Ultrahigh dimensional precision matrix estimation via refitted cross validation

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ABSTRACT

This paper develops a new estimation procedure for ultrahigh dimensional sparse precision matrix, the inverse of covariance matrix. Regularization methods have been proposed for sparse precision matrix estimation, but they may not perform well with ultrahigh dimensional data due to the spurious correlation. We propose a refitted cross validation (RCV) method for sparse precision matrix estimation based on its Cholesky decomposition, which does not require the Gaussian assumption. The proposed RCV procedure can be easily implemented with existing software for ultrahigh dimensional linear regression. We establish the consistency of the proposed RCV estimation and show that the rate of convergence of the RCV estimation without assuming banded structure is the same as that of those assuming the banded structure in Bickel and Levina (2008b). Monte Carlo studies were conducted to access the finite sample performance of the RCV estimation. Our numerical comparison shows that the RCV estimation outperforms the existing ones in various scenarios. We further apply the RCV estimation for an empirical analysis of asset allocation.

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1. Introduction

Precision matrix, the inverse of covariance matrix, plays an important role in statistical inference and statistical learning such as one sample mean test, graphical modeling and linear discrimination analysis. This study is motivated by an empirical analysis of portfolio allocation, in which the estimation of precision matrix is required (see Section 3.2 for more details). In the classic multivariate statistical analysis, the inverse of sample covariance matrix is a natural and consistent estimator of precision matrix. Sample covariance matrix, however, becomes singular for high dimensional data when the sample size \( n \) is less than the dimension \( p \) of data. Thus, it is of interest to develop alternative methods to estimate precision matrix. Bickel and Levina (2008a), Rothman et al. (2009) and Cai and Liu (2011) approached the estimation problem by applying different thresholding methods to the sample covariance matrix. These methods may not perform well for ultrahigh dimensional data. Yuan and Lin (2007) considered a LASSO-type method for precision matrix and solved the optimization problem by a maxdet algorithm. Due to the important role of precision matrix in Gaussian graphic model, Ren et al. (2015), Fan and Lv (2016) and Ren et al. (2019) studied the estimation and statistical
inference of sparse precision matrix by assuming conditional normal distribution. However, it is challenging in verifying the normality of high dimensional data in practice. People recast the problem of precision matrix estimation into a high dimensional regularization regression. A constrained $\ell_1$ minimization procedure was proposed by Cai et al. (2011) for sparse precision matrix estimation. Rothman et al. (2008) and Lam and Fan (2009) used the penalized likelihood method to estimate sparse precision matrix with $\ell_1$ and nonconvex penalty functions, respectively. These methods are easy to implement, and their statistical properties have been studied. However, when $p \gg n$, these methods may perform poorly due to severe spurious correlations. A comprehensive investigation of the impact of spurious collections on error variance estimation in ultrahigh dimensional data can be found in Fan et al. (2012). This paper aims to develop a consistent and easy-to-compute estimator of precision matrix for ultrahigh dimensional data without the Gaussian assumption.

In this paper, we propose an estimation procedure for the ultrahigh dimensional precision matrix by using refitted cross-validation (RCV, Fan et al., 2012). We first parameterize the precision matrix using the modified Cholesky decomposition. The decomposition is valid for any positive definite matrices and does not require the Gaussian assumption. It has been used in some literature for the estimation of covariance or precision matrix. Specifically, the modified Cholesky decomposition has been adopted for estimating the finite-dimensional error covariance matrix in longitudinal regression model (Pourahmadi, 1999, 2000), for the low-dimensional covariance estimation in Huang et al. (2006), and for the high-dimensional precision matrix estimation in Rothman et al. (2010). It is challenging to employ the Cholesky decomposition for ultrahigh dimensional covariance matrix estimation since it leads to ultrahigh dimensional linear regression problems. And the theoretical analyses of these problems are much more difficult than that in finite- or high-dimensional settings. To overcome the difficulty, we proposed the RCV estimation. The RCV estimation procedure for the ultrahigh dimensional covariance matrix can be carried out easily with the existing software for ultrahigh dimensional linear regressions. The resulting estimator of precision matrix is symmetric and positive definite. We show that the estimator is consistent. Moreover, we provide the convergence rate of the newly proposed estimator without assuming banded structure. The convergence rate is the same as the rate of those methods assuming the banded structure in Bickel and Levina (2008b).

The rest of this paper is organized as follows. In Section 2, we propose the RCV estimation procedure for the precision matrix and study its theoretical properties. In Section 3, we present a numerical comparison and an empirical study. Technical conditions and proofs are given in Appendix.

2. An RCV estimator for precision matrix

Let $x = (X_1, \ldots, X_p)^T$ be a $p$-dimensional random vector with mean vector $\mu$ and covariance matrix $\Sigma$. Of interest is to estimate the precision matrix $\Omega = \Sigma^{-1}$. Without loss of generality, assume $\mu = 0$ throughout this paper. As in Pourahmadi (1999), the modified Cholesky’s decomposition of $\Omega$ is defined as follows.

$$\Omega = L^T D^{-1} L,$$

(2.1)

where $L$ is unitriangular matrix (i.e. a lower triangular matrix satisfying that each diagonal element $l_{tt} = 1, \ t = 1, \ldots, p$) and $D = \text{diag}(\sigma_1^2, \ldots, \sigma_p^2)$ is a diagonal matrix. Regardless of the population distribution, the matrix decomposition is always true for the precision matrix. Since $\Omega = \Sigma^{-1}$, (2.1) is equivalent to

$$L \Sigma L^T = D,$$

which provides us a nice interpretation of the elements in $L$ and $D$. Define

$$e = Lx.$$

(2.2)

Let $-\beta_{tk} = l_{tk}$, the $(t, k)$-th entry of $L$ for $t < k$ and $e = (e_1, \ldots, e_p)^T$. Then (2.2) is equivalent to that $X_t = e_t$, and for $1 < t \leq p$,

$$X_t = \sum_{k=1}^{t-1} \beta_{tk} X_k + e_t,$$

(2.3)

with $E(e_t) = 0$ and $\text{var}(e_t) = \sigma_t^2$. Thus, the purpose of estimating $\Sigma$ and $\Omega$ can be achieved by estimating the regression coefficients and error variances in (2.2). Based on this idea, Pourahmadi (1999, 2000) proposed parametrization of error covariance matrix in the analysis of longitudinal data. In the presence of sparsity on $\beta_{tk}$, Huang et al. (2006) developed an estimation procedure for covariance matrix $\Sigma$ using penalized least squares method in the fixed and finite dimensional setting. Rothman et al. (2010) assumed that $L$ is banded and then applied LASSO-type estimation to estimate the regression coefficients $\beta$’s under the high dimensional setting.

2.1. A naive estimator

To illustrate the challenge of estimating ultrahigh dimensional precision matrix, let us start with a natural extension of the estimation procedure proposed by Huang et al. (2006). Since we will deal with ultrahigh dimensional regressions in (2.3), we impose sparsity assumption on $\beta_t = (\beta_{t1}, \ldots, \beta_{t(t-1)})^T$, and apply penalized least squares (PLS) method to
estimate $\beta_t$. After estimating $\beta_t$, we use the mean squared errors (MSE) of the corresponding linear models to estimate $\sigma_t^2$, the $t$th diagonal element of $D$.

Suppose that $x_i = (X_{i1}, \ldots, X_{ip})^T$, $i = 1, \ldots, n$ are the independent and identically distributed samples from the population of $x$. The PLS estimator of $\beta_t$ is

$$\tilde{\beta}_{t}^\text{LS} = \arg\min_{\beta_t} \frac{1}{2n} \sum_{i=1}^{n} \left( X_{it} - \sum_{k=1}^{t-1} \beta_{tk} X_{ik} \right)^2 + \sum_{k=1}^{t-1} p_{\lambda k}(|\beta_{tk}|).$$

(2.4)

where $p_{\lambda k}(\cdot)$ is a penalty function with tuning parameter $\lambda_k$. Various penalty functions can be applied, such as $\ell_1$-penalty (Tibshirani, 1996), SCAD penalty (Fan and Li, 2001) and MCP penalty (Zhang, 2010). With the estimate of $\beta_t$, we can estimate $\sigma_t^2$ by the MSE

$$\tilde{\sigma}_t^2 = \frac{1}{n - \|\tilde{\beta}_t\|_0} \sum_{i=1}^{n} \left( X_{it} - \sum_{k=1}^{t-1} \tilde{\beta}_{tik} X_{ik} \right)^2.$$

(2.5)

where $\|\cdot\|_0$ denotes the $\ell_0$-norm (i.e., the number of nonzero elements in the vector), and $\tilde{\beta}_{tik}$ is the least squares estimator under the model selected by the PLS method.

As the number of parameters increases, the computation consumption increases, while algorithm stability for Eq. (2.4) may decrease, especially for the ultrahigh dimensional cases ($p = O(e^\gamma n)$, $0 \leq \gamma < 1$). When $\log(t) = O(n^\alpha)$ for some constant $\alpha \in (0, 1)$, the penalized least squares (2.4) is difficult to carry out because of computational cost and algorithm stability. To tackle this problem, Fan and Lv (2008) proposed the sure independence screening (SIS) procedure to quickly reduce the ultrahigh dimension to a moderate dimension $d$ such as $d = [n/\log n]$. Consequently, we integrate SIS and PLS procedures estimate the precision matrix, and obtain a naive estimator of $\Omega$ as follows.

$$\hat{\Omega} = \tilde{L}^T \tilde{D}^{-1} \tilde{L},$$

(2.6)

where $\tilde{D}$ and $\tilde{L}$ are obtained from $\tilde{\sigma}_t^2$, $t = 1, \ldots, p$ and $\tilde{\beta}_t$, $t = 1, \ldots, p$, respectively.

The naive estimator $\hat{\Omega}$ may perform well in low dimensional case. However, it may not work well under ultrahigh dimensional setting due to the bias of $\tilde{\sigma}_t^2$. Let $M_t$ be the model selected by the PLS (2.4), $X_{M_t}$ be the design matrix for the selected model, and $P_{M_t} = X_{M_t}(X_{M_t}^T X_{M_t})^{-1} X_{M_t}^T$ be the projection matrix of the linear space spanned by the columns of $X_{M_t}$. Define $e_{M_t}^2 = e_{t}^T P_{M_t} e_{t}/e_{t}^T e_{t}$, where $e_{t} = (e_{t1}, \ldots, e_{tn})^T$. Under certain regularity conditions, Fan et al. (2012) showed that when $\log t$ and $n$ have the same order, $\tilde{\sigma}_t^2/(1 - \zeta_{M_t}^2) \rightarrow \sigma_t^2$ in probability and

$$\sqrt{n}(\tilde{\sigma}_t^2/(1 - \zeta_{M_t}^2) - \sigma_t^2) \rightarrow N(0, \text{var}(\tilde{\sigma}_t^2)).$$

Thus, $\tilde{\sigma}_t^2$ shrinks $\sigma_t^2$ by the factor $(1 - \zeta_{M_t}^2)$. Fan et al. (2012) demonstrated that $\zeta_{M_t}^2$ may not tend to zero under ultrahigh dimensional setting. This confirms that $\tilde{\sigma}_t^2$ can be an inconsistent estimator of $\sigma_t^2$ for large $t$. As a result, $\tilde{D}$ may not be a consistent estimator of $D$, and $\hat{\Omega}$ may not perform well.

Example 2.1. Before we pursue further, let us illustrate the bias of $\tilde{D}$ by running a small simulation study. Take $D = 0.01 I$ and $l_{.,k} = 0.6$ if $k = t - 1$, and $0$ otherwise. The corresponding covariance matrix $\Sigma$ satisfies the AR(1) structure. Set $n = 160$ and $p = 1000$. We run 200 simulations. For each simulation, we obtain $\tilde{\beta}_t$, the PLS with SCAD penalty while the tuning parameter is selected by cross-validation, and then obtain $\tilde{\sigma}_t^2$. The left panel of Fig. 1 presents the means of naive estimators $\tilde{D}_{it}$ for $t = 1, \ldots, p$ over 200 replications. The plot clearly shows a significant declining trend as the dimension increases. This implies that the bias of $\tilde{\sigma}_t^2$ becomes more severe as $t$ gets larger. Therefore, we need to develop a new statistical method to eliminate the bias.
2.2. RCV estimator

As shown in Fan et al. (2012), the bias of the naive estimator \( \hat{\sigma}^2_t \) comes from the spuriously correlated covariates in the selected models. To deal with the bias issue, we propose using the RCV method for the estimation of precision matrix. The RCV method has been used to estimate error variance in the ultrahigh dimensional linear regression (Fan et al., 2012) and in ultrahigh dimensional additive models (Chen et al., 2018), but the RCV method is new for estimation of ultrahigh dimensional precision matrix.

RCV procedure for precision matrix estimation consists of the following steps. In the first step, we randomly partition the \( n \) samples into two equally sized parts \( X^{(1)} \) and \( X^{(2)} \). In the second step, we conduct feature screening and variable selection procedures by using dataset \( X^{(1)} \) for the linear model (2.3) and obtain the selected model \( D_{1,t} \), for each \( t = 2, \ldots, p \). We fit the data \( X^{(2)} \) to the selected model \( D_{1,t} \) and obtain the regression coefficients \( \hat{\beta}^{(1)}_t \) and the mean squared errors \( \hat{\sigma}^2_t^{(1)} \) for \( t = 2, \ldots, p \). Based on these estimators \( \hat{\beta}^{(1)}_t \) and \( \hat{\sigma}^2_t^{(1)} \), we can obtain the estimators \( \hat{L}^{(1)} \) and \( \hat{D}^{(1)} \) for \( L \) and \( D \), respectively. In the third step, we switch the roles of \( X^{(1)} \) and \( X^{(2)} \). We use \( X^{(2)} \) to select model \( D_{2,t} \) and use \( X^{(1)} \) to refit \( D_{2,t} \) to obtain \( \hat{\beta}^{(2)}_t \) and \( \hat{\sigma}^2_t^{(2)} \) for \( t = 2, \ldots, p \). We further obtain estimators \( \hat{L}^{(2)} \) and \( \hat{D}^{(2)} \). Finally, the RCV estimators of \( D \) and \( L \) are defined as

\[
\hat{D} = \frac{\hat{D}^{(1)} + \hat{D}^{(2)}}{2}, \quad \hat{L} = \frac{\hat{L}^{(1)} + \hat{L}^{(2)}}{2},
\]

and the RCV estimator for \( \Omega \) is

\[
\hat{\Omega} = \hat{L}^{-1} \hat{D}.
\]

Since \( \hat{L} \) is still a lower triangular matrix, the estimator of precision matrix \( \hat{\Omega} \) is positive definite. If we partition the data into two parts with different sample sizes \( n_1 \) and \( n_2 \), \( n_1 \neq n_2 \), then the RCV estimators of \( D \) and \( L \) are accordingly adjusted to

\[
\hat{L} = \frac{n_1 \hat{L}^{(1)} + n_2 \hat{L}^{(2)}}{n_1 + n_2}, \quad \text{and} \quad \hat{D} = \frac{n_1 \hat{D}^{(1)} + n_2 \hat{D}^{(2)}}{n_1 + n_2}.
\]

As a direct comparison, we apply the RCV estimator \( \hat{D} \) for Example 2.1 and depicts the results in the right panel of Fig. 1, which shows that the two methods present different trends and the proposed method indeed eliminates the bias due to the spurious correlation. Fig. 2 demonstrates the flowchart of the RCV procedure for the precision matrix estimation.

2.3. Asymptotic properties

In this section, we study the asymptotic behavior of the proposed RCV estimator. We first introduce some notation. For a \( p \)-dimensional vector \( v = (v_1, \ldots, v_p)^T \), \( \|v\|_\infty \) stands for the \( \ell_\infty \)-norm of \( v \). In particular, \( v_0 \) is the number of nonzero elements of \( v \), and \( \|v\|_\infty = \max_i |v_i| \). For a \( p \)-by-\( q \) matrix \( A = (a_{ij}) \), the operator norm is defined by

\[
\|A\|_{(\alpha, \beta)} = \sup_{\|v\|_\alpha = 1} \|Av\|_\beta.
\]
where \( x \) is any \( p \)-dimensional unit vector and \( 0 \leq \alpha, \beta \leq \infty \). In particular,
\[
\| \Sigma \|_{(1,1)} = \max_j \sum_i |\sigma_{ij}|, \quad \| \Sigma \|_{(\infty, \infty)} = \max_i \sum_j |\sigma_{ij}|.
\]
Define the spectral norm \( \| A \| = \| A \|_{(2,2)} \). For a \( p \times p \) symmetric matrix \( \Sigma \), it follows that
\[
\| \Sigma \| = \lambda_{\text{max}}(\Sigma) \leq (\| \Sigma \|_{(1,1)}\| \Sigma \|_{(\infty, \infty)})^{1/2} = \max_j \sum_i |\sigma_{ij}|,
\]
where \( \lambda_{\text{max}}(\Sigma) \) and \( \lambda_{\text{min}}(\Sigma) \) stand for the largest and smallest eigenvalues of \( \Sigma \), respectively. (See Corollary 2.3.2, Golub and Van Loan (2012)). With a slight abuse of notation, let \( \| \Sigma \|_{\infty} = \max_{ij} |\sigma_{ij}| \) be the maximum norm of a matrix and \( |\mathcal{M}| \) stands for its cardinality of the set \( \mathcal{M} \).

Denote \( X = [x_1, \ldots, x_n]^T \) and \( X_t \) be a submatrix consisting of the first \( (t-1) \) columns of \( X \). Let \( \mathcal{M} \subseteq \{1, \ldots, p\} \) be the active index set. \( X_{\mathcal{M}, t} \) is the submatrix of \( X_t \), with the columns indexed by \( \mathcal{M} \). Define the \( m \)-sparse minimal eigenvalue and \( m \)-sparse maximal eigenvalue (Meinshausen and Yu, 2009) as
\[
\phi_{\text{min}}(m, t) = \min_{\mathcal{M} : |\mathcal{M}| \leq m} \lambda_{\text{min}}(n^{-1}X_{\mathcal{M}, t}^T X_{\mathcal{M}, t}),
\]
and
\[
\phi_{\text{max}}(m, t) = \max_{\mathcal{M} : |\mathcal{M}| \leq m} \lambda_{\text{max}}(n^{-1}X_{\mathcal{M}, t}^T X_{\mathcal{M}, t}).
\]

We need the following technical conditions to facilitate the proofs, although they may not be the weakest.

(A1) (Exponential tail conditions) \( X_t \) and \( e_t \) defined in (2.2), \( j = 1, \ldots, p \) satisfy the exponential condition. That is, there exist positive constants \( c_0 \) and \( T \) such that for \( 0 < |t| < T \),
\[
E e^{t^2/2} \leq c_0, \quad \text{and} \quad E e^{t^4} \leq c_0, \quad j = 1, \ldots, p.
\]

(A2) There exist a constant \( \lambda_0 > 0 \) and constant sequence \( \eta_n \) such that \( \eta_n = o(n) \) and \( P(\phi_{\text{min}}(\eta_n, t) \geq \lambda_0) = 1 \), for all \( n \) and \( t = 2, \ldots, p \).

(A3) Assume that there exists a constant \( \varepsilon_0 > 0 \) such that \( \varepsilon_0^{-1} \leq \lambda_{\text{min}}(\Omega) \leq \lambda_{\text{max}}(\Omega) \leq \varepsilon_0 \).

When \( x \) follows a multivariate normal distribution or \( X_t \)'s are bounded, the exponential tail condition (A1) is satisfied. This assumption is commonly adopted in the literature of high dimensional data modeling (Bickel and Levina, 2008a; Yuan, 2010; Cai et al., 2011; Cai and Liu, 2011).

Condition (A2) is a common assumption for the high dimensional variable selection, which is necessary for LASSO-type method to select all important variables. Here (A2) implies that the selected variables in stage one are not highly correlated. Condition (A3) implies that both the smallest and largest eigenvalues of the precision matrix (and the covariance matrix) are bounded by a constant. Recall notation \( \hat{\Sigma}^j \), \( j = 1, 2 \) in Section 2.2, and let \( \hat{\Sigma}^j_1 \) and \( \hat{\Sigma}^j_2 \) be the \( st \) column and \( th \) row of \( \hat{\Sigma}^j \), respectively. Denote \( k_n = \max_{1 \leq s \leq 5p} \max_{1 \leq t \leq 2} \max \{\| \hat{\Sigma}^j_1 \|_1, \| \hat{\Sigma}^j_2 \|_1\} \).

**Theorem 1.** Denote \( D_x^{\gamma} = \{k : \beta_k \neq 0\} \) and assume that \( D_x^{k_0} \) satisfies \( P(\mathcal{D}_k \supseteq D_x^{k_0}) = 1 \) for \( k = 1 \) and \( 2 \), and \( t = 2, \ldots, p \). Suppose that \( \log p = O(n^{\gamma}) \), \( \gamma \in (0, 1) \) and \( k_n = O((n/\log p)^{1/(3\alpha + 1)}) \), for any constant \( \alpha > 0 \). Under Condition (A1), (A2) and (A3), it follows that
\[
\| \hat{\Omega} - \Omega \| = O_p \left( \sqrt{\frac{k_n^2 \log p}{n}} \right).
\]

The condition \( \log p = O(n^\gamma) \), \( \gamma \in (0, 1) \) is usually adopted by the literature of ultrahigh dimensional data analysis. The order imposed on \( k_n \) implies that the number of nonzero elements in each row and each column of \( \hat{\Sigma}^j \) cannot diverge to infinite too fast. For banded \( \hat{\Sigma}^j \), we can simply control the number of nonzero elements in \( \hat{\Sigma}^j \). The convergence rate in Theorem 1 is quite similar to that of banded matrices estimation in Bickel and Levina (2008b). Note the fact that \( L \) of a banded matrix is still banded. If \( k_0 = (n/\log p)^{1/(3\alpha + 1)} \), for some constant \( \alpha > 0 \), then the convergence rate of (2.11) becomes \( O_p((\log p/n)^{\alpha/(2\alpha + 1)}) \), which is exactly the same as that in Bickel and Levina (2008b). However, the assumption of \( L \) in this paper is more flexible than that in Bickel and Levina (2008b) and Levina et al. (2008). Levina et al. (2008) directly assume the banded structure, and there is no selection procedure involved in their procedure, so the spurious correlation has much less impact on their estimation. In contrast, the proposed method assumes sparsity only rather than the banded structure.

3. Numerical studies

We investigate the finite sample performance of the proposed estimation procedure via Monte Carlo simulation study. We compare the proposed procedure with exiting ones in Section 3.1. We illustrate the proposed methodology by an empirical analysis of real financial market data.
Table 1
Simulation results for precision matrix $\Omega_1$ and $\Omega_2$.

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<th>p</th>
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<th>100</th>
<th>200</th>
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<td>60.9</td>
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<td>–</td>
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<tr>
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<td>74.44</td>
<td>193.0</td>
<td>555.8</td>
<td>1185.5</td>
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<tr>
<td>Naive-LASSO-AIC</td>
<td>1.44</td>
<td>6.37</td>
<td>18.74</td>
<td>54.7</td>
<td>233.0</td>
<td>714.6</td>
<td>2338.5</td>
<td>–</td>
</tr>
<tr>
<td>Naive-LASSO-BIC</td>
<td>0.71</td>
<td>2.06</td>
<td>4.62</td>
<td>10.1</td>
<td>28.9</td>
<td>62.7</td>
<td>137.9</td>
<td>–</td>
</tr>
<tr>
<td>Naive-SCAD-AIC</td>
<td>1.14</td>
<td>5.27</td>
<td>16.35</td>
<td>78.9</td>
<td>1170.3</td>
<td>8270.2</td>
<td>43641.6</td>
<td>–</td>
</tr>
<tr>
<td>Naive-SCAD-BIC</td>
<td>0.60</td>
<td>1.73</td>
<td>4.07</td>
<td>10.1</td>
<td>32.9</td>
<td>78.8</td>
<td>185.7</td>
<td>–</td>
</tr>
<tr>
<td>RCV-LASSO-AIC</td>
<td>0.79</td>
<td>2.43</td>
<td>5.67</td>
<td>13.4</td>
<td>49.2</td>
<td>142.5</td>
<td>433.2</td>
<td>–</td>
</tr>
<tr>
<td>RCV-LASSO-BIC</td>
<td>0.62</td>
<td>1.69</td>
<td>3.84</td>
<td>8.4</td>
<td>26.2</td>
<td>60.3</td>
<td>141.5</td>
<td>–</td>
</tr>
<tr>
<td>RCV-SCAD-AIC</td>
<td>0.71</td>
<td>2.46</td>
<td>7.88</td>
<td>31.1</td>
<td>182.7</td>
<td>575.9</td>
<td>1593.4</td>
<td>–</td>
</tr>
<tr>
<td>RCV-SCAD-BIC</td>
<td>0.62</td>
<td>1.69</td>
<td>3.86</td>
<td>8.5</td>
<td>26.1</td>
<td>60.0</td>
<td>143.2</td>
<td>–</td>
</tr>
</tbody>
</table>

3.1. Monte Carlo simulation

In this section, we conduct Monte Carlo simulation to examine the finite sample performance of the newly proposed method. We generate random samples from a $p$-dimensional multi-normal distribution $\mathcal{N}(\mathbf{0}, \Omega^{-1})$. We set the precision matrix $\Omega$ through its Cholesky decomposition $\Omega = \mathbf{L}^T \mathbf{D}^{-1} \mathbf{L}$, and consider four different scenarios.

(a) $\Omega_1 = \mathbf{L}^T \mathbf{D}^{-1} \mathbf{L}$ with $\mathbf{D} = \mathbf{I}$ and $l_{ij} = 1$, $l_{i+1,i} = \rho = 0.5$ with all other $l_{ij} = 0$. This is referred to as the AR(1) covariance matrix in the literature.

(b) $\Omega_2 = \mathbf{L}^T \mathbf{D}^{-1} \mathbf{L}$ with $\mathbf{D} = \mathbf{I}$. We randomly choose an entry $j, j < t$ in the $t$th row of $\mathbf{L}$, and set $l_{ij} = \rho$ and all other $l_{ij} = 0$.

(c) $\Omega_3 = \mathbf{L}^T \mathbf{D}^{-1} \mathbf{L}$ with $\mathbf{D} = \mathbf{I}$. We randomly choose two entries $j_1, j_2, j_1 < t, j_2 < t$ in the $t$th row of $\mathbf{L}$, and set $l_{j_1j_2} = l_{j_2j_1} = \rho$ and all other $l_{ij} = 0$.

(d) $\Omega_4 = \mathbf{L}^T \mathbf{D}^{-1} \mathbf{L}$ with $\mathbf{D} = \text{diag}(\sigma_1^2, \ldots, \sigma_p^2)$. $\mathbf{L}$ is set in the same way as in that in $\Omega_3$, but $\{\sigma_t^2, t = 1, \ldots, p\}$ are generated by the uniform distribution $\text{Unif}[1, 2]$.

In this simulation study, we evaluate estimation procedures by the following quadratic loss

$$\Delta(\hat{\Omega}, \Omega) = \text{tr}(\Omega^{-1} \hat{\Omega}^2 - 1)^2.$$ 

Note that the quadratic loss is zero if and only if $\hat{\Omega} = \Omega$. In our study, we also use different criteria such as the entropy loss (Muirhead, 1982), the Frobenius loss and the Kullback–Leibler divergence. Since the results of numerical comparison based on different measures are almost the same under all of our simulation settings, we only present those based on the quadratic loss to save space.

In our simulation, we set $p = 20, 50, 100, 200, 500, 1000$ and $2000$ to investigate the impact of dimensionality, and the sample size $n = 200$. For each combination, we conduct 100 replications. For the purpose of comparison, four existing methods are considered: (1) the inverse of sample covariance matrix when $p < n$, labeled by “Sample Covariance” in tables, (2) the regularized method in Bickel and Levina (2008b), denoted by “BL” for short, (3) the $\ell_1$ penalized normal likelihood method in Huang et al. (2006), denoted by “PNL”, and (4) the CLIME in Cai et al. (2011), denoted by “CLIME”. Since the computation cost of PNL and CLIME are much higher than others, we are able to get their results only when $p < n = 200$.

We also include the naive method in our numerical comparison. For both naive estimator and RCV method, we consider penalized least squares with the LASSO and SCAD penalties with tuning parameter chosen by Akaike information criterion (AIC) or Bayesian information criterion (BIC). This leads to 8 combinations (Naive/RCV-LASSO/SCAD-AIC/BIC).

The simulation results for $\Omega_1$ and $\Omega_2$ are depicted in the top and the bottom panels of Table 1. From the top panel of Table 1, we can see that the BL method preforms the best across all different $p$ dimensions. This is expected because $\Omega_1$ and...
Table 2
Simulation results for precision matrix $\Omega_3$ and $\Omega_4$.

<table>
<thead>
<tr>
<th>$\Omega$</th>
<th>p</th>
<th>20</th>
<th>50</th>
<th>100</th>
<th>200</th>
<th>500</th>
<th>1000</th>
<th>2000</th>
</tr>
</thead>
<tbody>
<tr>
<td>$\Omega_3$ Sample Covariance</td>
<td>3.18</td>
<td>36.65</td>
<td>530.51</td>
<td>–</td>
<td>–</td>
<td>–</td>
<td>–</td>
<td>–</td>
</tr>
<tr>
<td>PNL</td>
<td>1.53</td>
<td>5.85</td>
<td>16.12</td>
<td>–</td>
<td>–</td>
<td>–</td>
<td>–</td>
<td>–</td>
</tr>
<tr>
<td>CLIME</td>
<td>8.53</td>
<td>61.55</td>
<td>233.15</td>
<td>1431.8</td>
<td>3921.5</td>
<td>10462.0</td>
<td>–</td>
<td>–</td>
</tr>
<tr>
<td>BL</td>
<td>3.68</td>
<td>27.64</td>
<td>105.71</td>
<td>345.3</td>
<td>1439.0</td>
<td>4015.2</td>
<td>10562.0</td>
<td>–</td>
</tr>
<tr>
<td>Naive-LASSO-AIC</td>
<td>1.60</td>
<td>6.39</td>
<td>17.77</td>
<td>48.9</td>
<td>187.7</td>
<td>526.0</td>
<td>1443.1</td>
<td>–</td>
</tr>
<tr>
<td>Naive-LASSO-BIC</td>
<td>1.04</td>
<td>3.34</td>
<td>7.55</td>
<td>17.2</td>
<td>47.9</td>
<td>103.3</td>
<td>222.0</td>
<td>–</td>
</tr>
<tr>
<td>Naive-SCAD-AIC</td>
<td>1.65</td>
<td>3.82</td>
<td>9.93</td>
<td>25.9</td>
<td>106.9</td>
<td>347.0</td>
<td>1240.6</td>
<td>–</td>
</tr>
<tr>
<td>Naive-SCAD-BIC</td>
<td>0.73</td>
<td>2.20</td>
<td>5.01</td>
<td>11.8</td>
<td>37.5</td>
<td>91.9</td>
<td>232.0</td>
<td>–</td>
</tr>
<tr>
<td>RCV-LASSO-AIC</td>
<td>1.09</td>
<td>3.51</td>
<td>8.08</td>
<td>19.1</td>
<td>55.9</td>
<td>129.7</td>
<td>300.0</td>
<td>–</td>
</tr>
<tr>
<td>RCV-LASSO-BIC</td>
<td>1.14</td>
<td>3.34</td>
<td>7.00</td>
<td>15.2</td>
<td>38.4</td>
<td>83.3</td>
<td>179.7</td>
<td>–</td>
</tr>
<tr>
<td>RCV-SCAD-AIC</td>
<td>0.84</td>
<td>2.68</td>
<td>6.62</td>
<td>17.0</td>
<td>60.7</td>
<td>171.8</td>
<td>491.7</td>
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<tr>
<td>RCV-SCAD-BIC</td>
<td>0.87</td>
<td>2.49</td>
<td>5.57</td>
<td>12.5</td>
<td>34.8</td>
<td>93.4</td>
<td>233.1</td>
<td>–</td>
</tr>
<tr>
<td>$\Omega_4$ Sample Covariance</td>
<td>3.23</td>
<td>37.36</td>
<td>537.50</td>
<td>–</td>
<td>–</td>
<td>–</td>
<td>–</td>
<td>–</td>
</tr>
<tr>
<td>PNL</td>
<td>1.52</td>
<td>6.12</td>
<td>17.69</td>
<td>–</td>
<td>–</td>
<td>–</td>
<td>–</td>
<td>–</td>
</tr>
<tr>
<td>CLIME</td>
<td>38.25</td>
<td>213.24</td>
<td>710.07</td>
<td>1439.0</td>
<td>4015.2</td>
<td>10562.6</td>
<td>–</td>
<td>–</td>
</tr>
<tr>
<td>BL</td>
<td>3.71</td>
<td>28.60</td>
<td>107.97</td>
<td>341.2</td>
<td>1439.0</td>
<td>4015.2</td>
<td>10562.6</td>
<td>–</td>
</tr>
<tr>
<td>Naive-LASSO-AIC</td>
<td>1.68</td>
<td>6.65</td>
<td>17.77</td>
<td>48.5</td>
<td>187.5</td>
<td>528.5</td>
<td>1514.7</td>
<td>–</td>
</tr>
<tr>
<td>Naive-LASSO-BIC</td>
<td>1.10</td>
<td>3.42</td>
<td>7.78</td>
<td>17.0</td>
<td>47.7</td>
<td>103.4</td>
<td>222.3</td>
<td>–</td>
</tr>
<tr>
<td>Naive-SCAD-AIC</td>
<td>1.16</td>
<td>4.07</td>
<td>10.30</td>
<td>28.3</td>
<td>121.1</td>
<td>429.2</td>
<td>1710.4</td>
<td>–</td>
</tr>
<tr>
<td>Naive-SCAD-BIC</td>
<td>0.79</td>
<td>2.28</td>
<td>5.23</td>
<td>12.0</td>
<td>37.9</td>
<td>93.4</td>
<td>233.1</td>
<td>–</td>
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<tr>
<td>RCV-LASSO-AIC</td>
<td>1.21</td>
<td>3.68</td>
<td>8.40</td>
<td>18.8</td>
<td>57.1</td>
<td>134.4</td>
<td>323.9</td>
<td>–</td>
</tr>
<tr>
<td>RCV-LASSO-BIC</td>
<td>1.29</td>
<td>3.49</td>
<td>7.58</td>
<td>15.4</td>
<td>41.5</td>
<td>90.3</td>
<td>203.4</td>
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</tr>
<tr>
<td>RCV-SCAD-AIC</td>
<td>0.95</td>
<td>2.90</td>
<td>6.95</td>
<td>17.2</td>
<td>63.7</td>
<td>179.0</td>
<td>515.5</td>
<td>–</td>
</tr>
<tr>
<td>RCV-SCAD-BIC</td>
<td>0.99</td>
<td>2.68</td>
<td>6.00</td>
<td>12.8</td>
<td>37.2</td>
<td>84.5</td>
<td>194.1</td>
<td>–</td>
</tr>
</tbody>
</table>

its Cholesky factor $L$ are both with banded structure, and ‘BL’ method is designed for such a setting. Except for ‘BL’ method, RCV estimator is almost superior to all other estimators. For the naive method, the BIC tuning parameter selector results in significantly smaller quadratic loss than the AIC. This is because the AIC method tends to include more irrelevant variables, which are likely the spuriously correlated predictors. In general, for the naive method and RCV method, the SCAD also has smaller loss than the LASSO for most cases in Table 1. Comparing Naive-LASSO/SCAD-AIC with RCV-LASSO/SCAD-AIC, we can see that the RCV method can effectively eliminate the effect due to spurious correlation.

Comparing the top and bottom panels of Table 1, we can find that the BL method performs much worse for $\Omega_2$ than for $\Omega_1$, while other methods perform similarly for both $\Omega_1$ and $\Omega_2$. This is also expected since the Cholesky factor $L$ in $\Omega_2$ does not have regular sparse structure, and thus the banded structure of $\Omega_2$ cannot be guaranteed. This directly leads to the worse performance of BL method for $\Omega_2$. This also implies that the sparsity pattern has less influence on the results from methods other than BL.

Table 2 presents the results for $\Omega_3$ and $\Omega_4$. We can see from Table 2 that the BL method performs poorly for these two precision matrices. Overall patterns in Table 2 are similar to those in the bottom panel of Table 1. For cases in Table 2, the RCV estimators provide the best results, and can effectively improve the estimation accuracy of the naive method when $p > n$.

### 3.2. Real data analysis

In this section, we illustrate the proposed procedures by an empirical analysis. Portfolio allocation is of great interest in financial econometrics and quantitative finance. In the Markowitz’s portfolio theory (Markowitz, 1952), it considers the portfolio allocation with excess returns as an optimization problem

$$
\min \frac{1}{2} w^T \Sigma w
$$

subject to $w^T \mu \geq \mu_0$.

(3.1)

where $\mu = (\mu_1, \ldots, \mu_p)^T$ is the mean vector of excess returns of the $p$ assets, $\Sigma$ is the covariance matrix of assets returns, $\mu_0$ is the expected return. The optimal solution $(\mu^*, \sigma^*)$, the expected returns and the stand deviation, constitutes an efficient frontier (Markowitz, 1952), which theoretically clarifies that higher expected returns always comes with higher risks. Consequently, Sharpe (1966) and Sharpe (1994) proposed using the ratio of the return to the risk as a measure of portfolio allocation. The ratio is called the Sharpe ratio and represents the portfolio return per unit risk. Portfolio optimization still is an active research topic (Cai et al., 2019; Ao et al., 2018).

In this example, we collect monthly excess returns for stocks in S&P 500 index that have complete records from January 1980 to December 2012. The data, collected from the Center for Research in Security Prices (CRSP), contain the returns of 202 stocks with a time span of 396 months. We avoid the period of the subprime crisis starting in 2008 since the market performance was totally different during that period (Fig. 3). We set the returns of the last 6, 12 or 18 months as
Fig. 3. Equally weighted monthly Returns of portfolio (Jan. 1980–Dec. 2012).

the testing data and the remaining data as the training data to study the asset allocation. Denote the training data of the period from Jan. 1980 to Dec. 2006 by \( X^{(1)} = \{X_{ij}^{(1)}\}, i = 1, \ldots, 324, j = 1, \ldots, 202 \), and the corresponding testing data by \( Y^{(1)} = \{Y_{ij}^{(1)}\}, i = 1, \ldots, 12, j = 1, \ldots, 202 \).

It is well known that the optimal weights (i.e., the minimizer of (3.1)) is

\[
\mathbf{w}_{opt} = \mu_0 \Sigma^{-1} \mu = \mu_0 \Omega \mu,
\]

(3.2)

where \( \Omega = \Sigma^{-1} \). To illustrate the proposed procedure, we consider the annualized return rates (ARR) at 10\%, 20\% and 30\%, respectively. By using the theory of fixed income securities, the corresponding monthly return rates are 0.8\%, 1.53\% and 2.21\% respectively. To obtain the optimal weights, \( \Omega \) and \( \mu \) should be estimated through the training data. The mean vector \( \mu \) of the assets returns is estimated by the sample mean vector

\[
\hat{\mu} = \left( \hat{\mu}_1, \ldots, \hat{\mu}_p \right)^T = \left( n^{-1} \sum_{i=1}^{n} X_{1i}, \ldots, n^{-1} \sum_{i=1}^{n} X_{pi} \right)^T.
\]

The inverse of the sample covariance matrix \( S^{-1} \) does not perform very well due to the small samplesize. Thus, we apply the naive method and RCV method to estimate the portfolio allocation weights through the training data and to investigate the returns of the newly constructed portfolios through the testing data. The results are presented as the returns, the risks and the return–risk ratios. This ratio can be considered as the measure of excess return per unit of the investment risk, which is the realized Sharpe ratio (Sharpe, 1966, 1994).

For the purpose of comparison, we use the generalized inverse of the sample covariance matrix as the benchmark, and consider other estimates: (1) the \( \ell_1 \)-regularization. The tuning parameter are chosen by AIC or BIC; (2) The regularization with SCAD penalty. The tuning parameter is chosen by AIC or BIC; (3) The BL estimator (Bickel and Levina, 2008b). All regularization estimators are improved by RCV technique.

We first focus on the comparison between the naive method and the RCV method with different regularization methods. Table 3 presents the mean return, the risk and the Sharpe ratio. It clearly shows that, in terms of the Sharperatio, the RCV methods perform better than the naive methods when these two methods are implemented with the same regularization method for most cases. The RCV with SCAD-BIC performs the best among the four combinations of RCV methods (LASSO-AIC, SCAD-AIC, LASSO-BIC and SCAD-BIC).

Next, we compare the performance of the RCV-SCAD-BIC estimation, the generalized inverse of sample covariance matrix (labeled as Sample Cov in Table 3) and the BL estimation. From Table 3, the RCV method performs the best among these four methods in terms of the ratio.

4. Discussions

In this paper, we have proposed the RCV estimation for ultrahigh dimensional precision matrix, and studied the asymptotical properties of the RCV estimator. The RCV method may be used to estimate covariance matrix too since the Cholesky decomposition has been used to estimate covariance matrix (Pourahmadi, 1999, 2000; Huang et al., 2006). It will be of interest to examine the performance of the RCV estimation procedure for covariance matrix in the future study.

As one referee points out that the sparsity of linear regression models resulted from Cholesky decomposition depends on the order of variables. For this point, we regard it as a minor weakness of the proposed estimation procedure. The
Table 3
Estimated mean return and risk of the S&P 500 portfolios analysis by using excess returns.

<table>
<thead>
<tr>
<th></th>
<th>6-month prediction with 30% ARR</th>
<th>12-month prediction with 10% ARR</th>
<th>18-month prediction with 20% ARR</th>
</tr>
</thead>
<tbody>
<tr>
<td></td>
<td>Naive RCV</td>
<td>Risk</td>
<td>Naive RCV</td>
</tr>
<tr>
<td>LASSO-AIC</td>
<td>Return</td>
<td>0.0011</td>
<td>0.0062</td>
</tr>
<tr>
<td></td>
<td>Risk</td>
<td>0.0209</td>
<td>0.0254</td>
</tr>
<tr>
<td></td>
<td>Ratio</td>
<td>−0.0517</td>
<td>0.2456</td>
</tr>
<tr>
<td>SCAD-AIC</td>
<td>Return</td>
<td>0.0023</td>
<td>0.0086</td>
</tr>
<tr>
<td></td>
<td>Risk</td>
<td>0.0263</td>
<td>0.0265</td>
</tr>
<tr>
<td></td>
<td>Ratio</td>
<td>0.0884</td>
<td>0.3260</td>
</tr>
<tr>
<td>LASSO-BIC</td>
<td>Return</td>
<td>0.0022</td>
<td>0.0064</td>
</tr>
<tr>
<td></td>
<td>Risk</td>
<td>0.0172</td>
<td>0.0144</td>
</tr>
<tr>
<td></td>
<td>Ratio</td>
<td>0.1290</td>
<td>0.4640</td>
</tr>
<tr>
<td>SCAD-BIC</td>
<td>Return</td>
<td>0.0025</td>
<td>0.0071</td>
</tr>
<tr>
<td></td>
<td>Risk</td>
<td>0.0177</td>
<td>0.0154</td>
</tr>
<tr>
<td></td>
<td>Ratio</td>
<td>0.1421</td>
<td>0.4615</td>
</tr>
<tr>
<td>Sample Cov</td>
<td>Return</td>
<td>0.0029</td>
<td>0.0048</td>
</tr>
<tr>
<td></td>
<td>Risk</td>
<td>0.0317</td>
<td>0.0198</td>
</tr>
<tr>
<td></td>
<td>Ratio</td>
<td>0.0902</td>
<td>0.2450</td>
</tr>
<tr>
<td>BL</td>
<td>Return</td>
<td>0.0024</td>
<td>0.0017</td>
</tr>
<tr>
<td></td>
<td>Risk</td>
<td>0.0104</td>
<td>0.0098</td>
</tr>
<tr>
<td></td>
<td>Ratio</td>
<td>0.2302</td>
<td>0.1750</td>
</tr>
</tbody>
</table>

reason is that banded structure on covariance matrix or precision matrix has been assumed in the literature and the banded structure also depends on the order of variables. In addition, we can deal with this problem by sorting the variables according to the magnitude of their variance from the smallest to the largest in practice. This is implemented in our real data analysis.

Acknowledgments

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Appendix

For ease of presentation, let $y$ the $t$-column of $X$ defined in Section 2.3, and $Z = X_t$. Rewrite Eq. (2.3) in the matrix form:

$$y = Z\beta + \epsilon,$$

(A.1)

where we also compress the subscript in $\beta$. For an index set $\mathcal{M}$, denote $Z_M = X_{\mathcal{M},t}$, and $\phi_{\min}(m, t)$ and $\phi_{\max}(m, t)$ stand for $\phi_{\min}(m, t)$ and $\phi_{\max}(m, t)$, defined in Section 2.3, respectively.

Define the event $\mathcal{E}_{n,t} = \{\max_{1 \leq j < t} |Z_j^T \epsilon| \leq c \sqrt{n \log p}\}$, where $Z_j$ is the $j$th column of the sample matrix $Z$ and $c$ is a positive constant. To prove Theorem 1, we need the following two lemmas.

Lemma A.1. Suppose Condition (A1) hold. If $\log p/n = O(1)$, under model (A.1) it follows that $P(\mathcal{E}_{n,t}) \rightarrow 1$.

Proof. By using the exponential tail condition (A1), we have, for any $s \geq 1$,

$$P\left(|Z_0 \epsilon_t| \geq s\right) \leq P\left(|Z_0| \geq s^{1/2}\right) + P\left(|\epsilon_t| \geq s^{1/2}\right) \leq E\exp(Z_0^2) \exp(-s) + E\exp(\epsilon_t^2) \exp(-s) \leq c_0 \exp(-s).$$

(A.2)

By inequality (A.2) and using the integration by parts, there exists a constant $c_1 > 0$ such that $E\exp(2^{-1} |Z_k \epsilon_t|) \leq 2 + 2c_1$. For any $m \geq 2$, we have

$$E|Z_0 \epsilon_t|^m \leq 2^m m! E\exp\left\{2^{-1} |Z_k \epsilon_t|\right\} \leq 2^m m! (2 + 2c_1) \leq 2^{-1} (8(2 + 2c_1))^{2m-2}m!$$

(A.3)
From Bernstein's inequality (Theorem 2.2.11 in Vaart and Wellner (1996)),
\[ P(|Z_j^T e| \geq s) \leq 2 \exp \left\{ - \frac{s^2}{4(nL_1 + s)} \right\}, \] (A.4)
where \( L_1 = 4(2 + 2c_1) \). Thus, it follows that
\[
P \left( \max_{1 \leq i \leq p} |Z_i^T e| \geq c_2 \sqrt{n \log p} \right) \leq (t - 1) P \left( |Z_i^T e| \geq c_2 \sqrt{n \log p} \right)
\leq 2 \exp \left\{ \log(t - 1) \left[ 1 - \frac{1}{4(L_1 c_2^{-2} + c_2^{-1} \sqrt{\log p/n})} \right] \right\}. \] (A.5)
Since \( t < p, \log p/n = o(1) \), when \( c_2 \) is large enough, \( 4L_1 c_2^{-2} + 4c_2^{-1} \sqrt{\log p/n} < 1 \). It turns out that \( P(C_{n,t}^c) \to 0 \), that is \( P(C_{n,t}^c) \to 1 \).

**Lemma A.2.** Under Conditions in Theorem 1, it follows that
\[
\|\hat{L}^{(j)} - L\|_\infty = O_p \left( \sqrt{k_n \log p/n} \right),
\|\hat{D}^{(j)} - D\|_\infty = O_p \left( n^{-1/2} \log^{1/2} p \right),
\]
for \( j = 1, 2 \).

**Proof.** It is sufficient to show that (A.6) holds for \( j = 2 \). For model (A.1), we obtain by data splitting that \( y_t^{(j)}, Z_t^{(j)} \) and \( \epsilon_t^{(j)} \), \( j = 1 \) and 2, emphasize their dependence on \( t \). Consider the linear model
\[
y_t^{(2)} = Z_t^{(2)} \beta_t + \epsilon_t^{(2)}.
\] (A.7)
By their definitions, the \( t \)th row of \( L \) is \(-\beta_t^T\), and the \( t \)th diagonal element of \( D \) is \( \sigma_t^2 \), the error variance in (A.7). Let \( \mathcal{D}_t \) is the index of significant predictors that are selected from \( X_1, \ldots, X_{t-1} \) based on \( Z_t^{(1)} \), and consider
\[
y_t^{(2)} = Z_t^{(2)} \beta_{\mathcal{D}_t} + \epsilon_t^{(2)}.
\] (A.8)
Let \( \hat{\beta}_{\mathcal{D}_t}^{(2)} \) be the least squares estimate from (A.8), and \( \hat{\sigma}_t^2 \) is the mean squares errors from (A.8). We set the \( t \)-the diagonal element of \( \hat{D} \) to be \( \hat{\sigma}_t^2 \), and estimate \( \beta_{\mathcal{D}_t} \) by \( \hat{\beta}_{\mathcal{D}_t} \), and estimate \( \beta_{\mathcal{D}_t} \) by \( 0 \). To show (A.6), we establish the concentration inequality for \( \hat{\beta}_{\mathcal{D}_t}^{(2)} \) and \( \hat{\sigma}_t^2 \). To this end, we first derive the following tail probability
\[
P \left( \left\| \hat{\beta}_{\mathcal{D}_t}^{(2)} - \beta_{\mathcal{D}_t} \right\|_\infty \geq c_1 \sqrt{\log p/n} \right)
\]
for some constant \( c_1 \), and
\[
P \left( \max_{1 \leq i \leq p} |\hat{\sigma}_t^2 - \sigma_t^2| \geq c_2 \sqrt{\log p/n} \right)
\]
for some constant \( c_2 \).

By the assumption \( P(\mathcal{D}_t \supseteq \mathcal{D}^*) = 1 \) in Theorem 1 and using the proof of (A.5), it follows that
\[
P \left( \max_{k \in \mathcal{D}_t} \left\| Z_k^{(2)T} \epsilon_t^{(2)} \right\| \geq c_2 \sqrt{n \log k_n} \right) \leq 2 \exp \left\{ \log k_n \left[ 1 - \frac{1}{4(L_1 c_2^{-2} + c_2^{-1} \sqrt{\log k_n/n})} \right] \right\} \] (A.11)
When \( c \) is large enough, \( 4(L_1 c_2^{-2} + c_2^{-1} \sqrt{\log k_n/n}) < 1 \), and the exponential part of the last equation in (A.5) is negative. Note that
\[
\hat{\beta}_{\mathcal{D}_t}^{(2)} = \beta_{\mathcal{D}_t} + (Z_{\mathcal{D}_t}^{(2)T} Z_{\mathcal{D}_t}^{(2)})^{-1} Z_{\mathcal{D}_t}^{(2)T} \epsilon_t^{(2)}.
\] (A.12)
Then
\[
P \left( \left\| \hat{\beta}_{\mathcal{D}_t}^{(2)} - \beta_{\mathcal{D}_t} \right\|_\infty \geq c_1 \sqrt{\log p/n} \right)
= P \left( \left\| \frac{2}{n} Z_{\mathcal{D}_t}^{(2)T} Z_{\mathcal{D}_t}^{(2)} \epsilon_t^{(2)} \right\|_\infty \geq c_1 \sqrt{\log p/n} \right)
\]
Let \( A = (a_1, \ldots, a_p)^T \) be a \( p \)-by-\( p \) positive definite matrix and \( b \) be a \( p \)-dimension vector. Then
\[
\|Ab\|_\infty = \max_{k \in \{1, \ldots, p\}} |a_k^T b| \leq \|b\|_\infty \max_k \sum_{i=1}^{p} |a_{ki}| \leq \|b\|_\infty \|A\|_{(1,1)} \leq \sqrt{p} \|b\|_\infty \|A\|_.$$
By using Condition (A2) and (A.11), we have
\[
P\left( \left\| \tilde{\beta}_{\mathcal{D}_{1t}}^{(2)} - \beta_{\mathcal{D}_{1t}}^{(2)} \right\|_\infty \geq c_1 \sqrt{k_n \log k_n / n} \right) \\
\leq P\left( \frac{2}{n} \left\| Z_{\mathcal{D}_{1t}}^{(2)} \right\|_2^{2} \left\| Z_{\mathcal{D}_{1t}}^{(2)} e_t^{(2)} \right\|_\infty \geq c_1 k_n^{-1/2} \sqrt{k_n \log k_n / n} \right) \\
\leq P\left( \frac{2}{n} \left\| Z_{\mathcal{D}_{1t}}^{(2)} e_t^{(2)} \right\|_\infty \geq \varepsilon_0 c_1 \sqrt{\log k_n / n} \right) \\
= \mathbb{E} \left\{ \left\| Z_{\mathcal{D}_{1t}}^{(2)} e_t^{(2)} \right\|_\infty \geq \varepsilon_0 c_1 \sqrt{n \log k_n} \mathbf{x}_1^{(1)} \right\} \\
\leq 2 \exp \left\{ \log k_n \left[ 1 - \frac{1}{4(1 c_1^{-2} \varepsilon_0 + c_1^{-1} \varepsilon_0^{-1} \sqrt{\log k_n / n})} \right] \right\}.
\]
(A.13)

Since $k_n^2 = o(n)$, when $c_1$ is large enough, the right hand side of the above inequality goes to 0. Finally we have
\[
P\left( \left\| \mathbf{L}^{(2)} - \mathbf{L} \right\|_\infty \geq c_1 \sqrt{k_n \log k_n \log p / n} \right) \\
= \mathbb{P} \left( \max_{t \in \{1, \ldots, p\}} \left\| \tilde{\beta}_{\mathcal{D}_{1t}}^{(2)} - \beta_{\mathcal{D}_{1t}}^{(2)} \right\|_\infty \geq c_1 \sqrt{k_n \log k_n \log p / n} \right) \\
\leq \mathbb{P} \left( \left\| \tilde{\beta}_{\mathcal{D}_{1t}}^{(2)} - \beta_{\mathcal{D}_{1t}}^{(2)} \right\|_\infty \geq c_1 \sqrt{k_n \log p / n} \right) \\
\leq 2 \exp \left\{ \log p \left[ 1 - \frac{1}{4(1 c_1^{-2} \varepsilon_0 + c_1^{-1} \varepsilon_0^{-1} \sqrt{\log p / n})} \right] \right\} \to 0.
\]

Notice that $\log p = O(n^{\alpha_0})$, $k_n = O(n^{1-\alpha_0/2})$, $\alpha_0 \in [0, 1)$. Thus, the first equation in (A.6) holds.

Next we consider $\| D^{(2)} - D \|_\infty$. Denote $\tilde{d}_t^{(2)}$ and $d_t$ to be the $t$th element of $D^{(2)}$ and $D$. Thus $\tilde{d}_t^{(2)} = \tilde{d}_t^2$, the mean squared error of the least squares fit of model (A.7), and $d_t = \sigma_t^2$, the error variance of model (A.7). As a result, it can be written as, for $t = 1, \ldots, p$,
\[
\tilde{d}_t^{(2)} - d_t = (n/2)^{-1} (2^{(2)} - 2^{-1} nd_t) - d_t \\
= (n/2)^{-1} (2^{(2)} - 2^{-1} nd_t) - (n/2)^{-1} (2^{(2)} - 2^{-1} nd_t) - d_t
\]
(A.14)

where $p_{D_{1t}}^{(2)}$ is the projection matrix consisting of variables indexed by $D_{1t}$. Define the events $A_{n_1} = \{ D^* \subset D_{1t} \}$ and $A_{n_2} = \{ D^* \subset D_{2t} \}$, where $D^*$ is the true significant variable set of the $j$th regression model (A.14). For convenience, we still abbreviate the subscript $j$ in the above equations.

\[
P \left( \left| \tilde{d}_t^{(2)} - d_t \right| \geq v \right) \\
\leq P \left( (n/2)^{-1} (2^{(2)} - 2^{-1} nd_t) \geq v/2 \right) \\
+ P \left( (n/2)^{-1} (2^{(2)} - 2^{-1} nd_t) \leq -v/2 \right).
\]
(A.15)

Similar to Eq. (A.4), the first term follows that
\[
P \left( (n/2)^{-1} \left\| 2^{(2)} - 2^{-1} nd_t \right\| \geq v/2 \right) \leq 2 \exp \left\{ \frac{-n^2 v^2}{16(n \nu^2 + n \nu)} \right\},
\]
(A.16)

where $\kappa_t$ is the fourth moment of $\varepsilon_t$. For the second term, we have
\[
P \left( (n/2)^{-1} \left\| 2^{(2)} - 2^{-1} nd_t \right\| \leq v/2 \right) \leq P \left( (n/2)^{-1} \left\| 2^{(2)} - 2^{-1} nd_t \right\| \leq v/2 \right).
\]
(A.17)

By using the condition (A2), it follows that
\[
P \left( \left( Z_{D_{1t}}^{(2)} \right)^T Z_{D_{1t}}^{(2)} e_t^{(2)} \geq v \right) \\
\leq P \left( \left( Z_{D_{1t}}^{(2)} \right)^T Z_{D_{1t}}^{(2)} e_t^{(2)} \geq v \right) \\
= P \left( \left\| Z_{D_{1t}}^{(2)} e_t^{(2)} \right\|_2 \geq v \right) \leq P \left( \left\| Z_{D_{1t}}^{(2)} e_t^{(2)} \right\|_2^2 \geq \lambda_0^{-1} n v \right). 
\]
(A.18)

Since $D_{1t}$ comes from another dataset $I_1$, it is independent of $(Z_{i}^{(2)}, \varepsilon_{t}^{(2)})$. Recall the definition of $k_n$. With the probability one, we have
\[
\left\| Z_{D_{1t}}^{(2)} e_t^{(2)} \right\|_2^2 = \sum_{j \in D_{1t}} (x_j^{(2)} e_t)^2 \leq k_n \max_{j \notin D_{1t}} \left\| x_j^{(2)} e_t \right\|_2^2. 
\]
(A.19)
where $\mathbf{x}_j^{(2)}$ are the observations in $I_2$ of covariate $X_j$. Since $X_j$ and $\varepsilon_i$ are uncorrelated, $\mathbb{E}X_j^{(2)T}\varepsilon_i^{(2)} = 0$. Next consider the concentration inequality for the term $|X_j^{(2)T}\varepsilon_i^{(2)}|$. Provided that $\mathbf{x}_j$ and $\varepsilon_i$ satisfy the condition (A1) and are uncorrelated, we still have the Bernstein’s inequality

$$P\left(\max_{j\in\mathcal{D}_n} \left| X_j^{(2)T}\varepsilon_i^{(2)} \right| > v_1/2 \right) \leq k_n P\left( \left| \sum_{i=1}^{n} X_j^{(2)} \varepsilon_i^{(2)} \right| > v_1/2 \right) \leq 2k_n \exp \left\{ -\frac{v_1^2}{16(nL_1 + v_1/2)} \right\},$$

where $L_1$ is similarly defined in the lemma A.1. The last inequality holds due to Eq. (A.3) and (A.4). Consequently, we obtain that

$$P\left( \left| \varepsilon_i^{(2)T} P_{\mathcal{D}_n}^{(2)} \varepsilon_i^{(2)} \right| \geq v \right) \leq 2k_n \exp \left\{ -\frac{\lambda_0 k_n^{-1}v}{4nL_1 + 4\sqrt{\lambda_0 k_n^{-1}v}} \right\} \quad \text{(A.20)}$$

Taking $v = c_3 k_n \log k_n$, when $c_3$ is large enough, the last equation in (A.20) will exponentially converge to 0. Together with (A.15), (A.16) and (A.20), the leading term $\varepsilon_i^{(2)T} P_{\mathcal{D}_n}^{(2)} \varepsilon_i^{(2)}$ dominates the residual term $\varepsilon_i^{(2)T} P_{\mathcal{D}_n}^{(2)} \varepsilon_i^{(2)}$. Furthermore, we apply the similar arguments to study the diagonal matrix $\hat{D}$. Then under event $\mathcal{A}_n = \mathcal{A}_{n1} \cap \mathcal{A}_{n2}$, taking $v = c_4 \sqrt{\log p/n}$, provided $k_n = o(n^{1/2})$, we get that

$$P\left( \max_{t\in\{1, \ldots, p\}} \left| \hat{d}_t^{(2)} - d_t \right| \geq c_4 \sqrt{\log p/n} \right) \leq 2p \exp \left\{ -\frac{c_4^2 n \log p}{16(L_2 k_n + c_4 \sqrt{n \log p})} \right\} + 2pk_n \exp \left\{ -\frac{c_4^2 \lambda_0 k_n^{-1} n \sqrt{n \log p}}{4L_1 n + 4(c_4 \lambda_0 k_n^{-1} n \sqrt{n \log p})^{1/2}} \right\} \quad \text{(A.21)}$$

Since $k_n = O((n/\log p)^{1/(3(1+\alpha))})$, $\alpha > 0$ and $\log p = n^\nu$, $0 \leq \nu < 1$, the last two terms in (A.21) go to 0. Therefore,

$$\|\hat{D}^{(2)} - D\|_\infty = \max_{t\in\{1, \ldots, p\}} \left| \hat{d}_t^{(2)} - d_t \right| = O_p \left( n^{-1} \log p/n \right)^{1/2}. \quad \text{(A.22)}$$

**Proof of Theorem 1.** By assumption $0 < \epsilon_0^{-1} \leq \lambda_{\min}(\Omega) \leq \lambda_{\max}(\Omega) \leq \epsilon_0$, it follows that $\|L\| = \|D\| = O(1)$. Since $D$ is a diagonal matrix and RCV estimator $\hat{D} = (\hat{D}^{(1)} + \hat{D}^{(2)})/2$, it follows that

$$\|\hat{D} - D\|_\infty = O_p(\sqrt{\log p/n}). \quad \text{(A.23)}$$

Because of $\|\hat{d}_t^{-1} - d_t^{-1}\| = |\hat{d}_t - d_t| / |d_t d_t^{-1}|$ and $E\hat{d}_t = d_t$, $\|\hat{D}^{-1} - D^{-1}\| \asymp \|\hat{D} - D\|$. Using Hölder’s inequality for matrix norm, $\|A\| \leq \sqrt{\|A\|_{(1, 1)} \|A\|_{(\infty, \infty)}}$ and Lemma A.2., it follows that $\|\hat{L} - L\| \leq k_n \|\hat{L} - L\|_\infty = O_p \left( k_n \sqrt{k_n \log p/n} \right)$. Use an inequality in Bickel and Levina (2008a),

$$\|A_1 A_2 A_3 - B_1 B_2 B_3\| \leq \sum_{j=1}^{3} \|A_j - B_j\| \prod_{k \neq j} \|B_k\| + \sum_{j=1}^{3} \|B_j\| \prod_{k \neq j} \|A_k - B_k\| + \sum_{j=1}^{3} \|A_j - B_j\|. \quad \text{(A.24)}$$

Plugging $A_1 = A_1^{(3)} = \hat{L}$, $A_2 = \hat{D}^{-1}$, and $B_1 = B_1^{(3)} = L$, $B_2 = D^{-1}$ into (A.24). By Assumption (A3), $\|L\| = O(1)$ and $\|D^{-1}\| = O(1)$. Thus, $\|\hat{\Omega} - \Omega\| = \|\hat{L}^{-1} \hat{D}^{-1} \hat{L} - L^{-1} D^{-1} L\| = O_p(k_n \sqrt{k_n \log p/n})$. The proof of Theorem 1 is completed.

**References**


