We propose a Kronecker product model for correlation or covariance matrices in the large dimension case. The number of parameters of the model increases logarithmically with the dimension of the matrix. We propose a minimum distance (MD) estimator based on a log-linear property of the model, as well as a one-step estimator, which is a one-step approximation to the quasi-maximum likelihood estimator (QMLE). We establish the rate of convergence and a central limit theorem (CLT) for our estimators in the large dimensional case. A specification test and tools for Kronecker product model selection and inference are provided. In a Monte Carlo study where a Kronecker product model is correctly specified, our estimators exhibit superior performance. In an empirical application to portfolio choice for S&P500 daily returns, we demonstrate that certain Kronecker product models are good approximations to the general covariance matrix.
Estimation of a Multiplicative Correlation Structure in the Large Dimensional Case

Christian M. Hafner† Université catholique de Louvain
Oliver B. Linton‡ University of Cambridge
Haihan Tang§ Fudan University

September 28, 2018

Abstract

We propose a Kronecker product model for correlation or covariance matrices in the large dimension case. The number of parameters of the model increases logarithmically with the dimension of the matrix. We propose a minimum distance (MD) estimator based on a log-linear property of the model, as well as a one-step estimator, which is a one-step approximation to the quasi-maximum likelihood estimator (QMLE). We establish the rate of convergence and a central limit theorem (CLT) for our estimators in the large dimensional case. A specification test and tools for Kronecker product model selection and inference are provided. In an Monte Carlo study where a Kronecker product model is correctly specified, our estimators exhibit superior performance. In an empirical application to portfolio choice for S&P500 daily returns, we demonstrate that certain Kronecker product models are good approximations to the general covariance matrix.

Some key words: Correlation matrix; Kronecker product; Matrix logarithm; Multiway array data; Portfolio choice; Sparsity

JEL subject classification: C55, C58, G11

*We would like to thank Heather Battey, Michael I. Gil’, Liudas Giraitis, Bowei Guo, Xumin He, Chen Huang, Anders Bredahl Kock, Alexei Onatskiy, Qi-Man Shao, Richard Smith, Chen Wang, Tengyao Wang, Tom Wansbeck, Jianbin Wu, Qiwei Yao for useful comments. The paper has been presented at University of Köln, University of Messina, University of Bonn, University of Santander, SoFie 2017 in New York, and CFE’16 in Sevilla, and we are grateful for discussions and comments of seminar participants.

†Institut de statistique, biostatistique et sciences actuarielles, and CORE, Université catholique de Louvain, Louvain-la-Neuve, Belgium. Email: christian.hafner@uclouvain.be

‡Faculty of Economics, Austin Robinson Building, Sidgwick Avenue, Cambridge, CB3 9DD. Email: obl20@cam.ac.uk. Thanks to the ERC for financial support.

§Corresponding author. Fanhai International School of Finance and School of Economics, Fudan University. Email: hhtang@fudan.edu.cn.
1 Introduction

Covariance and correlation matrices are of great importance in many fields. In finance, they are a key element in portfolio choice and risk management. In psychology, scholars have long assumed that some observed variables are related to the key unobserved traits through a factor model, and then use the covariance matrix of the observed variables to deduce properties of the latent traits. \cite{??} is a classic statistical reference that studies the estimation of covariance matrices and hypotheses testing about them in the low dimensional case (i.e., the dimension of the covariance matrix, \( n \), is small compared with the sample size \( T \)).

More recent work has considered the case where \( n \) is large along with \( T \). This is because many datasets now used are large. For instance, as finance theory suggests that one should choose a well-diversified portfolio that perforce includes a large number of assets with non-zero weights, investors now consider many securities when forming a portfolio. The listed company Knight Capital Group claims to make markets in thousands of securities worldwide, and is constantly updating its inventories/portfolio weights to optimize its positions. If \( n/T \) is not negligible when compared to zero but still less than one, we call this the large dimensional case in this article. (We reserve the phrase "the high dimensional case" for \( n > T \).) The correct theoretical framework to study the large dimensional case is to use the joint asymptotics (i.e., both \( n \) and \( T \) diverge to infinity simultaneously albeit subject to some restriction on their relative growth rate), not the usual asymptotics (i.e., \( n \) fixed, \( T \) tends to infinity alone). Thus, standard statistical methods under the usual asymptotic framework, such as principal component analysis (PCA) and canonical-correlation analysis (CCA), do not directly generalise to the large dimension case; applications to, say, portfolio choice, face considerable difficulties (see \cite{??}).

There are many new methodological approaches for the large dimensional case, for example \cite{??}, \cite{??}, \cite{??}, \cite{??}, \cite{??}, \cite{??}, \cite{??}, \cite{??}. \cite{??} gave an excellent account of the recent developments in the theory and practice of estimating large dimensional covariance matrices. Generally speaking, the approach is either to impose some sparsity on the covariance matrix, meaning that many elements of the covariance matrix are assumed to be zero or small, thereby reducing the number of parameters in a model for the covariance matrix to be estimated, or to use some device, such as shrinkage or a factor model, to reduce dimension. Most of this literature assumes i.i.d. data.

We consider a parametric model for the covariance or correlation matrix - the Kronecker product model. For a real symmetric, positive definite \( n \times n \) matrix \( \Delta \), a Kronecker product model is a family of \( n \times n \) matrices \( \{ \Delta^* \} \), each of which has the following structure:

\[
\Delta^* = \Delta_1^* \otimes \Delta_2^* \otimes \cdots \otimes \Delta_v^*,
\]

(1.1)

where \( \Delta_j^* \) is an \( n_j \times n_j \) dimensional real symmetric, positive definite sub-matrix such that \( n = n_1 \times \cdots \times n_v \). We require that \( n_j \in \mathbb{Z} \) and \( n_j \geq 2 \) for all \( j \); the \( \{ n_j \}_j^v \) need not be distinct. We suppose that \( \Delta \) is the covariance or correlation matrix of an observable series with sample size \( T \) and \( \{ \Delta^* \} \) is a model for \( \Delta \).

We study the Kronecker product model in the large dimensional case. Since \( n \) tends to infinity in the joint asymptotics, there are two main cases: (1) \( n_j \to \infty \) for \( j = 1, \ldots, v \) and \( v \) is fixed; (2) \( \{ n_j \}_j^v \) are all fixed and \( v \to \infty \). We shall study case (2) in detail because of its dimensionality reduction property. In this case, the number of parameters of a Kronecker product model grows logarithmically with \( n \). In particular, we show that a Kronecker product model induces a type of sparsity on the covariance or correlation
matrix: The logarithm of a Kronecker product model has many zero elements, so that sparsity is explicitly imposed on the logarithm of the covariance or correlation matrix - we call this log sparsity.

The Kronecker product model has a number of intrinsic advantages for applications. The eigenvalues of a Kronecker product are products of the eigenvalues of its sub-matrices, which in the simplest case are obtainable in closed form. Compared with strict factor models whose eigenvalues have a spikedness property, Johnstone and Onatski (2018), the Kronecker model has more flexibility in the large dimensional case. The inverse covariance matrix, its determinant, and other key quantities are easily obtained from the corresponding quantities of the sub-matrices, which facilitates computation and analysis.

We focus on correlation matrices rather than covariance matrices. This is partly because the asymptotic theory for the correlation matrix model nests that for the covariance matrix model, and partly because this will allow us to adopt a more flexible approach to approximating a general covariance matrix: we can allow the diagonal elements of the covariance matrix to be unrestricted (and they can be estimated by other well-understood methods). In practice, fitting a correlation matrix with a Kronecker product model tends to perform better than doing so for its corresponding covariance matrix. To avoid confusion, we would like to remark that if a Kronecker product model is correctly specified for a correlation matrix, its corresponding covariance matrix need not have a Kronecker product structure, and vice versa. In other words, log sparsity on a correlation matrix does not necessarily imply that its corresponding covariance matrix has log sparsity, and vice versa.

We show that the logarithm of a Kronecker product model is linear in its unknown parameters, and use this as a basis to propose a minimum distance (MD) estimator. We establish a crude upper bound rate of convergence for the MD estimator under the joint asymptotics, but we anticipate that this bound could be improved with better technology and we leave this for future research. There is a large literature on the optimal rate of convergence for estimation of high-dimensional covariance matrices and inverse (i.e., precision) matrices (see ? and ?). ? gave a nice review on those recent results. However their optimal rates are not applicable to our setting because here sparsity is not imposed on the covariance or correlation matrix, but on its logarithm. In addition, we allow for weakly dependent data, whereas the above cited papers all assume i.i.d. structures.

Next, we discuss a quasi-maximum likelihood estimator (QMLE) and a one-step estimator, which is an approximate QMLE. Under the joint asymptotics, we provide feasible central limit theorems (CLT) for the MD and one-step estimators, the latter of which is shown to achieve the parametric efficiency bound (Cramer-Rao lower bound) in the fixed \( n \) case. When choosing the weighting matrix optimally, we also show that the optimally-weighted MD and one-step estimators have the same asymptotic distribution. These CLTs are of independent interest and contribute to the literature on the large dimensional CLTs (see ?, ?, ?, ?, ?, ?, ?, and ?). Last, we give a specification test which allows us to test whether a Kronecker product model is correctly specified.

We discuss in Section 2 what kind of data give rise to a Kronecker product model. However, a given covariance or correlation matrix might not exactly correspond to a Kronecker product; in which case a Kronecker product model is misspecified, so \( \Delta \notin \{\Delta^*\} \). The previous literature on Kronecker product models did not touch this, but we shall demonstrate in this article that a Kronecker product model is a very good approximating device to general covariance or correlation matrices, by trading off variance with bias. We show that for a given Kronecker product model there always exists a
member in it that is closest to the covariance or correlation matrix in some sense to be made precise shortly.

We provide some numerical evidence that the Kronecker product model works very well when it is correctly specified. In the empirical study, we apply the Kronecker product model to S&P500 daily stock returns and compare it with ?’s linear shrinkage estimator as well as ?’s direct nonlinear shrinkage estimator. We find that the minimum variance portfolio implied by a Kronecker product model is almost as good as that constructed from ?’s linear shrinkage estimator.

1.1 Literature Review

The Kronecker product model has been previously considered in the psychometric literature (see ?, ?, ?, ? etc). In a multitrait-multimethod (MTMM) context, "multi-mode" data give rise to a Kronecker product model naturally (we will further discuss this in Section 2). ? outlined several estimation methods of the model based on the least squares and maximum likelihood principles, and provided large sample variances under assumptions of Gaussianity and fixed $n$. There is a growing Bayesian and Frequentist literature on multiway array or tensor datasets, where a Kronecker product model is commonly employed. See for example ?, ?, ?, ?, ?, ?, ?, ?, ?, ?, ?, ?, ?, and ?. In this literature, they also work with fixed $n$.

In the spatial literature, there are a number of studies that consider a Kronecker product structure for the correlation matrix of a random field, see for example ?.

This article is the first one studying Kronecker product models in the large dimensional case. Our work is also among the first exploiting log sparsity; the other is ?, although there are a few differences. First, their log sparsity is an assumption from the onset, in a similar spirit as ?, whereas our log sparsity is induced by a Kronecker product model. Second, they work with covariance matrices while we shall focus on correlation matrices. Even if we look at covariance matrices for the purpose of comparison, the Kronecker product model imposes different sparsity restrictions - as compared to those imposed by ? - on the elements of the logarithm of the covariance matrix. Third and perhaps most important, we look at different estimators.

1.2 Roadmap

The rest of the article is structured as follows. In Section 2 we lay out the Kronecker product model in detail. Section 3 introduces the MD estimator, gives its asymptotic properties, and includes a specification test, while Section 4 discusses the QMLE and one-step estimator, and provides the asymptotic properties of the one-step estimator. Section 5 examines the issue of model selection. Section 6 provides numerical evidence for the model as well as an empirical application. Section 7 concludes. Major proofs are to be found in Appendix; the remaining proofs are put in Supplementary Material (SM in what follows).

1.3 Notation

Let $A$ be an $m \times n$ matrix. Let $\text{vec} A$ denote the vector obtained by stacking the columns of $A$ one underneath the other. The commutation matrix $K_{m,n}$ is an $mn \times mn$ orthogonal matrix which translates $\text{vec} A$ to $\text{vec}(A^\top)$, i.e., $\text{vec}(A^\top) = K_{m,n} \text{vec}(A)$. If $A$ is a symmetric
$n \times n$ matrix, its $n(n-1)/2$ supradiagonal elements are redundant in the sense that they can be deduced from symmetry. If we eliminate these redundant elements from vec $A$, we obtain a new $n(n+1)/2 \times 1$ vector, denoted vec $A$. They are related by the full-column-rank, $n^2 \times n(n+1)/2$ duplication matrix $D_n$: vec $A = D_n$ vec $A$. Conversely, vec $A = D_n^\top$ vec $A$, where $D_n^\top$ is $n(n+1)/2 \times n^2$ and the Moore-Penrose generalised inverse of $D_n$. In particular, $D_n^+ = (D_n^\top D_n)^{-1} D_n^\top$ because $D_n$ is full-column rank.

For $x \in \mathbb{R}^n$, let $\|x\|_2 := \sqrt{\sum_{i=1}^n x_i^2}$ and $\|x\|_\infty := \max_{1 \leq i \leq n} |x_i|$ denote the Euclidean norm and the element-wise maximum norm, respectively. Notation diag($x$) gives an $n \times n$ diagonal matrix with the diagonal being the elements of $x$. Let maxeval($\cdot$) and mineval($\cdot$) denote the maximum and minimum eigenvalues of some real symmetric matrix, respectively. For any real $m \times n$ matrix $A = (a_{i,j})_{1 \leq i \leq m, 1 \leq j \leq n}$, let $\|A\|_F := [\text{tr}(A^T A)]^{1/2}$, $\|A\|_{\ell_2} := \max_{\|x\|_2 = 1} \|Ax\|_2 \equiv \sqrt{\text{maxeval}(A^T A)}$, and $\|A\|_{\ell_\infty} := \max_{1 \leq i \leq m} \sum_{j=1}^n |a_{i,j}|$ denote the Frobenius norm, spectral norm ($\ell_2$ operator norm) and maximum row sum matrix norm ($\ell_\infty$ operator norm) of $A$, respectively. Note that $\|\cdot\|_\infty$ can also be applied to matrix $A$, i.e., $\|A\|_\infty = \max_{1 \leq i \leq m, 1 \leq j \leq n} |a_{i,j}|$; however $\|\cdot\|_\infty$ is not a matrix norm so it does not have the submultiplicativity property of a matrix norm.

Consider two sequences of $n \times n$ real random matrices $X_T$ and $Y_T$. Notation $X_T = O_p(\|Y_T\|)$, where $\|\cdot\|$ is some matrix norm, means that for every $\varepsilon > 0$, there exist $M_\varepsilon > 0$, $N_\varepsilon > 0$ and $T_\varepsilon > 0$ such that for all $n > N_\varepsilon$ and $T > T_\varepsilon$, $P(\|X_T\|/\|Y_T\| > M_\varepsilon) < \varepsilon$. Notation $X_T = o_p(\|Y_T\|)$, where $\|\cdot\|$ is some matrix norm, means that $\|X_T\|/\|Y_T\| \xrightarrow{p} 0$ as $n,T \rightarrow \infty$ simultaneously.

Let $a \lor b$ and $a \land b$ denote $\max(a,b)$ and $\min(a,b)$, respectively. For $x \in \mathbb{R}$, let $\lfloor x \rfloor$ denote the greatest integer strictly less than $x$ and $\lceil x \rceil$ denote the smallest integer greater than or equal to $x$. Notation $\sigma(\cdot)$ defines sigma algebra.

For matrix calculus, what we adopt is called the numerator layout or Jacobian formulation; that is, the derivative of a scalar with respect to a column vector is a row vector.

## 2 The Kronecker Product Model

In this section we provide more details on the model. Consider an $n$-dimensional weakly stationary time series vector $\{y_t\}_{t=1}^T$ where $\mu := \mathbb{E}y_t$ and covariance matrix $\Sigma := \mathbb{E}((y_t - \mu)(y_t - \mu)^\top)$. Let $D$ be the diagonal matrix containing the diagonal entries of $\Sigma^{1/2}$. The correlation matrix of $y_t$, denoted $\Theta$, is $\Theta := D^{-1/2} \Sigma D^{-1/2}$. A Kronecker product model for the covariance or correlation matrix is given by $1.1$.

Let $\Delta$ denote $\Sigma$ or $\Theta$ according to the modelling purpose. Given a factorization $n = n_1 \times \cdots \times n_g$, if $\Delta \in \{\Delta^*\}$, we say that the Kronecker product model $\{\Delta^*\}$ is correctly specified. Otherwise the Kronecker product model $\{\Delta^*\}$ is misspecified. We first make clearer when a Kronecker product model is correctly specified (see $?$ and $?$ for more discussion). A Kronecker product arises when data have some multiplicative array structure. For example, suppose that $u_{j,k}$ are error terms in a panel regression model with $j = 1, \ldots, J$ and $k = 1, \ldots, K$. The interactive effects model of $?$ is that $u_{j,k} = \gamma_j f_k$, which implies that $u = \gamma \otimes f$, where $u$ is the $JK \times 1$ vector containing all the elements of $u_{j,k}$, $\gamma = (\gamma_1, \ldots, \gamma_J)^\top$, and $f = (f_1, \ldots, f_K)^\top$. Suppose that $\gamma$, $f$ are random, where $\gamma$
is independent of \( f \), and both vectors have mean zero. Then,
\[
E[uu^\top] = E[\gamma\gamma^\top] \otimes E[ff^\top].
\]

In this case the covariance matrix of \( u \) is a Kronecker product of two sub-matrices. If one dimension were time and the other firm, then this implies that the variance matrix of \( u \) is the product of a covariance matrix representing cross-sectional dependence and a covariance matrix representing the time series dependence.

We can think of our more general model (1.1) arising from *multi-index data with \( v \) multiplicative factors*. Multiway arrays are one such example as each observation has \( v \) different indices (see ?). Suppose that \( u_{i_1,i_2,\ldots,i_v} = \epsilon_{i_1,i_2,\ldots,i_v} \cdot \epsilon_{v,i_v} \), \( i_j = 1,\ldots,n_j \) for \( j = 1,\ldots,v \), or in vector form
\[
u = (u_{1,1,\ldots,1}, \ldots, u_{n_1,n_2,\ldots,n_v})^\top = \epsilon_1 \otimes \epsilon_2 \otimes \cdots \otimes \epsilon_v,
\]
where the factor \( \epsilon_j = (\epsilon_{j,1}, \ldots, \epsilon_{j,n_j})^\top \) is a mean zero random vector of length \( n_j \) with covariance matrix \( \Sigma_j \) for \( j = 1,\ldots,v \), and in addition the factors \( \epsilon_1,\ldots,\epsilon_v \) are mutually independent. Then
\[
\Sigma = E[uu^\top] = \Sigma_1 \otimes \Sigma_2 \otimes \cdots \otimes \Sigma_v.
\]
We hence see that the covariance matrix is a Kronecker product of \( v \) sub-matrices. Often such multiplicative effects may be a valid description of a covariance or correlation structure.

In earlier versions of this article we emphasized the Kronecker product model for the covariance matrix. We now use it to model the correlation matrix for the reasons mentioned in the introduction and leave the diagonal variance matrix \( D \) unrestricted. For the present discussion we assume that \( D \) (as well as \( \mu \)) is known. A Kronecker product model for \( \Theta \) is given by (1.1) with \( \Delta^* \) and \( \{\Delta_j^*\}_{j=1}^v \) replaced by \( \Theta^* \) and \( \{\Theta_j^*\}_{j=1}^v \), respectively\(^4\). Since \( \Theta \) is a correlation matrix, this implies that the diagonal entries of \( \Theta_j^* \) must be the same, although this diagonal entry could differ as \( j \) varies. Without loss of generality, we shall impose a normalisation constraint that all these \( v \) diagonal entries of \( \{\Theta_j^*\}_{j=1}^v \) are equal to one.

The Kronecker product model substantially reduces the number of parameters to estimate. In an unrestricted correlation matrix, there are \( n(n-1)/2 \) parameters, while

\(^2\)For example, in portfolio choice, one might consider, say, 250 equity portfolios constructed by intersections of 5 size groups (quintiles), 5 book-to-market equity ratio groups (quintiles) and 10 industry groups, in the spirit of ?. For example, one equity portfolio might consist of stocks which are in the smallest size quintile, largest book-to-market equity ratio quintile, and construction industry simultaneously. Then a Kronecker product model is applicable either directly to the covariance matrix of returns of these 250 equity portfolios or to the covariance matrix of the residuals after purging other common risk factors such as momentum.

\(^3\)Note that if \( n \) is not composite, one can add a vector of pseudo variables to the system until the final dimension is composite. It is recommended to add a vector of independent variables \( z_1 \sim N(0,I_k) \) such that \( (y_1^\top, z_1^\top)^\top \) is an \( n' \times 1 \) random vector with \( n' \times n' \) correlation matrix
\[
\Theta = \begin{bmatrix} \Theta_y & 0 \\ 0 & I_k \end{bmatrix}.
\]

\(^4\)The Kronecker product model is invariant under the Lie group of transformations \( \mathcal{G} \) generated by \( A_1 \otimes A_2 \otimes \cdots \otimes A_v \), where \( A_j \) are \( n_j \times n_j \) nonsingular matrices (see ?). This structure can be used to characterise the tangent space \( \mathcal{T} \) of \( \mathcal{G} \) and to define a relevant equivariance concept for restricting the class of estimators for optimality considerations.
a Kronecker product model has only $\sum_{j=1}^{v} n_j(n_j - 1)/2$ parameters. As an extreme illustration, when $n = 256$, the unrestricted correlation matrix has 32,640 parameters while a Kronecker product model of factorization $256 = 2^8$ has only 8 parameters! Since we do not restrict the diagonal matrix we have an additional $n$ variance parameters\footnote{These parameters can be estimated in a first step by standard methods} so overall the correlation matrix version of the model has more parameters and more flexibility than the covariance Kronecker model. The Kronecker part of the model induces sparsity. Specifically, although $\Theta^*$ is not sparse, the matrix $\log \Theta^*$ is sparse, where $\log$ denotes the (principal) matrix logarithm defined through the eigendecomposition of a real symmetric, positive definite matrix (see $?$ p20 for a definition). This is due to a property of Kronecker products (see Lemma $??$ in SM $??$ for derivation), that

$$\log \Theta^* = \log \Theta^*_1 \otimes I_{n_1} \otimes \cdots \otimes I_{n_v} + I_{n_1} \otimes \log \Theta^*_2 \otimes I_{n_3} \otimes \cdots \otimes I_{n_v} + \cdots + I_{n_1} \otimes I_{n_2} \otimes \cdots \otimes \log \Theta^*_v,$$

whence we see that $\log \Theta^*$ has many zero elements, generated by identity sub-matrices.

We next discuss some further identification/parameterization issues. Even after the normalisation of diagonal entries of $\Theta^*_j$ to be 1 for all $j$, the choice of parameters in $\Theta^*_j$ still warrants some discussion. As an illustration, suppose

$$\Theta^*_1 = \begin{pmatrix} 1 & 0.8 & 0.5 \\ 0.8 & 1 & 0.2 \\ 0.5 & 0.2 & 1 \end{pmatrix},$$

and then one can compute that

$$\log \Theta^*_1 = \begin{pmatrix} -0.75 & 1.18 & 0.64 \\ 1.18 & -0.55 & -0.07 \\ 0.64 & -0.07 & -0.17 \end{pmatrix}.$$

Thus there are $\sum_{j=1}^{v} n_j(n_j + 1)/2$ parameters in $\{ \log \Theta^*_j \}_{j=1}^{v}$; we call these log parameters of some member $\Theta^*$ of the Kronecker product model. On the other hand, there are only $\sum_{j=1}^{v} n_j(n_j - 1)/2$ parameters in $\{ \Theta^*_j \}_{j=1}^{v}$; we call these original parameters of some member $\Theta^*$ of the Kronecker product model. These $n_j(n_j - 1)/2$ original parameters completely pin down those $n_j(n_j + 1)/2$ log parameters. In other words, there exists a function $f : \mathbb{R}^{n_j(n_j-1)/2} \rightarrow \mathbb{R}^{n_j(n_j+1)/2}$ which maps original parameters to log parameters. However, when $n_j > 4$, $f$ does not have a closed form because when $n_j > 4$ the continuous functions which map elements of a matrix to its eigenvalues have no closed form. When $n_j = 2$, we can solve $f$ by hand (see Example 2.1).

**Example 2.1.** Suppose

$$\Theta^*_1 = \begin{pmatrix} 1 & \rho^*_1 \\ \rho^*_1 & 1 \end{pmatrix}.$$  

The eigenvalues of $\Theta^*_1$ are $1 + \rho^*_1$ and $1 - \rho^*_1$, respectively. The corresponding eigenvectors are $(1, 1)^T/\sqrt{2}$ and $(1, -1)^T/\sqrt{2}$, respectively. Therefore

$$\log \Theta^*_1 = \begin{pmatrix} \frac{1}{2} \log(1 + |\rho^*_1|^2) & \frac{1}{2} \log(1 + |\rho^*_1|^2) \\ \frac{1}{2} \log(1 - |\rho^*_1|^2) & \frac{1}{2} \log(1 - |\rho^*_1|^2) \end{pmatrix},$$

and

$$\log \Theta^*_1 = \begin{pmatrix} 1 & 0 \\ 0 & \log(1 - |\rho^*_1|^2) \end{pmatrix}.$$
Thus
\[ f(\rho) = \left[3\right] \frac{1}{2} \log(1 - \rho^2), \frac{1}{2} \log[2] \frac{1 + \rho}{1 - \rho}, \frac{1}{2} \log(1 - \rho^2) \right]. \]

To separately identify the log parameters, we can set the first diagonal entry of \( \log \Theta^*_j \) to be zero for \( j = 1, \ldots, v-1 \). This is just one possible identification scheme; see Examples ?? and ?? in SM ?? for illustration of the necessity of an identification restriction in order to separately identify log parameters. In total there are
\[ s := \sum_{j=1}^{v} n_j(n_j + 1) - (v - 1) = O(\log n) \]
(identifiable) log parameters; let \( \theta^* \in \mathbb{R}^s \) denote these. On the other hand, to separately identify the original parameters, no additional identification restriction is needed.

To estimate a correlation matrix using a Kronecker product model, there are two approaches. First, one can estimate the original parameters using Gaussian quasi-maximum likelihood estimation (see Section 4.1) and form a direct estimate of the correlation matrix. Second, one can estimate the log parameters \( \theta^* \) using the principle of minimum distance or Gaussian quasi-maximum likelihood estimation (see Section 3 and Section 4.1); form an estimate of the logarithm of the correlation matrix and then recover the estimated correlation matrix via matrix exponential. To study the theoretical properties of a Kronecker product model, the second approach is more appealing as log parameters are additive in nature while original parameters are multiplicative in nature; additive objects are easier to analyse theoretically than multiplicative objects.

3 Minimum Distance Estimator

In this section, we define a class of estimators of the log parameters \( \theta^* \) of the Kronecker product model (1.1) and give its asymptotic properties.

3.1 Estimation

We first give the main useful model property that delivers a simple estimation strategy. In Theorem A.1 in Appendix A.1 we show that there exists an \( n(n + 1)/2 \times s \) full column rank, deterministic matrix \( E \) such that
\[ \text{vech}(\log \Theta^*) = E\theta^*. \]

Given a factorization \( n = n_1 \times \cdots \times n_v \), if the Kronecker product model \( \{\Theta^*\} \) is correctly specified (i.e., \( \Theta \in \{\Theta^*\} \)), then we necessarily have \( \text{vech}(\log \Theta) = E\theta \) for some \( \theta \in \mathbb{R}^s \).

Define the sample covariance matrix and sample correlation matrix
\[ \hat{\Sigma}_T := \frac{1}{T} \sum_{t=1}^{T} (y_t - \bar{y})(y_t - \bar{x})^\top, \quad \hat{\Theta}_T := \hat{D}_T^{-1/2}\hat{\Sigma}_T\hat{D}_T^{-1/2}, \]
where
\[ \hat{D}_T := \text{diag}(\hat{\theta}_1, \ldots, \hat{\theta}_s). \]

When one adopts the second approach, the diagonal of the estimated correlation matrix need not have exact ones. In this case, it is tempting to replace these diagonal estimates with 1. However, experience shows that when one needs to invert the estimated correlation matrix, not replacing the diagonal with 1 is numerically more stable.
where \( \bar{y} := (1/T) \sum_{t=1}^{T} y_t \) and \( \hat{D}_T \) is a diagonal matrix whose diagonal elements are diagonal elements of \( \Sigma_T \).

We show in Appendix A.2 that in the Kronecker product model \( \{\Theta^*\} \) there exists a unique member, denoted by \( \Theta^0 \), which is closest to the correlation matrix \( \Theta \) in the following sense:

\[
\theta^0 = \theta^0(W) := \arg \min_{\theta^* \in \mathbb{R}^*} [ \text{vech}(\log \Theta) - E\theta^*]^\top W [\text{vech}(\log \Theta) - E\theta^*],
\]

where \( W \) is a \( n(n+1)/2 \times n(n+1)/2 \) positive definite weighting matrix which is free to choose. Clearly, \( \theta^0 \) has the closed form solution

\[
\theta^0 = (E^\top WE)^{-1} E^\top W \text{vech}(\log \Theta).
\]

The population objective function (3.1) allows us to define a minimum distance (MD) estimator:

\[
\hat{\theta}_T = \hat{\theta}_T(W) := \arg \min_{b \in \mathbb{R}^*} [\text{vech}(\log \hat{\Theta}_T) - Eb]^\top W [\text{vech}(\log \hat{\Theta}_T) - Eb],
\]

whence we can solve

\[
\hat{\theta}_T = (E^\top WE)^{-1} E^\top W \text{vech}(\log \hat{\Theta}_T).
\]

Note that \( \theta^0 \) is the quantity which one should expect \( \hat{\theta}_T \) to converge to in some probabilistic sense regardless of whether the Kronecker product model \( \{\Theta^*\} \) is correctly specified or not. When \( \{\Theta^*\} \) is correctly specified, we have \( \theta^0 = (E^\top WE)^{-1} E^\top W \text{vech}(\log \Theta) = (E^\top WE)^{-1} E^\top WE\theta = \theta \). In this case, \( \hat{\theta}_T \) is indeed estimating the elements of the correlation matrix \( \Theta \).

In practice the MD estimator is fast and easy to compute. One only needs a user-defined function in some software to generate this matrix \( E \) before one can use formula (3.3) to compute the MD estimator.\(^7\)

### 3.2 Rate of Convergence

We now introduce some assumptions for our theoretical analysis. These conditions are sufficient but far from necessary.

**Assumption 3.1.**

(i) For all \( t \), for every \( a \in \mathbb{R}^n \) with \( ||a||_2 = 1 \), there exist absolute constants \( K_1 > 1, K_2 > 0, r_1 > 0 \) such that

\[
\mathbb{E}[2] \exp[1] K_2 |a^\top y_t|^{r_1} \leq K_1.
\]

(ii) The time series \( \{y_t\}_{t=1}^{T} \) are normally distributed.

**Assumption 3.2.** There exist absolute constants \( K_3 > 0 \) and \( r_2 > 0 \) such that for all \( h \in \mathbb{N} \)

\[
\alpha(h) \leq \exp[1] - K_3 h^{r_2},
\]

where \( \alpha(h) \) is the \( \alpha \)-mixing (i.e., strong mixing) coefficients of \( y_t \), which are defined by \( \alpha(0) = 1/2 \) and for \( h \in \mathbb{N} \)

\[
\alpha(h) := 2 \sup_t \sup_{A \in \sigma(\ldots, y_{t+1}, y_t)} \sup_{B \in \sigma(\ldots, y_{t+1}, y_t, y_{t+1}, \ldots)} [1] \mathbb{P}(A \cap B) - \mathbb{P}(A) \mathbb{P}(B),
\]

where \( \sigma(\cdot) \) defines sigma algebra.

\(^7\)We have written a user-defined function in Matlab which can return \( E \) within a few seconds for fairly large \( n \), say, \( n = 625 \). It is available upon request.

\(^8\)“Absolute constants” mean constants that are independent of both \( n \) and \( T \).
Assumption 3.3.

(i) Suppose $n, T \to \infty$ simultaneously, and $n/T \to 0$.

(ii) Suppose $n, T \to \infty$ simultaneously, and

$$\frac{n^4 \omega^4 \kappa(W) (\log^5 n) \log^2 (1 + T)}{T} = o(1)$$

where $\kappa(W)$ is the condition number of $W$ for matrix inversion with respect to the spectral norm, i.e., $\kappa(W) := \|W^{-1}\|_{\ell_2} \|W\|_{\ell_2}$ and $\omega$ is defined in Assumption 3.4(ii).

(iii) Suppose $n, T \to \infty$ simultaneously, and

(a) $$\frac{n^4 \omega^4 \kappa(W) (\log^5 n) \log^2 (1 + T)}{T} = o(1),$$

(b) $$\frac{\omega^2 \log n}{n} = o(1),$$

where $\kappa(W)$ is the condition number of $W$ for matrix inversion with respect to the spectral norm, i.e., $\kappa(W) := \|W^{-1}\|_{\ell_2} \|W\|_{\ell_2}$ and $\omega$ is defined in Assumption 3.4(ii).

Assumption 3.4.

(i) The minimum eigenvalue of $\Sigma$ is bounded away from zero by an absolute constant.

(ii) Suppose

$$\min\text{val}[3] \frac{1}{n} E^T E \geq \frac{1}{\omega} > 0.$$

(At most $\omega = o(n)$.)

Assumption 3.1(i) is standard in high-dimensional theoretical work (e.g., ?, ? etc). In essence it assumes that a random vector has some exponential-type tail probability (c.f. Lemma A.2 in Appendix A.3), which allows us to invoke some concentration inequality such as a version of the Bernstein’s inequality (e.g., Theorem A.2 in Appendix A.5). The parameter $r_1$ restricts the size of the tail of $y_t$ - the smaller $r_1$, the heavier the tail. When $r_1 = 2$, $y_t$ is said to be subgaussian, when $r_1 = 1$, $y_t$ is said to be subexponential, and when $0 < r_1 < 1$, $y_t$ is said to be semieponential.

Needless to say, Assumption 3.1(i) is stronger than a finite polynomial moment assumption as it assumes the existence of some exponential moment. In a setting of independent observations, ? replaced Assumption 3.1(i) with a finite polynomial moment condition and established a rate of convergence for covariance matrices, which is slightly worse than what we have in Theorem 3.1(i) for correlation matrices. For dependent data, relaxation of the subgaussian assumption is currently an active research area in probability theory and statistics. One of the recent work is ? in which they relaxed subgaussianity to a finite polynomial moment condition in high-dimensional linear models with help of Nagaev-type inequalities. Thus Assumption 3.1(i) is likely to be relaxed when new probabilistic tools become available.
Assumption 3.1(ii), which will only be used in Section 4 for one-step estimation, implies Assumption 3.1(i) with $0 < r_1 \leq 2$. Assumption 3.1(ii) is not needed for the minimum distance estimation (Theorem 3.2 or 3.3) though.

Assumption 3.2 assumes that $\{y_t\}_{t=1}^{\infty}$ is alpha mixing (i.e., strong mixing) because $\alpha(h) \to 0$ as $h \to \infty$. In fact, we require it to decrease at an exponential rate. The bigger $r_2$ gets, the faster the decay rate and the less dependence $\{y_t\}_{t=1}^{\infty}$ exhibits. This assumption covers a wide range of time series. It is well known that both classical ARMA and GARCH processes are strong mixing with mixing coefficients which decrease to zero at an exponential rate (see Section 2.6.1 of [8] and the references therein).

Assumption 3.3(i) is for the derivation of the rate of convergence of $\hat{\Theta}_T - \Theta$ in terms of spectral norm. To establish the same rate of convergence of $\hat{\Sigma}_T - \Sigma$ in terms of spectral norm, one only needs $n/T \to c \in [0, 1]$. However, for correlation matrices, we need $n/T \to 0$. This is because a correlation matrix involves inverses of standard deviations (see Lemma A.14 in Appendix A.5).

Assumptions 3.3(ii) and (iii) are sufficient conditions for the asymptotic normality of the minimum distance estimators (Theorems 3.2 and 3.3) and of the one-step estimator (Theorem 4.2), respectively. Assumption 3.3(ii) or (iii) necessarily requires $n^4/T \to 0$. At first glance, it looks restrictive, but we would like to emphasize that this is only a sufficient condition. We will have more to say on this assumption in the discussions following Theorem 3.2.

Assumption 3.4(i) is also standard. This ensures that $\Theta$ is positive definite with the minimum eigenvalue bounded away from 0 by an absolute positive constant (see Lemma A.7(i) in Appendix A.4) and its logarithm is well-defined. Assumption 3.4(ii) postulates a lower bound for the minimum eigenvalue of $E^T E/n$; that is

$$
\frac{1}{\sqrt{\text{minvec}[\frac{1}{n} E^T E]}} = O(\sqrt{\omega}).
$$

We divide $E^T E$ by $n$ because all the non-zero elements of $E^T E$ are a multiple of $n$ (see Lemma A.2 in Appendix A.1). In words, Assumption 3.4(ii) says that the minimum eigenvalue of $E^T E/n$ is allowed to slowly drift to zero.

The following theorem establishes an upper bound on the rate of convergence for the minimum distance estimator $\hat{\Theta}_T$. To arrive at this, we restrict $r_1$ and $r_2$ such that $1/r_1 + 1/r_2 > 1$. However, this is not a necessary condition.

**Theorem 3.1.**

(i) Suppose Assumptions 3.1(i), 3.2, 3.3(i) and 3.4 hold with $1/r_1 + 1/r_2 > 1$. Then

$$
\|\hat{\Theta}_T - \Theta\|_{\ell_2} = O_p[3] \sqrt{\frac{n}{T}},
$$

where $\|\cdot\|_{\ell_2}$ is the spectral norm.

(ii) Suppose that $\|\hat{\Theta}_T - \Theta\|_{\ell_2} < A$ with probability approaching 1 for some absolute constant $A > 1$, then we have

$$
\|\log \hat{\Theta}_T - \log \Theta\|_{\ell_2} = O_p(\|\hat{\Theta}_T - \Theta\|_{\ell_2}).
$$
(iii) Suppose Assumptions 3.1(i), 3.2, 3.3(i) and 3.4 hold with $1/r_1 + 1/r_2 > 1$. Then

$$\|\hat{\Theta}_T - \Theta^0\|_2 = O_p[3]\sqrt{\frac{n\varpi\kappa(W)}{T}}.$$ 

where $\| \cdot \|_2$ is the Euclidean norm, $\kappa(W)$ is the condition number of $W$ for matrix inversion with respect to the spectral norm, i.e., $\kappa(W) := \|W^{-1}\|_2\|W\|_2$, and $\varpi$ is defined in Assumption 3.3(ii).

Proof. See Appendix A.3.

Theorem 3.1(i) provides a rate of convergence of the spectral norm of $\hat{\Theta}_T - \Theta$, which is a stepping stone for the rest of theoretical results. This rate is the same as that of $\|\Sigma_T - \Sigma\|_2$. The rate $\sqrt{n/T}$ is optimal in the sense that it cannot be improved without a further structural assumption on $\Theta$ or $\Sigma$.

Theorem 3.1(ii) is of independent interest as it relates $\|\log \hat{\Theta}_T - \log \Theta\|_2$ to $\|\hat{\Theta}_T - \Theta\|_2$. It is due to $\varpi$.

Theorem 3.1(iii) gives a rate of convergence of the minimum distance estimator $\hat{\Theta}_T$. $\Theta^0$ are log parameters of the member in the Kronecker product model, which is closest to $\Theta$ in the sense discussed earlier. For sample correlation matrix $\hat{\Theta}_T$, the rate of convergence of $\|\text{vec}(\hat{\Theta}_T - \Theta)\|_2$ is $\sqrt{n^2/T}$ (square root of a sum of $n^2$ terms each of which has a convergence rate $1/T$). Thus the minimum distance estimator $\hat{\Theta}_T$ of the Kronecker product model converges faster provided $\varpi\kappa(W)$ is not too large, in line with the principle of dimension reduction. However, given that the dimension of $\Theta^0$ is $s = O(\log n)$, one would conjecture that the optimal rate of convergence for $\hat{\Theta}_T$ should be $\sqrt{\log n/T}$. In this sense, Theorem 3.1(iii) does not demonstrate the full advantages of a Kronecker product model. Because of the severe non-linearity introduced by the matrix logarithm it is a challenging problem to prove a faster rate of convergence for $\|\hat{\Theta}_T - \Theta^0\|_2$.

3.3 Asymptotic Normality

We define $y_t$’s natural filtration $\mathcal{F}_t := \sigma(y_t, y_{t-1}, \ldots, y_1)$ and $\mathcal{F}_0 = \emptyset$.

Assumption 3.5.

(i) Suppose that $\{y_t - \mu, \mathcal{F}_t\}$ is a martingale difference sequence; that is $\mathbb{E}[y_t - \mu|\mathcal{F}_{t-1}] = 0$ for all $t = 1, \ldots, T$.

(ii) Suppose that $\{y_t'y_t' - \mathbb{E}[y_t'y_t'], \mathcal{F}_t\}$ is a martingale difference sequence; that is $\mathbb{E}[1|y_t'i'y_t'j] = \mathbb{E}[y_t'i'y_t'j]|\mathcal{F}_t = 0$ for all $i, j = 1, \ldots, n$, $t = 1, \ldots, T$.

Assumption 3.5 allows us to establish inference results within a martingale framework. Outside this martingale framework, one encounters the issue of long-run variance whenever one tries to get some inference result. This is particularly challenging in the large dimensional case and we hence shall not consider it in this article.

To derive the asymptotic normality of the minimum distance estimator, we consider two cases

(i) $\mu$ is unknown but $D$ is known;

(ii) both $\mu$ and $D$ are unknown.
We will derive the asymptotic normality of the minimum distance estimator for both these cases. We first comment on the plausibility or relevance of case (i). We present five situations/arguments to show that case (i) is relevant and these are by no means exhaustive. First, one could use higher frequency data to estimate the individual variances and thereby utilise a very large sample size. But that is not an option for estimating correlations because of the non-synchronicity problem, which is acute in the large dimensional case, Park, Hong, and Linton (2016). Second, one could have unbalanced low frequency data meaning that each firm has a long time series but they start and finish at different times such that the overlap, which is relevant for estimation of correlations, can be quite a bit smaller. In that situation one might standardise marginally using all the individual time series data and then estimate pairwise correlations using the smaller overlapping data. Third, we could have a global parametric model for $D$ and $\mu$, but a local (in time) Kronecker product model for correlations, i.e., $\Theta(u)$ varies with rescaled time $u = t/T$. In this situation, the initial estimation of $D$ and $\mu$ can be done at a faster rate than estimation of the time varying correlation $\Theta(u)$, so effectively $D$ and $\mu$ could be treated as known quantities. Fourth, case (i) reflects our two-step estimation procedure where variances are estimated first without imposing any model structure. This is a common approach in dynamic volatility model estimation such as the DCC model of $\rightarrow$ and the GO-GARCH model $\rightarrow$. Indeed, in many of the early articles in that literature standard errors for dynamic parameters of the correlation process were constructed without regard to the effect of the initial procedure. Finally, we note that theoretically estimation of $\mu$ and $D$ is well understood even in the high dimensional case, so in keeping with much practice in the literature we do not emphasize estimation of $\mu$ and $D$ again.

Define the following $n^2 \times n^2$ dimensional matrix $H$:

$$H := \int_0^1 [t(\Theta - I) + I]^{-1} \otimes [t(\Theta - I) + I]^{-1} dt.$$  

Define also the $n \times n$ and $n^2 \times n^2$ matrices, respectively:

$$\tilde{\Sigma}_T := \frac{1}{T} \sum_{t=1}^T (y_t - \mu)(y_t - \mu)^\top.$$  

$$V := \text{var}[2:\text{vec}[1\vert(y_t - \mu)(y_t - \mu)^\top = \mathbb{E}[1\vert(y_t - \mu)(y_t - \mu)^\top \otimes (y_t - \mu)(y_t - \mu)^\top - \mathbb{E}[1\vert(y_t - \mu) \otimes (y_t - \mu)]\mathbb{E}[1\vert(y_t - \mu)^\top \otimes (y_t - \mu)^\top].$$

Since $x \mapsto ([x/n], x-[x/n]n)$ is a bijection from $\{1, \ldots, n^2\}$ to $\{1, \ldots, n\} \times \{1, \ldots, n\}$, it is easy to show that the $(a, b)$th entry of $V$ is

$$V_{a,b} \equiv V_{i,j,k,\ell} = \mathbb{E}[(y_{i,j} - \mu_i)(y_{i,j} - \mu_j)(y_{i,j} - \mu_k)(y_{i,j} - \mu_\ell)] - \mathbb{E}[(y_{i,j} - \mu_i)(y_{i,j} - \mu_j)]\mathbb{E}[(y_{i,j} - \mu_k)(y_{i,j} - \mu_\ell)],$$

where $\mu_i = \mathbb{E}y_{i,j}$ (similarly for $\mu_j, \mu_k, \mu_\ell$), $a, b \in \{1, \ldots, n^2\}$ and $i, j, k, \ell \in \{1, \ldots, n\}$. In the special case of normality, $V = 2D_nD_n^\top(\Sigma \otimes \Sigma)$ (Lemma 9).

**Assumption 3.6.** Suppose that $V$ is positive definite for all $n$, with its minimum eigenvalue bounded away from zero by an absolute constant and maximum eigenvalue bounded from above by an absolute constant.

Assumption 3.6 is also a standard regularity condition. It is automatically satisfied under normality given Assumptions 3.3(i) and 3.4(i) (via Lemma A.4(vi) in Appendix A.3).
Assumption 3.6 could be relaxed to the case where the minimum (maximum) eigenvalue of \( V \) is slowly drifting towards zero (infinity) at certain rate. The proofs for Theorem 3.2 and Theorem 3.3 remain unchanged, but this rate will need to be incorporated in Assumption 3.3(ii).

### 3.3.1 When \( \mu \) Is Unknown But \( D \) Is Known

In this case, \( \hat{\Theta}_T \) simplifies into \( \hat{\Theta}_{T,D} := D^{-1/2} \hat{\Sigma}_T D^{-1/2} \). Similarly, the minimum distance estimator \( \hat{\theta}_T \) simplifies into \( \hat{\theta}_{T,D} := (E^tWE)^{-1}E^tW \text{vech}(\log \hat{\Theta}_{T,D}) \). Let \( \hat{H}_{T,D} \) denote the \( n^2 \times n^2 \) matrix

\[
\hat{H}_{T,D} := \int_0^1 [t(\hat{\Theta}_{T,D} - I) + I]^{-1} \otimes [t(\hat{\Theta}_{T,D} - I) + I]^{-1} dt.
\]

Define \( V \)'s sample analogue \( \hat{V}_T \) whose \( (a,b) \)th entry is

\[
\hat{V}_{T,a,b} \equiv \hat{V}_{T,i,j,k,l} := \frac{1}{T} \sum_{t=1}^{T} (y_{t,i} - \bar{y}_i)(y_{t,j} - \bar{y}_j)(y_{t,k} - \bar{y}_k)(y_{t,l} - \bar{y}_l) - [2] \frac{1}{T} \sum_{t=1}^{T} (y_{t,i} - \bar{y}_i)(y_{t,j} - \bar{y}_j)[2] \frac{1}{T} \sum_{t=1}^{T} (y_{t,k} - \bar{y}_k)(y_{t,l} - \bar{y}_l),
\]

where \( \bar{y}_i := \frac{1}{T} \sum_{t=1}^{T} y_{t,i} \) (similarly for \( \bar{y}_j, \bar{y}_k \) and \( \bar{y}_l \)), \( a, b \in \{1, \ldots, n^2\} \) and \( i, j, k, l \in \{1, \ldots, n\} \).

For any \( c \in \mathbb{R}^s \) define the scalar

\[
c^tJ_{DC} := c^t(E^tWE)^{-1}E^tWD_n^+H(D^{-1/2} \otimes D^{-1/2})V(D^{-1/2} \otimes D^{-1/2})HD_n^{+1}WE(E^tWE)^{-1}c.
\]

In the special case of normality, \( c^tJ_{DC} \) could be simplified into (see Example ?? in SM ?? for details): \( 2c^t(E^tWE)^{-1}E^tWD_n^+H(\Theta \otimes \Theta)HD_n^{+1}WE(E^tWE)^{-1}c \). We also define the estimate \( c^t\hat{J}_{T,DC} \):

\[
c^t\hat{J}_{T,DC}
:= c^t(E^tWE)^{-1}E^tWD_n^+\hat{H}_{T,D}(D^{-1/2} \otimes D^{-1/2})\hat{V}_T(D^{-1/2} \otimes D^{-1/2})\hat{H}_{T,D}D_n^{+1}WE(E^tWE)^{-1}c.
\]

**Theorem 3.2.** Let Assumptions 3.1(i), 3.2, 3.3(ii), 3.4, 3.5 and 3.6 be satisfied with \( 1/r_1 + 1/r_2 > 1 \). In particular we set \( r_1 = 2 \). Then

\[
\frac{\sqrt{T}c^t(\hat{\Theta}_{T,D} - \Theta^0)}{\sqrt{c^t\hat{J}_{T,DC}}} \overset{d}{\to} N(0, 1),
\]

for any \( s \times 1 \) non-zero vector \( c \) with \( ||c||_2 = 1 \).

**Proof.** See Appendix A.4

Theorem 3.2 is a version of the large-dimensional CLT, whose proof is mathematically non-trivial. To simplify the technicality, we assume subgaussianity \( (r_1 = 2) \). Because the dimension of \( \Theta^0 \) is growing with the sample size, for a CLT to make sense, we need to
transform $\hat{\theta}_{T,D} - \theta^0$ to a univariate quantity by pre-multiplying $c^\top$. The magnitudes of the elements of $c$ are not important, so we normalize it to have unit Euclidean norm. What is important is whether the elements of $c$ are zero or not. The components of $\hat{\theta}_{T,D} - \theta^0$ whose positions correspond to the non-zero elements of $c$ are effectively entering the CLT.

We contribute to the literature on the large-dimensional CLT (see [7], [7], [7], [7], [7], [7], [7] and [7]). In this strand of literature, a distinct feature is that the dimension of parameter, say, $\theta^0$, is growing with the sample size, and at the same time we do not impose sparsity on $\theta^0$. As a result, the rate of growth of dimension of parameter has to be restricted by an assumption like Assumption 3.3(ii); in particular, the dimension of parameter cannot exceed the sample size. Assumption 3.3(ii) necessarily requires $n^4/T \rightarrow 0$. In [7], [7], [7], [7], [7], [7], they require $n^{3}/T \rightarrow 0$ for establishment of a CLT for an $n$-dimensional parameter. Hence there is much room of improvement for Assumption 3.3(ii) because the dimension of $\theta^0$ is $s = O(\log n)$. The difficulty for this relaxation is again, as we had mentioned when we discussed the rate of convergence of $\hat{\theta}_T$ (Theorem 3.1), due to the severe non-linearity introduced by matrix logarithm. In this sense Assumption 3.3(ii) is only a sufficient condition; the same reasoning applies to Assumption 3.3(iii).

Our approach is different from the recent literature on high-dimensional statistics such as Lasso, where one imposes sparsity on parameter to allow its dimension to exceed the sample size.

We also give a corollary which allows us to test multiple hypotheses like $H_0: A^\top \theta^0 = a$.

**Corollary 3.1.** Let Assumptions 3.1(i), 3.2, 3.3(ii), 3.4, 3.5 and 3.6 be satisfied with $1/r_1 + 1/r_2 > 1$. In particular we set $r_1 = 2$. Given a full-column-rank $s \times k$ matrix $A$ where $k$ is finite with $\|A\|_{\ell_2} = O(\sqrt{\log n \cdot n \kappa(W)})$, we have

$$\sqrt{T}(A^\top \hat{J}_{T,D}A)^{-1/2}A^\top(\hat{\theta}_{T,D} - \theta^0) \overset{d}{\rightarrow} N[1]0, I_k.$$

**Proof.** See SM ??.

Note that the condition $\|A\|_{\ell_2} = O(\sqrt{\log n \cdot n \kappa(W)})$ is trivial because the dimension of $A$ is only of order $O(\log n) \times O(1)$. Moreover we can always rescale $A$ when carrying out hypothesis testing.

If one chooses the weighting matrix $W$ optimally, albeit infeasibly,

$$W_{D,op} = [1]D_n^+ H(D^{-1/2} \otimes D^{-1/2})V(D^{-1/2} \otimes D^{-1/2})HD_n^+T^{-1},$$

the scalar $c^\top J_{D}c$ reduces to

$$c^\top [2]E^\top[1]D_n^+ H(D^{-1/2} \otimes D^{-1/2})V(D^{-1/2} \otimes D^{-1/2})HD_n^+T^{-1}E^{-1}c.$$

Under a further assumption of normality (i.e., $V = 2D_nD_n^+(\Sigma \otimes \Sigma)$), the preceding display further simplifies to

$$c^\top[3]E^\top D_n^+H^{-1}(\Theta^{-1} \otimes \Theta^{-1})H^{-1}D_nE^{-1}c,$$

by Lemmas 11 and 14 of [7]. We shall compare the preceding display with the variance of the asymptotic distribution of the one-step estimator in Section 4.
3.3.2 When Both \( \mu \) and \( D \) Are Unknown

The case where both \( \mu \) and \( D \) are unknown is considerably more difficult. If one simply recycles the proof for the case where only \( \mu \) is unknown and replaces \( D \) with its plug-in estimator \( \hat{D}_T \), it will not work.

Let \( \hat{H}_T \) denote the \( n^2 \times n^2 \) matrix

\[
\hat{H}_T := \int_0^1 [t(\hat{\Theta}_T - I) + I]^{-1} \otimes [t(\hat{\Theta}_T - I) + I]^{-1} dt.
\]

Define the \( n^2 \times n^2 \) matrix \( P \):

\[
P := I_{n^2} - D_n D_n^T (I_n \otimes \Theta) M_d, \quad M_d := \sum_{i=1}^n (F_i \otimes F_i),
\]

where \( F_i \) is an \( n \times n \) matrix with one in its \((i, i)\)th position and zeros elsewhere. Matrix \( M_d \) is an \( n^2 \times n^2 \) diagonal matrix with diagonal elements equal to 0 or 1; the positions of 1 in the diagonal of \( M_d \) correspond to the positions of diagonal entries of an arbitrary \( n \times n \) matrix \( A \) in \( \vec{A} \). Matrix \( P \) first appeared in (4.6) of [\textsuperscript{3}]. Note that for any correlation matrix \( \Theta \), matrix \( P \) is an idempotent matrix of rank \( n^2 - n \) and has \( n \) rows of zeros. [\textsuperscript{3}] proved that

\[
\frac{\partial \vec{\Theta}}{\partial \vec{\Sigma}} = P(D^{-1/2} \otimes D^{-1/2});
\]

that is, the derivative \( \frac{\partial \vec{\Theta}}{\partial \vec{\Sigma}} \) is a function of \( \Sigma \).

For any \( c \in \mathbb{R}^s \) define the scalar \( c^T J c \) and its estimate \( c^T \hat{J}_T c \):

\[
c^T J c := c^T (E^T W E)^{-1} E^T W D_n^+ H P(D^{-1/2} \otimes D^{-1/2}) V(D^{-1/2} \otimes D^{-1/2}) P^T H D_n^{+T} W E (E^T W E)^{-1} c.
\]

\[
c^T \hat{J}_T c := c^T (E^T W E)^{-1} E^T W D_n^+ \hat{H}_T \hat{P}_T (\hat{D}_T^{-1/2} \otimes \hat{D}_T^{-1/2}) \hat{V}_T (\hat{D}_T^{-1/2} \otimes \hat{D}_T^{-1/2}) \hat{P}_T^T \hat{H}_T D_n^{+T} W E (E^T W E)^{-1} c,
\]

where \( \hat{P}_T := I_{n^2} - D_n D_n^+ (I_n \otimes \hat{\Theta}_T) M_d \).

**Assumption 3.7.**

(i) For every positive constant \( C \)

\[
\sup_{\Sigma^* : \|\Sigma - \Sigma^*\|_F \leq C \sqrt{\frac{n}{T}}} \left| \frac{\partial \vec{\Theta}}{\partial \vec{\Sigma}} \right|_{\Sigma = \Sigma^*} - P(D^{-1/2} \otimes D^{-1/2}) \leq O[3] \sqrt{\frac{n}{T}};
\]

where \( |\Sigma = \Sigma^*\) means "evaluate the argument \( \Sigma \) at \( \Sigma^* \)".

(ii) The \( s \times s \) matrix

\[
E^T W D_n^+ H P(D^{-1/2} \otimes D^{-1/2}) V(D^{-1/2} \otimes D^{-1/2}) P^T H D_n^{+T} W E
\]

has full rank \( s \) (i.e., being positive definite). Moreover,

\[
\minval_2[E^T W D_n^+ H P(D^{-1/2} \otimes D^{-1/2}) V(D^{-1/2} \otimes D^{-1/2}) P^T H D_n^{+T} W E] \geq \frac{n}{\omega} \minval_2(W).
\]
Assumption 3.7(i) characterises some sort of uniform rate of convergence in terms of spectral norm of the Jacobian matrix $\left| \frac{\partial \text{vec} \Theta}{\partial \text{vec} \Sigma} \right|_{\Sigma = \hat{\Sigma}_T}$. This type of assumption is usually made when one wants to stop Taylor expansion, say, of vec $\hat{\Theta}_T$, at first order. If one goes into the second-order expansion (a tedious route), Assumption 3.7(i) can be completely dropped at some expense of further restricting the relative growth rate between $n$ and $T$. The radius of the shrinking neighbourhood $\sqrt{n^2/T}$ is determined by the rate of convergence in terms of the Frobenius norm of the sample covariance matrix $\hat{\Sigma}_T$. The rate on the right side of Assumption 3.7(i) is chosen to be $p n/T$ because it is the rate of convergence of $\left| \frac{\partial \text{vec} \Theta}{\partial \text{vec} \Sigma} \right|_{\Sigma = \hat{\Sigma}_T}^2$ which could be easily deduced from the proof of Theorem 3.3. This rate $p n/T$ could even be relaxed to $p n^2/T$ as the part of the proof of Theorem 3.3 which requires Assumption 3.7(i) is not the "binding" part of the whole proof.

We now examine Assumption 3.7(ii). The matrix $E|WD + nHP(D_1 = 2D = 2)$ is symmetric and positive semidefinite. By Observation 7.1.8 of ?, its rank is equal to rank($E|WD + nHP$) if $(D_1 = 2D = 2)$ is positive definite. In other words, Assumption 3.7(ii) is assuming rank($E|WD + nHP$) = $s$, provided $(D_1 = 2D = 2)$ is positive definite. Even though $P$ has only rank $n^2/n$, in general the rank condition does hold except in a special case. The special case is $\Theta = P_{n \& W} = I_{n(n+1)/2}$. In this special case

$$\text{rank}(E^T WD_n^+ HP(D_1 = 2D = 2)V(D_1 = 2D = 2)P^T HD_n^T W E)$$

is symmetric and positive semidefinite. By Observation 7.1.8 of ?, its rank is equal to rank($E^T WD_n^+ HP$), if $(D_1 = 2D = 2)V(D_1 = 2D = 2)$ is positive definite. In other words, Assumption 3.7(ii) is assuming rank($E^T WD_n^+ HP$) = $s$, provided $(D_1 = 2D = 2)V(D_1 = 2D = 2)$ is positive definite. Even though $P$ has only rank $n^2/n$, in general the rank condition does hold except in a special case. The special case is $\Theta = P_{n \& W} = I_{n(n+1)/2}$. In this special case

$$\text{rank}(E^T WD_n^+ HP) = \sum_{j=1}^v n_j(n_j - 1) < s.$$

The second part of Assumption 3.7(ii) postulates a lower bound for its minimum eigenvalue. The rate mineval($W(n)/n$) is specified as such because of Assumption 3.4(ii). Other magnitudes of the rate are also possible as long as the proof of Theorem 3.3 goes through.

**Theorem 3.3.** Let Assumptions 3.1(i), 3.2, 3.3(ii), 3.4, 3.5, 3.6 and 3.7 be satisfied with $1/r_1 + 1/r_2 > 1$. In particular we set $r_1 = 2$. Then

$$\frac{\sqrt{Tc^T(\hat{\Theta}_T - \Theta^0)}}{\sqrt{c^T J_T c}} \overset{d}{\rightarrow} N(0, 1),$$

for any $s \times 1$ non-zero vector $c$ with $\|c\|_2 = 1$.

**Proof.** See SM ??.
Corollary 3.2. Let Assumptions $3.1(i), 3.2, 3.3(ii), 3.4, 3.5, 3.6$ and $3.7$ be satisfied with $1/r_1 + 1/r_2 > 1$. In particular we set $r_1 = 2$. Given a full-column-rank $s \times k$ matrix $A$ where $k$ is finite with $\|A\|_{\ell_2} = O(\sqrt{\log^2 n \cdot n^2 (W)})$, we have
\[
\sqrt{T}(A^T \hat{J}_T A)^{-1/2} A^T (\hat{\theta}_T - \theta^0) \xrightarrow{d} N[1,0, I_k].
\]
Proof. Essentially the same as that of Corollary 3.1.

The condition $\|A\|_{\ell_2} = O(\sqrt{\log^2 n \cdot n^2 (W)})$ is trivial because the dimension of $A$ is only of order $O(\log n) \times O(1)$. Moreover we can always rescale $A$ when carrying out hypothesis testing. In the case of both $\mu$ and $D$ unknown, the infeasible optimal weighting matrix will be
\[
W_{op} = [1]D_n^+ HP(D^{-1/2} \otimes D^{-1/2}) V(D^{-1/2} \otimes D^{-1/2}) P^T H D_n^{+\top}^{-1}.
\]

3.4 Specification Test

We give a specification test (also known as an over-identification test) based on the minimum distance objective function in (3.2). Suppose we want to test whether the Kronecker product model $\{\Theta^*\}$ is correctly specified given the factorization $n = n_1 \times \cdots \times n_v$. That is,
\[
H_0 : \Theta \in \{\Theta^*\} \quad (i.e., \text{vech}(\log \Theta) = E \theta), \quad H_1 : \Theta \notin \{\Theta^*\}.
\]
We first fix $n$ (and hence $v$ and $s$). Recall (3.2):
\[
\hat{\theta}_T = \hat{\theta}_T(W) := \arg \min_{b \in \mathbb{R}^s} \text{vech}(\log \hat{\Theta}_T) - Eb]^T W \text{vech}(\log \hat{\Theta}_T) - Eb] =: \arg \min_{b \in \mathbb{R}^s} g_T(b)^T W g_T(b).
\]

Theorem 3.4. Fix $n$ (and hence $v$ and $s$).

(i) Suppose $\mu$ is unknown but $D$ is known. Let Assumptions $3.1(i), 3.2, 3.4, 3.5$ and $3.6$ be satisfied with $1/r_1 + 1/r_2 > 1$. In particular we set $r_1 = 2$. Thus, under $H_0$,
\[
T g_T(D) \hat{\theta}_T D)^\top \hat{S}_{T,D}^{-1} g_T(d) \hat{\theta}_T D)^\top \xrightarrow{d} \chi^2_{n(n+1)/2-s}, \quad (3.5)
\]
where
\[
g_T(D) \hat{\theta}_T D) := \text{vech}(\log \hat{\Theta}_T) - Eb
\]
\[
\hat{S}_{T,D} := D_n^+ \hat{H}_{T,D} (D^{-1/2} \otimes D^{-1/2}) \hat{V}_T (D^{-1/2} \otimes D^{-1/2}) \hat{H}_{T,D} D_n^{+\top}.
\]

(ii) Suppose both $\mu$ and $D$ are unknown. Let Assumptions $3.1(i), 3.2, 3.4, 3.5, 3.6$ and $3.7$ be satisfied with $1/r_1 + 1/r_2 > 1$. In particular we set $r_1 = 2$. Thus, under $H_0$,
\[
T g_T(d) \hat{\theta}_T D)^\top \hat{S}_{T,D}^{-1} g_T(d) \hat{\theta}_T D)^\top \xrightarrow{d} \chi^2_{n(n+1)/2-s}, \quad (3.5)
\]
where
\[
\hat{S}_{T} := D_n^+ \hat{H}_T \hat{P}_T (\hat{D}_T^{-1/2} \otimes \hat{D}_T^{-1/2}) \hat{V}_T (\hat{D}_T^{-1/2} \otimes \hat{D}_T^{-1/2}) \hat{P}_T^\top \hat{H}_T D_n^{+\top}.
\]

Proof. See SM ??.
Note that $\hat{S}^{-1}_{T,D}$ and $\hat{S}^{-1}_T$ are the feasible versions of optimal weighting matrices $W_{D,op}$ and $W_{op}$, respectively. From Theorem 3.4, we can easily get the following result of the diagonal path asymptotics, which is more general than the sequential asymptotics but less general than the joint asymptotics (see §).

Corollary 3.3.

(i) Suppose $\mu$ is unknown but $D$ is known. Let Assumptions 3.1(i), 3.2, 3.4, 3.5 and 3.6 be satisfied with $1/r_1 + 1/r_2 > 1$. In particular we set \( r_1 = 2 \). Under $H_0$,
\[
\frac{T g_{T,n,D}(\hat{\theta}_{T,n,D})\hat{S}^{-1}_{T,n,D} d g_{T,n,D}(\hat{\theta}_{T,n,D}) - [1]^n(n+1)}{[1]^n(n+1) - 2s^{1/2}} \xrightarrow[]{d} N(0,1),
\]
where $n = n_T$ as $T \to \infty$.

(ii) Suppose both $\mu$ and $D$ are unknown. Let Assumptions 3.1(i), 3.2, 3.4, 3.5, 3.6, and 3.7 be satisfied with $1/r_1 + 1/r_2 > 1$. In particular we set \( r_1 = 2 \). Under $H_0$,
\[
\frac{T g_{T,n}(\hat{\theta}_{T,n})\hat{S}^{-1}_{T,n} d g_{T,n}(\hat{\theta}_{T,n}) - [1]^n(n+1)}{[1]^n(n+1) - 2s^{1/2}} \xrightarrow[]{d} N(0,1),
\]
where $n = n_T$ as $T \to \infty$.

Proof. See SM ??.

4 QMLE and One-Step Estimator

4.1 QMLE

In the context of Gaussian quasi-maximum likelihood estimation (QMLE), given a factorization $n = n_1 \times \cdots \times n_v$, we shall additionally assume that the Kronecker product model $\{\Theta^*\}$ is correctly specified (i.e. $\text{vech}(\log \Theta) = E\theta$). Let $\rho \in [-1,1]^{n_\rho}$ be original parameters of some member of the Kronecker product model; we have mentioned that $s_\rho = \sum_{j=1}^v n_j(n_j - 1)/2$. Given Assumption 3.5, the log likelihood function in terms of original parameters $\rho$ for a sample \( \{y_1, y_2, \ldots, y_T\} \) is given by
\[
\ell_T(\mu, D, \rho) = -\frac{Tn}{2} \log(2\pi) - \frac{T}{2} \log |D^{1/2}\Theta(\rho)D^{1/2}| - \frac{1}{2} \sum_{t=1}^{T} (y_t - \mu)^T D^{-1/2} \Theta(\rho)^{-1} D^{-1/2} (y_t - \mu).
\]
(4.1)

Write $\Omega = \Omega(\theta) := \log \Theta = \log \Theta(\rho)$. Given Assumption 3.5, the log likelihood function in terms of log parameters $\theta$ for a sample \( \{y_1, y_2, \ldots, y_T\} \) is given by
\[
\ell_T(\mu, D, \theta)
= -\frac{Tn}{2} \log(2\pi) - \frac{T}{2} \log |D^{1/2} \exp(\Omega(\theta)) D^{1/2}| - \frac{1}{2} \sum_{t=1}^{T} (y_t - \mu)^T D^{-1/2} [\exp(\Omega(\theta))]^{-1} D^{-1/2} (y_t - \mu).
\]
(4.2)

In practice, conditional on some estimates of $\mu$ and $D$, we use an iterative algorithm based on the derivatives of $\ell_T$ with respect to either $\rho$ or $\theta$ to compute the QMLE of either
\( \rho \) or \( \theta \). Theorem 4.1 below provides formulas for the derivatives of \( \ell_T \) with respect to \( \theta \). The computations required are typically not too onerous, since for example the Hessian matrix is of an order \( \log n \) by \( \log n \). See \( ? \) and \( ? \) for a discussion of estimation algorithms in the case where the data are multiway array and \( v \) is of low dimension. Nevertheless since there is quite complicated non-linearity involved in the definition of the QMLE, it is not so easy to directly analyse QMLE.

Instead we shall consider a one-step estimator that uses the minimum distance estimator in Section 3 to provide a starting value and then takes a Newton-Raphson step towards the QMLE of \( \theta \). In the fixed \( n \) it is known that the one-step estimator is equivalent to the QMLE in the sense that it shares its asymptotic distribution (\( \mathcal{N} \)).

Below, for slightly abuse of notation, we shall use \( \rho;D;\) to denote the true parameter (i.e., characterising the data generating process) as well as the generic parameter of the likelihood function; we will be more specific whenever any confusion is likely to arise.

4.2 One-Step Estimator

Here we only examine the one-step estimator when \( \mu \) is unknown but \( D \) is known. When neither \( \mu \) nor \( D \) is known, one has to differentiate (4.2) with respect to both \( \theta \) and \( D \). The analysis becomes considerably more involved and we leave it for future work. Suppose \( D \) is known, the likelihood function (4.2) reduces to

\[
\ell_{T,D}(\theta, \mu) = -\frac{T}{2} n \log(2\pi) - \frac{T}{2} \log \left| D^{1/2} \exp(\Omega(\theta)) D^{1/2} \right| + \frac{1}{2} \sum_{t=1}^{T} (y_t - \mu)^T D^{-1/2} \left[ \exp(\Omega(\theta)) \right]^{-1} D^{-1/2} (y_t - \mu). 
\]

These well-known that for any choice of \( \Sigma \) (i.e., \( D \) and \( \theta \)), the QMLE for \( \mu \) is \( \hat{\mu} \). Hence we may define

\[
\hat{\theta}_{QMLE,D} = \arg \max_{\theta} \ell_{T,D}(\theta, \hat{\mu}).
\]

**Theorem 4.1.** The \( s \times 1 \) score function of (4.3) with respect to \( \theta \) takes the following form

\[
\frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta^T} = \frac{T}{2} E^T D_n^T \int_0^1 e^{t \Omega} \otimes e^{(1-t)\Omega} dt [2] \text{vec}[2] e^{-\Omega} D^{-1/2} \hat{\Sigma} T D^{-1/2} e^{-\Omega} - e^{-\Omega}.
\]

The \( s \times s \) block of the Hessian matrix of (4.3) corresponding to \( \theta \) takes the following form

\[
\frac{\partial^2 \ell_{T,D}(\theta, \mu)}{\partial \theta \partial \theta^T} = -\frac{T}{2} E^T D_n^T \Psi[1] e^{-\Omega} D^{-1/2} \hat{\Sigma} T D^{-1/2} \otimes I_n + I_n \otimes e^{-\Omega} D^{-1/2} \hat{\Sigma} T D^{-1/2} - I_n e^{\Omega} \otimes e^{-\Omega} \Psi[1] D_n E
\]

\[
+ \frac{T}{2} (\Psi[1] \otimes E^T D_n^T) \int_0^1 P_K[1] \text{vec} e^{(1-t)\Omega} \otimes \text{vec} e^{(1-s)\Omega} ds \cdot t dt D_n E
\]

\[
+ \frac{T}{2} (\Psi[1] \otimes E^T D_n^T) \int_0^1 P_K[1] \text{vec} e^{(1-t)\Omega} \otimes I_n e^{(1-s)(1-t)\Omega} ds \cdot (1 - t) dt D_n E,
\]

\[19\]
where \( P_K := I_n \otimes K_{n,n} \otimes I_n \), \( \hat{\Sigma} \) is defined in (3.4), and
\[
\Psi_1 = \Psi_1(\theta) := \int_0^1 e^{\Omega(\theta) t} \otimes e^{(1-t)\Omega(\theta)} dt,
\]
\[
\Psi_2 = \Psi_2(\theta) := \text{vec}[1] e^{-\Omega(\theta)} D^{-1/2} \hat{\Sigma}_T D^{-1/2} e^{-\Omega(\theta)} - e^{-\Omega(\theta)}.
\]

**Proof.** See SM ??.

Since \( \mathbb{E}\Psi_2(\theta) = 0 \), where \( \theta \) denotes the true parameter, so the negative normalized expected Hessian matrix evaluated at the true parameter \( \theta \) takes the following form
\[
\Upsilon_D := \mathbb{E}[3] - \frac{1}{T} \frac{\partial^2 \ell_{T,D}(\theta, \mu)}{\partial \theta \partial \theta^\top} = \frac{1}{2} E^T D_n^\top \Psi_1(\theta)[1] e^{-\Omega(\theta)} \otimes e^{-\Omega(\theta)} \Psi_1(\theta) D_n E
\]
\[
= \frac{1}{2} E^T D_n^\top \int_0^1 \int_0^1 e^{(t+s-1)\Omega} \otimes e^{(1-t-s)\Omega} dt ds D_n E
\]
\[
= \frac{1}{2} E^T D_n^\top [3] \int_0^1 \int_0^1 \Theta^{s-1} \otimes \Theta^{1-t-s} dt ds D_n E =: \frac{1}{2} E^T D_n^\top \Sigma D_n E. \tag{4.4}
\]

It can be shown that under normality (i.e., \( V = 2D_n D_n^\top (\Sigma \otimes \Sigma) \)), \( \Upsilon_D = \mathbb{E}[1] \frac{1}{T} \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta \partial \theta^\top} \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta \partial \theta^\top} \) (see Lemma ?? in SM ??). We then propose the following one-step estimator in the spirit of ?? or ?? p2150:
\[
\hat{\theta}_{T,D} := \hat{\theta}_{T,D} - \frac{1}{T} \frac{\partial \ell_{T,D}(\hat{\theta}_{T,D}, \hat{y})}{\partial \theta^\top} / T, \tag{4.5}
\]
where \( \hat{\theta}_{T,D} := \frac{1}{2} E^T D_n^\top [1] \int_0^1 \int_0^1 \hat{\Theta}^{s-1} \otimes \hat{\Theta}^{1-t-s} dt ds D_n E =: \frac{1}{2} E^T D_n^\top \hat{\Sigma}_T D_n E \). (We show in SM ?? that \( \hat{\theta}_{T,D} \) is invertible with probability approaching 1.) We did not use the plain vanilla one-step estimator because the Hessian matrix \( \frac{\partial^2 \ell_{T,D}(\theta, \mu)}{\partial \theta \partial \theta^\top} \) is rather complicated to analyse.

### 4.3 Large Sample Properties

To provide the large sample theory for the one-step estimator \( \hat{\theta}_{T,D} \), we make the following assumption.

**Assumption 4.1.** For every positive constant \( M \) and uniformly in \( b \in \mathbb{R}^s \) with \( \|b\|_2 = 1 \),
\[
\sup_{\theta^* : \|\theta^* - \theta\|_2 \leq M} \frac{[4] \sqrt{T} b^\top [3] \frac{1}{T} \frac{\partial \ell_{T,D}(\theta^*, \hat{y})}{\partial \theta^\top} - \frac{1}{T} \frac{\partial \ell_{T,D}(\theta, \hat{y})}{\partial \theta^\top} - \Upsilon_D (\theta^* - \theta) = o_p(1).
\]

Assumption 4.1 is one of the sufficient conditions needed for the asymptotic normality of \( \hat{\theta}_{T,D} \) (Theorem 4.2). This kind of assumption is standard in the asymptotics of one-step estimators (see (5.44) of ?? p71) or of M-estimation (see (C3) of ??). Assumption 4.1 implies that \( \frac{1}{T} \frac{\partial \ell_{T,D}(\theta, \hat{y})}{\partial \theta^\top} \) is differentiable at the true parameter \( \theta \), with derivative tending to \( \Upsilon_D \) in probability. The radius of the shrinking neighbourhood \( \sqrt{n \kappa(W) / T} \) is determined by the rate of convergence of any preliminary estimator, say, \( \hat{\theta}_{T,D} \) in our case. It is possible to relax the \( o_p(1) \) on the right side of the display in Assumption 4.1 to \( o_p(\sqrt{n/(\kappa^2 \log n)}) \) by examining the proof of Theorem 4.2.
Theorem 4.2. Suppose that the Kronecker product model $\{\Theta^*\}$ is correctly specified. Let Assumptions 3.1(ii), 3.3(iii), 3.4, 3.5 and 4.1 be satisfied with $1/r_1 + 1/r_2 > 1$ and $r_1 = 2$. Then
\[
\sqrt{T}c^T(\hat{\theta}_{T,D} - \theta) \overset{d}{\to} N(0,1)
\]
for any $s \times 1$ vector $c$ with $\|c\|_2 = 1$.

Proof. See SM ??.

Theorem 4.2 is a version of the large-dimensional CLT, whose proof is mathematically non-trivial. It has the same structure as that of Theorem 3.2 or Theorem 3.3. Note that under Assumption 3.1(ii), the QMLE is actually the maximum likelihood estimator (MLE). If we replace normality (Assumption 3.1(ii)) with the subgaussian assumption (Assumption 3.1(i) with $r_1 = 2$) - that is the Gaussian likelihood is not correctly specified - although the norm consistency of $\hat{\theta}_{T,D}$ should still hold, the asymptotic variance in Theorem 4.2 needs to be changed to have a sandwich formula. Theorem 4.2 says that $\sqrt{T}c^T(\hat{\theta}_{T,D} - \theta) \overset{d}{\to} N[1]0, c^T \hat{\Sigma}_{T,D} c$. In the fixed $n$ case, this estimator achieves the parametric efficiency bound by recognising a well-known result
\[
\sqrt{T}c^T(\hat{\theta}_{T,D} - \theta) \overset{d}{\to} N[1]0, c^T \hat{\Sigma}_{T,D} c.
\]
This shows that our one-step estimator $\hat{\theta}_{T,D}$ is efficient when $D$ (the variances) is known.

By recognising that $H^{-1} = \int_0^1 e^{t \log \Theta} \otimes e^{(1-t) \log \Theta} dt = \Psi_1$ (see Lemma ?? in SM ??), we see that, when $D$ is known, under normality and correct specification of the Kronecker product model, $\hat{\theta}_{T,D}$ and the optimal minimum distance estimator $\hat{\theta}_{T,D}(W_{D,op})$ have the same asymptotic variance, i.e., $[1]1/2 E^T D_n^H H^{-1}(\Theta^{-1} \otimes \Theta^{-1}) H^{-1} D_n E^{-1}$.

We also give the following corollary which allows us to test multiple hypotheses like $H_0 : A^T \theta = a$.

Corollary 4.1. Suppose the Kronecker product model $\{\Theta^*\}$ is correctly specified. Let Assumptions 3.1(ii), 3.3(iii), 3.4, 3.5 and 4.1 be satisfied with $1/r_1 + 1/r_2 > 1$ and $r_1 = 2$. Given a full-column-rank $s \times k$ matrix $A$ where $k$ is finite with $\|A\|_\ell_2 = O(\sqrt{\log n \cdot n})$, we have
\[
\sqrt{T}(A^T \hat{\Sigma}_{T,D}^{-1} A)^{-1/2} A^T(\hat{\theta}_{T,D} - \theta) \overset{d}{\to} N[1]0, I_k.
\]

Proof. Essentially the same as that of Corollary 3.1.

The condition $\|A\|_\ell_2 = O(\sqrt{\log n \cdot n})$ is trivial because the dimension of $A$ is only of order $O(\log n) \times O(1)$. Moreover we can always rescale $A$ when carrying out hypothesis testing.

5 Model Selection

We briefly discuss the issue of model selection here. One shall not worry about this if the data are in the multi-index format with $v$ multiplicative factors. This is because in this setting a Kronecker product model is pinned down by the structure of multiway arrays - the Kronecker product model is correctly specified. This issue will pop up when one uses Kronecker product models to approximate a general covariance or correlation matrix - all
Kronecker product models are then misspecified. The rest of discussions in this section will be based on this approximation framework.

First, if one permutes the data, the performance of a given Kronecker product model is likely to change. However, based on our experience, the performance of a Kronecker product model is not that sensitive to the ordering of the data. We will illustrate this in the empirical study. Moreover, usually one fixes the ordering of the data before considering the issue of covariance matrix estimation. Thus, Kronecker product models have a second-mover advantage: the choice of a Kronecker product model depends on the ordering of the data.

Second, if one fixes the ordering of the data as well as a factorization \( n = n_1 \times \cdots \times n_v \), but permutes \( \Theta_j \)'s, one obtains a different \( \Theta^* \) (i.e., a different Kronecker product model). Although the eigenvalues of these two Kronecker product models are the same, the eigenvectors of them are not.

Third, if one fixes the ordering of the data, but uses a different factorization of \( n \), one also obtains a different Kronecker product model. Suppose that \( n \) has the prime factorization \( n = p_1 \times p_2 \times \cdots \times p_v \) for some positive integer \( v \) and primes \( p_j \) for \( j = 1, \ldots, v \). Then there exist several different Kronecker product models, each of which is indexed by the dimensions of the sub-matrices. The base model has dimensions \( (p_1, p_2, \ldots, p_v) \), but there are many possible aggregations of this, for example, \( [1](p_1 \times p_2), \ldots, (p_{v-1} \times p_v) \).

To address the second and third issues, we might choose among Kronecker product models using some model selection criterion which penalizes models with more parameters. For example, we may define the Bayesian Information Criterion (BIC) in terms of the original parameters \( \rho \):

\[
BIC(\rho) = -\frac{2}{T} \ell_T(\mu, D, \rho) + \frac{\log T}{T} s_\rho,
\]

where \( \ell_T \) is the log likelihood function defined in (4.1), and \( s_\rho \) is the dimension of \( \rho \). We seek the Kronecker product model with the minimum preceding display. Typically there are not so many factorizations to consider, so this is not too computationally burdensome.

### 6 Monte Carlo Simulations and an Application

In this section, we first provide a set of Monte Carlo simulations that evaluate the performance of the QMLE and MD estimator, and then give a small application of our Kronecker product model to daily stock returns.

#### 6.1 Monte Carlo Simulations

We simulate \( T \) random vectors \( y_t \) of dimension \( n \) according to

\[
y_t = \Sigma^{1/2} z_t, \quad z_t \sim N(0, I_n) \quad \Sigma = \Sigma_1 \otimes \Sigma_2 \otimes \cdots \otimes \Sigma_v, \tag{6.1}
\]

where \( n = 2^v \) and \( v \in \mathbb{N} \). That is, the sub-matrices \( \Sigma_i \) are \( 2 \times 2 \) for \( i = 1, \ldots, v \). These sub-matrices \( \Sigma_j \) are generated with unit variances and off-diagonal elements drawn randomly from a uniform distribution on \((-1, 1)\). This ensures positive definiteness of \( \Sigma \). Note that we have two sources of randomness in this data generating process: random innovations \( z_t \) and random off-diagonal elements of the \( \Sigma_i \) for \( i = 1, \ldots, v \). Due to the unit variances, \( \Sigma \) is also the correlation matrix \( \Theta \) of \( y_t \), but the econometrician is unaware of this: He
applies a Kronecker product model to the correlation matrix $\Theta$. We consider the correctly specified case, i.e., the Kronecker product model has a factorization $n = 2^v$. The sample size is set to $T = 300$ while we vary $v$ (hence $n$). We set the Monte Carlo simulations to 1000.

We shall consider the QMLE and the MD estimator. For the QMLE, we estimate the original parameters $\rho$ and obtain an estimator for $\Theta$ (and hence $\Sigma$) directly. Recalling (4.1), we could use $\ell_T(\Bar{y}, \Hat{D}_T, \rho)$ to optimise $\rho$. For the MD estimator, we estimate log parameters $\theta^0$ via formula (3.3), obtain an estimator for log $\Theta$ and finally obtain an estimator for $\Theta$ (and hence $\Sigma$) via matrix exponential. In the MD cases, we need to specify a choice of the weighting matrix $W$. Given its sheer dimension $(n(n+1)/2 \times n(n+1)/2)$, any non-sparse $W$ will be a huge computational burden in terms of memory for the MD estimator. Hence we consider two diagonal weighting matrices

$$W_1 = I_{n(n+1)/2}, \quad W_2 = [2] D_n^+ [1] \Hat{D}_T \otimes \Hat{D}_T D_n^+ \tau^{-1}.$$ 

In the latter case, the MD estimator is inversely weighted by the sample variances. Weighting matrix $W_2$ resembles, but is not the same as, a feasible version of the optimal weighting matrix $W_{op}$. The choice of $W_2$ is based on heuristics. In an unreported simulation, we also consider the optimally weighted MD estimator. The optimally weighted MD estimator is extremely computationally intensive and its finite sample performance is not as good as those weighted by $W_1$ or $W_2$. This is probably because a data-driven, large-dimensional weighting matrix introduces additional sizeable estimation errors in small samples - such a phenomenon has been well documented in the GMM framework by $\text{?}$.

We compare our estimators with $\text{?}$ direct nonlinear shrinkage estimator (the LW2017 estimator hereafter).$^9$

Given a generic estimator $\Hat{\Sigma}$ of the covariance matrix $\Sigma$ and in each simulation, we can compute

$$1 - \frac{\|\Hat{\Sigma} - \Sigma\|_F^2}{\|\Sigma_T - \Sigma\|_F^2}.$$ 

The median of the preceding display is calculated among all the simulations and denoted RI in terms of $\Sigma$. Criterion RI is closely related to the percentage relative improvement in average loss (PRIAL) criterion in $\text{?}$.$^{10}$ As PRIAL, RI measures the performance of the estimator $\Hat{\Sigma}$ with respect to the sample covariance estimator $\Hat{\Sigma}_T$. Note that $\text{RI} \in (-\infty, 1]$: A negative value means $\Hat{\Sigma}$ performs worse than $\Hat{\Sigma}_T$ while a positive value means otherwise. RI is more robust to outliers than PRIAL.

Often an estimator of the precision matrix $\Sigma^{-1}$ is of more interest than that of $\Sigma$ itself, so we also compute RI for the inverse covariance matrix; that is, we compute the median of

$$1 - \frac{\|\Hat{\Sigma}^{-1} - \Sigma^{-1}\|_F^2}{\|\Sigma_T^{-1} - \Sigma^{-1}\|_F^2}.$$ 

---

$^9$The Matlab code for the direct nonlinear shrinkage estimator is downloaded from the website of Professor Michael Wolf from the Department of Economics at the University of Zurich. We are grateful for this.

$^{10}$It is defined as

$$\text{PRIAL} = 1 - \frac{\mathbb{E}\|\Hat{\Sigma} - \Sigma\|_F^2}{\mathbb{E}\|\Sigma_T - \Sigma\|_F^2}.$$ 

23
Note that this requires invertibility of the sample covariance matrix \( \hat{\Sigma}_T \) and therefore can only be calculated for \( n < T \).

Our final criterion is the minimum variance portfolio (MVP) constructed from an estimator of the covariance matrix. The weights of the minimum variance portfolio are given by

\[
w_{\text{MVP}} = \frac{\Sigma^{-1} t_n}{t_n^\top \Sigma^{-1} t_n},
\]

where \( t_n = (1, 1, \ldots, 1)^\top \) of dimension \( n \) (see ?, ?, etc). The first MVP weights are constructed using the sample covariance matrix \( \hat{\Sigma}_T \) while the second MVP weights are constructed using a generic estimator of \( \Sigma \). These two minimum variance portfolios are then evaluated by calculating their standard deviations in the out-of-sample data \((y_t)\) generated using the same mechanism. The out-of-sample size is set to \( T' = 21 \). The ratio of the standard deviation of the minimum variance portfolio constructed from \( \Sigma \) over that of the minimum variance portfolio constructed from \( \hat{\Sigma}_T \) is calculated. We report its median (VR) over Monte Carlo simulations. Note that VR \( \in [0, +\infty) \): A value greater than one means \( \Sigma \) performs worse than \( \hat{\Sigma}_T \) while a value less than one means otherwise.

Table 1 reports RI-1 (RI in terms of \( \Sigma \)), RI-2 (RI in terms of \( \Sigma^{-1} \)) and VR for various \( n \). We observe the following patterns. First, we see that all our estimators QMLE, MD1, MD2 outperform the sample covariance matrix in all dimensional cases including both the small-dimensional cases (e.g., \( n = 4 \)) and the large-dimensional cases (e.g., \( n = 256 \)). Note that in the large dimensional case like \( n = 256, T = 300 \), the ratio \( n/T \) is close to 1 - a case not really covered by Assumption 3.3. This perhaps illustrates that Assumption 3.3 is a sufficient but not necessary condition for theoretical analysis of our proposed methodology. Second, such a phenomenon holds in terms of RI-1, RI-2 and VR. The superiority of our estimators over the sample covariance matrix increases when \( n/T \) increases. Third, the QMLE outperforms the MD estimators whenever \( n/T \) is close to one, while the opposite holds when \( n/T \) is small. Fourth, the LW2017 estimator also beats the sample covariance matrix but its RI-1 margin is thin. This is perhaps not surprising as the LW2017 estimator does not utilise the Kronecker product structure of the data generating process. Overall, the QMLE is the best estimator in this baseline setting.

As robustness checks, we consider two modifications of our baseline data generating process:

(i) Time series \( y_t \) is still generated as in (6.1) but the actual data are \( w_t \):

\[
w_1 = y_1 \\
w_t = a_w w_{t-1} + \sqrt{1 - a_w^2} y_t, \quad t = 2, \ldots, T.
\]

The parameter \( a_w \) is set to be 0.5 to capture the temporal dependence.

(ii) Same as modification (i), but \( y_t \) is drawn from a multivariate \( t \) distribution of 5 degrees of freedom with \( \Sigma \) as its correlation matrix.

In modification (i), \( w_t \) is serially correlated given any non-zero autoregressive scalar \( a_w \) but its covariance matrix is still \( \Sigma \). A choice of \( a_w = 0.5 \) is consistent with Assumption 3.2. Our simulation results are reasonably robust to the choice of \( a_w \). In modification (ii), in addition to the serial dependence, we add heavy-tailed features to the data which
<table>
<thead>
<tr>
<th></th>
<th>n</th>
<th>4</th>
<th>8</th>
<th>16</th>
<th>32</th>
<th>64</th>
<th>128</th>
<th>256</th>
</tr>
</thead>
<tbody>
<tr>
<td>4* RI-1 QMLE</td>
<td>0.227</td>
<td>0.529</td>
<td>0.714</td>
<td>0.820</td>
<td>0.892</td>
<td>0.929</td>
<td>0.950</td>
<td></td>
</tr>
<tr>
<td></td>
<td>MD1</td>
<td>0.345</td>
<td>0.632</td>
<td>0.789</td>
<td>0.862</td>
<td>0.897</td>
<td>0.909</td>
<td>0.618</td>
</tr>
<tr>
<td></td>
<td>MD2</td>
<td>0.339</td>
<td>0.631</td>
<td>0.785</td>
<td>0.858</td>
<td>0.896</td>
<td>0.908</td>
<td>0.616</td>
</tr>
<tr>
<td></td>
<td>LW2017</td>
<td>0.020</td>
<td>0.027</td>
<td>0.046</td>
<td>0.063</td>
<td>0.087</td>
<td>0.106</td>
<td>0.127</td>
</tr>
<tr>
<td>4* RI-2 QMLE</td>
<td>0.323</td>
<td>0.615</td>
<td>0.805</td>
<td>0.914</td>
<td>0.973</td>
<td>0.995</td>
<td>1.000</td>
<td></td>
</tr>
<tr>
<td></td>
<td>MD1</td>
<td>0.354</td>
<td>0.632</td>
<td>0.771</td>
<td>0.752</td>
<td>0.665</td>
<td>0.588</td>
<td>0.837</td>
</tr>
<tr>
<td></td>
<td>MD2</td>
<td>0.344</td>
<td>0.643</td>
<td>0.790</td>
<td>0.796</td>
<td>0.714</td>
<td>0.628</td>
<td>0.846</td>
</tr>
<tr>
<td></td>
<td>LW2017</td>
<td>0.136</td>
<td>0.181</td>
<td>0.235</td>
<td>0.351</td>
<td>0.521</td>
<td>0.756</td>
<td>0.991</td>
</tr>
<tr>
<td>4* VR QMLE</td>
<td>0.999</td>
<td>0.995</td>
<td>0.980</td>
<td>0.953</td>
<td>0.899</td>
<td>0.770</td>
<td>0.389</td>
<td></td>
</tr>
<tr>
<td></td>
<td>MD1</td>
<td>0.999</td>
<td>0.993</td>
<td>0.979</td>
<td>0.953</td>
<td>0.900</td>
<td>0.774</td>
<td>0.401</td>
</tr>
<tr>
<td></td>
<td>MD2</td>
<td>0.999</td>
<td>0.993</td>
<td>0.979</td>
<td>0.954</td>
<td>0.899</td>
<td>0.774</td>
<td>0.400</td>
</tr>
<tr>
<td></td>
<td>LW2017</td>
<td>1.000</td>
<td>0.999</td>
<td>0.998</td>
<td>0.993</td>
<td>0.975</td>
<td>0.912</td>
<td>0.544</td>
</tr>
</tbody>
</table>

Table 1: The baseline setting. QMLE, MD1, MD2 and LW2017 stand for the quasi-maximum likelihood estimator of the Kronecker product model, the minimum distance estimator (weighted by \( W_1 \)) of the Kronecker product model, the minimum distance estimator (weighted by \( W_2 \)) of the Kronecker product model, and the direct nonlinear shrinkage estimator, respectively. RI-1 and RI-2 are RI criteria in terms of \( \Sigma \) and \( \Sigma^{-1} \), respectively. VR is the median of the ratio of the standard deviation of the MVP using a generic estimator to that using the sample covariance matrix out of sample. The sample size is fixed at \( T = 300 \).

might be a better reflection of reality. Heavy-tailed data are not covered by Assumption 3.1 so this modification serves as a robustness check for our theoretical findings.

The results of modification (i) are reported in Table 2. Those four observations we made from the baseline setting (Table 1) still hold when we relax the independence assumption of the data. Modification (ii) are reported in Table 3. When we switch on both temporal dependence and heavy tails, all estimators - ours and the LW2017 estimator - are adversely affected to a certain extent. In particular, in terms of RI-2, both the QMLE and LW2017 estimators fare worse than the sample covariance matrix in small dimensions. Overall, the identity weighted MD estimator is the best estimator in modification (ii). That the MD estimator trumps the QMLE in heavy-tailed data is intuitive because the MD estimator is derived not based on a particular distributional assumption.

### 6.2 An Application

We now consider estimation of the covariance matrix of \( n' = 441 \) stock returns \((y_t)\) in the S&P 500 index. We have daily observations from January 3, 2005 to November 6, 2015. The number of trading days is \( T = 2732 \). Since the underlying data might not have a multiplicative structure giving rise to a Kronecker product - or if they do but we are unaware of it - a Kronecker product model in this application is inherently misspecified. In other words, we are exploiting Kronecker product models’ approximating feature to a general covariance matrix.

We have proved in Appendix A.2 that in a given Kronecker product model there exists a member which is closest to the true covariance matrix. However, in order for this closest "distance" to be small, the chosen Kronecker product model needs to be
<table>
<thead>
<tr>
<th></th>
<th>4</th>
<th>8</th>
<th>16</th>
<th>32</th>
<th>64</th>
<th>128</th>
<th>256</th>
</tr>
</thead>
<tbody>
<tr>
<td>$n$</td>
<td>4</td>
<td>8</td>
<td>16</td>
<td>32</td>
<td>64</td>
<td>128</td>
<td>256</td>
</tr>
<tr>
<td>4*RI-1</td>
<td>QMLE</td>
<td>0.219</td>
<td>0.514</td>
<td>0.712</td>
<td>0.821</td>
<td>0.887</td>
<td>0.928</td>
</tr>
<tr>
<td></td>
<td>MD1</td>
<td>0.321</td>
<td>0.611</td>
<td>0.775</td>
<td>0.849</td>
<td>0.880</td>
<td>0.889</td>
</tr>
<tr>
<td></td>
<td>MD2</td>
<td>0.310</td>
<td>0.611</td>
<td>0.770</td>
<td>0.844</td>
<td>0.877</td>
<td>0.890</td>
</tr>
<tr>
<td></td>
<td>LW2017</td>
<td>0.025</td>
<td>0.032</td>
<td>0.049</td>
<td>0.065</td>
<td>0.093</td>
<td>0.117</td>
</tr>
<tr>
<td>4*RI-2</td>
<td>QMLE</td>
<td>0.320</td>
<td>0.654</td>
<td>0.824</td>
<td>0.932</td>
<td>0.980</td>
<td>0.997</td>
</tr>
<tr>
<td></td>
<td>MD1</td>
<td>0.338</td>
<td>0.639</td>
<td>0.737</td>
<td>0.691</td>
<td>0.593</td>
<td>0.517</td>
</tr>
<tr>
<td></td>
<td>MD2</td>
<td>0.347</td>
<td>0.652</td>
<td>0.770</td>
<td>0.844</td>
<td>0.877</td>
<td>0.890</td>
</tr>
<tr>
<td></td>
<td>LW2017</td>
<td>0.220</td>
<td>0.292</td>
<td>0.429</td>
<td>0.634</td>
<td>0.818</td>
<td>0.939</td>
</tr>
<tr>
<td>4*VR</td>
<td>QMLE</td>
<td>0.998</td>
<td>0.988</td>
<td>0.975</td>
<td>0.927</td>
<td>0.860</td>
<td>0.728</td>
</tr>
<tr>
<td></td>
<td>MD1</td>
<td>0.997</td>
<td>0.987</td>
<td>0.973</td>
<td>0.925</td>
<td>0.862</td>
<td>0.733</td>
</tr>
<tr>
<td></td>
<td>MD2</td>
<td>0.997</td>
<td>0.987</td>
<td>0.973</td>
<td>0.924</td>
<td>0.862</td>
<td>0.732</td>
</tr>
<tr>
<td></td>
<td>LW2017</td>
<td>1.000</td>
<td>0.999</td>
<td>0.997</td>
<td>0.990</td>
<td>0.970</td>
<td>0.907</td>
</tr>
</tbody>
</table>

Table 2: Modification (i). QMLE, MD1, MD2 and LW2017 stand for the quasi-maximum likelihood estimator of the Kronecker product model, the minimum distance estimator (weighted by $W_1$) of the Kronecker product model, the minimum distance estimator (weighted by $W_2$) of the Kronecker product model, and the direct nonlinear shrinkage estimator, respectively. RI-1 and RI-2 are RI criteria in terms of $\Sigma$ and $\Sigma^{-1}$, respectively. VR is the median of the ratio of the standard deviation of the MVP using a generic estimator to that using the sample covariance matrix out of sample. The sample size is fixed at $T = 300$.

<table>
<thead>
<tr>
<th></th>
<th>4</th>
<th>8</th>
<th>16</th>
<th>32</th>
<th>64</th>
<th>128</th>
<th>256</th>
</tr>
</thead>
<tbody>
<tr>
<td>$n$</td>
<td>4</td>
<td>8</td>
<td>16</td>
<td>32</td>
<td>64</td>
<td>128</td>
<td>256</td>
</tr>
<tr>
<td>4*RI-1</td>
<td>QMLE</td>
<td>0.021</td>
<td>0.105</td>
<td>0.203</td>
<td>0.320</td>
<td>0.442</td>
<td>0.564</td>
</tr>
<tr>
<td></td>
<td>MD1</td>
<td>0.071</td>
<td>0.211</td>
<td>0.348</td>
<td>0.492</td>
<td>0.621</td>
<td>0.719</td>
</tr>
<tr>
<td></td>
<td>MD2</td>
<td>0.084</td>
<td>0.242</td>
<td>0.378</td>
<td>0.510</td>
<td>0.626</td>
<td>0.712</td>
</tr>
<tr>
<td></td>
<td>LW2017</td>
<td>-0.023</td>
<td>-0.001</td>
<td>0.029</td>
<td>0.069</td>
<td>0.118</td>
<td>0.158</td>
</tr>
<tr>
<td>4*RI-2</td>
<td>QMLE</td>
<td>-0.035</td>
<td>-0.139</td>
<td>-0.357</td>
<td>-0.831</td>
<td>-0.202</td>
<td>0.867</td>
</tr>
<tr>
<td></td>
<td>MD1</td>
<td>0.006</td>
<td>0.035</td>
<td>0.111</td>
<td>0.385</td>
<td>0.896</td>
<td>0.636</td>
</tr>
<tr>
<td></td>
<td>MD2</td>
<td>-0.006</td>
<td>-0.009</td>
<td>0.032</td>
<td>0.255</td>
<td>0.894</td>
<td>0.724</td>
</tr>
<tr>
<td></td>
<td>LW2017</td>
<td>-0.103</td>
<td>-0.206</td>
<td>-0.428</td>
<td>-0.847</td>
<td>-0.279</td>
<td>0.825</td>
</tr>
<tr>
<td>4*VR</td>
<td>QMLE</td>
<td>0.996</td>
<td>0.982</td>
<td>0.956</td>
<td>0.923</td>
<td>0.842</td>
<td>0.708</td>
</tr>
<tr>
<td></td>
<td>MD1</td>
<td>0.994</td>
<td>0.982</td>
<td>0.955</td>
<td>0.921</td>
<td>0.840</td>
<td>0.720</td>
</tr>
<tr>
<td></td>
<td>MD2</td>
<td>0.994</td>
<td>0.982</td>
<td>0.955</td>
<td>0.920</td>
<td>0.840</td>
<td>0.719</td>
</tr>
<tr>
<td></td>
<td>LW2017</td>
<td>1.000</td>
<td>0.999</td>
<td>0.995</td>
<td>0.989</td>
<td>0.968</td>
<td>0.906</td>
</tr>
</tbody>
</table>

Table 3: Modification (ii). QMLE, MD1, MD2 and LW2017 stand for the quasi-maximum likelihood estimator of the Kronecker product model, the minimum distance estimator (weighted by $W_1$) of the Kronecker product model, the minimum distance estimator (weighted by $W_2$) of the Kronecker product model, and the direct nonlinear shrinkage estimator, respectively. RI-1 and RI-2 are RI criteria in terms of $\Sigma$ and $\Sigma^{-1}$, respectively. VR is the median of the ratio of the standard deviation of the MVP using a generic estimator to that using the sample covariance matrix out of sample. The sample size is fixed at $T = 300$. 

26
Table 4: MD, LW2004 and LW2017 stand for the (identity matrix weighted) minimum distance estimators of the Kronecker product models (factorisations given in parentheses), the linear shrinkage estimator, and the direct nonlinear shrinkage estimator, respectively. Impr is the median of the 106 quantities calculated based on (6.3) and Prop is the proportion of the times (out of 106) that a competitor MVP outperforms the sample covariance MVP (i.e., the proportion of the times when (6.3) is positive). A random permutation of the data means the 441 stocks are randomly reshuffled.

<table>
<thead>
<tr>
<th></th>
<th>MD</th>
<th>MD</th>
<th>MD</th>
<th>MD</th>
<th>2* LW2004</th>
<th>2* LW2017</th>
</tr>
</thead>
<tbody>
<tr>
<td></td>
<td>(2 × 2 × 3 × 37)</td>
<td>(4 × 111)</td>
<td>(3 × 148)</td>
<td>(2 × 222)</td>
<td></td>
<td></td>
</tr>
<tr>
<td>Impr</td>
<td>0.265</td>
<td>0.379</td>
<td>0.394</td>
<td>0.440</td>
<td>0.459</td>
<td>0.518</td>
</tr>
<tr>
<td>Prop</td>
<td>0.811</td>
<td>0.896</td>
<td>0.915</td>
<td>0.953</td>
<td>0.991</td>
<td>0.981</td>
</tr>
<tr>
<td>(lr)2-7</td>
<td>a random permutation of the data</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>Impr</td>
<td>0.259</td>
<td>0.364</td>
<td>0.404</td>
<td>0.431</td>
<td>0.459</td>
<td>0.518</td>
</tr>
<tr>
<td>Prop</td>
<td>0.811</td>
<td>0.887</td>
<td>0.915</td>
<td>0.943</td>
<td>0.991</td>
<td>0.981</td>
</tr>
<tr>
<td>(lr)2-7</td>
<td>a random permutation of the data</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>Impr</td>
<td>0.263</td>
<td>0.351</td>
<td>0.366</td>
<td>0.436</td>
<td>0.459</td>
<td>0.518</td>
</tr>
<tr>
<td>Prop</td>
<td>0.811</td>
<td>0.887</td>
<td>0.906</td>
<td>0.943</td>
<td>0.991</td>
<td>0.981</td>
</tr>
</tbody>
</table>

The table shows the performance of different estimators of the covariance matrix of the final system. The original ordering of the data was used in the first row, and a random permutation of the data was used in the second row. The results show that the MD estimator performed better than the sample covariance estimator, with Impr and Prop values indicating a higher proportion of times that the MD estimator outperformed the sample covariance estimator.

We add an $3 \times 1$ dimensional pseudo random vector $z_t$ which are $N(0, I_3)$ distributed and independent over $t$. The dimension of the final system is $n = 441 + 3 = 444$. Again we fit Kronecker product models to the correlation matrix of the final system and recover an estimator for the covariance matrix of the final system via left and right multiplication of the estimated correlation matrix of the final system by $\hat{D}_t^{1/2}$. Last, we extract the $441 \times 441$ upper-left block of the estimated covariance matrix of the final system to form our Kronecker product estimator of the covariance matrix of $y_t$. The dimension of the added pseudo random vector should not be too large to avoid introducing additional noise, which could adversely affect the performance of the Kronecker product models. We choose the dimension of the final system to be 444 because its prime factorization is $2 \times 2 \times 3 \times 37$, and we experiment with several Kronecker product models. We did try other dimensions for the final system and the pattern discussed below remains generally the same.

As we are considering less parsimonious models, the QMLE is computationally intensive and found to perform worse than the MD estimator in preliminary investigations, so we only use the MD estimator. The MD estimator is extremely fast because its formula is just (3.3). We choose the weighting matrix to be the identity matrix.

We follow the approach of and estimate our model on windows of size 504 days (equal to two years’ trading days) that are shifted from the beginning to the end of the sample. The Kronecker product estimator of the covariance matrix of $y_t$ is used to construct the minimum variance portfolio (MVP) weights as in (6.2). We also compute the MVP weights using the sample covariance matrix of $y_t$. These two minimum variance portfolios are then evaluated using the next 21 days (equal to one month’s trading days).
out-of-sample. In particular, we calculate

$$1 - \frac{\text{sd(a competitor MVP)}}{\text{sd(sample covariance MVP)}},$$

(6.3)

where \(\text{sd}(\cdot)\) computes standard deviation. Then the estimation window of 504 days is shifted forward by 21 days. This procedure is repeated until we reach the end of the sample; the total number of out-of-sample evaluations is 106. We consider two evaluation criteria of performance: \(\text{Impr}\) and \(\text{Prop}\). \(\text{Impr}\) is the median of the 106 quantities calculated based on (6.3). Note that \(\text{Impr} \in (-\infty, 1]\): A negative value means a competitor MVP performs worse than the sample covariance MVP while a positive value means otherwise. \(\text{Prop}\) is the proportion of the times (out of 106) that a competitor MVP outperforms the sample covariance MVP (i.e., the proportion of the times when (6.3) is positive).

For comparison, we will consider the linear shrinkage estimator of \(\gamma\) and the direct nonlinear shrinkage estimator of \(\gamma\). The results are reported in Table 4. We first use the original ordering of the data, i.e., alphabetical, and have the following observations. First, all the Kronecker product MVPs outperform the sample covariance MVP. Second, as we move from the most parsimonious factorisation \((444 = 2 \times 2 \times 3 \times 37)\) to the least parsimonious factorisation \((444 = 2 \times 222)\), the performance of Kronecker product MVPs monotonically improves. This is intuitive: Since we are using Kronecker product models to approximate a general covariance matrix, a more flexible Kronecker product model could fit the data better. There is no over-fitting at least in this application as we consider out-of-sample evaluation. Third, the performance of the \((2 \times 222)\) Kronecker product MVP is very close to that of a sophisticated estimator like \(\gamma\)’s linear shrinkage estimator. This is commendable because here a Kronecker product model is a misspecified parametric model for a general covariance matrix while the linear shrinkage estimator is in essence a data-driven, nonparametric estimator.

We next randomly reshuffle the 441 stocks two times and use the same Kronecker product models. In these two cases, the rows and columns of the true covariance matrix also get reshuffled. We see that the performances of those Kronecker product models are marginally affected by the reshuffle. \(\gamma\)’s and \(\gamma\)’s shrinkage estimators are, as expected, not affected by the ordering of the data.

7 Conclusions

We have established the large sample properties of estimators of Kronecker product models in the large dimensional case. In particular, we obtained norm consistency and the large dimensional CLTs for the MD and one-step estimators. Kronecker product models outperform the sample covariance matrix theoretically, in Monte Carlo simulations, and in an application to portfolio choice. When a Kronecker product model is correctly specified, Monte Carlo simulations show that estimators of it can beat \(\gamma\)’s direct non-linear shrinkage estimator. In the application, when one uses Kronecker product models as an approximating device to a general covariance matrix, a less parsimonious one can perform almost as good as \(\gamma\)’s linear shrinkage estimator. It is possible to extend the framework in various directions to improve performance.

A final motivation for the Kronecker product structure is that it can be used as a component of a super model consisting of several components. For instance, the idea of
the decomposition in (1.1) could be applied to components of dynamic models such as multivariate GARCH, an area in which Luc Bauwens has contributed significantly over the recent years, see also his highly cited review paper ?. For example, the dynamic conditional correlation (DCC) model of ?, or the BEKK model of ? both have intercept matrices that are required to be positive definite and suffer from the curse of dimensionality, for which model (1.1) would be helpful. Also, parameter matrices associated with the dynamic terms in the model could be equipped with a Kronecker product, similar to a suggestion by ? for vector autoregressions.

A Appendix

This appendix is organised as follows: Appendix A.1 further discusses this matrix $E$ of the minimum distance estimator in Section 3. Appendix A.2 shows that a Kronecker product model has a best approximation to a general covariance or correlation matrix. Appendix A.3 and A.4 contain proofs of Theorem 3.1 and of Theorem 3.2 respectively. Appendix A.5 contains auxiliary lemmas used in various places of this appendix.

A.1 Matrix $E$

The proof of the following theorem gives a concrete formula for the matrix $E$ of the minimum distance estimator.

**Theorem A.1.** Suppose that

$$\Theta^* = \Theta_1^* \otimes \Theta_2^* \otimes \cdots \otimes \Theta_v^*,$$

where $\Theta_j^*$ is $n_j \times n_j$ dimensional such that $n = n_1 \cdot n_2 \cdots n_v$. Taking the logarithm on both sides gives

$$\log \Theta^* = \log \Theta_1^* \otimes I_{n_2} \otimes \cdots \otimes I_{n_v} + I_{n_1} \otimes \log \Theta_2^* \otimes I_{n_3} \otimes \cdots \otimes I_{n_v} + \cdots + I_{n_1} \otimes I_{n_2} \otimes \cdots \otimes \log \Theta_v^*.$$

For identification we set the first diagonal entry of $\log \Theta_j^*$ to be 0 for $j = 1, \ldots, v - 1$. In total there are

$$s := \sum_{j=1}^v \frac{n_j(n_j + 1)}{2} - (v - 1)$$

(identifiable) log parameters in $\Theta_1^*, \ldots, \Theta_v^*$; let $\theta^* \in \mathbb{R}^s$ denote these. Then there exists a $n(n + 1)/2 \times s$ full column rank matrix $E$ such that

$$\text{vech}(\log \Theta^*) = E\theta^*.$$

**Proof.** Note that

$$\text{vec}(\log \Theta^*) = \text{vec}(\log \Theta_1^* \otimes I_{n_2} \otimes \cdots \otimes I_{n_v}) + \text{vec}(I_{n_1} \otimes \log \Theta_2^* \otimes I_{n_3} \otimes \cdots \otimes I_{n_v}) + \cdots + \text{vec}(I_{n_1} \otimes I_{n_2} \otimes \cdots \otimes \log \Theta_v^*).$$

If

$$\text{vec}(I_{n_1} \otimes \log \Theta_1^* \otimes I_{n_3} \otimes \cdots \otimes I_{n_v}) = E_i \text{vec}(\log \Theta_1^*)$$
for some $n^2 \times n_i(n_i + 1)/2$ matrix $E_i$ for $i = 1, \ldots, v$, then we have

$$\text{vech}(\log \Theta^*) = D_n^+ \cdot \text{vec}(\log \Theta^*) = D_n^+ \cdot \begin{bmatrix} E_1 & E_2 & \cdots & E_v \end{bmatrix} \begin{bmatrix} \text{vech}(\log \Theta_1^*) \\ \text{vech}(\log \Theta_2^*) \\ \vdots \\ \text{vech}(\log \Theta_v^*) \end{bmatrix}.$$ 

For identification we set the first diagonal entry of $\log \Theta_j^*$ to be 0 for $j = 1, \ldots, v - 1$. In total there are

$$s := \sum_{j=1}^v \frac{n_j(n_j + 1)}{2} - (v - 1)$$

( identifiable) log parameters in $\Theta_1^*, \ldots, \Theta_v^*$; let $\theta^* \in \mathbb{R}^s$ denote these. Then there exists a $n(n+1)/2 \times s$ full column rank matrix $E$ such that

$$\text{vech}(\log \Theta^*) = E\theta^*,$$

where

$$E := D_n^+ \cdot \begin{bmatrix} E_{1,(-1)} & E_{2,(-1)} & \cdots & E_{v-1,(-1)} & E_v \end{bmatrix}$$

and $E_{i,(-1)}$ stands for matrix $E_i$ with its first column removed. We now determine the formula for $E_i$. We first consider $\text{vec}(\log \Theta_1^* \otimes I_{n_2} \otimes \cdots \otimes I_{n_v})$.

$$\text{vec}(\log \Theta_1^* \otimes I_{n_2} \otimes \cdots \otimes I_{n_v}) = \text{vec}(\log \Theta_1^* \otimes I_{n_1}) = [1]I_{n_1} \otimes K_{n_1,n_1} \otimes I_{n_1} \cdot \text{vec}(\log \Theta_1^*) \otimes \text{vec} I_{n_1}$$

$$= [1]I_{n_1} \otimes K_{n_1,n_1} \otimes I_{n_1}[1]I_{n_1} \cdot \text{vec}(\log \Theta_1^*) \otimes \text{vec} I_{n_1}$$

$$= [1]I_{n_1} \otimes K_{n_1,n_1} \otimes I_{n_1}[1]I_{n_1} \cdot \text{vec} I_{n_1} D_{n_1} \cdot \text{vech}(\log \Theta_1^*),$$

where the second equality is due to ? Theorem 3.10 p55. Thus,

$$E_1 := [1]I_{n_1} \otimes K_{n_1,n_1} \otimes I_{n_1}[1]I_{n_1} \otimes \text{vec} I_{n_1} D_{n_1}.$$

We now consider $\text{vec}(I_{n_1} \otimes \cdots \otimes \log \Theta_1^* \otimes \cdots \otimes I_{n_v})$.

$$\text{vec}(I_{n_1} \otimes \cdots \otimes \log \Theta_1^* \otimes \cdots \otimes I_{n_v}) = \text{vec}[2]K_{n_1 \cdots n_v,n_1 \cdots n_v} \cdot [1] \otimes \log \Theta_1^* \otimes I_{n_1} \cdot K_{n_1 \cdots n_v,n_1 \cdots n_v} \cdot [1] \cdot \text{vec}[1] \otimes \log \Theta_1^* \otimes I_{n_1}$$

$$= [1]K_{n_1,n_1} \cdots n_v,n_1 \cdots n_v \otimes \cdots \otimes K_{n_1,n_1} \cdots n_v,n_1 \cdots n_v \cdot \text{vec}[1] \text{vec}[1] \otimes \log \Theta_1^* \otimes I_{n_1}$$

$$= [1]K_{n_1 \cdots n_v,n_1 \cdots n_v} \otimes K_{n_1 \cdots n_v,n_1 \cdots n_v} \cdot [1] I_{n_1} \otimes K_{n_1,n_1} \otimes I_{n_1}[1]I_{n_1} \otimes \text{vec} I_{n_1} D_{n_1} \cdot \text{vech}(\log \Theta_1^*),$$

where the first equality is due to the identity $B \otimes A = K_{p,m}(A \otimes B)K_{m,p}$ for $A (m \times m)$ and $B (p \times p)$. Thus

$$E_i := [1]K_{n_1 \cdots n_v,n_1 \cdots n_v} \otimes K_{n_1 \cdots n_v,n_1 \cdots n_v} \cdot [1] I_{n_i} \otimes K_{n_i,n_i} \otimes I_{n_i}[1]I_{n_i} \otimes \text{vec} I_{n_i} D_{n_i},$$

for $i = 2, \ldots, v$.

\[ \square \]

**Lemma A.1.** Given that $n = n_1 \cdot n_2 \cdots n_v$, the $s \times s$ matrix $E^TE$ takes the following form:
(i) For $i = 1, \ldots, s$, the $i$th diagonal entry of $E^T E$ records how many times the $i$th parameter in $\theta^0$ has appeared in $\text{vech}(\log \Theta^0)$. The value depends on to which $\log \Theta^0_j$ the $i$th parameter in $\theta^0, \theta_i^0$, belongs to. For instance, suppose $\theta_i^0$ is a parameter belonging to $\log \Theta^0_3$, then

$$(E^T E)_{i,i} = n/n_3.$$  

(ii) For $i, k = 1, \ldots, s$ ($i \neq k$), the $(i, k)$ entry of $E^T E$ (or the $(k, i)$ entry of $E^T E$ by symmetry) records how many times the $i$th parameter in $\theta^0, \theta_i^0$, and $k$th parameter in $\theta^0, \theta_k^0$, have appeared together (as summands) in an entry of $\text{vech}(\log \Theta^0)$. The value depends on to which $\log \Theta^0_j$ the $i$th parameter in $\theta^0, \theta_i^0$, and $k$th parameter in $\theta^0, \theta_k^0$, belong to. For instance, suppose $\theta_i^0$ is a parameter belonging to $\log \Theta^0_3$ and $\theta_k^0$ is a parameter belonging to $\log \Theta^0_3$, then

$$(E^T E)_{i,k} = (E^T E)_{k,i} = n/(n_3 \cdot n_5).$$

Note that if both $\theta_i^0$ and $\theta_k^0$ belong to the same $\log \Theta^0_j$, then $(E^T E)_{i,k} = (E^T E)_{k,i} = 0$. Also note that when $\theta_i^0$ is an off-diagonal entry of some $\log \Theta^0_j$, then

$$(E^T E)_{i,k} = (E^T E)_{k,i} = 0$$

for any $k = 1, \ldots, s$ ($i \neq k$).

Proof. Proof by spotting the pattern. \hfill \Box

We here give a concrete example to illustrate Lemma A.1.

**Example A.1.** Suppose that $n_1 = 3, n_2 = 2, n_3 = 2$. We have

\[
\log \Theta^0_1 = \begin{pmatrix} 0 & a_{12} & a_{13} \\ a_{12} & a_{22} & a_{23} \\ a_{13} & a_{23} & a_{33} \end{pmatrix}, \quad \log \Theta^0_2 = \begin{pmatrix} 0 & b_{12} \\ b_{12} & b_{22} \end{pmatrix}, \quad \log \Theta^0_3 = \begin{pmatrix} c_{11} & c_{12} \\ c_{12} & c_{22} \end{pmatrix}
\]

The leading diagonals of $\log \Theta^0_1$ and $\log \Theta^0_2$ are set to zero for identification as explained before. Thus

$$\theta^0 = (a_{1,2}, a_{1,3}, a_{2,2}, a_{2,3}, a_{3,3}, b_{1,2}, b_{2,2}, c_{1,1}, c_{1,2}, c_{2,2})^T.$$  

Then we can invoke Lemma A.1 to write down $E^T E$ without even using Matlab to compute $E$; that is,

\[
E^T E = \begin{pmatrix} 4 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\ 0 & 4 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\ 0 & 0 & 4 & 0 & 0 & 2 & 2 & 0 & 2 \\ 0 & 0 & 0 & 4 & 0 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 & 4 & 0 & 2 & 2 & 0 \\ 0 & 0 & 0 & 0 & 0 & 6 & 0 & 0 & 0 \\ 0 & 0 & 2 & 0 & 2 & 0 & 6 & 3 & 0 \\ 0 & 0 & 2 & 0 & 2 & 0 & 3 & 6 & 0 \\ 0 & 0 & 0 & 0 & 0 & 0 & 0 & 6 & 0 \\ 0 & 0 & 2 & 0 & 2 & 0 & 3 & 0 & 6 \end{pmatrix}
\]
A.2 Best Approximation

In this section of the appendix, we show that for any given $n \times n$ real symmetric, positive definite covariance matrix (or correlation matrix), there is a uniquely defined member of the Kronecker product model that is closest to the covariance matrix (or correlation matrix) in some sense in terms of the log parameter space, once a factorization $n = n_1 \cdots n_v$ is specified.

Let $\mathcal{M}_n$ denote the set of all $n \times n$ real symmetric matrices. For any $n(n+1)/2 \times n(n+1)/2$ known, deterministic, positive definite matrix $W$, define a map

$$\langle A, B \rangle_W := (\text{vech } A)^\top W \text{vech } B \quad A, B \in \mathcal{M}_n.$$ 

It is easy to show that $\langle \cdot, \cdot \rangle_W$ is an inner product. $\mathcal{M}_n$ with inner product $\langle \cdot, \cdot \rangle_W$ can be identified by $\mathbb{R}^{n(n+1)/2}$ with the usual Euclidean inner product. Since for finite $n \mathbb{R}^{n(n+1)/2}$ with the usual Euclidean inner product is a Hilbert space, so is $\mathcal{M}_n$. The inner product $\langle \cdot, \cdot \rangle_W$ induces the following norm

$$\| A \|_W := \sqrt{\langle A, A \rangle_W} = \sqrt{(\text{vech } A)^\top W \text{vech } A}.$$ 

Let $\mathcal{D}_n$ denote the set of matrices of the form

$$\Omega_1 \otimes I_{n_1} \otimes \cdots \otimes I_{n_v} + I_{n_1} \otimes \Omega_2 \otimes \cdots \otimes I_{n_v} + \cdots + I_{n_1} \otimes \cdots \otimes \Omega_v,$$

where $\Omega_j$ are $n_j \times n_j$ real symmetric matrices for $j = 1, \ldots, v$. Note that $\mathcal{D}_n$ is a (linear) subspace of $\mathcal{M}_n$ as, for $\alpha, \beta \in \mathbb{R},$

$$\alpha[1]\Omega_1 \otimes I_{n_1} \otimes \cdots \otimes I_{n_v} + I_{n_1} \otimes \Omega_2 \otimes \cdots \otimes I_{n_v} + \cdots + I_{n_1} \otimes \cdots \otimes \Omega_v +$$

$$\beta[1]\Xi_1 \otimes I_{n_1} \otimes \cdots \otimes I_{n_v} + I_{n_1} \otimes \Xi_2 \otimes \cdots \otimes I_{n_v} + \cdots + I_{n_1} \otimes \cdots \otimes \Xi_v = (\alpha\Omega_1 + \beta\Xi_1) \otimes I_{n_1} \cdots \otimes I_{n_v} + I_{n_1} \otimes (\alpha\Omega_2 + \beta\Xi_2) \otimes \cdots \otimes I_{n_v} + \cdots + I_{n_1} \otimes \cdots \otimes (\alpha\Omega_v + \beta\Xi_v) \in \mathcal{D}_n.$$ 

For finite $n$, $\mathcal{D}_n$ is also closed.

Consider a real symmetric, positive definite covariance matrix $\Sigma$. We have $\log \Sigma \in \mathcal{M}_n$. By the projection theorem of the Hilbert space, there exists a unique matrix $L^0 \in \mathcal{D}_n$ such that

$$\| \log \Sigma - L^0 \|_W = \min_{L \in \mathcal{D}_n} \| \log \Sigma - L \|_W.$$ 

(Note also that $\log \Sigma^{-1} = -\log \Sigma$, so that $-L^0$ simultaneously approximates the precision matrix $\Sigma^{-1}$ in the same norm.)

This says that any real symmetric, positive definite covariance matrix $\Sigma$ has a closest approximating matrix $\Sigma^0$ in a sense that

$$\| \log \Sigma - \log \Sigma^0 \|_W = \min_{L \in \mathcal{D}_n} \| \log \Sigma - L \|_W,$$

where $\Sigma^0 := \exp L^0$. Since $L^0 \in \mathcal{D}_n$, we can write

$$L^0 = L^0_1 \otimes I_{n_1} \otimes \cdots \otimes I_{n_v} + I_{n_1} \otimes L^0_2 \otimes \cdots \otimes I_{n_v} + \cdots + I_{n_1} \otimes \cdots \otimes L^0_v,$$
where \( L^0_j \) are \( n_j \times n_j \) real symmetric matrices for \( j = 1, \ldots, v \). Then

\[
\Sigma^0 = \exp L^0 = \exp[1]L^0_1 \otimes I_{n_1} \otimes \cdots \otimes I_{n_v} + I_{n_1} \otimes L^0_2 \otimes \cdots \otimes I_{n_v} + \cdots + I_{n_1} \otimes \cdots \otimes L^0_v
\]

where the third equality is due to Theorem 10.2 in p235 and the fact that \( L^0_1 \otimes I_{n_1} \otimes \cdots \otimes I_{n_v} \) and \( I_{n_1} \otimes L^0_2 \otimes \cdots \otimes I_{n_v} \) commute, the fourth equality is due to \( f(A) \otimes I = f(A \otimes I) \) for any matrix function \( f \) (e.g., Theorem 1.13 in p10), the fifth equality is due to a property of Kronecker products. Note that \( \Sigma^0_j \) is real symmetric, positive definite, and \( I_{n_j} \) are real symmetric matrices for \( j = 1, \ldots, v \).

We thus see that \( \Sigma^0 \) is of the Kronecker product form, and that the precision matrix \( \Sigma^{-1} \) has a closest approximating matrix \( (\Sigma^0)^{-1} \). This reasoning provides a justification (i.e., interpretation) for using \( \Sigma^0 \) even when the Kronecker product model is misspecified for the covariance matrix. The same reasoning applies to any real symmetric, positive definite correlation matrix \( \Theta \).

? and ? also considered this nearest approximation involving one Kronecker product only and in the original parameter space (not in the log parameter space). In that simplified problem, they showed that the optimisation problem could be solved by the singular value decomposition.

### A.3 The Proof of Theorem 3.1

In this subsection, we give a proof for Theorem 3.1. We will first give some preliminary lemmas leading to the proof of this theorem.

The following lemma characterises the relationship between an exponential-type moment assumption and an exponential tail probability.

**Lemma A.2.** Suppose that a random variable \( X \) satisfies the exponential-type tail condition, i.e., there exist absolute constants \( K_1 > 1, K_2 > 0, r_1 > 0 \) such that

\[
\mathbb{E}[\exp[1]K_2|X|^{r_1}] \leq K_1.
\]

(i) Then for every \( \epsilon \geq 0 \), there exists an absolute constant \( b_1 > 0 \) such that

\[
\mathbb{P}(|X| \geq \epsilon) \leq \exp[1]1 - (\epsilon/b_1)^{r_1}.
\]

(ii) We have \( \mathbb{E}[|X|] \leq \infty \).

(iii) Then for every \( \epsilon \geq 0 \), there exists an absolute constant \( c_1 > 0 \) such that

\[
\mathbb{P}(|X - \mathbb{E}X| \geq \epsilon) \leq \exp[1]1 - (\epsilon/c_1)^{r_1}.
\]

(iv) Suppose that another random variable \( Y \) satisfies \( \mathbb{E}[\exp[1]K_2^*|Y|^{r_1}] \leq K_1^* \) for some absolute constants \( K_1^* > 1, K_2^* > 0, r_1^* > 0 \). Then for every \( \epsilon \geq 0 \), there exists an absolute constant \( b_2 > 0 \) such that

\[
\mathbb{P}(|XY| \geq \epsilon) \leq \exp[1]1 - (\epsilon/b_2)^{r_2},
\]

where \( r_2 \in \left(0, \frac{r_1 r_1^*}{r_1 + r_1^*}\right) \).

33
Proof. For part (i), choose \( C := \log K_1 \lor 1 \) and \( b_1 := (C/K_2)^{1/r_1} \). If \( \epsilon > b_1 \), we have
\[
\mathbb{P}(|X| \geq \epsilon) \leq \frac{\mathbb{E}[1]\exp(K_2|X|^r_1)}{\exp(K_2^\epsilon r_1^\epsilon)} \leq K_1 e^{-K_2\epsilon^r_1} = e^{\log K_1 - K_2\epsilon^r_1} = e^{\log K_1 - C(\epsilon/b_1)^{r_1}}
\]
where the first inequality is due to a variant of Markov’s inequality. If \( \epsilon \leq b_1 \), we have
\[
\mathbb{P}(|X| \geq \epsilon) \leq 1 \leq e^{1-(\epsilon/b_1)^{r_1}}.
\]

For part (ii),
\[
\mathbb{E}|X| = \int_0^\infty \mathbb{P}(|X| \geq t)dt \leq \int_0^\infty e^{-t/b_1}\epsilon^r_1 dt = \int_0^\infty e^{-t/b_1}\epsilon^r_1 dt = \frac{eb_1}{r_1} \int_0^\infty y^{r_1-1}e^{-y}dy = \frac{eb_1}{r_1} \Gamma(r_1^{-1}) \leq \infty,
\]
where the first inequality is due to part (i), the third equality is due to change of variable \( y = (t/b_1)^{r_1} \), and the last equality is due to the recognition of \( \int_0^\infty \Gamma(r_1^{-1})^{-1}y^{r_1-1}e^{-y}dy = 1 \) using Gamma distribution. For part (iii),
\[
\mathbb{P}(|X - \mathbb{E}X| \geq \epsilon) \leq \mathbb{P}(|X| \geq \epsilon - \mathbb{E}|X|) = \mathbb{P}(|X| \geq \epsilon - \mathbb{E}|X| \wedge \epsilon) \leq \exp[3]\frac{1}{b_1^r} \frac{(\epsilon - \mathbb{E}|X| \wedge \epsilon)^{r_1}}{b_1^r}
\]
where the second inequality is due to part (i). First consider the case \( 0 < r_1 < 1 \).
\[
\exp[3]\frac{1}{b_1^r} - \frac{(\epsilon - \mathbb{E}|X| \wedge \epsilon)^{r_1}}{b_1^r} \leq \exp[3]\frac{1}{b_1^r} - \frac{\epsilon_1 - (\mathbb{E}|X| \wedge \epsilon)^{r_1}}{b_1^r} = \exp[3]\frac{1}{b_1^r} - \frac{\epsilon_1}{b_1^r} + \frac{(\mathbb{E}|X| \wedge \epsilon)^{r_1}}{b_1^r}
\]
\[
\leq \exp[3]\frac{1}{b_1^r} - \frac{\epsilon_1}{b_1^r} + \frac{(\mathbb{E}|X|)^{r_1}}{b_1^r} \leq \exp[3]C - \frac{\epsilon_1}{b_1^r} = \exp[3]C[3]\frac{1}{b_1^r} - \frac{\epsilon_1}{b_1^r}
\]
\[
=: \exp[3]C[3]\frac{1}{b_1^r} - \frac{\epsilon_1}{C_1^r}
\]

where the first inequality is due to subadditivity of the concave function: \( (x+y)^{r_1} - x^{r_1} \leq y^{r_1} \) for \( x, y \geq 0 \). If \( \epsilon > c_1 \), we have, via recognizing \( C > 1 \),
\[
\mathbb{P}(|X - \mathbb{E}X| \geq \epsilon) \leq \exp[3]C[3]\frac{1}{b_1^r} - \frac{\epsilon_1}{C_1^r} \leq \exp[3]\frac{1}{b_1^r} - \frac{\epsilon_1}{C_1^r}.
\]
If \( \epsilon \leq c_1 \), we have
\[
\mathbb{P}(|X - \mathbb{E}X| \geq \epsilon) \leq 1 \leq \exp[3]\frac{1}{b_1^r} - \frac{\epsilon_1}{C_1^r}.
\]

We now consider the case \( r_1 \geq 1 \). The proof is almost the same: Instead of relying on subadditivity of the concave function, we rely on Loève’s \( c_r \) inequality: \( |x + y|^{r_1} \leq 2^{r_1-1}(|x|^{r_1} + |y|^{r_1}) \) for \( r_1 \geq 1 \) to get \( 2^{1-r_1}\epsilon^{r_1} - (\mathbb{E}|X| \wedge \epsilon)^{r_1} \leq (\epsilon - \mathbb{E}|X| \wedge \epsilon)^{r_1} \). \( c_1 \) is now defined as \( C_1^r b_1^2 2^{r_1-1} \). For part (iv), an original proof could be found in p3338. Invoke part (i), \( \mathbb{P}(|Y| \geq \epsilon) \leq \exp[1]1 - (\epsilon/b_1^r)^{r_1} \). We have, for any \( \epsilon \geq 0 \), \( M := [1]\frac{(b_1^r)^{r_1+1}}{b_1^r}, b := b_1^r, r := \frac{r_1+1}{r_1+1}, \)
\[
\mathbb{P}(|XY| \geq \epsilon) \leq \mathbb{P}(|X| \geq \epsilon/M) + \mathbb{P}(|Y| \geq M) \leq \exp[1]1 - \frac{\epsilon/M^{r_1}}{b_1^r} + \exp[1]1 - \frac{M^{r_1}}{b_1^r}
\]
\[
= 2\exp[1]1 - (\epsilon/b)^r.
\]
Pick an $r_2 \in (0, r]$ and $b_2 > (1 + \log 2)^{1/r} b$. We consider the case $\epsilon \leq b_2$ first.

$$\mathbb{P}(|XY| \geq \epsilon) \leq 1 \leq \exp[1]1 - (\epsilon/b_2)^{r_2}.$$

We now consider the case $\epsilon > b_2$. Define a function $F(\epsilon) := (\epsilon/b)^r - (\epsilon/b_2)^{r_2}$. Using the definition of $b_2$, we have $F(b_2) > \log 2$. It is also not difficult to show that $F'(\epsilon) > 0$ when $\epsilon > b_2$. Thus we have $F(\epsilon) > F(b_2) > \log 2$ when $\epsilon > b_2$. Thus,

$$\mathbb{P}(|XY| \geq \epsilon) \leq 2 \exp[1]1 - (\epsilon/b)^r = \exp[1] \log 2 + 1 - (\epsilon/b)^r.$$

$$\leq \exp[(\epsilon/b)^r - (\epsilon/b_2)^{r_2} + 1 - (\epsilon/b)^r] = \exp[1 - (\epsilon/b_2)^{r_2}].$$

This following lemma gives the rate of convergence in terms of spectral norm for the sample covariance matrix.

**Lemma A.3.** Assume $n, T \to \infty$ simultaneously and $n/T \leq 1$. Suppose Assumptions 3.1(i) and 3.2 hold with $1/r_1 + 1/r_2 > 1$. Then

$$\|\hat{\Sigma}_T - \Sigma\|_{\ell_2} = O_p[3] \sqrt{\frac{n}{T}}.$$

**Proof.** Write $\hat{\Sigma}_T = \frac{1}{T} \sum_{t=1}^{T} y_t y_t^\top - \bar{y}\bar{y}^\top$. We have

$$\|\hat{\Sigma}_T - \Sigma\|_{\ell_2} \leq [3] \frac{1}{T} \sum_{t=1}^{T} y_t y_t^\top - \mathbb{E}y_t y_t^\top + \|\bar{y}\bar{y}^\top - \mu\mu^\top\|_{\ell_2}. \quad (A.1)$$

We consider the first term on the right hand side of (A.1) first. Invoke Lemma A.11 in Appendix A.5 with $\epsilon = 1/4$:

$$[3] \frac{1}{T} \sum_{t=1}^{T} y_t y_t^\top - \mathbb{E}y_t y_t^\top \leq 2 \max_{a \in \mathcal{N}_1} [3] a^\top [3] \frac{1}{T} \sum_{t=1}^{T} y_t y_t^\top - \mathbb{E}y_t y_t^\top : = 2 \max_{a \in \mathcal{N}_1} [3] \frac{1}{T} \sum_{t=1}^{T} (z_{a,t}^2 - \mathbb{E}z_{a,t}^2),$$

where $z_{a,t} := y_t^\top a$. First, invoke Lemma A.2(i) and (iv): For every $\epsilon \geq 0$, there exists an absolute constant $b_2 > 0$ such that

$$\mathbb{P}(|z_{a,t}^2| \geq \epsilon) \leq \exp[1]1 - (\epsilon/b_2)^{r_2}.$$

Next, invoke Lemma A.2(iii): For every $\epsilon > 0$, there exists an absolute constant $c_2 > 0$ such that

$$\mathbb{P}(|z_{a,t}^2 - \mathbb{E}z_{a,t}^2| \geq \epsilon) \leq \exp[1]1 - (\epsilon/c_2)^{r_2}.$$

Given Assumption 3.2 and the fact that mixing properties are hereditary in the sense that for any measurable function $m(\cdot)$, the process $\{m(y_t)\}$ possesses the mixing property of $\{y_t\}$ (7 p69), $z_{a,t}^2 - \mathbb{E}z_{a,t}^2$ is strong mixing with the same coefficient: $\alpha(h) \leq \exp[1] - K_3 h^{r_2}$. Define $r$ by $1/r := 2/r_1 + 1/r_2$. Using the fact that $2/r_1 + 1/r_2 > 1$, we can invoke a version of Bernstein’s inequality for strong mixing time series (Theorem A.2 in Appendix A.5), followed by Lemma A.12 in Appendix A.5:

$$2 \max_{a \in \mathcal{N}_1} [3] \frac{1}{T} \sum_{t=1}^{T} (z_{a,t}^2 - \mathbb{E}z_{a,t}^2) = O_p[3] \sqrt{\frac{\log |\mathcal{N}_1|}{T}}.$$
Invoking Lemma A.10 in Appendix A.5, we have \(|N_{1/4}| \leq 9^n\). Thus we have

\[
[3] \frac{1}{T} \sum_{t=1}^{T} y_t y_t^\top - \mathbb{E} y_t y_t^\top \leq 2 \max_{a \in N_{1/4}} [3] \frac{1}{T} \sum_{t=1}^{T} (z_{a,t}^2 - \mathbb{E} z_{a,t}^2) = O_p[3] \sqrt{\frac{n}{T}}.
\]

We now consider the second term on the right hand side of (A.1).

\[
\| \bar{y} y^\top - \mu \mu^\top \|_{\ell_2} = \| \bar{y} y^\top - \mu y^\top + \mu y^\top - \mu \mu^\top \|_{\ell_2} \leq 2 \max_{a \in N_{1/4}} [3] a^\top [3] \bar{y} y^\top - \mu y^\top + \mu y^\top - \mu \mu^\top a
\]

\[
= 2 \max_{a \in N_{1/4}} [3] a^\top [3] (\bar{y} - \mu) y^\top + \mu (\bar{y} - \mu)^\top a \leq 2 \max_{a \in N_{1/4}} [1] a^\top (\bar{y} - \mu) y^\top a + 2 \max_{a \in N_{1/4}} [1] a^\top \mu (\bar{y} - \mu)^\top a
\]

\[
\leq 2 \max_{a \in N_{1/4}} [1] a^\top (\bar{y} - \mu) \max_{a \in N_{1/4}} [1] \bar{y} \top a + 2 \max_{a \in N_{1/4}} [1] a^\top \mu \max_{a \in N_{1/4}} [1] (\bar{y} - \mu)^\top a.
\]

We consider \(\max_{a \in N_{1/4}} [1] (\bar{y} - \mu)^\top a\) first.

\[
\max_{a \in N_{1/4}} [1] (\bar{y} - \mu)^\top a = \max_{a \in N_{1/4}} [3] \frac{1}{T} \sum_{t=1}^{T} (y_t a - \mathbb{E}[y_t a]) =: \max_{a \in N_{1/4}} [3] \frac{1}{T} \sum_{t=1}^{T} (z_{a,t} - \mathbb{E} z_{a,t}).
\]

Recycling the proof for \(\max_{a \in N_{1/4}} [1] \frac{1}{T} \sum_{t=1}^{T} (z_{a,t} - \mathbb{E} z_{a,t}) = O_p[1] \sqrt{\frac{n}{T}}\) but with \(1/r := 1/r_1 + 1/r_2 > 1\), we have

\[
\max_{a \in N_{1/4}} [1] (\bar{y} - \mu)^\top a = \max_{a \in N_{1/4}} [3] \frac{1}{T} \sum_{t=1}^{T} (z_{a,t} - \mathbb{E} z_{a,t}) = O_p[3] \sqrt{\frac{\log |N_{1/4}|}{T}} = O_p[3] \sqrt{\frac{n}{T}}.
\]

(A.2)

Now let’s consider \(\max_{a \in N_{1/4}} [1] a^\top \mu\).

\[
\max_{a \in N_{1/4}} [1] a^\top \mu := \max_{a \in N_{1/4}} [1] \mathbb{E} a^\top y_t \leq \max_{a \in N_{1/4}} [1] \mathbb{E} a \mathbb{E} z_{a,t} = O(1),
\]

(A.3)

where the last equality is due to Lemma A.2(ii). Next we consider \(\max_{a \in N_{1/4}} [1] a^\top \bar{y}\).

\[
\max_{a \in N_{1/4}} [1] a^\top \bar{y} = \max_{a \in N_{1/4}} [1] a^\top (\bar{y} - \mu + \mu) \leq \max_{a \in N_{1/4}} [1] a^\top (\bar{y} - \mu) + \max_{a \in N_{1/4}} [1] a^\top \mu = O_p[3] \sqrt{\frac{n}{T}} + O(1)
\]

(A.4)

where the last equality is due to \(n \leq T\). Combining (A.2), (A.3) and (A.4), we have

\[
\| \bar{y} y^\top - \mu \mu^\top \|_{\ell_2} = O_p[3] \sqrt{\frac{n}{T}}.
\]

The following lemma gives the rate of convergence in terms of spectral norm for various quantities involving variances of \(y_t\). The rate \(\sqrt{n/T}\) is suboptimal, but there is no need improving it further as these quantities will not be the dominant terms in the proof of Theorem 3.1.

**Lemma A.4.** Suppose Assumptions 3.1(i), 3.2 3.3(i) and 3.4(i) hold with \(1/r_1 + 1/r_2 > 1\). Then
\( \| \hat{D}_T - D \|_{\ell_2} = O_p[3] \sqrt{\frac{n}{T}}. \)

**(ii)** The minimum eigenvalue of \( D \) is bounded away from zero by an absolute positive constant (i.e., \( \| D^{-1} \|_{\ell_2} = O(1) \)), so is the minimum eigenvalue of \( D^{1/2} \) (i.e., \( \| D^{-1/2} \|_{\ell_2} = O(1) \)).

\[ \| \hat{D}_T^{1/2} - D^{1/2} \|_{\ell_2} = O_p[3] \sqrt{\frac{n}{T}}. \]

**(iv)** \( \| \hat{D}_T^{-1/2} - D^{-1/2} \|_{\ell_2} = O_p[3] \sqrt{\frac{n}{T}}. \)

**(v)** \( \| \hat{D}_T^{-1/2} \|_{\ell_2} = O_p(1). \)

**(vi)** The maximum eigenvalue of \( \Sigma \) is bounded from the above by an absolute constant (i.e., \( \| \Sigma \|_{\ell_2} = O(1) \)). The maximum eigenvalue of \( D \) is bounded from the above by an absolute constant (i.e., \( \| D \|_{\ell_2} = O(1) \)).

\[ \| \hat{D}_T^{-1/2} \otimes \hat{D}_T^{-1/2} - D^{-1/2} \otimes D^{-1/2} \|_{\ell_2} = O_p[3] \sqrt{\frac{n}{T}}. \]

**Proof.** Define \( \sigma_i^2 := \mathbb{E}(y_{t,i} - \hat{y}_i)^2 \) and \( \hat{\sigma}_i^2 := \frac{1}{T} \sum_{t=1}^T (y_{t,i} - \bar{\hat{y}}_i)^2 \), where the subscript \( i \) denotes the \( i \)th component of the corresponding vector. For part (i),

\[
\| \hat{D}_T - D \|_{\ell_2} = \max_{1 \leq t \leq n} |\hat{\sigma}_i^2 - \sigma_i^2| = \max_{1 \leq t \leq n} |e_i^T(\hat{\Sigma}_T - \Sigma)e_i| \\
= \| \hat{\Sigma}_T - \Sigma \|_{\ell_2},
\]

where \( e_i \) denotes a unit vector whose \( i \)th component is 1. Now invoke Lemma \[A.13\] to get the result. For part (ii),

\[
\text{mineval}(D) = \min_{1 \leq i \leq n} \sigma_i^2 = \min_{1 \leq i \leq n} e_i^T \Sigma e_i \geq \min_{\|a\|_2=1} a^T \Sigma a = \text{mineval}(\Sigma) > 0
\]

where the last inequality is due to Assumption \[3.4(i)\]. For part (iii), invoking Lemma \[A.14\] in Appendix \[A.5\] gives

\[
\| \hat{D}_T^{1/2} - D^{1/2} \|_{\ell_2} \leq \frac{\| \hat{D}_T - D \|_{\ell_2}}{\text{mineval}(\hat{D}_T^{1/2}) + \text{mineval}(D^{1/2})} = O_p(1) \| \hat{D}_T - D \|_{\ell_2} = O_p[3] \sqrt{\frac{n}{T}},
\]

where the first and second equalities are due to parts (ii) and (i), respectively. Part (iv) follows from Lemma \[A.14\] in Appendix \[A.5\] via parts (ii) and (iii). For part (v),

\[
\| \hat{D}_T^{-1/2} \|_{\ell_2} = \| \hat{D}_T^{-1/2} - D^{-1/2} + D^{-1/2} \|_{\ell_2} \leq \| \hat{D}_T^{-1/2} - D^{-1/2} \|_{\ell_2} + \| D^{-1/2} \|_{\ell_2} \\
= O_p[3] \sqrt{\frac{n}{T}} + O(1) = O_p(1).
\]

37
For part (vi), we have

$$\|\Sigma\|_{\ell_2} = \max_{\|a\|_2 = 1} [1] a^T [1] \mathbb{E}[y_i y_i^T] - \mu \mu^T a \leq \max_{\|a\|_2 = 1} \mathbb{E} z_{a,t}^2 + \max_{\|a\|_2 = 1} (\mathbb{E} z_{a,t})^2 \leq 2 \max_{\|a\|_2 = 1} \mathbb{E} z_{a,t}^2.$$  

We have shown that in the proof of Lemma A.3 that $z_{a,t}^2$ has an exponential tail for any $\|a\|_2 = 1$. This says that $\mathbb{E} z_{a,t}^2$ is bounded for any $\|a\|_2 = 1$ via Lemma A.2(ii), so the result follows. Next we consider

$$\|D\|_{\ell_2} = \max_{1 \leq i \leq n} \sigma_i^2 = \max_{1 \leq i \leq n} e_i^T \Sigma e_i \leq \max_{\|a\|_2 = 1} \|a^T \Sigma a\| = \text{maxeval}(\Sigma) < \infty.$$  

For part (vii),

$$\|\hat{D}_{T^{-1/2}} \otimes \hat{D}_{T^{-1/2}} - D_{T^{-1/2}} \otimes D_{T^{-1/2}}\|_{\ell_2}$$  

$$= \|\hat{D}_{T^{-1/2}} \otimes \hat{D}_{T^{-1/2}} - \hat{D}_{T^{-1/2}} \otimes \hat{D}_{T^{-1/2}} + \hat{D}_{T^{-1/2}} \otimes \hat{D}_{T^{-1/2}} - D_{T^{-1/2}} \otimes D_{T^{-1/2}}\|_{\ell_2}$$  

$$\leq \|\hat{D}_{T^{-1/2}} \otimes \hat{D}_{T^{-1/2}}\|_{\ell_2} + \|\hat{D}_{T^{-1/2}} \otimes \hat{D}_{T^{-1/2}}\|_{\ell_2}$$  

$$= [1]\|\hat{D}_{T^{-1/2}}\|_{\ell_2} + [1]\|D_{T^{-1/2}}\|_{\ell_2} = O_p\left(\frac{n}{T}\right),$$  

where the second equality is due to Lemma A.16 in Appendix A.5.

To prove part (ii) of Theorem 3.1 we shall use Lemma 4.1 of ?. That lemma will further simplify when we consider real symmetric, positive definite matrices. For the ease of reference, we state this simplified version of Lemma 4.1 of ? here.

**Lemma A.5 (Simplified from Lemma 4.1 of ?).** For $n \times n$ real symmetric, positive definite matrices $A, B$, if

$$\|A - B\|_{\ell_2} < a,$$

for some absolute constant $a > 1$, then

$$\|\log A - \log B\|_{\ell_2} \leq C\|A - B\|_{\ell_2},$$

for some positive constant $C$.

**Proof.** First note that for any real symmetric, positive definite matrix $A$, $p(A, x) = x$ for any $x > 0$ in Lemma 4.1 of ?. Since $A$ is real symmetric and positive definite, all its eigenvalues lie in the region $|\arg(z - a)| \leq \pi/2$. Then according to ? p11, we have for any $t \geq 0$ not coinciding with eigenvalues of $A$

$$\rho(A, -t) \geq (a + t) \sin(\pi/2) = a + t$$

$$\rho(A, -t) - \delta \geq a + t - \delta,$$

where

$$\delta := \left\{\begin{array}{ll}
\|A - B\|_{\ell_2}^{1/n} & \text{if } \|A - B\|_{\ell_2} \leq 1 \\
\|A - B\|_{\ell_2} & \text{if } \|A - B\|_{\ell_2} \geq 1
\end{array}\right.$$  

and $\rho(A, -t)$ is defined in ? p3. Then the condition of Lemma A.5 allows one to invoke Lemma 4.1 of ? as

$$\rho(A, -t) \geq a + t \geq a > \delta.$$  

38
Lemma 4.1 of \cite{2020} says
\[
\| \log A - \log B \|_{\ell_2} \leq \| A - B \|_{\ell_2} \int_{0}^{\infty} \frac{1}{\rho(A,-t)} \frac{1}{\rho(A,-t) - \delta} \, dt
\]
\[
= \| A - B \|_{\ell_2} \int_{0}^{\infty} \frac{1}{\rho(A,-t) \rho(A,-t) - \delta} \, dt \leq \| A - B \|_{\ell_2} \int_{0}^{\infty} \frac{1}{(a+t)(a+t-\delta)} \, dt
\]
\[
\leq \| A - B \|_{\ell_2} \int_{0}^{\infty} \frac{1}{(a+t-\delta)^2} \, dt = \| A - B \|_{\ell_2} \frac{1}{a-\delta} =: C \| A - B \|_{\ell_2}.
\]

We are now ready to give a proof for Theorem 3.1

\textbf{Proof of Theorem 3.1} For part (i), recall that
\[
\hat{\Theta}_T = \tilde{D}_T^{-1/2} \hat{\Sigma}_T \tilde{D}_T^{-1/2}, \quad \Theta = D^{-1/2} \Sigma D^{-1/2}.
\]
Then we have
\[
\| \hat{\Theta}_T - \Theta \|_{\ell_2} = \| \tilde{D}_T^{-1/2} \hat{\Sigma}_T \tilde{D}_T^{-1/2} - \tilde{D}_T^{-1/2} \Sigma \tilde{D}_T^{-1/2} + \tilde{D}_T^{-1/2} \Sigma \tilde{D}_T^{-1/2} - D^{-1/2} \Sigma D^{-1/2} \|_{\ell_2}
\]
\[\leq \| \tilde{D}_T^{-1/2} \hat{\Sigma}_T - \Sigma \|_{\ell_2} + \| \tilde{D}_T^{-1/2} \Sigma \tilde{D}_T^{-1/2} - D^{-1/2} \Sigma D^{-1/2} \|_{\ell_2}. \tag{A.5}\]
Invoking Lemmas A.3 and A.4(v), we conclude that the first term of (A.5) is \(O_p(\sqrt{n/T})\). Let’s consider the second term of (A.5). Write
\[
\| \tilde{D}_T^{-1/2} \Sigma \tilde{D}_T^{-1/2} - D^{-1/2} \Sigma D^{-1/2} \|_{\ell_2}
\]
\[\leq \| (\tilde{D}_T^{-1/2} - D^{-1/2}) \Sigma \tilde{D}_T^{-1/2} \|_{\ell_2} + \| D^{-1/2} \Sigma (\tilde{D}_T^{-1/2} - D^{-1/2}) \|_{\ell_2}
\]
\[\leq \| \tilde{D}_T^{-1/2} \|_{\ell_2} \| \Sigma \|_{\ell_2} \| \tilde{D}_T^{-1/2} - D^{-1/2} \|_{\ell_2} + \| D^{-1/2} \|_{\ell_2} \| \Sigma \|_{\ell_2} \| \tilde{D}_T^{-1/2} - D^{-1/2} \|_{\ell_2}.
\]
Invoking Lemma A.4(ii), (iv), (v) and (vi), we conclude that the second term of (A.5) is \(O_p(\sqrt{n/T})\). For part (ii), it follows trivially from Lemma A.5 For part (iii), we have
\[
\| \hat{\theta}_T - \theta \|_{2} = \| (E^T W E)^{-1} E^T W \|_{\ell_2} \| D_+^T \|_{\ell_2} \log \hat{\Theta}_T - \log \Theta \|_F
\]
\[\leq \| (E^T W E)^{-1} E^T W \|_{\ell_2} \sqrt{n} \| \log \hat{\Theta}_T - \log \Theta \|_{\ell_2} = O(\sqrt{\log(W)/n}) \sqrt{n} O_p(\sqrt{n/T})
\]
\[= O_p[3] \sqrt{\frac{n \log(W)}{T}},
\]
where the first inequality is due to (A.8), and the second equality is due to (A.14) and parts (i)-(ii) of this theorem.

\textbf{A.4 The Proof of Theorem 3.2}

In this subsection, we give a proof for Theorem 3.2. We will first give some preliminary lemmas leading to the proof of this theorem.

The following lemma linearizes the matrix logarithm.
Lemma A.6. Suppose both $n \times n$ matrices $A+B$ and $A$ are real, symmetric, and positive definite for all $n$ with the minimum eigenvalues bounded away from zero by absolute constants. Suppose the maximum eigenvalue of $A$ is bounded from above by an absolute constant. Further suppose

$$[1]|t(A - I) + I|^{-1}tB_{t_2} \leq C < 1$$

(A.6)

for all $t \in [0,1]$ and some constant $C$. Then

$$\log(A + B) - \log A = \int_0^1 [t(A - I) + I]^{-1}B[t(A - I) + I]^{-1}dt + O(\|B\|^2_{l_2} \vee \|B\|^3_{l_2}).$$

The conditions of the preceding lemma implies that for every $t \in [0,1]$, $t(A - I) + I$ is positive definite for all $n$ with the minimum eigenvalue bounded away from zero by an absolute constant (see p181). Lemma A.6 has a flavour of Frechet derivative because $\int_0^1 [t(A - I) + I]^{-1}B[t(A - I) + I]^{-1}dt$ is the Frechet derivative of matrix logarithm at $A$ in the direction $B$ (see p272); however, this lemma is slightly stronger in the sense of a sharper bound on the remainder.

Proof. Since both $A + B$ and $A$ are positive definite for all $n$, with minimum eigenvalues real and bounded away from zero by absolute constants, by Theorem A.3 in Appendix A.5, we have

$$\log(A + B) = \int_0^1 (A + B - I)[t(A + B - I) + I]^{-1}dt, \quad \log A = \int_0^1 (A - I)[t(A - I) + I]^{-1}dt.$$

Use (A.6) to invoke Lemma A.15 in Appendix A.5 to expand $[t(A - I) + I + tB]^{-1}$ to get

$$[t(A - I) + I + tB]^{-1} = [t(A - I) + I]^{-1} - [t(A - I) + I]^{-1}tB[t(A - I) + I]^{-1} + O(\|B\|^2_{l_2})$$

and substitute into the expression of $\log(A + B)$

$$\log(A + B) = \int_0^1 (A + B - I) \left\{ [t(A - I) + I]^{-1} - [t(A - I) + I]^{-1}tB[t(A - I) + I]^{-1} + O(\|B\|^2_{l_2}) \right\} dt$$

$$= \log A + \int_0^1 B[t(A - I) + I]^{-1}dt - \int_0^1 (A + B - I)[t(A - I) + I]^{-1}B[t(A - I) + I]^{-1}dt$$

$$+ (A + B - I)O(\|B\|^2_{l_2})$$

$$= \log A + \int_0^1 [t(A - I) + I]^{-1}B[t(A - I) + I]^{-1}dt - \int_0^1 tB[t(A - I) + I]^{-1}B[t(A - I) + I]^{-1}dt$$

$$+ (A + B - I)O(\|B\|^2_{l_2})$$

$$= \log A + \int_0^1 [t(A - I) + I]^{-1}B[t(A - I) + I]^{-1}dt + O(\|B\|^2_{l_2} \vee \|B\|^3_{l_2}),$$

where the last equality follows from $\maxeval(A) < C < \infty$ and $\minival[t(A - I) + I] > C' > 0$. 

40
**Lemma A.7.** Suppose Assumptions 3.1(i), 3.2, 3.3(i) and 3.4(i) hold with $1/r_1 + 1/r_2 > 1$.

(i) Then $\Theta$ has minimum eigenvalue bounded away from zero by an absolute constant and maximum eigenvalue bounded from above by an absolute constant.

(ii) Then $\hat{\Theta}_T$ has minimum eigenvalue bounded away from zero by an absolute constant and maximum eigenvalue bounded from above by an absolute constant with probability approaching 1.

**Proof.** For part (i), the maximum eigenvalue of $\Theta$ is its spectral norm, i.e., $\|\Theta\|_{\ell_2}$.

$$\|\Theta\|_{\ell_2} = \|D^{-1/2} \Sigma D^{-1/2}\|_{\ell_2} \leq \|D^{-1/2}\|_{\ell_2}^2 \Sigma\|\ell_2 < C,$$

where the last inequality is due to Lemma A.4(ii) and (vi). Now let’s consider the minimum eigenvalue of $\Theta$.

$$\text{mineval} (\Theta) = \text{mineval}(D^{-1/2} \Sigma D^{-1/2}) = \min_{\|a\|_2=1} a^\top D^{-1/2} \Sigma D^{-1/2} a \geq \min_{\|a\|_2=1} \text{mineval}(\Sigma) \|D^{-1/2} a\|_2^2$$

$$= \text{mineval}(\Sigma) \min_{\|a\|_2=1} a^\top D^{-1} a = \text{mineval}(\Sigma) \text{mineval}(D^{-1}) = \frac{\text{mineval}(\Sigma)}{\text{maxeval}(D)} > 0,$$

where the second equality is due to Rayleigh-Ritz theorem, and the last inequality is due to Assumption 3.4(i) and Lemma A.4(vi). For part (ii), the maximum eigenvalue of $\hat{\Theta}_T$ is its spectral norm, i.e., $\|\hat{\Theta}\|_{\ell_2}$.

$$\|\hat{\Theta}_T\|_{\ell_2} \leq \|\hat{\Theta}_T - \Theta\|_{\ell_2} + \|\Theta\|_{\ell_2} = O_p[3] \sqrt{\frac{n}{T}} + \|\Theta\|_{\ell_2} = O_p(1)$$

where the first equality is due to Theorem 3.1(i) and the last equality is due to part (i). The minimum eigenvalue of $\hat{\Theta}_T$ is $1/\text{maxeval}(\hat{\Theta}_T^{-1})$. Since $\|\Theta^{-1}\|_{\ell_2} = \text{maxeval}(\Theta^{-1}) = 1/\text{mineval}(\Theta) = O(1)$ by part (i) and $\|\hat{\Theta}_T - \Theta\|_{\ell_2} = O_p(\sqrt{n/T})$ by Theorem 3.1(i), we can invoke Lemma A.14 in Appendix A.5 to get

$$\|\hat{\Theta}_T^{-1} - \Theta^{-1}\|_{\ell_2} = O_p(\sqrt{n/T}),$$

whence we have

$$\|\hat{\Theta}_T^{-1}\|_{\ell_2} \leq \|\hat{\Theta}_T^{-1} - \Theta^{-1}\|_{\ell_2} + \|\Theta^{-1}\|_{\ell_2} = O_p(1).$$

Thus the minimum eigenvalue of $\hat{\Theta}_T$ is bounded away from zero by an absolute constant. \hfill $\Box$

Define

$$\hat{H}_T := \int_0^1 [t(\hat{\Theta}_T - I) + I]^{-1} \otimes [t(\hat{\Theta}_T - I) + I]^{-1} dt.$$

The following lemma gives the rate of convergence for $\hat{H}_T$. The following lemma is also true when one replaces $H_T$ with $H_{T,D}$.

**Lemma A.8.** Let Assumptions 3.1(i), 3.2, 3.3(i) and 3.4(i) be satisfied with $1/r_1 + 1/r_2 > 1$. Then we have

$$\|H\|_{\ell_2} = O(1), \quad \|\hat{H}_T\|_{\ell_2} = O_p(1), \quad \|\hat{H}_T - H\|_{\ell_2} = O_p[3] \sqrt{\frac{n}{T}}. \quad (A.7)$$
Proof. The proofs for \(|H|l_2 = O(1)\) and \(|\hat{H}_T|l_2 = O_p(1)\) are exactly the same, so we only give the proof for the latter. Define \(A_t := [t(\hat{\Theta}_T - I) + I]^{-1}\) and \(B_t := [t(\Theta - I) + I]^{-1}\).

\[
\|\hat{H}_T\|l_2 = [3] \int_0^1 A_t \otimes A_t dt \leq \int_0^1 [1] A_t \otimes A_t dt \leq \max_{t \in [0, 1]} [1] A_t \otimes A_t \leq \max_{t \in [0, 1]} [0] A_t^2 \leq \max_{t \in [0, 1]} \{\text{maxeval}([t(\hat{\Theta}_T - I) + I]^{-1})\}^2 = \max_{t \in [0, 1]} [3] \frac{1}{\text{mineval}(t(\hat{\Theta}_T - I) + I)}^2 = O_p(1),
\]

where the second equality is due to Lemma A.14 in Appendix A.5, and the last equality is due to Lemma A.7(ii). Now,

\[
\|\hat{H}_T - H\|l_2 = [3] \int_0^1 A_t \otimes A_t - B_t \otimes B_t dt \leq \int_0^1 \|A_t \otimes A_t - B_t \otimes B_t\|l_2 dt \leq \max_{t \in [0, 1]} \|A_t \otimes A_t - B_t \otimes B_t\|l_2 \leq \max_{t \in [0, 1]} \|A_t \otimes (A_t - B_t)\|l_2 + \|B_t \otimes B_t\|l_2 = \max_{t \in [0, 1]} \|A_t \otimes (A_t - B_t)\|l_2 + \|B_t \otimes B_t\|l_2 = O_p(1) \max_{t \in [0, 1]} \|t(\hat{\Theta}_T - I) + I^{-1} - [t(\Theta - I) + I^{-1}]\|l_2 \]

where the first inequality is due to Jensen’s inequality, the third equality is due to special properties of Kronecker product, the fourth equality is due to Lemma A.16 in Appendix A.5, and the last equality is because Lemma A.7 implies

\[
\|t(\hat{\Theta}_T - I) + I^{-1}\|l_2 = O_p(1) \quad \|t(\Theta - I) + I^{-1}\|l_2 = O(1).
\]

Now

\[
\|t(\hat{\Theta}_T - I) + I - [t(\Theta - I) + I]\|l_2 = t\|\hat{\Theta}_T - \Theta\|l_2 = O_p(\sqrt{n/T}),
\]

where the last equality is due to Theorem 3.1(i). The lemma then follows after invoking Lemma A.14 in Appendix A.5.

\[
\boxed{\text{Lemma A.9. Given the } n^2 \times n(n + 1)/2 \text{ duplication matrix } D_n \text{ and its Moore-Penrose generalised inverse } D_n^+ = (D_n^T D_n)^{-1} D_n^T \text{ (i.e., } D_n \text{ is full-column rank), we have}}
\]

\[
\|D_n^+\|l_2 = \|D_n^{+T}\|l_2 = 1, \quad \|D_n\|l_2 = \|D_n^T\|l_2 = 2.
\]

\[
\boxed{\text{A.8}}
\]

\[
\text{Proof. First note that } D_n^T D_n \text{ is a diagonal matrix with diagonal entries either 1 or 2. Using the fact that for any matrix } A, AA^T \text{ and } A^TA \text{ have the same non-zero eigenvalues, we have}
\]

\[
\|D_n^{+T}\|l_2 = \text{maxeval}(D_n^+ D_n^{+T}) = \text{maxeval}((D_n^T D_n)^{-1}) = 1
\]

\[
\|D_n^T\|l_2 = \text{maxeval}(D_n^{T+} D_n^T) = \text{maxeval}(D_n^T D_n^T) = \text{maxeval}((D_n^T D_n)^{-1}) = 1
\]

\[
\|D_n\|l_2 = \text{maxeval}(D_n^T D_n) = 2
\]

\[
\|D_n^+\|l_2 = \text{maxeval}(D_n^+ D_n^T) = \text{maxeval}(D_n^T D_n) = 2
\]
We are now ready to give a proof for Theorem 3.2.

**Proof of Theorem 3.2.** We first show that (A.6) is satisfied with probability approaching 1 for \( A = \Theta \) and \( B = \hat{\Theta}_T - \Theta \). That is,

\[
\| [t(\Theta - I) + I]^{-1} t(\hat{\Theta}_T - \Theta) \|_{\ell_2} \leq C < 1 \quad \text{with probability approaching 1},
\]

for some constant \( C \).

\[
\| [t(\Theta - I) + I]^{-1} t(\hat{\Theta}_T - \Theta) \|_{\ell_2} \leq \| [t(\Theta - I) + I]^{-1} \|_{\ell_2} \| \hat{\Theta}_T - \Theta \|_{\ell_2} \leq t \| [t(\Theta - I) + I]^{-1} \|_{\ell_2} O_p(\sqrt{n/T}) = O_p(\sqrt{n/T})/\text{mineval}(t(\Theta - I) + I) = o_p(1),
\]

where the first equality is due to Theorem 3.1(i), and the last equality is due to mineval(\( t(\Theta - I) + I \)) \( > C > 0 \) for some absolute constant \( C \) (implied by Lemma A.7(i)) and Assumption 3.3(i). Together with Lemma A.7(ii) and Lemma 2.12 in ?, we can invoke Lemma A.6 stochastically with \( A = \Theta \) and \( B = \hat{\Theta}_T - \Theta \):

\[
\log \hat{\Theta}_T - \log \Theta = \int_0^1 [t(\Theta - I) + I]^{-1}(\hat{\Theta}_T - \Theta)[t(\Theta - I) + I]^{-1} dt + O_p(\| \hat{\Theta}_T - \Theta \|_{\ell_2}^2). \tag{A.9}
\]

(We can invoke Lemma A.6 stochastically because the remainder of the log linearization is zero when the perturbation is zero. Moreover, we have \( \| \hat{\Theta}_T - \Theta \|_{\ell_2} \overset{p}{\to} 0 \) under Assumption 3.3(i).) Note that (A.9) also holds with \( \hat{\Theta}_T \) replaced by \( \hat{\Theta}_{T,D} \) by repeating the same argument. That is,

\[
\log \hat{\Theta}_{T,D} - \log \Theta = \int_0^1 [t(\Theta - I) + I]^{-1}(\hat{\Theta}_{T,D} - \Theta)[t(\Theta - I) + I]^{-1} dt + O_p(\| \hat{\Theta}_{T,D} - \Theta \|_{\ell_2}^2).
\]

Now we can write

\[
\sqrt{T} c^\top (\hat{\Theta}_{T,D} - \Theta^\top) = \sqrt{T} c^\top (E^\top WE)^{-1} E^\top WD_n^+ HD_{n}^1 \otimes D_{n}^{-1/2} \text{vec}(\hat{\Sigma}_T - \Sigma) \sqrt{c^\top J_{T,D} C} + \frac{\sqrt{T} c^\top (E^\top WE)^{-1} E^\top WD_n^+ \text{vec} O_p(\| \hat{\Theta}_{T,D} - \Theta \|_{\ell_2}^2)}{\sqrt{c^\top J_{T,D} C}} = \hat{t}_{D,1} + \hat{t}_{D,2}.
\]

Define

\[
t_{D,1} := \frac{\sqrt{T} c^\top (E^\top WE)^{-1} E^\top WD_n^+ HD_{n}^1 \otimes D_{n}^{-1/2} \text{vec}(\hat{\Sigma}_T - \Sigma)}{\sqrt{c^\top J_{D} C}},
\]

To prove Theorem 3.2 it suffices to show \( t_{D,1} \overset{d}{\to} N(0,1) \), \( t_{D,1} - \hat{t}_{D,1} = o_p(1) \), and \( \hat{t}_{D,2} = o_p(1) \).
A.4.1 \( t_{D,1} \overset{d}{\to} N(0,1) \)

We now prove that \( t_{D,1} \) is asymptotically distributed as a standard normal.

\[
  t_{D,1} = \frac{\sqrt{T}c^\top (E^\top WE)^{-1}E^\top WD_n^{+}H(D^{-1/2} \otimes D^{-1/2})vec[2]^\top \sum_{t=1}^{T} [y_t - \mu](y_t - \mu)^\top - \mathbb{E}(y_t - \mu)(y_t - \mu)^\top}{\sqrt{c^\top J_D c}}
\]

\[
  = \sum_{t=1}^{T} \frac{T^{-1/2}c^\top (E^\top WE)^{-1}E^\top WD_n^{+}H(D^{-1/2} \otimes D^{-1/2})vec[1](y_t - \mu)(y_t - \mu)^\top - \mathbb{E}(y_t - \mu)(y_t - \mu)^\top}{\sqrt{c^\top J_D c}}
\]

\[
  =: \sum_{t=1}^{T} U_{D,T,n,t}.
\]

Define a triangular array of sigma algebras \( \{\mathcal{F}_{T,n,t}, t = 0, 1, 2, \ldots, T\} \) by \( \mathcal{F}_{T,n,t} := \mathcal{F}_t \)

(only non-standard thing is that this triangular array has one more subscript \( n \)). It is easy to see that \( U_{D,T,n,t} \) is \( \mathcal{F}_{T,n,t} \)-measurable. We now show that \( \{U_{D,T,n,t}, \mathcal{F}_{T,n,t}\} \)

is a martingale difference sequence (i.e., \( \mathbb{E}[U_{D,T,n,t} | \mathcal{F}_{T,n,t-1}] = 0 \) almost surely for \( t = 1, \ldots, T \)). It suffices to show for all \( t \)

\[
  \mathbb{E}[1](y_t - \mu)(y_t - \mu)^\top - \mathbb{E}[(y_t - \mu)(y_t - \mu)^\top] | \mathcal{F}_{T,n,t-1} = 0 \quad a.s. \quad (A.10)
\]

This is straightforward via Assumption 3.5. Now we check conditions (i)-(iii) of Theorem A.5 in Appendix A.5. We first investigate at what rate the denominator \( \sqrt{c^\top J_D c} \) goes to zero:

\[
  c^\top J_D c = c^\top (E^\top WE)^{-1}E^\top WD_n^{+}H(D^{-1/2} \otimes D^{-1/2})V(D^{-1/2} \otimes D^{-1/2})HD_n^{+\top}WE(E^\top WE)^{-1}c
\]

\[
  \geq \text{mineval}(V)\text{mineval}(D^{-1} \otimes D^{-1})\text{mineval}(H^2)\text{mineval}(D_n^{+}D_n^{+\top})\text{mineval}(W)\text{mineval}((E^\top WE)^{-1})
\]

\[
  = \frac{\text{mineval}(V)\text{mineval}^2(H)}{\text{maxeval}(D \otimes D)\text{maxeval}(D_n^{+}D_n^{+\top})\text{maxeval}(W^{-1})\text{maxeval}(E^\top WE)}
\]

\[
  \geq \frac{\text{mineval}(V)\text{mineval}^2(H)}{\text{maxeval}(D \otimes D)\text{maxeval}(D_n^{+}D_n^{+\top})\text{maxeval}(W^{-1})\text{maxeval}(W)\text{maxeval}(E^\top E)}
\]

where the first and third inequalities are true by repeatedly invoking the Rayleigh-Ritz theorem. Note that

\[
  \text{maxeval}(E^\top E) \leq \text{tr}(E^\top E) \leq s \cdot n, \quad (A.11)
\]

where the last inequality is due to Lemma A.1. For future reference

\[
  \|E\|_{\ell_2} = \|E^\top\|_{\ell_2} = \sqrt{\text{maxeval}(E^\top E)} \leq \sqrt{sn}. \quad (A.12)
\]

Since the minimum eigenvalue of \( H \) is bounded away from zero by an absolute constant by Lemma A.7(i), the maximum eigenvalue of \( D \) is bounded from above by an absolute constant (Lemma A.4(vi)), and \( \text{maxeval}[D_n^{+}D_n^{+\top}] \) is bounded from above since \( D_n^{+}D_n^{+\top} \) is a diagonal matrix with diagonal entries either 1 or 2, we have, via Assumption 3.6

\[
  \frac{1}{\sqrt{c^\top J_D c}} = O(\sqrt{s \cdot n \cdot \kappa(W)}). \quad (A.13)
\]
Also note that
\[
\|(E^TW^E)^{-1}E^TW^{1/2}\|_2 = \sqrt{\maxval(1)[(E^TW^E)^{-1}E^TW^{1/2} ]}
\]
\[
= \sqrt{\maxval(1)(E^TW^E)^{-1}E^TW^{1/2}}
\]
\[
= \sqrt{\maxval(1)(E^TW^E)^{-1}E^TW^{1/2}W^{1/2}E(E^TW^E)^{-1}}
\]
\[
= \sqrt{\maxval(1)(E^TW^E)^{-1}} = \sqrt{\frac{1}{\minval(E^TW^E)}} \leq \sqrt{\frac{1}{\minval(E^TW^E)\minval(W)}}
\]
\[
= O[1] \sqrt{\frac{\kappa(W)}{n} \sqrt{\|W^{-1}\|_2}}
\]
where the second equality is due to the fact that for any matrix $A$, $AA^T$ and $A^TA$ have
the same non-zero eigenvalues, the third equality is due to $(A^T)^{-1} = (A^{-1})^T$, and the last
equality is due to Assumption 3.4(ii). Thus
\[
\|(E^TW^E)^{-1}E^TW\|_2 = O(\sqrt{\kappa(W)/n}),
\]
whence we have
\[
[1]c^T(E^TW^E)^{-1}E^TW^2D_n^+H(D^{-1/2} \otimes D^{-1/2})_2 = O(\sqrt{\kappa(W)/n}),
\]
via (A.7) and Lemma A.4(ii). We now verify (i) and (ii) of Theorem A.5 in Appendix A.5.
We shall use Orlicz norms as defined in \[?\]: Let $\psi: \mathbb{R}^+ \to \mathbb{R}^+$ be a non-decreasing, convex
function with $\psi(0) = 0$ and $\lim_{x \to \infty} \psi(x) = \infty$, where $\mathbb{R}^+$ denotes the set of nonnegative
real numbers. Then, the Orlicz norm of a random variable $X$ is given by
\[
X_\psi = \inf \{C > 0 : \mathbb{E}\psi(|X|/C) \leq 1\},
\]
where $\inf \emptyset = \infty$. We shall use Orlicz norms for $\psi(x) = \psi_p(x) = e^{px} - 1$ for $p = 1, 2$ in
this article. We consider $|U_{D,T,n,t}|$ first.
\[
|U_{D,T,n,t}| = [3] \sqrt{T^{-1/2}c^T(E^TW^E)^{-1}E^TW^2D_n^+H(D^{-1/2} \otimes D^{-1/2})_2 \vec{1}[1](y_t - \mu)(y_t - \mu)^T - \mathbb{E}(y_t - \mu)(y_t - \mu)^T}
\]
\[
\leq \sqrt{T^{-1/2}c^T(E^TW^E)^{-1}E^TW^2D_n^+H(D^{-1/2} \otimes D^{-1/2})_2 \vec{1}[1](y_t - \mu)(y_t - \mu)^T - \mathbb{E}(y_t - \mu)(y_t - \mu)^T}_F
\]
\[
\leq O[3] \sqrt{\frac{\kappa^2(W)}{T}} [1](y_t - \mu)(y_t - \mu)^T - \mathbb{E}(y_t - \mu)(y_t - \mu)^T}_\infty
\]
where the second equality is due to (A.13) and (A.15). Consider
\[
[2][1](y_t - \mu)(y_t - \mu)^T - \mathbb{E}(y_t - \mu)(y_t - \mu)^T_\infty \psi_1
\]
\[
= [2] \max_{1 \leq i < j \leq n} [1](y_{t,i} - \mu_i)(y_{t,j} - \mu_j) - \mathbb{E}(y_{t,i} - \mu_i)(y_{t,j} - \mu_j)\psi_1
\]
\[
\leq \log(1 + n^2) \max_{1 \leq i < j \leq n} [2](y_{t,i} - \mu_i)(y_{t,j} - \mu_j) - \mathbb{E}(y_{t,i} - \mu_i)(y_{t,j} - \mu_j)\psi_1
\]
\[
\leq 2 \log(1 + n^2) \max_{1 \leq i < j \leq n} [2](y_{t,i} - \mu_i)(y_{t,j} - \mu_j)\psi_1
\]
where the first inequality is due to Lemma 2.2.2 in ?. Assumption 3.1(i) with $r_1 = 2$ gives $\mathbb{E}[1]\exp(K_2|y_{t,i}|^2) \leq K_1$ for all $i$. Then

$$
P[1](y_{t,i} - \mu_i)(y_{t,j} - \mu_j) \geq \varepsilon \leq \mathbb{P}[1]|y_{t,i} - \mu_i| \geq \sqrt{\varepsilon} + \mathbb{P}[1]|y_{t,j} - \mu_j| \geq \sqrt{\varepsilon}$$

$$
\leq 2\exp[1] - (\sqrt{\varepsilon}/c_1)^2 =: Ke^{-\varepsilon}
$$

where the second inequality is due to Lemma A.2(iii). It follows from Lemma 2.2.1 in ? that $\|y_{t,i} - \mu_i)(y_{t,j} - \mu_j)||_\psi \leq (1 + K)/C$ for all $i, j, t$. Thus

$$
[2]\max_{1 \leq t \leq T} |U_{D,T,n,t}| \leq \log(1 + T) \max_{1 \leq t \leq T} [2]U_{D,T,n,t}\psi_1
$$

where the first inequality is due to Assumption 3.3(ii). Since $\|U\|_L \leq r\|U\|_\psi$ for any random variable $U$ (? , p95), we conclude that (i) and (ii) of Theorem A.5 in Appendix A.5 are satisfied. We now verify condition (iii) of Theorem A.5 in Appendix A.5. Since we have already shown in (A.13) that $sn\kappa(W)c^nJ_{DC}$ is bounded away from zero by an absolute constant, it suffices to show

$$
sn\kappa(W) \cdot [3] \frac{1}{T} \sum_{t=1}^{T} [2]c^T(E^NWE)^{-1}E^NWD_n^+H(D^{-1/2} \otimes D^{-1/2})u_t^2 - c^TJ_{DC} = o_p(1),
$$

where $u_t := \text{vec}[1](y_{t}-\mu)(y_{t}-\mu)^T - \mathbb{E}(y_{t}-\mu)(y_{t}-\mu)^T$. Note that

$$
sn\kappa(W) \cdot [3] \frac{1}{T} \sum_{t=1}^{T} [2]c^T(E^NWE)^{-1}E^NWD_n^+H(D^{-1/2} \otimes D^{-1/2})u_t^2 - c^TJ_{DC}
$$

$$
\leq sn\kappa(W)[3] \frac{1}{T} \sum_{t=1}^{T} u_tu_t^T - V \|c^T(E^NWE)^{-1}E^NWD_n^+H(D^{-1/2} \otimes D^{-1/2})\|_1^2
$$

$$
\leq sn^3\kappa(W)[3] \frac{1}{T} \sum_{t=1}^{T} u_tu_t^T - V \|c^T(E^NWE)^{-1}E^NWD_n^+H(D^{-1/2} \otimes D^{-1/2})\|_1^2
$$

$$
\leq sn^3\kappa(W)[3] \frac{1}{T} \sum_{t=1}^{T} u_tu_t^T - V \|(E^NWE)^{-1}E^NWD_n^+H(D^{-1/2} \otimes D^{-1/2})\|_1^2 \|D_n^+\|_1^2 \|H\|_1^2 \|D^{-1/2} \otimes D^{-1/2}\|_2^2
$$

$$
= O_p(sn^3\kappa(W)) \frac{\log n}{T} \cdot \frac{\kappa(W)}{n} = O_p[3] \frac{\sqrt{s^2n^4\kappa^4(W)} \log n \cdot \kappa(W)}{T} = o_p(1)
$$

where the first equality is due to Lemma A.4(ii), Lemma A.16 in Appendix A.5 (A.7), (A.14), (A.8), and the fact that $[1]^T\sum_{t=1}^{T} u_tu_t^T - V = O_p(\frac{\log n}{T})$, which can be deduced from the proof of Lemma ?? in SM ??, the last equality is due to Assumption 3.3(ii). Thus condition (iii) of Theorem A.5 in Appendix A.5 is verified and $t_{D,1} \overset{d}{\rightarrow} N(0, 1)$. 

46
A.4.2 $t_{D,1} - \hat{t}_{D,1} = o_p(1)$

We now show that $t_{D,1} - \hat{t}_{D,1} = o_p(1)$. Let $A_D$ and $\hat{A}_D$ denote the numerators of $t_{D,1}$ and $\hat{t}_{D,1}$, respectively.

$$t_{D,1} - \hat{t}_{D,1} = \frac{A_D}{\sqrt{c^T J_{DC}}} - \frac{\hat{A}_D}{\sqrt{c^T \hat{J}_{T,DC}}} = \frac{\sqrt{\text{snk}(W)} A_D}{\sqrt{\text{snk}(W)c^T J_{DC}}} - \frac{\sqrt{\text{snk}(W)} \hat{A}_D}{\sqrt{\text{snk}(W)c^T \hat{J}_{T,DC}}}.$$ 

Since we have already shown in (A.13) that $\text{snk}(W)c^T J_{DC}$ is bounded away from zero by an absolute constant, it suffices to show the denominators as well as numerators of $t_{D,1}$ and $\hat{t}_{D,1}$ are asymptotically equivalent.

A.4.3 Denominators of $t_{D,1}$ and $\hat{t}_{D,1}$

We first show that the denominators of $t_{D,1}$ and $\hat{t}_{D,1}$ are asymptotically equivalent, i.e.,

$$\text{snk}(W)|c^T \hat{J}_{T,DC} - c^T J_{DC}| = o_p(1).$$

Define

$$c^T \hat{J}_{T,DC} := c^T (E^T WE)^{-1} E^T WD_n^+ \hat{H}_{T,D}(D^{-1/2} \otimes D^{-1/2})V(D^{-1/2} \otimes D^{-1/2})\hat{H}_{T,D}D_n^+ WE(E^T WE)^{-1}c.$$ 

By the triangular inequality:

$$|\text{snk}(W)c^T \hat{J}_{T,DC} - \text{snk}(W)c^T J_{DC}| \leq |\text{snk}(W)c^T \hat{J}_{T,DC} - \text{snk}(W)c^T \hat{J}_{T,DC}| + |\text{snk}(W)c^T \hat{J}_{T,DC} - \text{snk}(W)c^T J_{DC}|.$$ 

First, we prove $|\text{snk}(W)c^T \hat{J}_{T,DC} - \text{snk}(W)c^T J_{DC}| = o_p(1)$.

$$\text{snk}(W)|c^T \hat{J}_{T,DC} - c^T J_{DC}| = \text{snk}(W)|c^T (E^T WE)^{-1} E^T WD_n^+ \hat{H}_{T,D}(D^{-1/2} \otimes D^{-1/2})V(D^{-1/2} \otimes D^{-1/2})\hat{H}_{T,D}D_n^+ WE(E^T WE)^{-1}c - c^T (E^T WE)^{-1} E^T WD_n^+ \hat{H}_{T,D}(D^{-1/2} \otimes D^{-1/2})V(D^{-1/2} \otimes D^{-1/2})\hat{H}_{T,D}D_n^+ WE(E^T WE)^{-1}c|$$

$$\leq \text{snk}(W)||\hat{V}_T - V||_\infty||(D^{-1/2} \otimes D^{-1/2})\hat{H}_{T,D}D_n^+ WE(E^T WE)^{-1}c||_2^2$$

$$\leq \text{snk}(W)||\hat{V}_T - V||_\infty||(D^{-1/2} \otimes D^{-1/2})||_{\ell_2}^2 \|\hat{H}_{T,D}||_{\ell_2}^2 \|D_n^+||_{\ell_2}$$

$$= O_p(\text{snk}(W)^2\|\hat{V}_T - V\|_\infty = O_p[3]\sqrt{\frac{n^4 k^4(W) s^2 \omega^2 \log n}{T}} = o_p(1),$$

where $|| \cdot ||_{\ell_\infty}$ denotes the absolute elementwise maximum, the third equality is due to Lemma A.4(ii), Lemma A.16 in Appendix A.5, (A.7), (A.14), and (A.8), the second last equality is due to Lemma ?? in SM ??, and the last equality is due to Assumption 3.3(ii). We now prove $|\text{snk}(W)c^T \hat{J}_{T,DC} - c^T J_{DC}| = o_p(1)$.

$$\text{snk}(W)|c^T \hat{J}_{T,DC} - c^T J_{DC}| = \text{snk}(W)|c^T (E^T WE)^{-1} E^T WD_n^+ \hat{H}_{T,D}(D^{-1/2} \otimes D^{-1/2})V(D^{-1/2} \otimes D^{-1/2})\hat{H}_{T,D}D_n^+ WE(E^T WE)^{-1}c - c^T (E^T WE)^{-1} E^T WD_n^+ \hat{H}_{T,D}(D^{-1/2} \otimes D^{-1/2})V(D^{-1/2} \otimes D^{-1/2})HD_n^+ WE(E^T WE)^{-1}c|$$

$$\leq \text{snk}(W)[1] \max\{1\}|(D^{-1/2} \otimes D^{-1/2})V(D^{-1/2} \otimes D^{-1/2})^2 ||(\hat{H}_{T,D} - H)D_n^+ WE(E^T WE)^{-1}c||_2^2$$

$$+ 2\text{snk}(W)||((D^{-1/2} \otimes D^{-1/2})V(D^{-1/2} \otimes D^{-1/2}) HD_n^+ WE(E^T WE)^{-1}c||_2$$

$$\cdot ||(\hat{H}_{T,D} - H)D_n^+ WE(E^T WE)^{-1}c||_2$$

(A.16)
where the inequality is due to Lemma A.17 in Appendix A.5. We consider the first term of (A.16) first.

\[ sn(W)[1] \max_{\|H\|_{\ell_2} \leq \epsilon} [D^{-1/2} \otimes D^{-1/2}] V (D^{-1/2} \otimes D^{-1/2})^2 \| (\hat{H}_{T, D} - H) D_n^{+T} W E (E^T W E)^{-1} c \|_2^2 \]

\[ = O(snk(W)) \| \hat{H}_{T, D} - H \|_{\ell_2}^2 \| D_n^{+T} \|_{\ell_2}^2 \| W E (E^T W E)^{-1} c \|_2^2 \]

\[ = O_p(snk^2(W) \omega/T) = o_p(1), \]

where the second last equality is due to (A.7), (A.8), and (A.14), and the last equality is due to Assumption 3.3(ii). We consider the first term of (A.16).

\[ 2snk(W)((D^{-1/2} \otimes D^{-1/2}) V (D^{-1/2} \otimes D^{-1/2}) H D_n^{+T} W E (E^T W E)^{-1} c \|_2 \]

\[ \cdot \| (\hat{H}_{T, D} - H) D_n^{+T} W E (E^T W E)^{-1} c \|_2 \]

\[ \leq O(snk(W)) \| H \|_{\ell_2} \| \hat{H}_{T, D} - H \|_{\ell_2} \| D_n^{+T} \|_{\ell_2} \| W E (E^T W E)^{-1} c \|_2^2 = O(\sqrt{nK^4(W)s^2\omega^2/T}) = o_p(1), \]

where the first equality is due to (A.7), (A.8), and (A.14), and the last equality is due to Assumption 3.3(ii). We have proved \(|snk(W)c^T \hat{J}_{T, DC} - snk(W)c^T J_{DC}| = o_p(1)| and hence \(|snk(W)c^T \hat{J}_{T, DC} - snk(W)c^T J_{DC}| = o_p(1). |

### A.4.4 Numerators of \(t_{D,1}\) and \(\hat{t}_{D,1}\)

We now show that numerators of \(t_{D,1}\) and \(\hat{t}_{D,1}\) are asymptotically equivalent, i.e.,

\[ \sqrt{snk(W)}|A_D - \hat{A}_D| = o_p(1). \]

This is relatively straightforward.

\[ \sqrt{Tsnk(W)[1]c^T (E^T W E)^{-1} E^T W D_n^{+T} H (D^{-1/2} \otimes D^{-1/2}) \text{vec}(\Sigma_T - \hat{\Sigma}_T + \Sigma) \]

\[ = \sqrt{Tsnk(W)[1]c^T (E^T W E)^{-1} E^T W D_n^{+T} H (D^{-1/2} \otimes D^{-1/2}) \text{vec}(\Sigma_T - \hat{\Sigma}_T) \]

\[ = \sqrt{Tsnk(W)[1]c^T (E^T W E)^{-1} E^T W D_n^{+T} H (D^{-1/2} \otimes D^{-1/2}) \text{vec}(1)(\bar{y} - \mu)(\bar{y} - \mu)^T \]

\[ \leq \sqrt{Tsnk(W)||E^T W E||_{\ell_2} D_n^{+T} \|H\|_{\ell_2} D^{-1/2} \otimes D^{-1/2} \|_{\ell_2} \text{vec}(1)(\bar{y} - \mu)(\bar{y} - \mu)^T \|_2 \]

\[ = O(\sqrt{Tsnk(W)}\sqrt{\pi\kappa(W)/n}\|\bar{y} - \mu\|\|\bar{y} - \mu\|^T \|_F \]

\[ \leq O(\sqrt{Tsnk(W)}\sqrt{\pi\kappa(W)/n}\|\bar{y} - \mu\|\|\bar{y} - \mu\|^T \|_\infty \]

\[ = O(\sqrt{Tsnk^2(W)\omega} \max_{1 \leq i \leq n} [1](\bar{y} - \mu)_i(\bar{y} - \mu)_j = O_p(\sqrt{Tsnk^2(W)\omega}) \log n/T \]

\[ = O_p[3\sqrt{\log^3 n \cdot n^2 \kappa^2(W)\omega T} = o_p(1), \]

where the third equality is due to (A.7), (A.8), and (A.14), the third last equality is due to (??), and the last equality is due to Assumption 3.3(ii).

### A.4.5 \(\hat{t}_{D,2} = o_p(1)\)

Write

\[ \hat{t}_{D,2} = \frac{\sqrt{T}\sqrt{\pi\kappa(W)c^T (E^T W E)^{-1} E^T W D_n^{+T} \text{vec} O_p(\|H_{T, D} - \hat{H}_{T, D}\|_{\ell_2}^2)} \sqrt{snk(W)c^T \hat{J}_{T, DC}}}{\sqrt{snk(W)c^T \hat{J}_{T, DC}}} = o_p(1). \]
Since the denominator of the preceding equation is bounded away from zero by an absolute constant with probability approaching one by (A.13) and that $|sk(W)c^TJ_{T,D}c - snk(W)c^TJ_{D}c| = o_p(1)$, it suffices to show

$$\sqrt{T}\sqrt{\frac{skn(W)}{T}}c^T(E^TW)E^{-1}E^WD_n^+\text{vec } O_p(\|\hat{\Theta}_{T,D} - \Theta\|_{l_2}^2) = o_p(1).$$

This is straightforward:

$$|\sqrt{T}\sqrt{skn(W)}c^T(E^TW)E^{-1}E^WD_n^+\text{vec } O_p(\|\hat{\Theta}_{T,D} - \Theta\|_{l_2}^2)|$$

$$\leq \sqrt{T}\sqrt{skn(W)}\|c^T(E^TW)E^{-1}E^WD_n^+\|_2\text{vec } O_p(\|\hat{\Theta}_{T,D} - \Theta\|_{l_2}^2)\|_F$$

$$= O(\sqrt{T}\sqrt{skn(W)})\|O_p(\|\hat{\Theta}_{T,D} - \Theta\|_{l_2}^2)\|_F = O(\sqrt{T}\sqrt{skn(W)})\|O_p(\|\hat{\Theta}_{T,D} - \Theta\|_{l_2}^2)\|_{l_2}$$

$$= O(\sqrt{T}\sqrt{skn(W)})O_p(\|\hat{\Theta}_{T,D} - \Theta\|_{l_2}^2) = O_p[3\frac{\kappa(W)\sqrt{T}\sqrt{skn}}{T}] = O_p[3\frac{\sqrt{skn^3}\kappa^2(W)}{T}] = o_p(1),$$

where the last equality is due to Assumption 3.3(ii). 

\section{A.5 Auxiliary Lemmas}

This subsection of Appendix contains auxiliary lemmas which have been used in other subsections of Appendix.

\begin{definition}[Nets and covering numbers] Let $(T,d)$ be a metric space and fix $\varepsilon > 0$.

(i) A subset $\mathcal{N}_\varepsilon$ of $T$ is called an $\varepsilon$-net of $T$ if every point $x \in T$ satisfies $d(x,y) \leq \varepsilon$ for some $y \in \mathcal{N}_\varepsilon$.

(ii) The minimal cardinality of an $\varepsilon$-net of $T$ is denote $\mathcal{N}(\varepsilon,d)$ and is called the covering number of $T$ (at scale $\varepsilon$). Equivalently, $\mathcal{N}(\varepsilon,d)$ is the minimal number of balls of radius $\varepsilon$ and with centers in $T$ needed to cover $T$.

\begin{lemma} The unit Euclidean sphere $\{x \in \mathbb{R}^n : \|x\|_2 = 1\}$ equipped with the Euclidean metric $d$ satisfies for every $\varepsilon > 0$ that

$$\mathcal{N}(\varepsilon,d) \leq \left(1 + \frac{2}{\varepsilon}\right)^n.$$\end{lemma}

\begin{proof} See ? p8. \end{proof}

Recall that for a symmetric $n \times n$ matrix $A$, its $\ell_2$ spectral norm can be written as: $\|A\|_{l_2} = \max_{\|x\|_2 = 1} |x^TAx|$.

\begin{lemma} Let $A$ be a symmetric $n \times n$ matrix, and let $\mathcal{N}_\varepsilon$ be an $\varepsilon$-net of the unit sphere $\{x \in \mathbb{R}^n : \|x\|_2 = 1\}$ for some $\varepsilon \in [0,1)$. Then

$$\|A\|_{l_2} \leq \frac{1}{1 - 2\varepsilon} \max_{x \in \mathcal{N}_\varepsilon} |x^TAx|.$$\end{lemma}

\begin{proof} See ? p8. \end{proof}

The following theorem is a version of Bernstein’s inequality which accommodates strong mixing time series.
Theorem A.2 (Theorem 1 of ). Let \( \{X_t\}_{t \in \mathbb{Z}} \) be a sequence of centered real-valued random variables. Suppose that for every \( \epsilon \geq 0 \), there exist absolute constants \( \gamma_2 \in (0, +\infty) \) and \( b \in (0, +\infty) \) such that

\[
\sup_{t \geq 1} \mathbb{P}(|X_t| \geq \epsilon) \leq \exp[1] 1 - (\epsilon/b)^{\gamma_2}.
\]

Moreover, assume its alpha mixing coefficient \( \alpha(h) \) satisfies

\[
\alpha(h) \leq \exp(-ch^{\gamma_1}), \quad h \in \mathbb{N}
\]

for absolute constants \( c > 0 \) and \( \gamma_1 > 0 \). Define \( \gamma \) by \( 1/\gamma := 1/\gamma_1 + 1/\gamma_2 \); constants \( \gamma_1 \) and \( \gamma_2 \) need to be restricted to make sure \( \gamma < 1 \). Then, for any \( T \geq 4 \), there exist positive constants \( C_1, C_2, C_3, C_4, C_5 \) depending only on \( b, c, \gamma_1, \gamma_2 \) such that, for every \( \epsilon \geq 0 \),

\[
\mathbb{P} \left[ \max_{1 \leq i \leq n} \frac{1}{T} \sum_{t=1}^{T} X_{t,i} \geq \epsilon \right] \leq T \exp[3] \frac{(T\epsilon)^{\gamma}}{C_1} + \exp[3] \frac{(T\epsilon)^2}{C_2(1 + C_3 T)} + \exp[3] \frac{(T\epsilon)^2}{C_4 T} \exp[3] \frac{(T\epsilon)^{\gamma(1-\gamma)}}{C_5 (\log(T\epsilon))^{\gamma}}.
\]

We can use the preceding theorem to establish a rate for the maximum.

Lemma A.12. Suppose that we have for \( 1 \leq i \leq n \), for every \( \epsilon \geq 0 \),

\[
\mathbb{P} \left[ \max_{1 \leq i \leq n} \frac{1}{T} \sum_{t=1}^{T} X_{t,i} \geq \epsilon \right] \leq T \exp[3] \frac{(T\epsilon)^{\gamma}}{C_1} + \exp[3] \frac{(T\epsilon)^2}{C_2(1 + C_3 T)} + \exp[3] \frac{(T\epsilon)^2}{C_4 T} \exp[3] \frac{(T\epsilon)^{\gamma(1-\gamma)}}{C_5 (\log(T\epsilon))^{\gamma}}.
\]

Suppose \( \log n = o(T^{\frac{1}{\gamma(1-\gamma)}}) \) if \( n > T \). Then

\[
\max_{1 \leq i \leq n} \frac{1}{T} \sum_{t=1}^{T} X_{t,i} = O_p[3] \sqrt{\frac{\log n}{T}}.
\]

Proof.

\[
\mathbb{P} \left[ \max_{1 \leq i \leq n} \frac{1}{T} \sum_{t=1}^{T} X_{t,i} \geq \epsilon \right] \leq \sum_{i=1}^{n} \mathbb{P} \left[ \frac{1}{T} \sum_{t=1}^{T} X_{t,i} \geq \epsilon \right] \leq n T \exp[3] \frac{(T\epsilon)^{\gamma}}{C_1} + n \exp[3] \frac{(T\epsilon)^2}{C_2(1 + C_3 T)} + n \exp[3] \frac{(T\epsilon)^2}{C_4 T} \exp[3] \frac{(T\epsilon)^{\gamma(1-\gamma)}}{C_5 (\log(T\epsilon))^{\gamma}}.
\]

We shall choose \( \epsilon = C \sqrt{\log n/T} \) for some \( C > 0 \) and consider the three terms on the right side of inequality separately. We consider the first term for the case \( n \leq T \)

\[
n T \exp[3] \frac{(T\epsilon)^{\gamma}}{C_1} = \exp[3] \log(n T) - \frac{C_1}{C_1} (T \log n)^{\gamma/2} = \exp[3] (T \log n)^{\gamma/2} (\frac{\log(n T)}{(T \log n)^{\gamma/2}} - \frac{C_1}{C_1})
\]

\[
\leq \exp[3] (T \log n)^{\gamma/2} (\frac{2 \log T}{(T \log n)^{\gamma/2}} - \frac{C_1}{C_1}) = \exp[3] (T \log n)^{\gamma/2} o(1) - \frac{C_1}{C_1} = o(1),
\]

50
for large enough \( C \). We next consider the first term for the case \( n > T \)

\[
nT \exp[3] - \frac{(Te)^\gamma}{C_1} = \exp[3] \log(nT) - \frac{C^\gamma}{C_1} (T \log n)^{\gamma/2} = \exp[3] (T \log n)^{\gamma/2}[3] \frac{\log(nT)}{(T \log n)^{\gamma/2}} - \frac{C^\gamma}{C_1} \leq \exp[3] (T \log n)^{\gamma/2}[3] \frac{2 \log n}{(T \log n)^{\gamma/2}} - \frac{C^\gamma}{C_1} = \exp[3] (T \log n)^{\gamma/2}[3] o(1) - \frac{C^\gamma}{C_1} = o(1),
\]

for large enough \( C \) given the assumption \( \log n = o(T^{\gamma/\gamma}) \). We consider the second term.

\[
n \exp[3] - \frac{(Te)^2}{C_2(1 + C_3T)} = \exp[3] \log n - \frac{C^2 \log n}{C_2/T + C_2C_3} = \exp[3] \log n[3] 1 - \frac{C^2}{C_2/T + C_2C_3} = o(1)
\]

for large enough \( C \). We consider the third term.

\[
n \exp[3] - \frac{(Te)^2}{C_4} \frac{\exp[3] (Te)^{(1-\gamma)}}{C_5 \log(Te)} \leq n \exp[3] - \frac{(Te)^2}{C_4} \frac{\exp[3] (Te)^{(1-\gamma)}}{C_5 (Te)^\gamma} \leq n \exp[3] - \frac{(Te)^2}{C_4} \frac{\exp[3]}{C_5 (Te)^\gamma} \leq n \exp[3] - \frac{(Te)^2}{C_4 T} (1 + o(1)) = \exp[3] \log n - \frac{C^2 \log n}{C_4} (1 + o(1)) = o(1),
\]

for large enough \( C \). This yields the result. \( \square \)

**Lemma A.13.** Let \( A, B \) be \( n \times n \) positive semidefinite matrices and not both singular. Then

\[ \| A - B \|_{\ell_2} \leq \frac{\| A^2 - B^2 \|_{\ell_2}}{\text{mineval}(A) + \text{mineval}(B)}. \]

**Proof.** See ? p410. \( \square \)

**Lemma A.14.** Let \( \hat{\Omega}_n \) and \( \Omega_n \) be invertible (both possibly stochastic) square matrices whose dimensions could be growing. Let \( T \) be the sample size. For any matrix norm, suppose that \( \| \Omega_n^{-1} \| = O_p(1) \) and \( \| \hat{\Omega}_n - \Omega_n \| = O_p(a_n,T) \) for some sequence \( a_n,T \) with \( a_n,T \to 0 \) as \( n \to \infty \), \( T \to \infty \) simultaneously (joint asymptotics). Then \( \| \hat{\Omega}_n^{-1} - \Omega_n^{-1} \| = O_p(a_n,T) \).

**Proof.** The original proof could be found in ? Lemma A.2.

\[ \| \hat{\Omega}_n^{-1} - \Omega_n^{-1} \| \leq \| \hat{\Omega}_n^{-1} \| \| \Omega_n - \hat{\Omega}_n \| \| \Omega_n^{-1} \| \leq [1] \| \Omega_n^{-1} \| + \| \hat{\Omega}_n^{-1} - \hat{\Omega}_n \| \| \hat{\Omega}_n^{-1} \| \| \Omega_n^{-1} \|. \]

Let \( v_n,T \), \( z_n,T \) and \( x_n,T \) denote \( \| \Omega_n^{-1} \| \), \( \| \hat{\Omega}_n^{-1} - \Omega_n^{-1} \| \) and \( \| \Omega_n - \hat{\Omega}_n \| \), respectively. From the preceding equation, we have

\[ w_n,T := \frac{z_n,T}{(v_n,T + z_n,T)v_n,T} \leq x_n,T = O_p(a_n,T) = o_p(1). \]

We now solve for \( z_n,T \):

\[ z_n,T = \frac{v_n,T^2 w_n,T}{1 - v_n,T w_n,T} = O_p(a_n,T). \]

\( \square \)
Theorem A.3 (p269; ?). For $A \in \mathbb{C}^{n \times n}$ with no eigenvalues lying on the closed negative real axis $(-\infty, 0)$,

$$\log A = \int_0^1 (A - I)[t(A - I) + I]^{-1} dt.$$ 

**Lemma A.15.** Let $A, B$ be $n \times n$ real matrices. Suppose that $A$ is symmetric, positive definite for all $n$ and its minimum eigenvalue is bounded away from zero by an absolute constant. Assume $\|A^{-1}B\|_{\ell_2} \leq C < 1$ for some constant $C$. Then $A + B$ is invertible for every $n$ and

$$(A + B)^{-1} = A^{-1} - A^{-1}BA^{-1} + O(\|B\|_{\ell_2}^2).$$

**Proof.** We write $A + B = A[I - (-A^{-1}B)]$. Since $\| - A^{-1}B\|_{\ell_2} \leq C < 1$, $I - (-A^{-1}B)$ and hence $A + B$ are invertible (p301). We then can expand

$$(A + B)^{-1} = \sum_{k=0}^{\infty} (-A^{-1}B)^k A^{-1} = A^{-1} - A^{-1}BA^{-1} + \sum_{k=2}^{\infty} (-A^{-1}B)^k A^{-1}.$$ 

Then

$$\sum_{k=2}^{\infty} (-A^{-1}B)^k A^{-1} \leq \sum_{k=2}^{\infty} (-A^{-1}B)^k \|A^{-1}\|_{\ell_2} \|A^{-1}\|_{\ell_2} \leq \sum_{k=2}^{\infty} \|(-A^{-1}B)^k\|_{\ell_2} \|A^{-1}\|_{\ell_2} \leq \frac{\|A^{-1}\|_{\ell_2}^2 \|B\|_{\ell_2}^2}{1 - \|A^{-1}\|_{\ell_2}^2}.$$ 

where the first and third inequalities are due to the submultiplicative property of a matrix norm, the second inequality is due to the triangular inequality. Since $A$ is real, symmetric, and positive definite with the minimum eigenvalue bounded away from zero by an absolute constant, $\|A^{-1}\|_{\ell_2} = \text{maxeval}(A^{-1}) = 1/\text{mineval}(A) < D < \infty$ for some absolute constant $D$. Hence the result follows. \hfill \square

**Lemma A.16.** Consider real matrices $A$ ($m \times n$) and $B$ ($p \times q$). Then

$$\|A \otimes B\|_{\ell_2} = \|A\|_{\ell_2} \|B\|_{\ell_2}.$$ 

**Proof.**

$$\|A \otimes B\|_{\ell_2} = \sqrt{\maxeval[(A \otimes B)^T(A \otimes B)]} = \sqrt{\maxeval[(A^T \otimes B^T)(A \otimes B)]}$$

$$= \sqrt{\maxeval[A^T A \otimes B^T B]} = \sqrt{\maxeval[A^T A] \maxeval[B^T B]} = \|A\|_{\ell_2} \|B\|_{\ell_2},$$

where the fourth equality is due to the fact that both $A^T A$ and $B^T B$ are symmetric, positive semidefinite. \hfill \square

**Lemma A.17.** Let $A$ be a $p \times p$ symmetric matrix and $\hat{v}, v \in \mathbb{R}^p$. Then

$$|\hat{v}^T A \hat{v} - v^T A v| \leq |\maxeval(A)|^2 \|\hat{v} - v\|_2^2 + 2(\|Av\|_2 \|\hat{v} - v\|_2).$$
Proof. See Lemma 3.1 in the supplementary material of ?.

Theorem A.4. For arbitrary $n \times n$ complex matrices $A$ and $E$, and for any matrix norm $\| \cdot \|$, 

$$ \| e^{A+E} - e^A \| \leq \| E \| \exp(\| E \|) \exp(\| A \|). $$

Proof. See ? p430.

Lemma A.18 (? p27).

$$ \frac{X_k^2 - k}{\sqrt{2k}} \xrightarrow{d} N(0, 1), $$

as $k \to \infty$.

Lemma A.19 (? p41). For $T, n \in \mathbb{N}$ let $X_{T,n}$ be random vectors such that $X_{T,n} \sim X_n$ as $T \to \infty$ for every fixed $n$ such that $X_n \sim X$ as $n \to \infty$. Then there exists a sequence $n_T \to \infty$ such that $X_{T,n_T} \sim X$ as $T \to \infty$.

Theorem A.5 (?). Let $\{X_{n,i}, i = 1, \ldots, k_n\}$ be a martingale difference array with respect to the triangular array of $\sigma$-algebras $\{\mathcal{F}_{n,i}, i = 0, \ldots, k_n\}$ (i.e., $X_{n,i}$ is $\mathcal{F}_{n,i}$-measurable and $\mathbb{E}[X_{n,i} | \mathcal{F}_{n,i-1}] = 0$ almost surely for all $n$ and $i$) satisfying $\mathcal{F}_{n,i-1} \subset \mathcal{F}_{n,i}$ for all $n \geq 1$. Assume,

(i) $\max_{i \leq k_n} |X_{n,i}|$ is uniformly (in $n$) bounded in $L_2$ norm,

(ii) $\max_{i \leq k_n} |X_{n,i}| \xrightarrow{p} 0$, and

(iii) $\sum_{i=1}^{k_n} X_{n,i}^2 \xrightarrow{p} 1$.

Then, $S_n = \sum_{i=1}^{k_n} X_{n,i} \xrightarrow{d} N(0, 1)$ as $n \to \infty$.


53
8 Supplementary Material

This section is the supplementary materials to the main article. SM 8.1 contains additional materials related to the Kronecker product (models). SM 8.2 gives a lemma characterising a rate for $\|\hat{V}_T - V\|_\infty$, which is used in the proofs of the limiting distributions of our estimators. SM 8.3, SM 8.4, and SM 8.5 provide proofs of Theorem 3.3, Theorem 4.1, and Theorem 4.2, respectively. SM 8.6 gives proofs of Theorem 3.4 and Corollary 3.3. SM 8.7 contains some miscellaneous results.

8.1 Additional Materials Related to the Kronecker Product

The following lemma proves a property of Kronecker products.

**Lemma 8.1.** Suppose $v = 2, 3, \ldots$ and that $A_1, A_2, \ldots, A_v$ are real symmetric and positive definite matrices of sizes $a_1 \times a_1, \ldots, a_v \times a_v$, respectively. Then

$$
\log(A_1 \otimes A_2 \otimes \cdots \otimes A_v) = \log(A_1 \otimes I_{a_2} \otimes \cdots \otimes I_{a_v}) + I_{a_1} \otimes \log(A_2 \otimes I_{a_3} \otimes \cdots \otimes I_{a_v}) + \cdots + I_{a_1} \otimes I_{a_2} \otimes \cdots \otimes \log(A_v).
$$

**Proof.** We prove by mathematical induction. We first give a proof for $v = 2$; that is,

$$
\log(A_1 \otimes A_2) = \log(A_1 \otimes I_{a_2}) + I_{a_1} \otimes \log(A_2).
$$

Since $A_1, A_2$ are real symmetric, they can be orthogonally diagonalized: for $i = 1, 2$, $A_i = U_i^T A_i U_i$,

where $U_i$ is orthogonal, and $\Lambda_i = \text{diag}(\lambda_{i,1}, \ldots, \lambda_{i, a_i})$ is a diagonal matrix containing the $a_i$ eigenvalues of $A_i$. Positive definiteness of $A_1, A_2$ ensures that their Kronecker product is positive definite. Then the logarithm of $A_1 \otimes A_2$ is:

$$
\log(A_1 \otimes A_2) = \log((U_1 \otimes U_2)^T (\Lambda_1 \otimes \Lambda_2)(U_1 \otimes U_2)) = (U_1 \otimes U_2)^T \log(\Lambda_1 \otimes \Lambda_2)(U_1 \otimes U_2),
$$

(8.1)
where the first equality is due to the mixed product property of the Kronecker product, and the second equality is due to a property of matrix functions. Next,

\[
\log(A_1 \otimes A_2) = \text{diag}(\log(\lambda_{1,1}A_2), \ldots, \log(\lambda_{1,a_1}A_2)) = \text{diag}(\log(\lambda_{1,1}I_{a_2}A_2), \ldots, \log(\lambda_{1,a_1}I_{a_2}A_2)) \\
= \text{diag}(\log(\lambda_{1,1}I_{a_2}) + \log(A_2), \ldots, \log(\lambda_{1,a_1}I_{a_2}) + \log(A_2)) \\
= \text{diag}(\log(\lambda_{1,1}I_{a_2}), \ldots, \log(\lambda_{1,a_1}I_{a_2})) + \text{diag}(\log(A_2), \ldots, \log(A_2)) \\
= \log(\Lambda_1) \otimes I_{a_2} + I_{a_1} \otimes \log(A_2),
\]

(8.2)

where the third equality holds only because \(\lambda_{1,1}I_{a_2}\) and \(A_2\) have real positive eigenvalues only and commute for all \(j = 1, \ldots, a_1\) [Higham (2008) p270 Theorem 11.3]. Substitute (8.2) into (8.1):

\[
\log(A_1 \otimes A_2) = (U_1 \otimes U_2)^T \log(\Lambda_1) \otimes \log(A_2)(U_1 \otimes U_2) = (U_1 \otimes U_2)^T(\log(A_1 \otimes I_{a_2} + I_{a_1} \otimes \log(A_2))(U_1 \otimes U_2) \\
= (U_1 \otimes U_2)^T(\log(A_1 \otimes I_{a_2})(U_1 \otimes U_2) + (U_1 \otimes U_2)^T(I_{a_1} \otimes \log(A_2))(U_1 \otimes U_2) \\
= \log A_1 \otimes I_{a_2} + I_{a_1} \otimes \log A_2.
\]

We now assume that this lemma is true for \(v = k\). That is,

\[
\log(A_1 \otimes A_2 \otimes \cdots \otimes A_k) = \log A_1 \otimes I_{a_2} \otimes \cdots \otimes I_{a_k} + I_{a_1} \otimes \log A_2 \otimes \cdots \otimes \log A_k + \cdots + I_{a_1} \otimes I_{a_2} \otimes \cdots \otimes \log A_k.
\]

(8.3)

We now prove that the lemma holds for \(v = k + 1\). Let \(A_{1-k} := A_1 \otimes \cdots \otimes A_k\) and \(I_{a_1-a_k} := I_{a_1} \otimes \cdots \otimes I_{a_k}\).

\[
\log(A_1 \otimes A_2 \otimes \cdots \otimes A_k \otimes A_{k+1}) = \log(A_{1-k} \otimes A_{k+1}) = \log(A_{1-k} \otimes I_{a_k-a_{k+1}} + I_{a_1-a_k} \otimes \log A_{k+1}) \\
= \log A_1 \otimes I_{a_2} \otimes \cdots \otimes I_{a_k} \otimes I_{a_{k+1}} + I_{a_1} \otimes \log A_2 \otimes \cdots \otimes I_{a_k} \otimes I_{a_{k+1}} + \cdots + \\
I_{a_1} \otimes I_{a_2} \otimes \cdots \otimes \log A_k \otimes I_{a_{k+1}} + I_{a_1} \otimes \cdots \otimes I_{a_k} \otimes \log A_{k+1}.
\]

Thus the lemma holds for \(v = k + 1\). By induction, the lemma is true for \(v = 2, 3, \ldots\) \(\square\)

Next we provide two examples to illustrate the necessity of an identification restriction in order to separately identify log parameters.

**Example 8.1.** Suppose that \(n_1, n_2 = 2\). We have

\[
\log \Theta_1 = \begin{pmatrix} a_{11} & a_{12} \\ a_{12} & a_{22} \end{pmatrix}, \quad \log \Theta_2^* = \begin{pmatrix} b_{11} & b_{12} \\ b_{12} & b_{22} \end{pmatrix}
\]

Then we can calculate

\[
\log \Theta^* = \log \Theta_1^* \otimes I_2 + I_2 \otimes \log \Theta_2^* = \begin{pmatrix} a_{11} + b_{11} & b_{12} & a_{12} & 0 \\ b_{12} & a_{11} + b_{22} & 0 & a_{12} \\ a_{12} & 0 & a_{22} + b_{11} & b_{12} \\ 0 & a_{12} & b_{12} & a_{22} + b_{22} \end{pmatrix}
\]

Log parameters \(a_{12}, b_{12}\) can be separately identified from the off-diagonal entries of \(\log \Theta^*\) because they appear separately. We now examine whether log parameters \(a_{11}, b_{11}, a_{22}, b_{22}\) can be separately identified from diagonal entries of \(\log \Theta^*\). The answer is no. We have the following linear system

\[
Ax := \begin{pmatrix} 1 & 0 & 1 & 0 \\ 1 & 0 & 0 & 1 \\ 0 & 1 & 1 & 0 \\ 0 & 1 & 0 & 1 \end{pmatrix} \begin{pmatrix} a_{11} \\ a_{22} \\ b_{11} \\ b_{22} \end{pmatrix} = \begin{pmatrix} \log \Theta_1^* \\ \log \Theta_2^* \end{pmatrix} =: d.
\]
Note that the rank of $A$ is 3. There are three effective equations and four unknowns; the linear system has infinitely many solutions for $x$. Hence one identification restriction is needed to separately identify log parameters $a_{11}, b_{11}, a_{22}, b_{22}$. We choose to set $a_{11} = 0$.

**Example 8.2.** Suppose that $n_1, n_2, n_3 = 2$. We have

$$\log \Theta^*_1 = \begin{pmatrix} a_{11} \\ a_{12} \\ a_{22} \end{pmatrix} \quad \log \Theta^*_2 = \begin{pmatrix} b_{11} \\ b_{12} \\ b_{22} \end{pmatrix} \quad \log \Theta^*_3 = \begin{pmatrix} c_{11} \\ c_{12} \\ c_{22} \end{pmatrix}$$

Then we can calculate

$$\log \Theta^* = \log \Theta^*_1 \otimes I_2 \otimes I_2 + I_2 \otimes \log \Theta^*_2 \otimes I_2 + I_2 \otimes I_2 \otimes \log \Theta^*_3 =$$

$$\begin{pmatrix}
  a_{11} + b_{11} + c_{11} & a_{12} & \ldots & a_{12} \\
  a_{12} & b_{12} & \ldots & b_{12} \\
  \vdots & \vdots & \ddots & \vdots \\
  a_{12} & a_{12} & \ldots & c_{22} \\
  a_{22} & a_{22} & \ldots & a_{22} \\
  \vdots & \vdots & \ddots & \vdots \\
  a_{22} & a_{22} & \ldots & c_{22} \\
  \end{pmatrix}.$$

Log parameters $a_{12}, b_{12}, c_{12}$ can be separately identified from off-diagonal entries of $\log \Theta^*$ because they appear separately. We now examine whether log parameters $a_{11}, b_{11}, c_{11}, a_{22}, b_{22}, c_{22}$ can be separately identified from diagonal entries of $\log \Theta^*$. The answer is no. We have the following linear system

$$Ax := \begin{pmatrix}
  1 & 0 & 1 & 0 & 1 & 0 \\
  1 & 0 & 1 & 0 & 0 & 1 \\
  1 & 0 & 0 & 1 & 1 & 0 \\
  1 & 0 & 0 & 1 & 0 & 1 \\
  0 & 1 & 1 & 0 & 1 & 0 \\
  0 & 1 & 1 & 0 & 0 & 1 \\
  0 & 1 & 0 & 1 & 1 & 0 \\
  0 & 1 & 0 & 1 & 0 & 1 \\
  \end{pmatrix} \begin{pmatrix}
  a_{11} \\
  a_{22} \\
  b_{11} \\
  b_{22} \\
  c_{11} \\
  c_{22} \\
  \end{pmatrix} = \begin{pmatrix}
  \log \Theta^*_1 \\
  \log \Theta^*_2 \\
  \log \Theta^*_3 \\
  \log \Theta^*_4 \\
  \log \Theta^*_5 \\
  \log \Theta^*_6 \\
  \end{pmatrix} =: \mathbf{d}.$$

Note that the rank of $A$ is 4. There are four effective equations and six unknowns; the linear system has infinitely many solutions for $x$. Hence two identification restrictions are needed to separately identify log parameters $a_{11}, b_{11}, c_{11}, a_{22}, b_{22}, c_{22}$. We choose to set $a_{11} = b_{11} = 0$.

### 8.2 Lemma [8.2]

The following lemma characterises a rate for $\|\hat{V}_\gamma - V\|_\infty$, which is used in the proofs of the limiting distributions of our estimators.

**Lemma 8.2.** Let Assumptions 3.1(i) and 3.2 be satisfied with $1/\gamma := 1/r_1 + 1/r_2 > 1$. Suppose $\log n = o(T^{2/\gamma})$ if $n > T$. Then

$$\|\hat{V}_T - V\|_\infty = O_p\left(\sqrt{\frac{\log n}{T}}\right).$$

**Proof.** Let $\hat{y}_{t,i}$ denote $y_{t,i} - \bar{y}_i$, similarly for $\hat{y}_{t,j}, \hat{y}_{t,k}, \hat{y}_{t,t}$. Let $\bar{y}_{t,i}$ denote $y_{t,i} - \mu_i$, similarly for
\[ \dot{y}_{t,j}, \dot{y}_{t,k}, \dot{y}_{t,\ell}. \]

\[ \|\dot{V}_T - V\|_\infty := \max_{1 \leq x,y \leq n^2} |\dot{V}_{T,x,y} - V_{x,y}| = \max_{1 \leq i,j,k,\ell \leq n} |\dot{V}_{T,i,j,k,\ell} - V_{i,j,k,\ell}| \]

\[ \leq \max_{1 \leq i,j,k,\ell \leq n} \left| \frac{1}{T} \sum_{t=1}^{T} \dot{y}_{t,i} \dot{y}_{t,j} \dot{y}_{t,k} \dot{y}_{t,\ell} - \frac{1}{T} \sum_{t=1}^{T} \dot{y}_{t,i} \dot{y}_{t,j} \dot{y}_{t,k} \dot{y}_{t,\ell} \right| \]

\[ + \max_{1 \leq i,j,k,\ell \leq n} \left| \frac{1}{T} \sum_{t=1}^{T} \dot{y}_{t,i} \dot{y}_{t,j} \dot{y}_{t,k} \dot{y}_{t,\ell} - E[\dot{y}_{t,i} \dot{y}_{t,j} \dot{y}_{t,k} \dot{y}_{t,\ell}] \right| \]

\[ + \max_{1 \leq i,j,k,\ell \leq n} \left| \frac{1}{T} \sum_{t=1}^{T} \dot{y}_{t,j} \dot{y}_{t,k} \dot{y}_{t,\ell} - E[\dot{y}_{t,j} \dot{y}_{t,k} \dot{y}_{t,\ell}] \right| \]

(8.4)

Assumption 3.1(i) says that for all \( t \), there exist absolute constants \( K_1 > 1, K_2 > 0, r_1 > 0 \) such that

\[ E \left[ \exp \left( K_2 |y_{t,i}|^{r_1} \right) \right] \leq K_1 \quad \text{for all } i = 1, \ldots, n. \]

By repeated using Lemma A.2 in Appendix A.3 we have for all \( i, j, k, \ell = 1, 2, \ldots, n \), every \( \epsilon \geq 0 \), absolute constants \( b_1, c_1, b_2, c_2, b_3, c_3 > 0 \) such that

\[ \mathbb{P}(|y_{t,i}| \geq \epsilon) \leq \exp \left[ 1 - (\epsilon/b_1)^{r_1} \right] \]

\[ \mathbb{P}(|\dot{y}_{t,i}| \geq \epsilon) \leq \exp \left[ 1 - (\epsilon/c_1)^{r_1} \right] \]

\[ \mathbb{P}(|\dot{y}_{t,i} y_{t,j} - E[\dot{y}_{t,i} y_{t,j}]| \geq \epsilon) \leq \exp \left[ 1 - (\epsilon/b_2)^{r_2} \right] \]

\[ \mathbb{P}(|\dot{y}_{t,i} y_{t,j} - E[\dot{y}_{t,i} y_{t,j}] y_{t,k} y_{t,\ell}| \geq \epsilon) \leq \exp \left[ 1 - (\epsilon/b_3)^{r_2} \right] \]

where \( r_2 \in (0, r_1/2) \) and \( r_3 \in (0, r_1/4) \). Use the assumption \( 1/r_1 + 1/r_2 > 1 \) to invoke Lemma A.12 in Appendix A.5 to get

\[ \max_{1 \leq i,j,k,\ell \leq n} \left| \frac{1}{T} \sum_{t=1}^{T} \dot{y}_{t,i} \dot{y}_{t,j} \dot{y}_{t,k} \dot{y}_{t,\ell} - E[\dot{y}_{t,i} \dot{y}_{t,j} \dot{y}_{t,k} \dot{y}_{t,\ell}] \right| = O_p \left( \sqrt{\frac{\log n}{T}} \right). \]

(8.8)

(8.5)

We now consider (8.7).

\[ \max_{1 \leq i,j,k,\ell \leq n} \left| \frac{1}{T} \sum_{t=1}^{T} \dot{y}_{t,i} \dot{y}_{t,j} \right| \left( \frac{1}{T} \sum_{t=1}^{T} \dot{y}_{t,k} \dot{y}_{t,\ell} \right) - E[\dot{y}_{t,i} \dot{y}_{t,j}] E[\dot{y}_{t,k} \dot{y}_{t,\ell}] \]

\[ \leq \max_{1 \leq i,j,k,\ell \leq n} \left| \frac{1}{T} \sum_{t=1}^{T} \dot{y}_{t,i} \dot{y}_{t,j} \right| \left( \frac{1}{T} \sum_{t=1}^{T} \dot{y}_{t,k} \dot{y}_{t,\ell} - E[\dot{y}_{t,k} \dot{y}_{t,\ell}] \right) \]

\[ + \max_{1 \leq i,j,k,\ell \leq n} \left| E[\dot{y}_{t,k} \dot{y}_{t,\ell}] \right| \left( \frac{1}{T} \sum_{t=1}^{T} \dot{y}_{t,i} \dot{y}_{t,j} - E[\dot{y}_{t,i} \dot{y}_{t,j}] \right). \]

(8.9)
Consider (8.9).

\[
\max_{1 \leq i,j,k,\ell \leq n} \left( \frac{1}{T} \sum_{t=1}^{T} \tilde{y}_{t,i} \tilde{y}_{t,j} \right) \left( \frac{1}{T} \sum_{t=1}^{T} \tilde{y}_{t,k} \tilde{y}_{t,\ell} - E\tilde{y}_{t,k} \tilde{y}_{t,\ell} \right)
\]

\[
\leq \max_{1 \leq i,j \leq n} \left( \frac{1}{T} \sum_{t=1}^{T} \tilde{y}_{t,i} \tilde{y}_{t,j} - E\tilde{y}_{t,i} \tilde{y}_{t,j} \right) + \left( \frac{1}{T} \sum_{t=1}^{T} \tilde{y}_{t,k} \tilde{y}_{t,\ell} - E\tilde{y}_{t,k} \tilde{y}_{t,\ell} \right)
\]

\[
= \left( O_p \left( \sqrt{\frac{\log n}{T}} \right) + O(1) \right) O_p \left( \sqrt{\frac{\log n}{T}} \right) = O_p \left( \sqrt{\frac{\log n}{T}} \right)
\]

where the first equality is due to Lemma A.2 (ii) in Appendix A.3 and Lemma A.12 in Appendix A.5. Now consider (8.10).

\[
\max_{1 \leq i,j,k,\ell \leq n} \left| \mathbb{E}[\tilde{y}_{t,k} \tilde{y}_{t,\ell}] \left( \frac{1}{T} \sum_{t=1}^{T} \tilde{y}_{t,i} \tilde{y}_{t,j} - E[\tilde{y}_{t,i} \tilde{y}_{t,j}] \right) \right|
\]

\[
\leq \max_{1 \leq k,\ell \leq n} \left| \mathbb{E}[\tilde{y}_{t,k} \tilde{y}_{t,\ell}] \right| \max_{1 \leq i,j \leq n} \left| \frac{1}{T} \sum_{t=1}^{T} \tilde{y}_{t,i} \tilde{y}_{t,j} - E\tilde{y}_{t,i} \tilde{y}_{t,j} \right| = O_p \left( \sqrt{\frac{\log n}{T}} \right)
\]

where the equality is due to Lemma A.12 in Appendix A.5. Thus

\[
\max_{1 \leq i,j,k,\ell \leq n} \left| \left( \frac{1}{T} \sum_{t=1}^{T} \tilde{y}_{t,i} \tilde{y}_{t,j} \right) \left( \frac{1}{T} \sum_{t=1}^{T} \tilde{y}_{t,k} \tilde{y}_{t,\ell} \right) - \mathbb{E}[\tilde{y}_{t,i} \tilde{y}_{t,j}]\mathbb{E}[\tilde{y}_{t,k} \tilde{y}_{t,\ell}] \right| = O_p \left( \sqrt{\frac{\log n}{T}} \right). \tag{8.11}
\]

We first give a rate for \(\max_{1 \leq i \leq n} |\tilde{y}_i - \mu_i|\). The index \(i\) is arbitrary and could be replaced with \(j, k, \ell\). Invoking Lemma A.12 in Appendix A.5, we have

\[
\max_{1 \leq i \leq n} |\tilde{y}_i - \mu_i| = \max_{1 \leq i \leq n} \left| \frac{1}{T} \sum_{t=1}^{T} (y_{t,i} - E\tilde{y}_{t,i}) \right| = O_p \left( \sqrt{\frac{\log n}{T}} \right). \tag{8.12}
\]

Then we also have

\[
\max_{1 \leq i \leq n} |\tilde{y}_i| = \max_{1 \leq i \leq n} |\tilde{y}_i - \mu_i + \mu_i| \leq \max_{1 \leq i \leq n} |\tilde{y}_i - \mu_i| + \max_{1 \leq i \leq n} |\mu_i| = O_p \left( \sqrt{\frac{\log n}{T}} \right) + O(1) = O_p(1). \tag{8.13}
\]

We now consider (8.4):

\[
\max_{1 \leq i,j,k,\ell \leq n} \left| \frac{1}{T} \sum_{t=1}^{T} \tilde{y}_{t,i} \tilde{y}_{t,j} \tilde{y}_{t,k} \tilde{y}_{t,\ell} - \frac{1}{T} \sum_{t=1}^{T} \tilde{y}_{t,i} \tilde{y}_{t,j} \tilde{y}_{t,k} \tilde{y}_{t,\ell} \right|.
\]

With tedious expansion, simplification and recognition the indices \(i, j, k, \ell\) are completely symmetric, we can bound (8.4) by

\[
\max_{1 \leq i,j,k,\ell \leq n} \left| \tilde{y}_{i} \tilde{y}_{j} \tilde{y}_{k} \tilde{y}_{\ell} - \mu_i \mu_j \mu_k \mu_\ell \right| \tag{8.14}
\]

\[
+ 4 \max_{1 \leq i,j,k,\ell \leq n} \left| \tilde{y}_{i} \left( \tilde{y}_{j} \tilde{y}_{k} \tilde{y}_{\ell} - \mu_j \mu_k \mu_\ell \right) \right| \tag{8.15}
\]

\[
+ 6 \max_{1 \leq i,j,k,\ell \leq n} \left| \left( \frac{1}{T} \sum_{t=1}^{T} y_{t,i} y_{t,j} \right) \left( \tilde{y}_{k} \tilde{y}_{\ell} - \mu_k \mu_\ell \right) \right| \tag{8.16}
\]

\[
+ 4 \max_{1 \leq i,j,k,\ell \leq n} \left| \left( \frac{1}{T} \sum_{t=1}^{T} y_{t,i} y_{t,j} y_{t,k} \right) \left( \tilde{y}_{\ell} - \mu_\ell \right) \right|. \tag{8.17}
\]
We consider (8.14) first. (8.14) can be bounded by repeatedly invoking triangular inequalities (e.g., inserting terms like \( \mu i,j \tilde{y} i,j \)) using Lemma A.2(ii) in Appendix A.3 (8.13) and (8.12). (8.14) is of order \( O_p(\sqrt{\log n/T}) \). (8.15) is of order \( O_p(\sqrt{\log n/T}) \) by a similar argument. (8.16) and (8.17) are of the same order \( O_p(\sqrt{\log n/T}) \) using a similar argument provided that both \( \max_{1 \leq i,j \leq n} [\sum_{t=1}^{T} y_t i,j y_t j]/T \) and \( \max_{1 \leq i,j,k \leq n} [\sum_{t=1}^{T} y_t i,j y_t j y_t k]/T \) are \( O_p(1) \); these follow from Lemma A.2(ii) in Appendix A.3 and Lemma A.12 in Appendix A.5. Thus

\[
\max_{1 \leq i,j,k,\ell \leq n} \left| \frac{1}{T} \sum_{t=1}^{T} y_t i,j \tilde{y} t,j y_t k,\ell - \frac{1}{T} \sum_{t=1}^{T} y_t i,j \tilde{y} t,j y_t k,\ell \right| = O_p(\sqrt{\log n/T}).
\]  

(8.18)

We now consider (8.6).

\[
\max_{1 \leq i,j,k,\ell \leq n} \left| \left( \frac{1}{T} \sum_{t=1}^{T} y_t i,j \tilde{y} t,j \right) \left( \frac{1}{T} \sum_{t=1}^{T} \tilde{y} t,k y_t \ell \right) - \left( \frac{1}{T} \sum_{t=1}^{T} \tilde{y} t,i \tilde{y} t,j \right) \left( \frac{1}{T} \sum_{t=1}^{T} \tilde{y} t,k y_t \ell \right) \right|
\leq \max_{1 \leq i,j,k,\ell \leq n} \left| \left( \frac{1}{T} \sum_{t=1}^{T} \tilde{y} t,i \tilde{y} t,j \right) \left( \frac{1}{T} \sum_{t=1}^{T} \tilde{y} t,k y_t \ell - \tilde{y} t,k \tilde{y} t \ell \right) \right|
+ \max_{1 \leq i,j,k,\ell \leq n} \left| \left( \frac{1}{T} \sum_{t=1}^{T} y_t k y_t \ell \right) \left( \frac{1}{T} \sum_{t=1}^{T} \tilde{y} t,i \tilde{y} t,j \right) \right|\]  

(8.19)

It suffices to give a bound for (8.19) as the bound for (8.20) is of the same order and follows through similarly. First, it is easy to show that \( \max_{1 \leq i,j,k \leq n} [\sum_{t=1}^{T} \tilde{y} t,i \tilde{y} t,j] = \max_{1 \leq i,j \leq n} [\frac{1}{T} \sum_{t=1}^{T} y_t i,y_t j - \tilde{y} t i \tilde{y} t j] = O_p(1) \) (using Lemma A.2(ii) in Appendix A.3 and Lemma A.12 in Appendix A.5). Next

\[
\max_{1 \leq i,j,k,\ell \leq n} \left| \frac{1}{T} \sum_{t=1}^{T} \left( \tilde{y} t,k y_t \ell - \tilde{y} t,k \tilde{y} t \ell \right) \right| = \max_{1 \leq k,\ell \leq n} \left| -(\tilde{y} k - \mu k)(\tilde{y} \ell - \mu \ell) \right| = O_p\left( \frac{\log n}{T} \right).
\]  

(8.21)

The lemma follows after summing up the rates for (8.8), (8.11), (8.18) and (8.21). \( \square \)

8.3 Proof of Theorem 3.3

In this subsection, we give a proof for Theorem 3.3. We will first give a preliminary lemma leading to the proof of this theorem.

Lemma 8.3. Let Assumptions 3.1(i), 3.2, 3.3(i) and 3.4(i) hold with \( 1/r_1 + 1/r_2 > 1 \). Then we have

\[
\| P \|_{\ell_2} = O(1), \quad \| \hat{P} \|_{\ell_2} = O_p(1), \quad \| \hat{P} - P \|_{\ell_2} = O_p\left( \sqrt{\frac{n}{T}} \right).
\]  

(8.22)

Proof. The proofs for \( \| P \|_{\ell_2} = O(1) \) and \( \| \hat{P} \|_{\ell_2} = O_p(1) \) are exactly the same, so we only give the proof for the latter.

\[
\| \hat{P} \|_{\ell_2} = \| I_n - D_n D_n^+ (I_n \otimes \hat{\Theta}) M_d \|_{\ell_2} \leq 1 + \| D_n D_n^+ (I_n \otimes \hat{\Theta}) M_d \|_{\ell_2}
\leq 1 + \| D_n \|_{\ell_2} \| D_n^+ \|_{\ell_2} \| I_n \otimes \hat{\Theta} \|_{\ell_2} \| M_d \|_{\ell_2} = 1 + 2 \| I_n \|_{\ell_2} \| \hat{\Theta} \|_{\ell_2} = O_p(1)
\]  

where the second equality is due to (A.8) and Lemma A.16 in Appendix A.5, and last equality is due to Lemma A.7(ii). Now,

\[
\| \hat{P} - P \|_{\ell_2} = \| I_n - D_n D_n^+ (I_n \otimes \hat{\Theta}) M_d - (I_n - D_n D_n^+ (I_n \otimes \Theta) M_d) \|_{\ell_2}
= \| D_n D_n^+ (I_n \otimes \hat{\Theta}) M_d - D_n D_n^+ (I_n \otimes \Theta) M_d \|_{\ell_2} = \| D_n D_n^+ (I_n \otimes (\hat{\Theta} - \Theta)) M_d \|_{\ell_2}
= O_p(\sqrt{n/T}),
\]  

where the last equality is due to Theorem 3.1(i). \( \square \)
We are now ready to give a proof for Theorem 3.3.

**Proof of Theorem 3.3.** We write

\[
\frac{\sqrt{T}c^t(\hat{\theta}_T - \theta^0)}{\sqrt{c^t Jc}} = \frac{\sqrt{T}c^t(E^tWE)^{-1}E^tWD_n^+Hvec(\hat{\Theta}_T - \Theta)}{\sqrt{c^t Jc}} + \frac{\sqrt{T}c^t(E^tWE)^{-1}E^tWD_n^+vec O_p(\|\hat{\Theta}_T - \Theta\|_{l_2})}{\sqrt{c^t Jc}}
\]

\[
= \frac{\sqrt{T}c^t(E^tWE)^{-1}E^tWD_n^+Hvec \left( \frac{\partial vec \Theta}{\partial vec \Sigma} \right) \bigg|_{\Sigma = \hat{\Sigma}_T(j)} (\hat{\Sigma}_T - \Sigma)}{\sqrt{c^t Jc}}
\]

\[
+ \frac{\sqrt{T}c^t(E^tWE)^{-1}E^tWD_n^+vec O_p(\|\hat{\Theta}_T - \Theta\|_{l_2})}{\sqrt{c^t Jc}}
\]

\[= \hat{t}_1 + \hat{t}_2,
\]

where \( \frac{\partial vec \Theta}{\partial vec \Sigma} \bigg|_{\Sigma = \hat{\Sigma}_T(j)} \) denotes the matrix whose jth row is the jth row of the Jacobian matrix \( \frac{\partial vec \Theta}{\partial vec \Sigma} \) evaluated at \( \hat{\Sigma}_T(j) \), which is a point between \( vec \Sigma \) and \( vec \hat{\Sigma}_T \), for \( j = 1, \ldots, n^2 \).

Define

\[ t_1 := \frac{\sqrt{T}c^t(E^tWE)^{-1}E^tWD_n^+HP(D^{-1/2} \otimes D^{-1/2})vec(\hat{\Theta}_T - \Sigma)}{\sqrt{c^t Jc}}. \]

To prove Theorem 3.3 it suffices to show \( t_1 \overset{d}{\rightarrow} N(0,1) \), \( t_1 - \hat{t}_1 = o_p(1) \), and \( \hat{t}_2 = o_p(1) \). The proof is similar to that of Theorem 3.2 so we will be concise for the parts which are almost identical to that of Theorem 3.2.

**8.3.1** \( t_1 \overset{d}{\rightarrow} N(0,1) \)

We now prove that \( t_1 \) is asymptotically distributed as a standard normal.

\[ t_1 = \frac{\sqrt{T}c^t(E^tWE)^{-1}E^tWD_n^+HP(D^{-1/2} \otimes D^{-1/2})vec \left( \frac{1}{T} \sum_{t=1}^{T} [(y_t - \mu)(y_t - \mu)^T - \mathbb{E}(y_t - \mu)(y_t - \mu)^T] \right)}{\sqrt{c^t Jc}} \]

\[= \sum_{t=1}^{T} \frac{T^{-1/2}c^t(E^tWE)^{-1}E^tWD_n^+HP(D^{-1/2} \otimes D^{-1/2})vec \left( [(y_t - \mu)(y_t - \mu)^T - \mathbb{E}(y_t - \mu)(y_t - \mu)^T] \right)}{\sqrt{c^t Jc}} \]

\[= \sum_{t=1}^{T} U_{T,n,t}. \]

Again it is straightforward to show that \( \{U_{T,n,t}, \mathcal{F}_{T,n,t}\} \) is a martingale difference sequence. We first investigate at what rate the denominator \( \sqrt{c^t Jc} \) goes to zero:

\[ c^t Jc = c^t(E^tWE)^{-1}E^tWD_n^+HP(D^{-1/2} \otimes D^{-1/2})V(D^{-1/2} \otimes D^{-1/2})P^t HD_n^+WE(E^tWE)^{-1}c \]

\[\geq \text{mineval} (E^tWD_n^+HP(D^{-1/2} \otimes D^{-1/2})V(D^{-1/2} \otimes D^{-1/2})P^t HD_n^+WE) \| (E^tWE)^{-1}c \|_2^2 \]

\[\geq \frac{n}{\omega} \text{mineval}^2(W)c^t(E^tWE)^{-2}c \geq \frac{n}{\omega} \text{mineval}^2(W)\text{mineval} \left( (E^tWE)^{-2} \right) \]

\[= \frac{n \cdot \text{mineval}^2(W)}{\omega \text{maxeval}^2(E^tWE)} \geq \frac{\omega}{\omega \kappa^2(W)} \text{maxeval}^2(E^tWE) \]

\[= \frac{n}{\omega \kappa^2(W)} \text{maxeval}^2(E^tWE) \]
where the second inequality is due to Assumption 3.7(ii). Using (A.11), we have
\[
\frac{1}{\sqrt{c^\top Jc}} = O(\sqrt{s^2 \cdot n \cdot \kappa^2(W) \cdot \omega}).
\]

(8.23)

Verification of conditions (i)-(iii) of Theorem A.5 in Appendix A.5 will be exactly the same as that in Section A.4.

8.3.2 \( t_1 - \hat{t}_1 = o_p(1) \)

We now show that \( t_1 - \hat{t}_1 = o_p(1) \). Let \( A \) and \( \hat{A} \) denote the numerators of \( t_1 \) and \( \hat{t}_1 \), respectively.
\[
t_1 - \hat{t}_1 = \frac{A}{\sqrt{c^\top Jc}} - \frac{\hat{A}}{\sqrt{c^\top J_{TC}}} = \frac{\sqrt{s^2 \cdot n \cdot \kappa^2(W) \cdot \omega} A}{\sqrt{s^2 \cdot n \cdot \kappa^2(W) \cdot \omega} c^\top Jc} - \frac{\sqrt{s^2 \cdot n \cdot \kappa^2(W) \cdot \omega} \hat{A}}{\sqrt{s^2 \cdot n \cdot \kappa^2(W) \cdot \omega} c^\top \hat{J}_{TC}}.
\]

Since we have already shown in (8.23) that \( s^2 \cdot n \cdot \kappa^2(W) \cdot \omega \) is bounded away from zero by an absolute constant, it suffices to show the denominators as well as numerators of \( t_1 \) and \( \hat{t}_1 \) are asymptotically equivalent.

8.3.3 Denominators of \( t_1 \) and \( \hat{t}_1 \)

We first show that the denominators of \( t_1 \) and \( \hat{t}_1 \) are asymptotically equivalent, i.e.,
\[
s^2 \cdot n \cdot \kappa^2(W) \cdot \omega |c^\top J_{TC} - c^\top Jc| = o_p(1).
\]

Define
\[
c^\top \hat{J}_{TC} = c^\top (E^\top WE)^{-1} E^\top WD_n^+ \hat{H}_T \hat{P}_T (\hat{D}_T^{-1/2} \otimes \hat{D}_T^{-1/2}) V (\hat{D}_T^{-1/2} \otimes \hat{D}_T^{-1/2}) \hat{P}_T^\top \hat{H}_T D_n^+ W (E^\top WE)^{-1} c.
\]

By the triangular inequality: \( s^2 \cdot n \cdot \kappa^2(W) \cdot \omega |c^\top J_{TC} - c^\top Jc| \leq s^2 \cdot n \cdot \kappa^2(W) \cdot \omega |c^\top \hat{J}_{TC} - c^\top J_{TC}| + s^2 \cdot n \cdot \kappa^2(W) \cdot \omega |c^\top \hat{J}_{TC} - c^\top Jc| \).

First, we prove \( s^2 \cdot n \cdot \kappa^2(W) \cdot \omega |c^\top \hat{J}_{TC} - c^\top J_{TC}| = o_p(1) \).
\[
s^2 \cdot n \cdot \kappa^2(W) \cdot \omega |c^\top \hat{J}_{TC} - c^\top J_{TC}|
\]

\[
= s^2 \cdot n \cdot \kappa^2(W) \cdot \omega |c^\top (E^\top WE)^{-1} E^\top WD_n^+ \hat{H}_T \hat{P}_T (\hat{D}_T^{-1/2} \otimes \hat{D}_T^{-1/2}) V (\hat{D}_T^{-1/2} \otimes \hat{D}_T^{-1/2}) \hat{P}_T^\top \hat{H}_T D_n^+ W (E^\top WE)^{-1} c
\]

\[
- c^\top (E^\top WE)^{-1} E^\top WD_n^+ \hat{H}_T \hat{P}_T (\hat{D}_T^{-1/2} \otimes \hat{D}_T^{-1/2}) V (\hat{D}_T^{-1/2} \otimes \hat{D}_T^{-1/2}) \hat{P}_T^\top \hat{H}_T D_n^+ W (E^\top WE)^{-1} c|
\]

\[
= s^2 \cdot n \cdot \kappa^2(W) \cdot \omega \cdot |c^\top (E^\top WE)^{-1} E^\top WD_n^+ \hat{H}_T \hat{P}_T (\hat{D}_T^{-1/2} \otimes \hat{D}_T^{-1/2}) V (\hat{D}_T^{-1/2} \otimes \hat{D}_T^{-1/2}) \hat{P}_T^\top \hat{H}_T D_n^+ W (E^\top WE)^{-1} c| |
\]

\[
\leq s^2 \cdot n \cdot \kappa^2(W) \cdot \omega \|V_T - V\|_\infty \|((\hat{D}_T^{-1/2} \otimes \hat{D}_T^{-1/2}) \hat{P}_T^\top \hat{H}_T D_n^+ W (E^\top WE)^{-1} c| |
\]

\[
\leq s^2 \cdot n \cdot \kappa^2(W) \cdot \omega \|V_T - V\|_\infty \|((\hat{D}_T^{-1/2} \otimes \hat{D}_T^{-1/2}) \hat{P}_T^\top \hat{H}_T D_n^+ W (E^\top WE)^{-1} c| |
\]

\[
\leq s^2 \cdot n \cdot \kappa^2(W) \cdot \omega \|V_T - V\|_\infty \|((\hat{D}_T^{-1/2} \otimes \hat{D}_T^{-1/2}) \hat{P}_T^\top \hat{H}_T D_n^+ W (E^\top WE)^{-1} c| |
\]

\[
= O_p(s^2 \cdot n \cdot \kappa^2(W) \cdot \omega^2) \|V_T - V\|_\infty = O_p\left(\sqrt{n^4 \kappa^4(W) \cdot s^4 \cdot \omega^4 \cdot \log n \cdot T}\right) = o_p(1),
\]

where \( \| \cdot \|_\infty \) denotes the absolute elementwise maximum, the third equality is due to Lemma A.4(v), Lemma A.16 in Appendix A.5, A.7, A.14, A.8 and 8.22, the second last equality is due to Lemma 8.2 in SM 8.3 and the last equality is due to Assumption 3.3(ii).

We now prove \( s^2 \cdot n \cdot \kappa^2(W) \cdot \omega |c^\top \hat{J}_{TC} - c^\top Jc| = o_p(1) \). Define
\[
c^\top \hat{J}_{a,c} := c^\top (E^\top WE)^{-1} E^\top WD_n^+ \hat{H}_T \hat{P}_T (D^{-1/2} \otimes D^{-1/2}) V (D^{-1/2} \otimes D^{-1/2}) \hat{P}_T^\top \hat{H}_T D_n^+ W (E^\top WE)^{-1} c
\]
\[
c^\top \hat{J}_{b,c} := c^\top (E^\top WE)^{-1} E^\top WD_n^+ \hat{H}_T P(D^{-1/2} \otimes D^{-1/2}) V (D^{-1/2} \otimes D^{-1/2}) P^\top \hat{H}_T D_n^+ W (E^\top WE)^{-1} c.
\]
We use triangular inequality again
\[ s^2n_k^2(W)\varpi|c^T\hat{J}_{T,C} - c^TJC| \leq s^2n_k^2(W)\varpi|c^T\hat{J}_{T,C} - c^T\hat{J}_{T,a,c} + s^2n_k^2(W)\varpi|c^T\hat{J}_{T,a,c} - c^T\hat{J}_{T,b,c} + s^2n_k^2(W)\varpi|c^T\hat{J}_{T,b,c} - c^T\hat{J}_{T,c}| \]
\[ \leq s^2n_k^2(W)\varpi|\max\{V\}|^2(\tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2})\hat{P}_{T}^{\dagger}\hat{H}_{T}D_{n}^+\varpi(\hat{E}^{T}WE)^{-1}c \]
\[ + 2s^2n_k^2(W)\varpi|\max\{V\}|^2(\tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2})\hat{P}_{T}^{\dagger}\hat{H}_{T}D_{n}^+\varpi(\hat{E}^{T}WE)^{-1}c \]
\[ \cdot(\hat{P}_{T}^{\dagger}\hat{H}_{T}D_{n}^+\varpi(\hat{E}^{T}WE)^{-1}c)_{2} \]
(8.25)
where the inequality is due to Lemma A.17 in Appendix A.5

We consider the first term of (8.24).
\[ s^2n_k^2(W)\varpi|c^T\hat{J}_{T,C} - c^T\hat{J}_{T,a,c}| = \]
\[ s^2n_k^2(W)\varpi|c^T(\hat{E}^{T}WE)^{-1}E^{T}WD_{n}^+\hat{H}_{T}\hat{P}_{T}(\tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2})V(\tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2})\hat{P}_{T}^{\dagger}\hat{H}_{T}D_{n}^+\varpi(\hat{E}^{T}WE)^{-1}c \]
\[ = O(s^2n_k^2(W)\varpi)\|\hat{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2}\|_{\ell_{2}}\|\hat{P}_{T}^{\dagger}\|_{2}\|\hat{H}_{T}\|_{2}\|D_{n}^+\|_{2}\|\varpi(\hat{E}^{T}WE)^{-1}c\|_{2} \]
\[ = O_p(s^2n_k^3(W)\varpi^{2}/T) = o_p(1), \]
where the second last equality is due to (A.7), (A.8), (A.14), (8.22) and Lemma A.4(vii), and the last equality is due to Assumption 3.3(ii).

We now consider the second term of (8.25).
\[ 2s^2n_k^2(W)\varpi|\max\{V\}|^2(\tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2})\hat{P}_{T}^{\dagger}\hat{H}_{T}D_{n}^+\varpi(\hat{E}^{T}WE)^{-1}c \]
\[ \cdot(\tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2})\hat{P}_{T}^{\dagger}\hat{H}_{T}D_{n}^+\varpi(\hat{E}^{T}WE)^{-1}c \]
\[ \leq O(s^2n_k^2(W)\varpi)\|\hat{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2}\|_{\ell_{2}}\|\hat{P}_{T}^{\dagger}\|_{2}\|\hat{H}_{T}\|_{2}\|D_{n}^+\|_{2}\|\varpi(\hat{E}^{T}WE)^{-1}c\|_{2} \]
\[ = O(\sqrt{s^4n_k^6(W)\varpi^{2}/T}) = o_p(1), \]
where the first equality is due to (A.7), (A.8), (A.14), (8.22) and Lemma A.4(vii), and the last equality is due to Assumption 3.3(ii). We have proved \( s^2n_k^2(W)\varpi|c^T\hat{J}_{T,C} - c^T\hat{J}_{T,a,c}| = o_p(1) \).

We consider the second term on the right hand side of (8.24).
\[ s^2n_k^2(W)\varpi|c^T\hat{J}_{T,a,c} - c^T\hat{J}_{T,b,c}| = \]
\[ s^2n_k^2(W)\varpi|c^T(\hat{E}^{T}WE)^{-1}E^{T}WD_{n}^+\hat{H}_{T}\hat{P}_{T}(\tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2})V(\tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2})\hat{P}_{T}^{\dagger}\hat{H}_{T}D_{n}^+\varpi(\hat{E}^{T}WE)^{-1}c \]
\[ - c^T(\hat{E}^{T}WE)^{-1}E^{T}WD_{n}^+\hat{H}_{T}\hat{P}_{T}(\tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2})V(\tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2})\hat{P}_{T}^{\dagger}\hat{H}_{T}D_{n}^+\varpi(\hat{E}^{T}WE)^{-1}c \]
\[ \leq s^2n_k^2(W)\varpi|\max\{V\}|^2(\tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2})\hat{P}_{T}^{\dagger}\hat{H}_{T}D_{n}^+\varpi(\hat{E}^{T}WE)^{-1}c \]
\[ + 2s^2n_k^2(W)\varpi(\tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2})\hat{P}_{T}^{\dagger}\hat{H}_{T}D_{n}^+\varpi(\hat{E}^{T}WE)^{-1}c \]
\[ \cdot(\hat{P}_{T}^{\dagger}\hat{H}_{T}D_{n}^+\varpi(\hat{E}^{T}WE)^{-1}c)_{2} \]
(8.26)
where the inequality is due to Lemma A.17 in Appendix A.5

We consider the first term of (8.26) first.
\[ s^2n_k^2(W)\varpi|\max\{V\}|^2(\tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2} - \tilde{D}_{T}^{-1/2})\hat{P}_{T}^{\dagger}\hat{H}_{T}D_{n}^+\varpi(\hat{E}^{T}WE)^{-1}c \]
\[ = O(s^2n_k^2(W)\varpi)\|\hat{P}_{T}^{\dagger}\|_{2}\|\hat{H}_{T}\|_{2}\|D_{n}^+\|_{2}\|\varpi(\hat{E}^{T}WE)^{-1}c\|_{2} \]
\[ = O_p(s^2n_k^3(W)\varpi^{2}/T) = o_p(1), \]
9
where the second last equality is due to (A.7), (A.8), (A.14), and (8.22), and the last equality is due to Assumption 3.3(ii).

We now consider the second term of (8.26).

$$
2s^2 n \kappa^2 (W) \omega \Vert (D^{-1/2} \otimes D^{-1/2}) V (D^{-1/2} \otimes D^{-1/2}) P^T \hat{H}_T D_n^+ W (E^T W E)^{-1} c \Vert_2
\cdot \Vert (\hat{H}_T - P)^T \hat{H}_T D_n^+ W (E^T W E)^{-1} c \Vert_2
\leq O(s^2 n \kappa^2 (W) \omega \Vert \hat{H}_T - P \Vert_2^2 \Vert \hat{H}_T \Vert_2^2 \Vert D_n^+ \Vert_2^2 \Vert W (E^T W E)^{-1} \Vert_2^2
= O(\sqrt{s^4 n \kappa^6 (W) \omega^4 / T}) = o_p(1),
$$

where the first equality is due to (A.7), (A.8), (A.14), and (8.22), and the last equality is due to Assumption 3.3(ii). We have proved $s^2 n \kappa^2 (W) \omega [c^T \hat{J}_{T,c} - c^T J_{T,c}] = o_p(1)$.

We consider the third term on the right hand side of (8.24).

$$
s^2 n \kappa^2 (W) \omega [c^T \hat{J}_{T,b,c} - c^T J_c] =
$$

where the inequality is due to Lemma A.17 in Appendix A.5. We consider the first term of (8.27) first.

$$
s^2 n \kappa^2 (W) \omega \max [P(D^{-1/2} \otimes D^{-1/2}) V (D^{-1/2} \otimes D^{-1/2}) P^T] \Vert (\hat{H}_T - H) D_n^+ W (E^T W E)^{-1} c \Vert_2^2
\leq s^2 n \kappa^2 (W) \omega \max [P(D^{-1/2} \otimes D^{-1/2}) V (D^{-1/2} \otimes D^{-1/2}) P^T] \Vert P(D^{-1/2} \otimes D^{-1/2}) P^T H D_n^+ W (E^T W E)^{-1} c \Vert_2^2
+ 2s^2 n \kappa^2 (W) \omega \Vert P(D^{-1/2} \otimes D^{-1/2}) V (D^{-1/2} \otimes D^{-1/2}) P^T H D_n^+ W (E^T W E)^{-1} c \Vert_2
\cdot \Vert (\hat{H}_T - H) D_n^+ W (E^T W E)^{-1} c \Vert_2
= O(s^2 n \kappa^3 (W) \omega^2 / T) = o_p(1),
$$

where the second last equality is due to (A.7), (A.8), and (A.14), and the last equality is due to Assumption 3.3(ii).

We now consider the second term of (8.27).

$$
2s^2 n \kappa^2 (W) \omega \Vert P(D^{-1/2} \otimes D^{-1/2}) V (D^{-1/2} \otimes D^{-1/2}) P^T H D_n^+ W (E^T W E)^{-1} c \Vert_2
\cdot \Vert (\hat{H}_T - H) D_n^+ W (E^T W E)^{-1} c \Vert_2
\leq O(s^2 n \kappa^2 (W) \omega \Vert \hat{H}_T - H \Vert_2^2 \Vert D_n^+ \Vert_2^2 \Vert W (E^T W E)^{-1} \Vert_2^2
= O(\sqrt{s^4 n \kappa^6 (W) \omega^4 / T}) = o_p(1),
$$

where the first equality is due to (A.7), (A.8), and (A.14), and the last equality is due to Assumption 3.3(ii). We have proved $s^2 n \kappa^2 (W) \omega [c^T \hat{J}_{T,b,c} - c^T J_{T,c}] = o_p(1)$. Hence we have proved $s^2 n \kappa^2 (W) \omega [c^T \hat{J}_{T,b,c} - c^T J_{T,c}] = o_p(1)$.

### 8.3.4 Numerators of $t_1$ and $\hat{t}_1$

We now show that numerators of $t_1$ and $\hat{t}_1$ are asymptotically equivalent, i.e.,

$$
\sqrt{s^2 n \kappa^2 (W) \omega} |A - \hat{A}| = o_p(1).
$$
To show the final equality is due to Assumption 3.3(ii), where the second equality is due to Assumption 3.7(i), the third equality is due to Lemma A.3, the fourth equality is due to Assumption 3.7(i). Note that

\[ √{\mathcal{N}} = O \]

\[ s \sqrt{nκ} \cdot \| \mathcal{D} \|_2 \parallel \vec{Θ} \parallel_2 H \left( \frac{\partial \vec{Θ}}{\partial \vec{Σ}} \right)_{\Sigma=\tilde{Σ}_T} = O(D^{-1/2} \otimes D^{-1/2}) \parallel \vec{Σ}_T - \Sigma \parallel_2 \]

where the second equality is due to Assumption 3.7(i), the third equality is due to Lemma A.3 and the fourth equality is due to Assumption 3.3(ii).

We now show that \( \sqrt{s nk^2(W)\bar{w}} | \hat{A}_n - A | = o_p(1) \).

\[ \sqrt{s n k^2(W)\bar{w}} | c^T (E^T W E)^{-1} E^T W D_n^\dagger H | \left| \frac{\partial \vec{Θ}}{\partial \vec{Σ}} \right|_{\Sigma=\tilde{Σ}_T} = O(D^{-1/2} \otimes D^{-1/2}) \parallel \vec{Σ}_T - \Sigma \parallel_2 \]

\[ \leq O(\sqrt{s n k^2(W)\bar{w}}) \sqrt{n} \parallel \vec{Σ}_T - \Sigma \parallel_F \leq O(\sqrt{s n k^2(W)\bar{w}}) \sqrt{n} \parallel \vec{Σ}_T - \Sigma \parallel_2 \]

\[ = O_p(\sqrt{\frac{n s^2 κ^2(W)\bar{w}^2}{T}}) = o_p(1), \]

where the third equality is due to \( (8.21) \), the last equality is due to Assumption 3.3(ii), and the second equality is due to \( (A.7) \), \( (A.8) \), \( (A.14) \), and the fact that

\[ \parallel \frac{\partial \vec{Θ}}{\partial \vec{Σ}} \parallel_{\Sigma=\tilde{Σ}_T} = O_p\left(\sqrt{\frac{n}{T}} + O(1) = O_p(1)\right). \]
8.3.5 \( \hat{t}_2 = o_p(1) \)

Write

\[
\hat{t}_2 = \frac{\sqrt{T} \sqrt{s^2 n^2 (W) \omega \epsilon^T (E^T W E)^{-1} E^T W D_n^+} \vec{O}_p(\|\hat{\Theta} - \Theta\|^2_{\ell_2})}{\sqrt{s^2 n^2 (W) \omega \epsilon^T J_{TC}}}
\]

Since the denominator of the preceding equation is bounded away from zero by an absolute constant with probability approaching one by (8.23) and that \( s^2 n^2 (W) \omega \epsilon^T |J_{TC} - \omega \epsilon^T Jc| = o_p(1) \), it suffices to show

\[
\sqrt{T} \sqrt{s^2 n^2 (W) \omega \epsilon^T (E^T W E)^{-1} E^T W D_n^+} \vec{O}_p(\|\hat{\Theta} - \Theta\|^2_{\ell_2}) = o_p(1).
\]

This is straightforward:

\[
|\sqrt{T} s^2 n^2 (W) \omega \epsilon^T (E^T W E)^{-1} E^T W D_n^+ \vec{O}_p(\|\hat{\Theta} - \Theta\|^2_{\ell_2})| \\
\leq \sqrt{T} s^2 n^2 (W) \omega \epsilon^T (E^T W E)^{-1} E^T W D_n^+ \vec{O}_p(\|\hat{\Theta} - \Theta\|^2_{\ell_2})_2 \\
= O(\sqrt{T} s^2 n^2 (W) \omega \epsilon^T) \vec{O}_p(\|\hat{\Theta} - \Theta\|^2_{\ell_2})_F = O(\sqrt{T} s^2 n^2 (W) \omega \epsilon^T) \vec{O}_p(\|\hat{\Theta} - \Theta\|^2_{\ell_2})_{\ell_2} \\
= O(\sqrt{s^2 n^2 (W) \omega \epsilon^T}) \vec{O}_p(\|\hat{\Theta} - \Theta\|^2_{\ell_2}) = O_p(1),
\]

where the last equality is due to Assumption 3.3 (ii).

8.4 Proof of Theorem 4.1

In this subsection, we give a proof for Theorem 4.1. We first give a useful lemma which is used in the proof of Theorem 4.1.

**Lemma 8.4.** Let \( A \) and \( B \) be \( m \times n \) and \( p \times q \) matrices, respectively. There exists a unique permutation matrix \( P_K := I_n \otimes K_{q,m} \otimes I_p \), where \( K_{q,m} \) is the commutation matrix, such that

\[ \vec{c}(A \otimes B) = P_K(\vec{c}A \otimes \vec{c}B). \]


**Proof of Theorem 4.1.** At each step, we take the symmetry of \( \Omega(\theta) \) into account.

\[
d\ell_{T,D(\theta, \mu)} = -\frac{T}{2} d\log|D^{1/2} \exp(\Omega) D^{1/2}| - \frac{T}{2} dtr\left(\frac{1}{T} \sum_{t=1}^{T} (y_t - \mu)^T D^{-1/2} [\exp(\Omega)]^{-1} D^{-1/2} (y_t - \mu)\right)
\]

\[
= -\frac{T}{2} d\log|D^{1/2} \exp(\Omega) D^{1/2}| - \frac{T}{2} dtr\left(D^{-1/2} \tilde{\Sigma}_T D^{-1/2} [\exp(\Omega)]^{-1}\right)
\]

\[
= -\frac{T}{2} tr\left([\exp(\Omega)]^{-1} d \exp(\Omega)\right) - \frac{T}{2} tr\left(D^{-1/2} \tilde{\Sigma}_T D^{-1/2} [\exp(\Omega)]^{-1}\right)
\]

\[
= -\frac{T}{2} tr\left([\exp(\Omega)]^{-1} d \exp(\Omega)\right) + \frac{T}{2} tr\left(D^{-1/2} \tilde{\Sigma}_T D^{-1/2} [\exp(\Omega)]^{-1}\right)
\]

\[
= -\frac{T}{2} \left(\vec{c}[\exp(\Omega)]^{-1} D^{-1/2} \tilde{\Sigma}_T D^{-1/2} [\exp(\Omega)]^{-1} - [\exp(\Omega)]^{-1}\right)_{\ell_2} d \exp(\Omega)
\]

\[
= \frac{T}{2} \left(\vec{c}[\exp(\Omega)]^{-1} D^{-1/2} \tilde{\Sigma}_T D^{-1/2} [\exp(\Omega)]^{-1} - [\exp(\Omega)]^{-1}\right)^T \vec{c} d \exp(\Omega)
\]

\[
= \frac{T}{2} \left(\vec{c}[\exp(\Omega)]^{-1} D^{-1/2} \tilde{\Sigma}_T D^{-1/2} [\exp(\Omega)]^{-1} - [\exp(\Omega)]^{-1}\right)^T \vec{c} d \exp(\Omega),
\]
where in the second equality we used the definition of $\tilde{\Sigma}_T$, the third equality is due to that $d \log |X| = \text{tr}(X^{-1} dX)$, the fifth equality is due to that $dX^{-1} = -X^{-1}(dX)X^{-1}$, the seventh equality is due to that $\text{tr}(AB) = (\text{vec}[A^T])^T \text{vec} B$, and the eighth equality is due to that matrix function preserves symmetry and we can interchange inverse and transpose operators.

The following differential of matrix exponential can be inferred from (10.15) in Higham (2008) p238:

\[ d \exp(\Omega) = \int_0^1 e^{(1-t)\Omega}(d\Omega)e^{t\Omega}dt. \]

Therefore,

\[ \text{vec} d \exp(\Omega) = \int_0^1 e^{t\Omega} \otimes e^{(1-t)\Omega} dt \text{vec}(d\Omega) = \int_0^1 e^{t\Omega} \otimes e^{(1-t)\Omega} dtD_n \text{vech}(d\Omega) \]

Hence,

\[ d\ell_{T,D}(\theta, \mu) = \frac{T}{2} \left[ \text{vec} \left( [\exp(\Omega)]^{-1/2} \tilde{\Sigma}_T D^{-1/2} [\exp(\Omega)]^{-1} \right) - [\exp(\Omega)]^{-1} \right]^T \int_0^1 e^{t\Omega} \otimes e^{(1-t)\Omega} dtD_n E d\theta. \]

and

\[ y := \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta^T} \]

\[ = \frac{T}{2} E^T D_n^T \int_0^1 e^{t\Omega} \otimes e^{(1-t)\Omega} dt \left[ \text{vec} \left( [\exp(\Omega)]^{-1/2} \tilde{\Sigma}_T D^{-1/2} [\exp(\Omega)]^{-1} \right) \right] \]

\[ =: \frac{T}{2} E^T D_n^T \Psi_1 \Psi_2. \]

Now we derive the Hessian matrix.

\[ dy = \frac{T}{2} E^T D_n^T (d\Psi_1) \Psi_1^T + \frac{T}{2} E^T D_n^T \Psi_2^T d\Psi_2 = \frac{T}{2} (\Psi_2^T \otimes E^T D_n^T) \text{vec} d\Psi_1 + \frac{T}{2} E^T D_n^T \Psi_1 d\Psi_2. \] (8.28)

Consider $d\Psi_1$ first.

\[ d\Psi_1 = d \int_0^1 e^{t\Omega} \otimes e^{(1-t)\Omega} dt = \int_0^1 de^{t\Omega} \otimes e^{(1-t)\Omega} dt + \int_0^1 e^{t\Omega} \otimes de^{(1-t)\Omega} dt \]

\[ =: \int_0^1 A \otimes e^{(1-t)\Omega} dt + \int_0^1 e^{t\Omega} \otimes B dt, \]

where

\[ A := \int_0^1 e^{(1-s)\Omega} d(t) e^{s\Omega} ds, \quad B := \int_0^1 e^{(1-s)(1-t)\Omega} d((1-t)\Omega) e^{s(1-t)\Omega} ds. \]
Therefore,
\[
\text{vec} \, d\Psi_1 = \int_0^1 \text{vec} \, (A \otimes e^{(1-t)\Omega}) \, dt + \int_0^1 \text{vec} \, (e^{t\Omega} \otimes B) \, dt \\
= \int_0^1 P_K \left( \text{vec} \, A \otimes \text{vec} \, e^{(1-t)\Omega} \right) \, dt + \int_0^1 P_K \left( \text{vec} \, e^{t\Omega} \otimes \text{vec} \, B \right) \, dt \\
= \int_0^1 P_K \left( I_{n^2} \otimes \text{vec} \, e^{(1-t)\Omega} \right) \, vec \, Adt + \int_0^1 P_K \left( \text{vec} \, e^{t\Omega} \otimes I_{n^2} \right) \, vec \, Bdt \\
= \int_0^1 P_K \left( I_{n^2} \otimes \text{vec} \, e^{(1-t)\Omega} \right) \, \int_0^1 e^{st\Omega} \otimes e^{(1-s)\Omega} \, ds \cdot vec \, d(t\Omega) \, dt \\
+ \int_0^1 P_K \left( \text{vec} \, e^{t\Omega} \otimes I_{n^2} \right) \, \int_0^1 e^{s(1-t)\Omega} \otimes e^{(1-s)(1-t)\Omega} \, ds \cdot \text{vec} \, d((1-t)\Omega) \, dt \\
= \int_0^1 P_K \left( I_{n^2} \otimes \text{vec} \, e^{(1-t)\Omega} \right) \, \int_0^1 e^{st\Omega} \otimes e^{(1-s)\Omega} \, ds \cdot tdtT,DnEdt \\
+ \int_0^1 P_K \left( \text{vec} \, e^{t\Omega} \otimes I_{n^2} \right) \, \int_0^1 e^{s(1-t)\Omega} \otimes e^{(1-s)(1-t)\Omega} \, ds \cdot (1-t)dtT,DnEdt \\
\tag{8.29}
\]
where \( P_K := I_n \otimes K_{n,n} \otimes I_n \), the second equality is due to Lemma 8.4. We now consider \( d\Psi_2 \).
\[
d\Psi_2 = d \text{vec} \left( [\exp(\Omega)]^{-1} D^{-1/2} \tilde{\Sigma}_T D^{-1/2} [\exp(\Omega)]^{-1} - [\exp(\Omega)]^{-1} \right) \\
= \text{vec} \left( d[\exp(\Omega)]^{-1} D^{-1/2} \tilde{\Sigma}_T D^{-1/2} [\exp(\Omega)]^{-1} \right) \\
+ \left( [\exp(\Omega)]^{-1} D^{-1/2} \tilde{\Sigma}_T D^{-1/2} d[\exp(\Omega)]^{-1} \right) - \text{vec} \left( d [\exp(\Omega)]^{-1} \right) \\
= \text{vec} \left( -[\exp(\Omega)]^{-1} d\exp(\Omega) [\exp(\Omega)]^{-1} D^{-1/2} \tilde{\Sigma}_T D^{-1/2} [\exp(\Omega)]^{-1} \right) \\
+ \text{vec} \left( [\exp(\Omega)]^{-1} D^{-1/2} \tilde{\Sigma}_T D^{-1/2} d\exp(\Omega) [\exp(\Omega)]^{-1} \right) \\
+ \text{vec} \left( [\exp(\Omega)]^{-1} \otimes [\exp(\Omega)]^{-1} \right) \text{vec} \, d\exp(\Omega) \\
- \left( [\exp(\Omega)]^{-1} \otimes [\exp(\Omega)]^{-1} D^{-1/2} \tilde{\Sigma}_T D^{-1/2} [\exp(\Omega)]^{-1} \right) \text{vec} \, d\exp(\Omega) \\
- \left( [\exp(\Omega)]^{-1} D^{-1/2} \tilde{\Sigma}_T D^{-1/2} [\exp(\Omega)]^{-1} \otimes [\exp(\Omega)]^{-1} \right) \text{vec} \, d\exp(\Omega) \\
\tag{8.30}
\]
Substituting \[8.29\] and \(8.30\) into \(8.28\) yields the result:
\[
\frac{\partial^2 \ell_{T,D}(\theta, \mu)}{\partial \theta \partial \theta^T} = \\
- \frac{T}{2} E^T D_n^T \Psi_1 \left( [\exp(\Omega)]^{-1} D^{-1/2} \tilde{\Sigma}_T D^{-1/2} \otimes I_n + I_n \otimes [\exp(\Omega)]^{-1} D^{-1/2} \tilde{\Sigma}_T D^{-1/2} - I_{n^2} \right) \cdot \\
\left( [\exp(\Omega)]^{-1} \otimes [\exp(\Omega)]^{-1} \right) \Psi_1 D_n E \\
+ \frac{T}{2} (\Psi_2 \otimes E^T D_n^T) \int_0^1 P_K \left( I_{n^2} \otimes \text{vec} \, e^{(1-t)\Omega} \right) \, \int_0^1 e^{s(1-t)\Omega} \otimes e^{(1-s)(1-t)\Omega} \, ds \cdot tdtT,DnE \\
+ \frac{T}{2} (\Psi_2 \otimes E^T D_n^T) \int_0^1 P_K \left( \text{vec} \, e^{t\Omega} \otimes I_{n^2} \right) \, \int_0^1 e^{s(1-t)\Omega} \otimes e^{(1-s)(1-t)\Omega} \, ds \cdot (1-t)dtT,DnE.
\]

\[8.5\] **Proof of Theorem 4.2**

In this subsection, we give a proof for Theorem 4.2. We will first give some preliminary lemmas leading to the proof of this theorem.
Lemma 8.5. Suppose Assumptions 3.1(i), 3.2, 3.3(i) and 3.4(i) hold with $1/r_1 + 1/r_2 > 1$. Then

(i) \[ \Xi = \int_0^1 \int_0^1 \Theta^{t+s-1} \otimes \Theta^{1-t-s} dtds \]
has minimum eigenvalue bounded away from zero by an absolute constant and maximum eigenvalue bounded from above by an absolute constant.

(ii) \[ \hat{\Xi}_T = \int_0^1 \int_0^1 \Theta_T^{t+s-1} \otimes \Theta_T^{1-t-s} dtds \]
has minimum eigenvalue bounded away from zero by an absolute constant and maximum eigenvalue bounded from above by an absolute constant with probability approaching 1.

(iii) \[ \|\hat{\Xi}_T - \Xi\|_{\ell_2} = O_p\left(\sqrt{\frac{n}{T}}\right) \]

(iv) \[ \|\Psi_1\|_{\ell_2} = \left\| \int_0^1 e^{\Omega(\theta)} \otimes e^{(1-t)\Omega(\theta)} dt \right\|_{\ell_2} = O(1) \]

Proof. The proofs for the first two parts are the same, so we only give one for part (i). Under assumptions of this lemma, we can invoke Lemma A.7(i) to have eigenvalues of $\Theta$ to be bounded away from zero and from above by absolute positive constants. Let $\lambda_1, \ldots, \lambda_n$ denote these eigenvalues. Suppose $\Theta = Q^T\text{diag}(\lambda_1, \ldots, \lambda_n)Q$ (orthogonal diagonalization). By definition of matrix function, we have

\[ \Theta^{(t+s-1)} = Q^T\text{diag}(\lambda_1^{(t+s-1)}, \ldots, \lambda_n^{(t+s-1)})Q \]
\[ \Theta^{(1-s-t)} = Q^T\text{diag}(\lambda_1^{(1-s-t)}, \ldots, \lambda_n^{(1-s-t)})Q \]
\[ \Theta^{(t+s-1)} \otimes \Theta^{(1-s-t)} = (Q \otimes Q)^T \left[ \text{diag}(\lambda_1^{(t+s-1)}, \ldots, \lambda_n^{(t+s-1)}) \otimes \text{diag}(\lambda_1^{(1-s-t)}, \ldots, \lambda_n^{(1-s-t)}) \right] (Q \otimes Q) \]
\[ =: (Q \otimes Q)^T M_2 (Q \otimes Q), \]

where $M_2$ is an $n^2 \times n^2$ diagonal matrix whose $[(i-1)n + j]^{th}$ diagonal entry is

\[ \begin{cases} 
1 & \text{if } i = j \\
\frac{1}{\log(\frac{\lambda_i}{\lambda_j})} \frac{\lambda_i}{\lambda_j} \left( \frac{\lambda_i}{\lambda_j} - 1 \right)^2 & \text{if } i \neq j, \lambda_i \neq \lambda_j 
\end{cases} \]

for $i, j = 1, \ldots, n$.

Thus

\[ \int_0^1 \int_0^1 \Theta^{t+s-1} \otimes \Theta^{1-t-s} dtds = (Q \otimes Q)^T \int_0^1 \int_0^1 M_2 dtds (Q \otimes Q) \]

where $\int_0^1 \int_0^1 M_2 dtds$ is an $n^2 \times n^2$ diagonal matrix whose $[(i-1)n + j]^{th}$ diagonal entry is

\[ \begin{cases} 
1 & \text{if } i = j \\
1 & \text{if } i \neq j, \lambda_i = \lambda_j \\
\frac{1}{\log(\frac{\lambda_i}{\lambda_j})} \frac{\lambda_i}{\lambda_j} \left( \frac{\lambda_i}{\lambda_j} - 1 \right)^2 & \text{if } i \neq j, \lambda_i \neq \lambda_j 
\end{cases} \]
To see this,

\[
\int_0^1 \int_0^1 \left( \frac{\lambda_i}{\lambda_j} \right)^{s+t-1} dtds = \frac{\lambda_j}{\lambda_i} \int_0^1 \left( \frac{\lambda_i}{\lambda_j} \right)^s ds \int_0^1 \left( \frac{\lambda_i}{\lambda_j} \right)^t dt \\
= \frac{\lambda_j}{\lambda_i} \left[ \int_0^1 \left( \frac{\lambda_i}{\lambda_j} \right)^s ds \right]^2 = \frac{\lambda_j}{\lambda_i} \left[ \frac{\left( \frac{\lambda_i}{\lambda_j} \right)^s}{\log \left( \frac{\lambda_i}{\lambda_j} \right)} \right]_0^1 = \frac{1}{\log \left( \frac{\lambda_i}{\lambda_j} \right)^2} \frac{\lambda_j}{\lambda_i} \left[ \frac{\lambda_i}{\lambda_j} - 1 \right]^2 .
\]

For part (iii), we have

\[
\left\| \int_0^1 \int_0^1 \hat{\Theta}^{t+s-1}_T \otimes \hat{\Theta}^{1-t-s}_T dtds - \int_0^1 \int_0^1 \Theta^{t+s-1} \otimes \Theta^{1-t-s} dtds \right\|_{\ell_2}
\]

\[
\leq \int_0^1 \int_0^1 \left\| \hat{\Theta}^{t+s-1}_T - \Theta^{t+s-1} \otimes \Theta^{1-t-s} \right\|_{\ell_2} dtds
\]

\[
= \int_0^1 \int_0^1 \left\| \hat{\Theta}^{t+s-1}_T - \hat{\Theta}^{1-t-s}_T - \Theta^{t+s-1} \otimes \Theta^{1-t-s} + \Theta^{t+s-1} \otimes \Theta^{1-t-s} - \Theta^{t+s-1} \otimes \Theta^{1-t-s} \right\|_{\ell_2} dtds
\]

\[
= \int_0^1 \int_0^1 \left[ \| \hat{\Theta}^{t+s-1}_T - \Theta^{1-t-s} \|_{\ell_2} + \| \hat{\Theta}^{t+s-1}_T - \Theta^{t+s-1} \|_{\ell_2} + \| \Theta^{t+s-1} - \Theta^{1-t-s} \|_{\ell_2} + \| \Theta^{t+s-1} - \Theta^{t+s-1} \|_{\ell_2} \right] dtds
\]

\[
\leq \max_{t,s \in [0,1]} \left[ \| \hat{\Theta}^{t+s-1}_T - \Theta^{1-t-s} \|_{\ell_2} + \| \hat{\Theta}^{t+s-1}_T - \Theta^{t+s-1} \|_{\ell_2} + \| \Theta^{t+s-1} - \Theta^{1-t-s} \|_{\ell_2} + \| \Theta^{t+s-1} - \Theta^{t+s-1} \|_{\ell_2} \right] .
\]

First, note that for any \( t, s \in [0,1] \), \( \| \hat{\Theta}^{t+s-1}_T \|_{\ell_2} \) and \( \| \Theta^{1-t-s} \|_{\ell_2} \) are \( O_p(1) \) and \( O(1) \), respectively. For example, diagonalize \( \Theta \), apply the function \( f(x) = x^{1-t-s} \), and take the spectral norm.

The theorem would then follow if we show that

\[
\max_{t,s \in [0,1]} \| \hat{\Theta}^{t+s-1}_T - \Theta^{1-t-s} \|_{\ell_2} = O_p(\sqrt{n/T}), \quad \max_{t,s \in [0,1]} \| \hat{\Theta}^{t+s-1}_T - \Theta^{t+s-1} \|_{\ell_2} = O_p(\sqrt{n/T}).
\]

It suffices to give a proof for the first equation, as the proof for the second is similar.

\[
\| \hat{\Theta}^{1-t-s}_T - \Theta^{1-t-s} \|_{\ell_2} = \| e^{(1-t-s) \log \hat{\Theta}_T} - e^{(1-t-s) \log \Theta} \|
\]

\[
\leq \| (1-t-s) (\log \hat{\Theta}_T - \log \Theta) \|_{\ell_2} \exp[(1-t-s) \| \log \hat{\Theta}_T - \log \Theta \|_{\ell_2}] + \| (1-t-s) \| \log \Theta \|_{\ell_2}] = \| (1-t-s) (\log \hat{\Theta}_T - \log \Theta) \|_{\ell_2} \exp[(1-t-s) \| \log \Theta \|_{\ell_2}] + \| (1-t-s) \| \log \Theta \|_{\ell_2}] \log \Theta \|_{\ell_2},
\]

where the first inequality is due to Theorem A.4 in Appendix A.5 and the second equality is due to the fact that all the eigenvalues of \( \Theta \) are bounded away from zero and infinity by absolute positive constants. Now use Theorem 3.1 to get

\[
\| \log \hat{\Theta}_T - \log \Theta \|_{\ell_2} = O_p \left( \sqrt{\frac{n}{T}} \right).
\]

The result follows after recognising \( \exp(o_p(1)) = O_p(1) \).

For part (iv), since \( \Theta = Q^T \text{diag}(\lambda_1, \ldots, \lambda_n)Q \), we have

\[
\Theta^t = Q^T \text{diag}(\lambda_1^t, \ldots, \lambda_n^t)Q, \quad \Theta^{1-t} = Q^T \text{diag}(\lambda_1^{1-t}, \ldots, \lambda_n^{1-t})Q.
\]

Then

\[
\Theta^t \otimes \Theta^{1-t} = (Q \otimes Q)^T [ \text{diag}(\lambda_1^t, \ldots, \lambda_n^t) \otimes \text{diag}(\lambda_1^{1-t}, \ldots, \lambda_n^{1-t}) ] (Q \otimes Q) =: (Q \otimes Q)^T M_3(Q \otimes Q),
\]

16
where $M_3$ is an $n^2 \times n^2$ diagonal matrix whose $[(i-1)n+j]$th diagonal entry is

$$
\begin{cases}
    1 & \text{if } i = j \\
    1 & \text{if } i \neq j, \lambda_i = \lambda_j \\
    \lambda_j \left( \frac{\lambda_i}{\lambda_j} \right)^t & \text{if } i \neq j, \lambda_i \neq \lambda_j
\end{cases}
$$

for $i, j = 1, \ldots, n$.

Thus

$$
\Psi_1 = \int_0^1 \Theta^t \otimes \Theta^{1-t} dt = (Q \otimes Q)^t \int_0^1 M_3 dt (Q \otimes Q)
$$

where $\int_0^1 M_3 dt$ is an $n^2 \times n^2$ diagonal matrix whose $[(i-1)n+j]$th diagonal entry is

$$
\begin{cases}
    1 & \text{if } i = j \\
    1 & \text{if } i \neq j, \lambda_i = \lambda_j \\
    \lambda_i - \lambda_j \log \lambda_i - \log \lambda_j & \text{if } i \neq j, \lambda_i \neq \lambda_j
\end{cases}
$$

To see this,

$$
\lambda_j \int_0^1 \left( \frac{\lambda_i}{\lambda_j} \right)^t dt = \lambda_j \int_0^1 \left( \frac{\lambda_i}{\lambda_j} \right)^t dt = \lambda_j \left[ \left( \frac{\lambda_i}{\lambda_j} \right)^t \right]_0^1 = \frac{1}{\log(\frac{\lambda_i}{\lambda_j})} \lambda_j \left[ \frac{\lambda_i}{\lambda_j} - 1 \right].
$$

Lemma 8.6. Suppose Assumptions 3.1(i), 3.2, 3.3(i) and 3.4 hold with $1/r_1 + 1/r_2 > 1$. Then

(i)

$$
\| \hat{\Upsilon}_{T,D} - \Upsilon_D \|_{\ell_2} = O_p \left( \frac{s n \sqrt{n}}{T} \right).
$$

(ii)

$$
\| \hat{\Upsilon}_{T,D}^{-1} - \Upsilon_D^{-1} \|_{\ell_2} = O_p \left( \frac{s^2 n \sqrt{1/nT}}{T} \right).
$$

Proof. For part (i),

$$
\| \hat{\Upsilon}_{T,D} - \Upsilon_D \|_{\ell_2} = \frac{1}{2} \| E^T D_n^\top (\hat{\Xi}_T - \Xi) D_n E \|_{\ell_2} \le \frac{1}{2} \| E^T \|_{\ell_2} \| D_n^\top \|_{\ell_2} \| \hat{\Xi}_T - \Xi \|_{\ell_2} \| D_n \|_{\ell_2} \| E \|_{\ell_2} = O(1) \| \hat{\Xi}_T - \Xi \|_{\ell_2} \| E \|_{\ell_2} = O_p \left( \frac{s n \sqrt{n}}{T} \right),
$$

where the second equality is due to (A.8), and the last equality is due to (A.12) and Lemma 8.5(iii).

For part (ii),

$$
\| \hat{\Upsilon}_{T,D}^{-1} - \Upsilon_D^{-1} \|_{\ell_2} = \| \hat{\Upsilon}_{T,D}^{-1} (\Upsilon_D - \hat{\Upsilon}_{T,D}) \Upsilon_D^{-1} \|_{\ell_2} \le \| \hat{\Upsilon}_{T,D}^{-1} \|_{\ell_2} \| \Upsilon_D - \hat{\Upsilon}_{T,D} \|_{\ell_2} \| \Upsilon_D^{-1} \|_{\ell_2} = O_p \left( \frac{s^2 n \sqrt{1/nT}}{T} \right),
$$

where the second last equality is due to (8.31).
Therefore without loss of generality, we have, for any $\parallel \cdot \parallel_{T,D}$ bounded away from zero by an absolute constant with probability approaching one. Also as a by-product
\[
\parallel \hat{\Upsilon}_{T,D}^{-1} \parallel_{\ell_2} = \frac{1}{\text{mineval}(\Upsilon_{T,D})} = O_p\left(\frac{\varpi}{n}\right) \quad \parallel \Upsilon_D^{-1} \parallel_{\ell_2} = \frac{1}{\text{mineval}(\Upsilon_D)} = O\left(\frac{\varpi}{n}\right). \tag{8.31}
\]
From the definition of $\hat{\theta}_{T,D}$, for any $b \in \mathbb{R}^s$ with $\parallel b \parallel_2 = 1$ we can write
\[
\sqrt{T}b^\top(\hat{\Upsilon}_{T,D})(\hat{\theta}_{T,D} - \theta) = \sqrt{T}\hat{b}^\top(\hat{\Upsilon}_{T,D})(\hat{\theta}_{T,D} - \theta) - \sqrt{T}\hat{b}^\top 1 T \frac{\partial \ell_{T,D}(\hat{\theta}_{T,D}, \bar{y})}{\partial \theta^\top} = \sqrt{T}b^\top(\hat{\Upsilon}_{T,D})(\hat{\theta}_{T,D} - \theta) - \sqrt{T}b^\top \Upsilon_D(\hat{\theta}_{T,D} - \theta) + o_p(1)
\]
where the second equality is due to Assumption 4.1 and the fact that $\hat{\theta}_{T,D}$ is $C_2\sqrt{n\log n/T}$-consistent. Defining $a^\top := b^\top(-\hat{\Upsilon}_{T,D})$, we write
\[
\sqrt{T}a^\top\parallel a \parallel_{\ell_2}(\hat{\theta}_{T,D} - \theta) = \sqrt{T}a^\top\hat{\Upsilon}_{T,D}^{-1}(\hat{\Upsilon}_{T,D} - \Upsilon)(\hat{\theta}_{T,D} - \theta) - a^\top\parallel a \parallel_{\ell_2}\hat{\Upsilon}_{T,D}^{-1}1 T \frac{\partial \ell_{T,D}(\theta, \bar{y})}{\partial \theta^\top} + o_p(1/\parallel a \parallel_{\ell_2}).
\]
By recognising that $\parallel a^\top \parallel_2 = \parallel b^\top \hat{\Upsilon}_{T,D} \parallel_2 \geq \text{mineval}(\hat{\Upsilon}_{T,D})$, we have
\[
\frac{1}{\parallel a \parallel_2} = O_p\left(\frac{\varpi}{n}\right).
\]
Thus without loss of generality, we have, for any $c \in \mathbb{R}^s$ with $\parallel c \parallel_2 = 1$,
\[
\sqrt{T}c^\top(\hat{\theta}_{T,D} - \theta) = \sqrt{T}c^\top\hat{\Upsilon}_{T,D}^{-1}(\Upsilon_{T,D} - \Upsilon)(\hat{\theta}_{T,D} - \theta) - c^\top\hat{\Upsilon}_{T,D}^{-1}1 T \frac{\partial \ell_{T,D}(\theta, \bar{y})}{\partial \theta^\top} + o_p(\varpi/n).
\]
We now determine a rate for the first term on the right side. This is straightforward
\[
\sqrt{T}\parallel c \parallel_{\ell_2}\parallel \hat{\Upsilon}_{T,D}^{-1}(\Upsilon_{T,D} - \Upsilon)(\hat{\theta}_{T,D} - \theta) \parallel_{\ell_2} \leq \sqrt{T}\parallel c \parallel_{\ell_2}\parallel \hat{\Upsilon}_{T,D}^{-1}_{\ell_2}\parallel_{\ell_2}\parallel \hat{\Upsilon}_{T,D} - \Upsilon \parallel_{\ell_2}\parallel \hat{\theta}_{T,D} - \theta \parallel_{\ell_2}
\]
We now determine a rate for the first term on the right side. This is straightforward
\[
\sqrt{T}\parallel c \parallel_{\ell_2}\parallel \hat{\Upsilon}_{T,D}^{-1}(\Upsilon_{T,D} - \Upsilon)(\hat{\theta}_{T,D} - \theta) \parallel_{\ell_2} \leq \sqrt{T}\parallel c \parallel_{\ell_2}\parallel \hat{\Upsilon}_{T,D}^{-1}_{\ell_2}\parallel_{\ell_2}\parallel \hat{\Upsilon}_{T,D} - \Upsilon \parallel_{\ell_2}\parallel \hat{\theta}_{T,D} - \theta \parallel_{\ell_2}
\]
We now determine a rate for the first term on the right side. This is straightforward
\[
\sqrt{T}\parallel c \parallel_{\ell_2}\parallel \hat{\Upsilon}_{T,D}^{-1}(\Upsilon_{T,D} - \Upsilon)(\hat{\theta}_{T,D} - \theta) \parallel_{\ell_2} \leq \sqrt{T}\parallel c \parallel_{\ell_2}\parallel \hat{\Upsilon}_{T,D}^{-1}_{\ell_2}\parallel_{\ell_2}\parallel \hat{\Upsilon}_{T,D} - \Upsilon \parallel_{\ell_2}\parallel \hat{\theta}_{T,D} - \theta \parallel_{\ell_2}
\]
where the first equality is due to (8.31), Lemma 8.6(i) and the rate of convergence for the minimum distance estimator $\hat{\theta}_T (\hat{\theta}_{T,D})$. Thus

$$
\sqrt{T} c^\top (\hat{\theta}_{T,D} - \theta) = -c^\top \hat{\Upsilon}_{T,D}^{-1} T \frac{\partial \ell_{T,D}(\theta, \tilde{y})}{\partial \theta} + \text{rem}
$$

$$
\text{rem} = O_p \left( \sqrt{n^2 \log^2 n \kappa(W)} / T \right) + o_p(\kappa/n)
$$

whence, if we divide by $\sqrt{c^\top \Upsilon_{T,D}^{-1} c}$, we have

$$
\frac{\sqrt{T} c^\top (\hat{\theta}_{T,D} - \theta)}{\sqrt{c^\top \Upsilon_{T,D}^{-1} c}} = -c^\top \hat{\Upsilon}_{T,D}^{-1} T \frac{\partial \ell_{T,D}(\theta, \bar{y})}{\partial \theta} + \frac{\text{rem}}{\sqrt{c^\top \Upsilon_{T,D}^{-1} c}} =: t_{os,D,1} + t_{os,D,2}.
$$

Define

$$
t_{os,D,1} := -c^\top \Upsilon_{D}^{-1} T \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta} / T, \\
t_{os,D,2} := \frac{\text{rem}}{\sqrt{c^\top \Upsilon_{D}^{-1} c}}.
$$

To prove Theorem 4.2 it suffices to show $t_{os,D,1} \xrightarrow{d} N(0,1)$, $t_{os,D,2} = o_p(1)$, and $t_{os,D,2} = o_p(1)$.}

**8.5.1** $t_{os,D,1} \xrightarrow{d} N(0,1)$

We now prove that $t_{os,D,1}$ is asymptotically distributed as a standard normal. Write

$$
t_{os,D,1} := -c^\top \Upsilon_{D}^{-1} T \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta} / T = \sqrt{c^\top \Upsilon_{D}^{-1} c}
$$

$$
\sum_{t=1}^T \frac{1}{2} c^\top \Upsilon_{D}^{-1} E^\dagger D_n^\dagger \Psi_1(\Theta^{-1} \otimes \Theta^{-1}) (D^{-1/2} \otimes \tilde{D}^{1/2} T^{-1/2} \text{vec} [(y_t - \mu)(y_t - \mu)^\dagger - \bar{E}(y_t - \mu)(y_t - \mu)^\dagger]}
$$

$$
\frac{\sqrt{c^\top \Upsilon_{D}^{-1} c}}{\sqrt{c^\top \Upsilon_{D}^{-1} c}} =: \sum_{t=1}^T U_{os,D,T,n,t}.
$$

The proof is very similar to that of $t_{D,1} \xrightarrow{d} N(0,1)$ in Section A.4.1. It is straightforward to show that $\{U_{os,D,T,n,t}, F_{T,n,t}\}$ is a martingale difference sequence. We first investigate that at what rate the denominator $\sqrt{c^\top \Upsilon_{D}^{-1} c}$ goes to zero.

$$
c^\top \Upsilon_{D}^{-1} c = 2c^\top (E^\dagger D_n^\dagger \Xi D_n E)^{-1} c \geq 2\text{meval} \left( (E^\dagger D_n^\dagger \Xi D_n E)^{-1} \right) = \frac{2}{\text{meval}(E^\dagger D_n^\dagger \Xi D_n E)}.
$$

Since,

$$
\text{meval} \left( E^\dagger D_n^\dagger \Xi D_n E \right) \leq \text{meval}(\Xi) \text{meval}(D_n^\dagger D_n) \text{meval}(E^\dagger E) \leq Csn,
$$

for some positive constant $C$. Thus we have

$$
\frac{1}{\sqrt{c^\top \Upsilon_{D}^{-1} c}} = O(\sqrt{sn}). \quad (8.32)
$$
We now verify (i) and (ii) of Theorem A.5 in Appendix A.5. We consider $|U_{os,D,T,n,t}|$ first.

$$
|U_{os,D,T,n,t}| = \frac{1}{2} c^\top Y_D^{-1} E^\top D_n^1 \Psi_1 (\Theta^{-1} \otimes \Theta^{-1}) (D^{-1/2} \otimes D^{-1/2}) T^{-1/2} \vec{\|(y_t - \mu)(y_t - \mu)\top - E(y_t - \mu)(y_t - \mu)\top\|_2}
$$

$$
\leq \frac{1}{2} T^{-1/2} \left\| c^\top Y_D^{-1} E^\top D_n^1 \Psi_1 (\Theta^{-1} \otimes \Theta^{-1}) (D^{-1/2} \otimes D^{-1/2}) \right\|_2 \left\| \vec{\{(y_t - \mu)(y_t - \mu)^\top - E(y_t - \mu)(y_t - \mu)^\top\}} \right\|_2
$$

$$
= O \left( \sqrt{\frac{s^2 \omega^2}{T}} \right) \left\| (y_t - \mu)(y_t - \mu)^\top - E(y_t - \mu)(y_t - \mu)^\top \right\|_F
$$

$$
\leq O \left( \sqrt{\frac{n^2 s^2 \omega^2}{T}} \right) \left\| (y_t - \mu)(y_t - \mu)^\top - E(y_t - \mu)(y_t - \mu)^\top \right\|_\infty,
$$

where the last equality is due to

$$
\left\| c^\top Y_D^{-1} E^\top D_n^1 \Psi_1 (\Theta^{-1} \otimes \Theta^{-1}) (D^{-1/2} \otimes D^{-1/2}) \right\|_2 \leq \left\| Y_D^{-1} \right\|_{\ell_2} \left\| E^\top \right\|_{\ell_2} \left\| D_n^1 \right\|_{\ell_2} \left\| \Psi_1 \right\|_{\ell_2} \left\| \Theta^{-1} \otimes \Theta^{-1} \right\|_{\ell_2} \left\| D^{-1/2} \otimes D^{-1/2} \right\|_{\ell_2} = O \left( \frac{\omega}{\sqrt{n}} \right) \sqrt{s m} = O \left( \sqrt{\frac{s^2 \omega^2}{n}} \right)
$$

via (8.31) and (A.12). Next,

$$
\left\| \max_{1 \leq t \leq T} \left| U_{os,D,T,n,t} \right| \right\|_{\psi_1} \leq \log(1 + T) \max_{1 \leq t \leq T} \left| U_{os,D,T,n,t} \right|_{\psi_1}
$$

$$
= \log(1 + T) O \left( \sqrt{\frac{n^2 s^2 \omega^2}{T}} \right) \left\| \max_{1 \leq t \leq T} \left\| (y_t - \mu)(y_t - \mu)^\top - E(y_t - \mu)(y_t - \mu)^\top \right\|_\infty \right\|_{\psi_1}
$$

$$
= \log(1 + T) \log(1 + n^2) O \left( \sqrt{\frac{n^2 s^2 \omega^2}{T}} \right) \max_{1 \leq t \leq T} \left\| (y_t - \mu)(y_t - \mu)^\top - E(y_t - \mu)(y_t - \mu)^\top \right\|_\infty
$$

$$
= O \left( \frac{n^2 s^2 \omega^2 \log^2 (1 + T) \log^2 (1 + n^2)}{T} \right) = o(1)
$$

where the last equality is due to Assumption 3.3 iii). Since $\|U\|_{L_r} \leq r!\|U\|_{\psi_1}$ for any random variable $U$ (van der Vaart and Wellner 1996, p95), we conclude that (i) and (ii) of Theorem A.5 in Appendix A.5 are satisfied.

We now verify condition (iii) of Theorem A.5 in Appendix A.5. Since we have already shown that $\text{sci} Y_D^{-1} c$ is bounded away from zero by an absolute constant, it suffices to show

$$
\text{sn} \left| \frac{1}{T} \sum_{t=1}^{T} \left( \frac{1}{2} c^\top Y_D^{-1} E^\top D_n^1 \Psi_1 (\Theta^{-1} \otimes \Theta^{-1}) (D^{-1/2} \otimes D^{-1/2}) u_t \right)^2 - c^\top Y_D^{-1} c \right| = o_p(1),
$$

where $u_t := \vec{\{(y_t - \mu)(y_t - \mu)^\top - E(y_t - \mu)(y_t - \mu)^\top\}}$. Recognise that

$$
c^\top Y_D^{-1} c = c^\top Y_D^{-1} Y_D Y_D^{-1} c = c^\top Y_D^{-1} \left( \frac{1}{2} E^\top D_n^1 \Psi_1 (\Theta^{-1} \otimes \Theta^{-1}) \right) Y_D^{-1} c
$$

$$
= c^\top Y_D^{-1} \left( \frac{1}{2} E^\top D_n^1 \Psi_1 (\Theta^{-1} \otimes \Theta^{-1})(\Theta \otimes \Theta)(\Theta^{-1} \otimes \Theta^{-1}) \right) Y_D^{-1} c
$$

$$
= c^\top Y_D^{-1} \left( \frac{1}{2} E^\top D_n^1 \Psi_1 (\Theta^{-1} \otimes \Theta^{-1})(\Theta^{-1} \otimes \Theta^{-1})(\Theta \otimes \Theta)(\Theta^{-1} \otimes \Theta^{-1}) D_n E \right) Y_D^{-1} c
$$

$$
= \frac{1}{4} c^\top Y_D^{-1} E^\top D_n^1 \Psi_1 (\Theta^{-1} \otimes \Theta^{-1})(\Theta^{-1} \otimes \Theta^{-1})(\Theta \otimes \Theta)(\Theta^{-1} \otimes \Theta^{-1}) D_n E Y_D^{-1} c
$$

$$
= \frac{1}{4} c^\top Y_D^{-1} E^\top D_n^1 \Psi_1 (\Theta^{-1} \otimes \Theta^{-1})(\Theta^{-1} \otimes \Theta^{-1}) V(D^{-1/2} \otimes D^{-1/2})(\Theta^{-1} \otimes \Theta^{-1}) D_n E Y_D^{-1} c
$$

$$
= \frac{1}{4} c^\top Y_D^{-1} E^\top D_n^1 \Psi_1 (\Theta^{-1} \otimes \Theta^{-1})(\Theta^{-1} \otimes \Theta^{-1}) V(D^{-1/2} \otimes D^{-1/2})(\Theta^{-1} \otimes \Theta^{-1}) D_n E Y_D^{-1} c,
$$

20
where the second last equality is due to Lemma 11 of [Magnus and Neudecker (1986)] and the last equality is due to Assumption 3.1(ii). Thus

\[
\begin{align*}
sn & \left| \frac{1}{T} \sum_{t=1}^{T} \left( \frac{1}{2} c^\top \hat{\varphi}_{t,D}^{-1} E^\top D_n \Psi_1(\Theta^{-1} \otimes \Theta^{-1})(D^{-1/2} \otimes D^{-1/2}) u_t \right)^2 - c^\top \varphi_{t,D}^{-1} c \right| \\
& \leq sn \left\| \frac{1}{T} \sum_{t=1}^{T} u_t u_t^\top - V \right\|_\infty \left\| (D^{-1/2} \otimes D^{-1/2})(\Theta^{-1} \otimes \Theta^{-1}) \Psi_1 D_n E \varphi_{t,D}^{-1} c \right\|_2^2 \\
& \leq sn^3 \left\| \frac{1}{T} \sum_{t=1}^{T} u_t u_t^\top - V \right\|_\infty \left\| D^{-1/2} \otimes D^{-1/2} \left\| \Theta^{-1} \otimes \Theta^{-1} \left\| \Psi_1 \right\|_2 \left\| D_n \right\|_2 \left\| E \right\|_2 \left\| \varphi_{t,D}^{-1} c \right\|_2^2 \right. \\
& \left. \leq sn^3 \left( \frac{1}{T} \sum_{t=1}^{T} \left\| u_t u_t^\top - V \right\|_\infty \left\| \theta^{-1} \otimes \theta^{-1} \left\| \Psi_1 \right\|_2 \left\| D_n \right\|_2 \left\| E \right\|_2 \left\| \varphi_{t,D}^{-1} c \right\|_2 \right. \right) \\
& = O_p(sn^3) \sqrt{\frac{\log n}{T}} \cdot sn \cdot \omega^2 = O_p \left( \sqrt{\frac{n^4 \cdot \log n \cdot \omega^4 \cdot \log^4 n}{T}} \right) = o_p(1)
\end{align*}
\]

where the first equality is due to (8.31), (A.12) and the proof of Lemma 8.2, and the last equality is due to Assumption 3.3(iii).

8.5.2 \( \hat{t}_{os,D,1} - t_{os,D,1} = o_p(1) \)

We now show that \( \hat{t}_{os,D,1} - t_{os,D,1} = o_p(1) \). Let \( A_{os,D} \) and \( \hat{A}_{os,D} \) denote the numerators of \( t_{os,D,1} \) and \( \hat{t}_{os,D,1} \), respectively.

\[
\hat{t}_{os,D,1} - t_{os,D,1} = \frac{\hat{A}_{os,D}}{\sqrt{c^\top \varphi_{t,D}^{-1} c}} - \frac{A_{os,D}}{\sqrt{c^\top \varphi_{t,D}^{-1} c}} = \frac{\sqrt{sn} \hat{A}_{os,D}}{\sqrt{sn} c^\top \varphi_{t,D}^{-1} c} - \frac{\sqrt{sn} A_{os,D}}{\sqrt{sn} c^\top \varphi_{t,D}^{-1} c}
\]

Since we have already shown in (8.32) that snc \( \varphi_{t,D}^{-1} c \) is bounded away from zero by an absolute constant, it suffices to show the denominators as well as numerators of \( \hat{t}_{os,D,1} \) and \( t_{os,D,1} \) are asymptotically equivalent.

8.5.3 Denominators of \( \hat{t}_{os,D,1} \) and \( t_{os,D,1} \)

We need to show

\[
\text{sn} \left| c^\top (\hat{\varphi}_{t,D}^{-1} - \varphi_{t,D}^{-1}) c \right| = o_p(1).
\]

This is straightforward.

\[
\text{sn} \left| c^\top (\hat{\varphi}_{t,D}^{-1} - \varphi_{t,D}^{-1}) c \right| \leq \text{sn} \left\| (\hat{\varphi}_{t,D}^{-1} - \varphi_{t,D}^{-1}) \right\|_2 = \text{sn} \left( s \omega^2 \sqrt{\frac{1}{nT}} \right) = O_p \left( s^2 \omega^2 \sqrt{\frac{n}{T}} \right) = o_p(1),
\]

where the last equality is due to Assumption 3.3(iii).

8.5.4 Numerators of \( \hat{t}_{os,D,1} \) and \( t_{os,D,1} \)

We now show

\[
\sqrt{sn} \left| c^\top \hat{\varphi}_{t,D}^{-1} \sqrt{T} \frac{\partial \ell_{T,D}(\theta, \bar{y})}{\partial \theta} / T - c^\top \varphi_{t,D}^{-1} \sqrt{T} \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta} / T \right| = o_p(1).
\]
Using triangular inequality, we have
\[
\sqrt{s_n} \left| c^T \hat{T}^1_{T,D} \sqrt{T} \frac{\partial \ell_{T,D}(\hat{\theta}, \hat{\mu})}{\partial \theta^T} / T - c^T Y^1_{D} \sqrt{T} \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta^T} / T \right| \\
\leq \sqrt{s_n} \left| c^T \hat{T}^1_{T,D} \sqrt{T} \frac{\partial \ell_{T,D}(\hat{\theta}, \hat{\mu})}{\partial \theta^T} / T \right| \\
+ \sqrt{\frac{s_n}{n}} \left| c^T Y^1_{D} \sqrt{T} \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta^T} / T \right| (8.33)
\]

We first show that the first term of (8.33) is \(o_p(1)\).

\[
\sqrt{s_n} \left| c^T \hat{T}^1_{T,D} - Y^1_{D} \right| \sqrt{T} \left| \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta^T} / T \right| \\
= \sqrt{s_n} \left| c^T \hat{T}^1_{T,D} - Y^1_{D} \right| \sqrt{T} \left| \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta^T} \right| / T \\
\leq \sqrt{s_n} \left| \hat{T}^1_{T,D} - Y^1_{D} \right| \sqrt{T} \left| \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta^T} \right| / T \\
\leq \sqrt{s_n} \frac{\Sigma}{\sqrt{n/T}} \sqrt{T} \left| \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta^T} \right| / T \\
= O_p \left( \frac{\log n}{\sqrt{T}} \right) = o_p(1),
\]

where the last equality is due to Assumption 3.3 iii).

We now show that the second term of (8.33) is \(o_p(1)\).

\[
\sqrt{s_n} \left| c^T Y^1_{D} \sqrt{T} \left( \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta^T} / T - \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta^T} / T \right) \right| \\
= \sqrt{s_n} \left| c^T Y^1_{D} \sqrt{T} \left( \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta^T} / T \right) \right| \\
= O(\sqrt{s_n}) \left| \frac{\Sigma}{\sqrt{n/T}} \sqrt{T} \left| \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta^T} \right| / T \\
= O_p \left( \frac{\log n}{\sqrt{T}} \right) = o_p(1),
\]

where the third last equality is due to (8.21), and the last equality is due to Assumption 3.3 iii).

8.5.5 \(t_{os, D, 2} = o_p(1)\)

To prove \(t_{os, D, 2} = o_p(1)\), it suffices to show that \(\sqrt{s_n} \text{rem} = o_p(1)\). This is delivered by Assumption 3.3 iii). \(\square\)

8.6 Proof of Theorem 3.4 and Corollary 3.3

In this subsection, we give proofs of Theorem 3.4 and Corollary 3.3.

**Proof of Theorem 3.4** We only give a proof for part (i), as that for part (ii) is similar. Note that under \(H_0\),
\[
\sqrt{T} g_{T,D}(\theta) = \sqrt{T} [\text{vech}(\log \hat{\Theta}_{T,D}) - E\theta] = \sqrt{T} [\text{vech}(\log \check{\Theta}_{T,D}) - \text{vech}(\log \Theta)] \\
= \sqrt{T} D^1_n \text{vec}(\log \check{\Theta}_{T,D} - \log \Theta).
\]

Thus we can adopt the same method as in Theorem 3.2 to establish the asymptotic distribution of \(\sqrt{T} g_{T,D}(\theta)\). In fact, it will be much simpler here because we fixed \(n\). We should have
\[
\sqrt{T} g_{T,D}(\theta) \overset{d}{\to} N(0, S), \quad S := D^1_n H(D^{-1/2} \otimes D^{-1/2}) V(D^{-1/2} \otimes D^{-1/2}) HD^1_n, \quad (8.34)
\]
where \( S \) is positive definite given the assumptions of this theorem. The closed-form solution for \( \hat{\theta}_{T,D} \) has been given in (3.3), but this is not important. We only need that \( \hat{\theta}_{T,D} \) sets the first derivative of the objective function to zero:

\[
E^T W g_{T,D}(\hat{\theta}_{T,D}) = 0. \tag{8.35}
\]

Notice that

\[
g_{T,D}(\hat{\theta}_{T,D}) - g_{T,D}(\theta) = -E(\hat{\theta}_{T,D} - \theta). \tag{8.36}
\]

Pre-multiply (8.36) by \( \frac{\partial g_{T,D}(\hat{\theta}_{T,D})}{\partial \theta} \) to give

\[
-E^T W [g_{T,D}(\hat{\theta}_{T,D}) - g_{T,D}(\theta)] = E^T W (\hat{\theta}_{T,D} - \theta),
\]

whence we obtain

\[
\hat{\theta}_{T,D} - \theta = -(E^T W E)^{-1} E^T W [g_{T,D}(\hat{\theta}_{T,D}) - g_{T,D}(\theta)]. \tag{8.37}
\]

Substitute (8.37) into (8.36)

\[
\sqrt{T} g_{T,D}(\hat{\theta}_{T,D}) = \left[ I_{n(n+1)/2} - E (E^T W E)^{-1} E^T W \right] \sqrt{T} g_{T,D}(\theta) + E (E^T W E)^{-1} \sqrt{T} E^T W g_{T,D}(\hat{\theta}_{T,D})
\]

\[
= \left[ I_{n(n+1)/2} - E (E^T W E)^{-1} E^T W \right] \sqrt{T} g_{T,D}(\theta),
\]

where the second equality is due to (8.35). Using (8.34), we have

\[
\sqrt{T} g_{T,D}(\hat{\theta}_{T,D}) \xrightarrow{d} N \left( 0, I_{n(n+1)/2} - E (E^T W E)^{-1} E^T W \right) S \left( I_{n(n+1)/2} - E (E^T W E)^{-1} E^T W \right)^T.
\]

Now choosing \( W = S^{-1} \), we can simplify the asymptotic covariance matrix in the preceding display to

\[
S^{1/2} \left( I_{n(n+1)/2} - S^{-1/2} E (E^T S^{-1} E)^{-1} E^T S^{-1/2} \right) S^{1/2}.
\]

Thus

\[
\sqrt{T} \hat{S}_{T,D}^{-1/2} g_{T,D}(\hat{\theta}_{T,D}) \xrightarrow{d} N \left( 0, I_{n(n+1)/2} - S^{-1/2} E (E^T S^{-1} E)^{-1} E^T S^{-1/2} \right),
\]

because \( \hat{S}_{T,D} \) is a consistent estimate of \( S \) given (A.7) and Lemma 8.2, which hold under the assumptions of this theorem. The asymptotic covariance matrix in the preceding display is idempotent and has rank \( n(n+1)/2 - s \). Thus, under \( H_0 \),

\[
T g_{T,D}(\hat{\theta}_{T,D})^T \hat{S}_{T,D}^{-1/2} g_{T,D}(\hat{\theta}_{T,D}) \xrightarrow{d} \chi^2_{n(n+1)/2-s}.
\]

\[
\square
\]

Proof of Corollary 3.3. We only give a proof for part (i), as that for part (ii) is similar. From (3.5) and the Slutsky lemma, we have for every fixed \( n \) (and hence \( v \) and \( s \))

\[
\frac{T g_{T,D}(\hat{\theta}_{T,D})^T \hat{S}_{T,D}^{-1/2} g_{T,D}(\hat{\theta}_{T,D}) - \left( \frac{n(n+1)}{2} - s \right)}{n(n+1) - 2s} \xrightarrow{d} \chi^2_{n(n+1)/2-s} - \left( \frac{n(n+1)}{2} - s \right),
\]

as \( T \to \infty \). Then invoke Lemma A.18 in Appendix A.5

\[
\frac{\chi^2_{n(n+1)/2-s} - \left( \frac{n(n+1)}{2} - s \right)}{n(n+1) - 2s} \xrightarrow{d} N(0,1),
\]

as \( n \to \infty \) under \( H_0 \). Next invoke Lemma A.19 in Appendix A.5 there exists a sequence \( n = n_T \)

such that

\[
\frac{T g_{T,n,D}(\hat{\theta}_{T,n,D})^T \hat{S}_{T,n,D}^{-1/2} g_{T,n,D}(\hat{\theta}_{T,n,D}) - \left( \frac{n(n+1)}{2} - s \right)}{n(n+1) - 2s} \xrightarrow{d} N(0,1), \quad \text{under } H_0
\]

as \( T \to \infty \).
8.7 Miscellaneous Results

This subsection contains miscellaneous results of the article.

*Proof of Corollary 3.1.* Theorem 3.2 and a result we proved before, namely,
\[ |c^T \hat{J}_{T,D} c - c^T J_D c| = |c^T \hat{J}_{T,D} c - c^T J_D c| = o_p \left( \frac{1}{snk(W)} \right), \tag{8.38} \]

imply
\[ \sqrt{T} c^T (\hat{\theta}_{T,D} - \theta^0) \xrightarrow{d} N(0, c^T J_D c). \tag{8.39} \]

Consider an arbitrary, non-zero vector \( b \in \mathbb{R}^k \). Then
\[ \left\| \frac{Ab}{\|Ab\|_2} \right\|_2 = 1, \]
so we can invoke (8.39) with \( c = Ab/\|Ab\|_2 \):
\[ \sqrt{T} \frac{1}{\|Ab\|_2} b^T A^T (\hat{\theta}_{T,D} - \theta^0) \xrightarrow{d} N \left( 0, \frac{b^T A^T J_D A}{\|Ab\|_2} \right), \]
which is equivalent to
\[ \sqrt{T} b^T A^T (\hat{\theta}_{T,D} - \theta^0) \xrightarrow{d} N \left( 0, b^T A^T J_D A \right). \]

Since \( b \in \mathbb{R}^k \) is non-zero and arbitrary, via the Cramer-Wold device, we have
\[ \sqrt{T} A^T (\hat{\theta}_{T,D} - \theta^0) \xrightarrow{d} N \left( 0, A^T J_D A \right). \]

Since we have shown in the mathematical display above (A.11) that \( J_D \) is positive definite and \( A \) has full-column rank, \( A^T J_D A \) is positive definite and its negative square root exists. Hence,
\[ \sqrt{T} (A^T J_D A)^{-1/2} A^T (\hat{\theta}_{T,D} - \theta^0) \xrightarrow{d} N \left( 0, I_k \right). \]

Next from (8.38),
\[ |b^T Bb| := |b^T A^T \hat{J}_{T,D} A b - b^T A^T J_D A b| = o_p \left( \frac{1}{snk(W)} \right) \|Ab\|_2 \leq o_p \left( \frac{1}{snk(W)} \right) \|A\|^2_{l_2} \|b\|^2_{l_2}. \]

By choosing \( b = e_j \) where \( e_j \) is a vector in \( \mathbb{R}^k \) with \( j \)th component being 1 and the rest of components being 0, we have for \( j = 1, \ldots, k \)
\[ |B_{jj}| \leq o_p \left( \frac{1}{snk(W)} \right) \|A\|^2_{l_2} = o_p(1), \]
where the equality is due to \( \|A\|_{l_2} = O(\sqrt{snk(W)}) \). By choosing \( b = e_{ij} \), where \( e_{ij} \) is a vector in \( \mathbb{R}^k \) with \( i \)th and \( j \)th components being \( 1/\sqrt{2} \) and the rest of components being 0, we have
\[ |B_{ij}/2 + B_{jj}/2 + B_{ij}| \leq o_p \left( \frac{1}{snk(W)} \right) \|A\|^2_{l_2} = o_p(1). \]

Then
\[ |B_{ij}| \leq |B_{ij} + B_{ii}/2 + B_{jj}/2| + |(B_{ii}/2 + B_{jj}/2)| = o_p(1). \]

Thus we proved
\[ B = A^T \hat{J}_{T,D} A - A^T J_D A = o_p(1), \]
because the dimension of the matrix \( B, k, \) is finite. By Slutsky’s lemma
\[ \sqrt{T} (A^T \hat{J}_{T,D} A)^{-1/2} A^T (\hat{\theta}_{T,D} - \theta^0) \xrightarrow{d} N \left( 0, I_k \right). \]

\( \square \)
Lemma 8.7. For any positive definite matrix $\Theta$,
\[
\left( \int_0^1 [t(\Theta - I) + I]^{-1} \otimes [t(\Theta - I) + I]^{-1}dt \right)^{-1} = \int_0^1 e^{t\log \Theta} \otimes e^{(1-t)\log \Theta} dt.
\]

Proof. (11.9) and (11.10) of Higham (2008) p272 give, respectively, that
\[
\text{vec } E = \int_0^1 e^{t\log \Theta} \otimes e^{(1-t)\log \Theta} dt \text{vec } L(\Theta, E),
\]
\[
\text{vec } L(\Theta, E) = \int_0^1 [t(\Theta - I) + I]^{-1} \otimes [t(\Theta - I) + I]^{-1}dt \text{vec } E.
\]
Substitute the preceding equation into the second last
\[
\text{vec } E = \int_0^1 e^{t\log \Theta} \otimes e^{(1-t)\log \Theta} dt \int_0^1 [t(\Theta - I) + I]^{-1} \otimes [t(\Theta - I) + I]^{-1}dt \text{vec } E.
\]
Since $E$ is arbitrary, the result follows. \(\square\)

Example 8.3. In the special case of normality, $V = 2D_n D_n^+ (\Sigma \otimes \Sigma)$ (Magnus and Neudecker [1986] Lemma 9). Then $c^T J_D c$ could be simplified into
\[
c^T J_D c =
2c^T (E^T WE)^{-1} E^T WD_n^+ H(D^{-1/2} \otimes D^{-1/2}) D_n D_n^+ (\Sigma \otimes \Sigma)(D^{-1/2} \otimes D^{-1/2}) HD_n^T WE (E^T WE)^{-1} c
\]
\[
= 2c^T (E^T WE)^{-1} E^T WD_n^+ H(D^{-1/2} \otimes D^{-1/2})(\Sigma \otimes \Sigma)(D^{-1/2} \otimes D^{-1/2}) HD_n^T WE (E^T WE)^{-1} c
\]
\[
= 2c^T (E^T WE)^{-1} E^T WD_n^+ H(D^{-1/2} \Sigma D^{-1/2} \otimes D^{-1/2} \Sigma D^{-1/2}) HD_n^T WE (E^T WE)^{-1} c
\]
\[
= 2c^T (E^T WE)^{-1} E^T WD_n^+ H(\Theta \otimes \Theta) HD_n^T WE (E^T WE)^{-1} c,
\]
where the second equality is true because, given the structure of $H$, via Lemma 11 of Magnus and Neudecker (1986), we have the following identity:
\[
D_n^+ H(D^{-1/2} \otimes D^{-1/2}) = D_n^+ H(D^{-1/2} \otimes D^{-1/2}) D_n D_n^+.
\]

Lemma 8.8. Let $\frac{\partial^2 \ell_{T,D}(\theta, \mu)}{\partial \theta \partial \mu}$ and $\frac{\partial^2 \ell_{T,D}(\theta, \mu)}{\partial \theta \partial \mu}$ be given as in Theorem 4.1. Then we have under normality (i.e., $V = 2D_n D_n^+ (\Sigma \otimes \Sigma)$)
\[
\mathbb{E} \left[ -\frac{1}{T} \frac{\partial^2 \ell_{T,D}(\theta, \mu)}{\partial \theta \partial \mu} \right] = \mathbb{E} \left[ \frac{1}{T} \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \theta} \frac{\partial \ell_{T,D}(\theta, \mu)}{\partial \mu} \right].
\]

Proof. We show in (4.4) that
\[
\mathbb{E} \left[ -\frac{1}{T} \frac{\partial^2 \ell_{T,D}(\theta, \mu)}{\partial \theta \partial \mu} \right] = \frac{1}{2} E^T D_n^T \mathcal{V}_1 (e^{-\Omega} \otimes e^{-\Omega}) \Psi_1 D_n E.
\]
Now using the expression for $\frac{\partial \ell_{T,D}(\theta,\mu)}{\partial \theta}$ and the fact that it has zero expectation, we have

$$
\mathbb{E} \left[ \frac{1}{T} \frac{\partial \ell_{T,D}(\theta,\mu)}{\partial \theta} \right] = \frac{T}{4} \mathbb{E} \mathbf{D}^\top \mathbf{n} \Psi_1 \text{var} \left( \text{vec} \left( e^{-\Omega} \mathbf{D}^{-1/2} \Sigma_T \mathbf{D}^{-1/2} e^{-\Omega} \right) \right) \Psi_1 \mathbf{D} \mathbb{E}
$$

$$
= \frac{T}{4} \mathbb{E} \mathbf{D}^\top \mathbf{n} \Psi_1 \left( e^{-\Omega} \otimes e^{-\Omega} \right) (D^{-1/2} \otimes D^{-1/2}) \text{var} \left( \text{vec} \left[ \frac{1}{T} \sum_{t=1}^T (y_t - \mu)(y_t - \mu)^\top \right] \right)
$$

$$
\cdot (D^{-1/2} \otimes D^{-1/2}) (e^{-\Omega} \otimes e^{-\Omega}) \Psi_1 \mathbf{D} \mathbb{E}
$$

$$
= \frac{1}{4} \mathbb{E} \mathbf{D}^\top \mathbf{n} \Psi_1 \left( e^{-\Omega} \otimes e^{-\Omega} \right) (D^{-1/2} \otimes D^{-1/2}) \text{var} \left( \text{vec} \left[ (y_t - \mu)(y_t - \mu)^\top \right] \right)
$$

$$
\cdot (D^{-1/2} \otimes D^{-1/2}) (e^{-\Omega} \otimes e^{-\Omega}) \Psi_1 \mathbf{D} \mathbb{E}
$$

$$
= \frac{1}{4} \mathbb{E} \mathbf{D}^\top \mathbf{n} \Psi_1 \left( e^{-\Omega} \otimes e^{-\Omega} \right) (D^{-1/2} \otimes D^{-1/2}) 2D_n \mathbf{D}^+ (\Sigma \otimes \Sigma)(D^{-1/2} \otimes D^{-1/2}) (e^{-\Omega} \otimes e^{-\Omega}) \Psi_1 \mathbf{D} \mathbb{E}
$$

$$
= \frac{1}{2} \mathbb{E} \mathbf{D}^\top \mathbf{n} \Psi_1 \left( e^{-\Omega} \otimes e^{-\Omega} \right) (D^{-1/2} \otimes D^{-1/2}) D_n \mathbf{D}^+ (\Sigma \otimes \Sigma)(D^{-1/2} \otimes D^{-1/2}) (e^{-\Omega} \otimes e^{-\Omega}) \Psi_1 \mathbf{D} \mathbb{E}
$$

where the third equality is due to weak stationarity of $y_t$ and (A.10) via Assumption 3.5. There seems to be a caveat. □

References


——— (2007): Matrix Differential Calculus with Applications in Statistics and Econometrics, John Wiley and Sons Ltd.