$

Now observe that $|[N] \setminus S'| \leq 2\epsilon N$ for all S' , and take $Y = \max_{S'} \{I_{3,S'}\}$, $X = \sum_{i \in [N]} (\tilde{Y}_i - \mu)$, $f(y) = g(y^\top (\mathbf{I}_d - X_k^\dagger) y)$ where $g(y) = \frac{|[N] \setminus S'|}{N} y + \sqrt{4y \cdot \epsilon N (\epsilon N \ln(\frac{e}{\epsilon}) + t)}$ and $a = g(Nk\sigma^2 + \sqrt{\frac{4k^*}{\alpha}} N\sigma^2)$. By the argument above and the fact that $g(\cdot)$ is an increasing function on the positive real line, we know

$$\mathbb{P}(f(X) \geq a) = \mathbb{P}\left(y^\top (\mathbf{I}_d - X_k^\dagger) y \geq Nk\sigma^2 + \sqrt{\frac{4k}{\alpha}} N\sigma^2\right) \leq \frac{\alpha}{2}.$$

From the maximal inequality above, we also have $\mathbb{E}[\mathbb{P}(Y \geq f(X) | X)] \leq \mathbb{E}e^{-t} = e^{-t}$. Taking $t = \ln(\frac{2}{\alpha})$, and combining the results above, we conclude that with probability greater than $1 - \alpha$ the following event holds for all S' .

$$\left\{ I_{3,S'} \leq |[N] \setminus S'| k\sigma^2 + |[N] \setminus S'| \sqrt{\frac{4k}{\alpha}} \sigma^2 + \sqrt{8 \left(Nk\sigma^2 + \sqrt{\frac{4k}{\alpha}} N\sigma^2 \right) \epsilon N \left(2\epsilon N \ln\left(\frac{e}{2\epsilon}\right) + \ln\left(\frac{2}{\alpha}\right) \right) \sigma^2} \right\}.$$

By the symmetrical nature of $I_{3,S'}$, same technique and argument can also be applied to the probability of the lower tail of $I_{3,S'}$.

Bound on II. It is part of the expectation, and there is no randomness.

Bound on III. By splitting $S' = [N] \setminus ([N] \setminus S')$, we have

$$\text{III} = \underbrace{|S'| \sum_{i \in [N]} (\tilde{Y}_i - \mu)^\top A_k^\dagger \mu}_{\text{III}_{1,S'}} - \underbrace{|S'| \sum_{i \in [N] \setminus S'} (\tilde{Y}_i - \mu)^\top A_k^\dagger \mu}_{\text{III}_{2,S'}}.$$

For $\text{III}_{1,S'}$, observe that $\text{III}_{1,S'} \sim \mathcal{N}\left(0, |S'|^2 N\sigma^2 \|A_k^\dagger \mu\|_2^2\right)$. By the Chebyshev's inequality and the fact $|S'| \leq N$, we have

$$\mathbb{P}\left(|\text{III}_{1,S'}| \leq \sqrt{\frac{N^3 \sigma^2 \|P^* \mu\|_2^2}{\alpha}}\right) \geq 1 - \alpha.$$

For $\text{III}_{2,S'}$, similarly, we have $\text{III}_{2,S'} \sim \mathcal{N}\left(0, |S'|^2 |[N] \setminus S'| \|P^* \mu\|_2^2 \sigma^2\right)$ for each fixed S' . By the Chebyshev's inequality and the union bound over S' (similar argument as before), we know that the following holds.

$$\mathbb{P}\left(|\text{III}_{2,S'}| \leq \sqrt{2\epsilon N^3 \|A_k^\dagger \mu\|_2^2 \sigma^2} \cdot \sqrt{2\epsilon N \ln\left(\frac{e}{2\epsilon}\right) + \ln\left(\frac{1}{\alpha}\right)}, \forall S' \subset [N]\right) \geq 1 - \alpha.$$

In summary, for each unknown term of I and III, we have the following expectations and concentrations on the variations.

I₁ : Expectation: $kN\sigma^2$,

Variation: $\sqrt{\frac{2k}{\alpha}}N\sigma^2$,

Techniques: the Hanson–Wright inequality.

I_{2,S'} : Expectation: $k|[N]\setminus S'|\sigma^2$,

Variation: $c \max \left\{ \sqrt{[\epsilon N \ln(\frac{\epsilon}{\epsilon}) + \ln(\frac{1}{\alpha})] k \epsilon^2 N^2 \sigma^4}, [\epsilon N \ln(\frac{1}{\epsilon}) + \ln(\frac{1}{\alpha})] \epsilon N \sigma^2 \right\}$,

Techniques: Lemma B.2, union bound over S' , and the Hanson–Wright inequality.

I_{3,S'} : Expectation: $k|[N]\setminus S'|\sigma^2$,

Variation: $\left[|[N]\setminus S'| \sqrt{\frac{4k}{\alpha}} + \sqrt{8 \left(Nk + \sqrt{\frac{4k}{\alpha}} N \right) \epsilon N \left(2\epsilon N \ln(\frac{\epsilon}{2\epsilon}) + \ln(\frac{2}{\alpha}) \right)} \right] \sigma^2$,

Techniques: conditional argument and (sub-)Gaussian maximal inequality.

III_{1,S'}: Expectation: 0,

Variation: $\sqrt{\frac{N^3 \sigma^2 \|A_k^\dagger \mu\|_2^2}{\alpha}}$,

Techniques: Chebyshev's inequality.

III_{2,S'}: Expectation: 0,

Variation: $\sqrt{2\epsilon N^3 \|A_k^\dagger \mu\|_2^2} \sigma^2 \cdot \sqrt{2\epsilon N \ln(\frac{\epsilon}{2\epsilon})}$,

Techniques: Lemma B.2, union bound over S' and Chebyshev's inequality.

The sum of the expectations for (B.2) is exactly $k|S'|\sigma^2$, which is the second term in the LHS. It remains is to require all the variations above to be asymptotically less than the term II, which is equivalent to the condition $\|A_k^\dagger \mu\|_2 \gtrsim E_{\text{raw}}(\epsilon, K, N, \sigma^2, k)$ with simple algebra.

Now consider the general corrupted observations \mathbf{Y} . From the arguments above, we know that the original samples $\tilde{\mathbf{Y}}$ is consistent with ϕ_e with high (constant) probability given the condition (A.3). Combining with the simple fact that if $\tilde{\mathbf{Y}}$ is consistent, then $\tilde{\mathbf{Y}}_S$ is consistent for any $S \subset [N]$ with $|S| \geq (1 - \epsilon)N$, we have, therefore, with high probability, there exists at least one consistent subset within \mathbf{Y} (at least $\mathbf{Y}_{[N]\setminus C} = \tilde{\mathbf{Y}}_{[N]\setminus C}$) for ϕ_e . The theoretical algorithm works by scanning through all the subsets of \mathbf{Y} with cardinality no less than $(1 - \epsilon)N$ to check their consistency. Denote any found consistent subset as Y_{S_0} . Then the testing result on $\tilde{\mathbf{Y}}$ is recovered through the connection that $\phi_e(\mathbf{Y}_{S_0}) = \phi_e(\mathbf{Y}_{S_0 \cap [N]\setminus C}) = \phi_e(\mathbf{Y}_{[N]\setminus C}) = \phi_e(\tilde{\mathbf{Y}})$ and the fact that $|S_0 \cap [N]\setminus C| \geq (1 - 2\epsilon)N$. Since ϕ_e itself is a valid test for the uncontaminated testing problem, the proof is completed. \square

B.5 Necessary Additional Concentration Inequalities for Theorem 4.11

Lemma B.3. *Assume $\mathbf{X} \in \mathbb{R}^{n \times d}$ and each row is independently drawn from $\mathcal{N}(\mathbf{0}, \Sigma)$ with $\mathbf{0} \preceq \Sigma \preceq \mathbf{I}_d$ and $\text{tr}(\Sigma) = k \leq d$, then for $\delta > 0$, with probability higher than $1 - 2\delta$, we have:*

$$(i), n \geq k: \quad \|\mathbf{X}^\top \mathbf{X} - n\Sigma\|_2 \lesssim \sqrt{nk} + \sqrt{n \ln(\frac{1}{\delta})} + \ln(\frac{1}{\delta}),$$

$$(ii), n < k: \quad \|\mathbf{X}\mathbf{X}^\top - k\mathbf{I}_n\|_2 \lesssim \sqrt{nk} + \sqrt{k \ln(\frac{1}{\delta})} + \ln(\frac{1}{\delta}).$$

Proof. For the concentration bound on the sample covariance matrix, the proof involves the Talagrand generic chaining method and Talagrand's majorizing measure theorem. We kindly refer

the reader to Corollary 2 of [Koltchinskii and Lounici \[2017\]](#) for such results and proof. Here we only prove (ii) via the Hanson–Wright inequality and standard ϵ -net argument.

From the definition of ℓ_2 operator norm, we know

$$\left\| \mathbf{X}\mathbf{X}^\top - k\mathbf{I}_n \right\|_2 = \max_{u \in \mathbb{S}^{n-1}} \left| \left\| \mathbf{X}^\top u \right\|_2^2 - k \right|.$$

Since $\|u\|_2 = 1$, we know $\mathbf{X}^\top u$ is also a Gaussian random variable with distribution $\mathcal{N}(\mathbf{0}, \Sigma)$. Let $\mathbf{X}^\top = \sqrt{\Sigma}\mathbf{Z}$ where $\mathbf{Z} \sim \mathcal{N}(\mathbf{0}, \mathbf{I}_d)$. For any fixed $u \in \mathbb{S}^{n-1}$, by the Hanson–Wright inequality, we have

$$\mathbb{P} \left(\left| \left\| \mathbf{X}^\top u \right\|_2^2 - k \right| \geq t \right) = \mathbb{P} \left(\left| \mathbf{Z}^\top \Sigma \mathbf{Z} - k \right| \geq t \right) \stackrel{(i)}{\leq} 2 \exp \left\{ -c \min \left\{ \frac{t^2}{k}, t \right\} \right\},$$

where (i) is from the fact that $\mathbf{0} \preceq \Sigma \preceq \mathbf{I}_d$ and therefore $\|\Sigma\|_F^2 \leq \|\Sigma\|_F = k$, $\|\Sigma\|_2 \leq 1$.

For \mathbb{S}^{n-1} , from standard ϵ -net argument, we know that there exists a set $\mathcal{S} = \{v_1, v_2, \dots\}$ with $|\mathcal{S}| \leq 9^n$ and $\|u - v_i\|_2 \leq \frac{1}{4}$ for any $u \in \mathbb{S}^{n-1}$ and some $v_i \in \mathcal{S}$. Denote $Q = \mathbf{X}\mathbf{X}^\top - k\mathbf{I}_n$, then for any $u \in \mathbb{S}^{n-1}$ and $v_u \in \mathcal{S}$ such that $\|u - v_u\|_2 \leq \frac{1}{4}$, we have

$$\begin{aligned} \left| u^\top Q u - v_u^\top Q v_u \right| &= \left| u^\top Q(u - v_u) + (u - v_u)^\top Q v_u \right| \\ &\leq \left| u^\top Q(u - v_u) \right| + \left| (u - v_u)^\top Q v_u \right| \\ &\leq \|u\|_2 \|Q\|_2 \|u - v_u\|_2 + \|u - v_u\|_2 \|M\|_2 \|v_u\|_2 \\ &\leq \frac{1}{2} \|Q\|_2. \end{aligned}$$

Therefore, we have

$$\left\| u^\top Q u \right\|_2 \leq \left| v_u^\top Q v_u \right|_2 + \frac{1}{2} \|Q\|_2 \leq \sup_{v \in \mathcal{S}} \left| v^\top Q v \right| + \frac{1}{2} \|Q\|_2.$$

Taking supreme over u , we have

$$\|Q\|_2 \leq 2 \sup_{v \in \mathcal{S}} \left| v^\top Q v \right|.$$

Finally, we have

$$\mathbb{P}(\|Q\|_2 \geq t) \leq \mathbb{P} \left(\sup_{v \in \mathcal{S}} \left| v^\top Q v \right| \geq \frac{t}{2} \right) \leq 9^n \cdot 2 \exp \left\{ -c \min \left\{ \frac{t^2}{k}, t \right\} \right\}.$$

Simple calculations and the condition that $n < k$ yield that with probability greater than $1 - 2\delta$ we have

$$\|Q\|_2 \lesssim \sqrt{kn} + \sqrt{k \ln \left(\frac{1}{\delta} \right)} + \ln \left(\frac{1}{\delta} \right).$$

The proof is completed. □

Distinct from Lemma [B.3](#), Lemma [B.4](#) below bounds the ℓ_2 -operator norm of the weighted empirical covariance $\sum_{i=1}^n \omega_i X_i X_i^\top$ under small weights ω . This result plays a crucial role in verifying the second condition of ω -regularity subsequent to Algorithm [5](#) (or [6](#)). While the proof builds upon Fact 4.2 in [Dong et al. \[2019\]](#), our analysis is carefully tailored to handle the specific structural assumptions imposed on Σ .

Lemma B.4. Assume $\mathbf{X} \in \mathbb{R}^{n \times d}$ and each row is independently drawn from $\mathcal{N}(0, \boldsymbol{\Sigma})$ where $\mathbf{0} \preceq \boldsymbol{\Sigma} \preceq \mathbf{I}_d$ and $\text{tr}(\boldsymbol{\Sigma}) = k$. Then, with some constants $c_1, c_2 > 0$, the following holds with probability at least $1 - \delta$ for any ω that $0 \leq \omega_i \leq 1$ and $\|\omega\|_1 \leq c_2 \cdot \epsilon n$:

$$\left\| \sum_{i=1}^n \omega_i X_i X_i^\top \right\|_2 \leq c_1 \left(\epsilon n \ln \left(\frac{1}{\epsilon} \right) + k + \ln \left(\frac{1}{\delta} \right) \right).$$

Proof. Without loss of generality, we assume $c_2 = 1$. For any fixed $v \in \mathbb{S}^{d-1}$, We have:

$$v^\top \left(\sum_{i=1}^n \omega_i X_i X_i^\top \right) v = \sum_{i=1}^n \omega_i (v^\top X_i)^2.$$

Since $(v^\top X_i)^2 \geq 0$, we further assume $\|\omega\|_1 = \epsilon n$. The set $\{\omega \mid \|\omega\|_1 = \epsilon n\}$ is convex, consequently the RHS above is maximized at some vertex $\omega_S := \mathbf{1}_S, |S| = \epsilon n$. Therefore, we have

$$\sum_{i=1}^n \omega_i (v^\top X_i)^2 = \sum_{i \in S} \left[(v^\top X_i)^2 - v^\top \boldsymbol{\Sigma} v \right] + \epsilon n (v^\top \boldsymbol{\Sigma} v) \leq \left\| \sum_{i \in S} X_i X_i^\top - \epsilon n \boldsymbol{\Sigma} \right\|_2 + \epsilon n.$$

From Corollary 2 of [Koltchinskii and Lounici \[2017\]](#) (also see Lemma B.3), for any fixed S with $|S| = \epsilon n$, we have

$$\left\| \sum_{i \in S} X_i X_i^\top - \epsilon n \boldsymbol{\Sigma} \right\|_2 \lesssim \sqrt{\epsilon n k} + k + \sqrt{\epsilon n t} + t$$

with probability greater than $1 - e^{-t}$. From Lemma B.2, the total amount of such S is at most $\left(\frac{\epsilon}{\delta}\right)^{\epsilon n}$. Setting $t = \delta \cdot \left(\frac{\epsilon}{\delta}\right)^{\epsilon n}$ in the inequality above and using the union bound, we have:

$$\mathbb{P} \left(\left\| \sum_{i \in S} X_i X_i^\top - \epsilon n \boldsymbol{\Sigma} \right\|_2 \lesssim \sqrt{\epsilon n k} + k + \sqrt{\epsilon n \ln \left(\frac{1}{\delta} \right)} + \ln \left(\frac{1}{\delta} \right) + \epsilon n \ln \left(\frac{1}{\epsilon} \right), \forall S, |S| = \epsilon n \right) \geq 1 - \delta.$$

Finally, we have

$$\sqrt{\epsilon n k} + k + \sqrt{\epsilon n \ln \left(\frac{1}{\delta} \right)} + \ln \left(\frac{1}{\delta} \right) + \epsilon n \ln \left(\frac{1}{\epsilon} \right) \lesssim \epsilon n \ln \left(\frac{1}{\epsilon} \right) + k + \ln \left(\frac{1}{\delta} \right).$$

The proof is completed by taking the supreme over $v \in \mathbb{S}^{d-1}$. □