On Capacity of Large-Scale MIMO Multiple Access Channels with Distributed Sets of Correlated Antennas *

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Abstract

In this paper, a deterministic equivalent of ergodic sum rate and an algorithm for evaluating the capacity-achieving input covariance matrices for the uplink large-scale multiple-input multiple-output (MIMO) antenna channels are proposed. We consider a large-scale MIMO system consisting of multiple users and one base station with several distributed antenna sets. Each link between a user and an antenna set forms a two-sided spatially correlated MIMO channel with line-of-sight (LOS) components. Our derivations are based on novel techniques from large dimensional random matrix theory (RMT) under the assumption that the numbers of antennas at the terminals approach to infinity with a fixed ratio. The deterministic equivalent results (the deterministic equivalent of ergodic sum rate and the capacity-achieving input covariance matrices) are easy to compute and shown to be accurate for realistic system dimensions. In addition, they are shown to be invariant to several types of fading distribution.

Index Terms—Deterministic equivalent, large dimensional RMT, large-scale MIMO, Stieltjes transform.

1 Introduction

To achieve higher rates, much efforts have been put to improving the spectral efficiency and data throughput of wireless communication systems. The multi-antenna technology is one key technology for wireless communication and is envisaged to be adopted ubiquitously. With the number of antennas at the base stations (BSs) and user equipments (UEs) being increased, communications systems will have better rate and link reliability [1, 2]. However, the actual achievable spectral efficiency could be greatly compromised by interference arising from simultaneous communications in neighboring areas.

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A promising solution to interference management is the large-scale multiple-input multiple-output (MIMO) technology, e.g., [3–7]. Figure 1 illustrates a possible scenario where the antenna array of a BS is composed of multiple geographically distributed low-power antenna sets, installed onto a ring of high-speed fibre-bus, and this BS is communicating with several multi-antenna UEs. The large-scale MIMO setting is beneficial not only in terms of communication performances (such as better coverage and efficient radio resource utilization) but also in terms of energy-saving.¹ In this complex system model, a number of practical factors such as correlation effects and line-of-sight (LOS) components need to be included, which occur due to the space limitation of UEs and the densification of the antenna arrays resulting in a visible propagation path from the UEs, respectively. For typical systems of tens of distributed antenna sets and hundreds of UEs, even computer simulations become challenging [9], which makes performance analysis of such large-scale MIMO systems an important and a new subject of research.

When a system is large, exact performance analysis is no longer suitable because an exact analytical expression would be too complex to appreciate. Hence, alternatives have emerged and the large dimensional random matrix theory (RMT) [6, 10–19] provides a powerful tool in dealing with large-scale MIMO systems. Utilizing the large dimensional RMT, this paper aims to derive information-theoretic results of the large-scale MIMO systems. In particular, our focus is on the uplink large-scale MIMO systems consisting of $K$ UEs and a BS with $L$ distributed sets of multiple antennas. Let $n_k$ and $N_l$ denote, respectively, the numbers of antennas at the $k$-th UE and the $l$-th antenna set of the BS receiver. The channel between the $k$-th UE and the $l$-th antenna set is modeled as the $N_l \times n_k$ complex matrix $H_{l,k} = R_{l,k}^{1/2} X_{l,k} T_{l,k}^{1/2} + \bar{H}_{l,k}$, where $X_{l,k}$’s are statistically independent random matrices of independent and identically distributed (i.i.d.) entries (but not necessary Gaussian), $\bar{H}_{l,k}$ is a deterministic matrix reflecting the LOS components of the

¹Using the setting, the number of BS can be greatly reduced. Note that the energy consumption for air conditioning for each BS is consuming up to 20,000 kWh each year on average which is sometimes higher than other equipments in a BS [8].

²Despite the Rayleigh or Rician distribution being the most popular distributions for small-scale amplitude fading, there
channel, and $R_{l,k}$ and $T_{l,k}$, respectively, characterize the spatial correlation structures at the receiver and transmitter sides separatively. Since the signals from multiple antenna sets are collected into a BS, the corresponding channel matrix of UE $k$ can be expressed as $H_k \triangleq [H^T_{T,1,k} \cdots H^T_{T,L,k}]^T$. An important objective of this study is to obtain a deterministic equivalent of the ergodic sum rate for the distributed uplink MIMO channel $\sum_{k=1}^K H_k H_k^H$ so that the system sum rate can be efficiently and accurately computed.

Although there have been quite many such results on MIMO capacity analysis utilizing large dimensional RMT [10, 12–16], the general model studied in this paper has not been addressed. To appreciate the objective of this paper, it is important to understand the limitations of the existing results. First, previous works in the large-scale MIMO systems usually assumed $n_k = 1$ and $N_l = 1$ for all $k, l$. That is, the UEs have only one antenna each and the BS is equipped with completely distributed antennas (i.e., one antenna in each antenna set). The elements of this channel matrix merely reflect the path loss differences between the links. Regarding the channel model (the channel with a variance profile), the most relevant work is [12] (or [11, Theorem 3.8] without the LOS components). In [12], a deterministic equivalent of the mutual information\(^3\) was derived based on the Bai-and-Silverstein method [22] (or [17, Chapter 6.2.1]). In fact, the results of [12] can be easily extended to the case with $n_k \geq 1$ and $N_l \geq 1$ but those spatial correlation matrices $T_{l,k}$’s and $R_{l,k}$’s are required to be diagonal.

The deterministic approximations of [11, 12] have found many applications in various system optimization designs such as scheduling [9, 23], training length designs [24], cell planning [25], and many others [11, 17]. This is because the designs based on the deterministic approximations not only can provide an efficient computation method but also give insight into what the optimal strategies look like. However, inheriting from the limitations of [11, 12], these results do not allow the UEs or each antenna set of the BS to be equipped with multiple spatially correlated antennas. Because of the potential applicability of deterministic equivalent results to system designs, there is a strong desire to deriving new deterministic equivalents as those given in [12] for the general model of our interest. However, even for an extension to the one-sided spatially correlated case, there will be several obstacles when one intends to get the deterministic equivalent of mutual information by using the Bai-and-Silverstein method and alike.\(^4\)

To date, there are only very few results dealing with random matrix models where the entries are correlated across both rows and columns. Most studies only considered random matrices with independent

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\(^3\)Formally, it should be read as the mutual information between the input and output over the channel with a variance profile. In this paper, we often simply refer to it as “the mutual information” if no confusion would occur.

\(^4\)If $H_{l,k} = 0$ $\forall l,k$, a partial generalization is possible by the Bai-and-Silverstein method. Specifically, with minor modifications for the case in [16, 26], the asymptotic mutual information can be obtained for the case that $R_{l,k}$’s were permitted to be nonnegative definite, while $H_{l,k} = 0$ and $T_{l,k}$’s are diagonal. If $T_{l,k}$’s are generally nonnegative definite, difficulties arise.
complex Gaussian random variables and used the fact that the correlated Gaussian random matrix can be transformed to an uncorrelated one with non-identically distributed entries without changing the concerned objects (e.g., the eigenvalue distribution and the mutual information). For convenience, we will refer to this transformation as the *decorrelation procedure*. Because of the assumption of Gaussianity, the entries are in fact uncorrelated, and so the Bai-and-Silverstein method can be used. For the latest results using this trick, refer to, e.g., [16]. Unfortunately, the channel model of our interest (i.e., $H_k$) cannot be transformed to a Gaussian random matrix with uncorrelated columns even if $X_{l,k}$'s are assumed to be Gaussian. For this to be possible, it would require that $T_{1,k}, \ldots, T_{L,k}$ be simultaneously unitarily diagonalizable for every $k$. Clearly, this restriction in the model does not permit UEs to have multiple spatially correlated antennas, which is unrealistic and greatly limits the significance of the model.

If the entries of the random matrices are Gaussian, then an alternative method, known as the Gaussian method [27] (the integration by part formula and Poincaré-Nash inequality), is much more useful. In this context, Hachem *et al.* [13, 15] have succeeded in obtaining the deterministic equivalent of mutual information for Kronecker (or separately) correlated Rayleigh and Rician MIMO channels. Compared to the Bai-and-Silverstein method, the Gaussian method is only suited to random matrices with Gaussian entries. However, one may extend the results obtained for matrices with Gaussian entries to any random matrices with independent entries following two recent developments, the Lindeberg principle [28] and the interpolation trick [29]. For the latest results, see, e.g., [19], where the Lindeberg principle is applied.

Early analyses using the Gaussian method were only for the typical Kronecker MIMO channel [13, 15]. In that case, the correlated Gaussian random matrix was transformed into an uncorrelated one, and the decorrelation procedure was employed. As such, the Gaussian method was merely an alternative tool to study large dimensional random matrices. Its superiority in dealing with random matrices with correlated pattern is largely unexplored until most recently, Dupuy and Loubaton in [18] derived the deterministic equivalent of average mutual information for a frequency selective MIMO channel, in which the decorrelation procedure could not be applied. We believe that the Gaussian method can be useful to treat other random matrices with involved correlation. With the aid of the Lindeberg principle, one may further extend the results obtained for matrices with Gaussian entries to any random matrices. Following this approach, this paper combines the two techniques to get the deterministic equivalents for the concerned channel model.

In particular, we first use the Gaussian method to derive the deterministic equivalent of ergodic sum rate for the large-scale MIMO multiple access channel (MAC) when $X_{l,k}$’s are Gaussian distributed. Our
results are much more general and can cope with several complex applications. As a special case, this contribution complements the results of [18] by extending the analysis to the case with LOS components. This extension is non-trivial.\footnote{Using the Gaussian tools, the asymptotic mutual information expressions for Rayleigh fading Kronecker MIMO channels were first proved by [13]. Two years later, the authors in the same group generalized the results to Rician fading channels [15]. This in some ways reflects the difficulty of such extension even for the typical Kronecker MIMO channel.} Next, by the generalized Lindeberg principle [28, 30], we generalize the deterministic equivalent for random matrices with Gaussian entries to those with non-Gaussian entries. Simulation results reveal that even for systems with realistic system dimensions, the deterministic approximation of ergodic sum rate provides reliable estimates to those obtained by Monte-Carlo simulations. Then, we apply the approximation to design the input covariances that tend to maximize the ergodic sum rate of the large-scale MIMO MAC, and provide an iterative water-filling optimization algorithm when only the statistical CSI at the transmitter (precisely, $T_{l,k}$’s, $R_{l,k}$’s, and $H_{l,k}$’s) is available. Finally, we conduct several simulations to confirm the comparability between results by our approach and those by the true (but time-consuming) optimization procedure under several types of fading distribution.

Notations—We use uppercase and lowercase boldface letters to denote matrices and vectors, respectively. $I_N$ denotes an $N \times N$ identity matrix while an all-zero matrix is denoted by $0$, and an all-one matrix is denoted by $1$. The matrix inequality $\succeq$ shows the positive semi-definiteness. The superscripts $(\cdot)^H$, $(\cdot)^T$, and $(\cdot)^*$ represent the conjugate-transpose, transpose, and conjugate operations, respectively. Also, we use $E\{\cdot\}$ to denote expectation with respect to all random variables within the brackets; $\log(\cdot)$ is the natural logarithm; $\rho(\cdot)$ denotes the spectral radius (i.e., the largest absolute value of the eigenvalues) of a matrix. $\|\cdot\|$ represents the Euclidean norm of an input vector or the spectral norm of an input matrix, while $\|\cdot\|_F$ denotes the Frobenius norm of a matrix, and $\|\cdot\|_\infty$ represents the maximum row sum matrix norm. The complex number field is denoted by $\mathbb{C}$. For any matrix $A \in \mathbb{C}^{N \times n}$, we use $[A]_{l,k}$, $[A]_{l,k}$ or $A_{kl}$ to denote the $(l,k)$-th entry, and $a_k$ denotes the $k$-th entry of the column vector $a$. The operators $(\cdot)^{\frac{1}{2}}$, $(\cdot)^{-1}$, $\text{tr}(\cdot)$ and $\text{det}(\cdot)$ represent the matrix principal square root, inverse, trace and determinant, respectively. In addition, $\text{diag}(x)$ denotes a diagonal matrix with an input vector $x$ representing its diagonal elements.

2 Channel Model and Problem Statement

2.1 Uplink Large MIMO

As shown in Figure 1, we consider the large-scale MIMO MAC with $K$ UEs, labeled as UE$_1$, . . . , UE$_K$, which are equipped with $n_1$, . . . , $n_K$ antennas, respectively. The $K$ UEs transmit simultaneously to a central coordinator with $L$ distributed antenna sets, labeled as BS$_1$, . . . , BS$_L$, which are equipped with
\( N_1, \ldots, N_L \) antennas, respectively. In this paper, we use the Kronecker model to characterize the spatial correlation of the MIMO channel for each MIMO link so that the correlation at an antenna set and a UE is modeled separately, as in [31]. Specifically, the channel from \( UE_k \) to \( BS_l \), \( H_{l,k} \in \mathbb{C}^{N_l \times n_k} \), can be written as

\[
H_{l,k} = \tilde{H}_{l,k} + \bar{H}_{l,k} = R_{l,k}^{\frac{1}{2}} X_{l,k} T_{l,k}^{\frac{1}{2}} + \bar{H}_{l,k}, \tag{1}
\]

where \( R_{l,k} \in \mathbb{C}^{N_l \times N_l} \) and \( T_{l,k} \in \mathbb{C}^{n_k \times n_k} \) are deterministic nonnegative definite matrices, characterizing the spatial correlation of the received signals across the antenna elements of \( BS_l \) and that of the transmitted signals across the antenna elements of \( UE_k \), respectively; \( X_{l,k} \equiv [\frac{1}{\sqrt{n_k}} X_{ij}^{(l,k)}] \in \mathbb{C}^{N_l \times n_k} \) consists of the random components of the channel in which the elements are i.i.d. complex random variables with zero mean and unit variance; and \( \tilde{H}_{l,k} \in \mathbb{C}^{N_l \times n_k} \) is a deterministic matrix corresponding to the channel LOS.

With the channel given above, we define the Rician factor between \( UE_k \) and \( BS_l \) as

\[
\kappa_{l,k} = \frac{\|\bar{H}_{l,k}\|_F^2}{\mathbb{E}\{\|\tilde{H}_{l,k}\|_F^2\}}. \tag{2}
\]

We also denote the distance-dependent pathloss of the \((l,k)\)-th pair by \( g_{l,k} = \mathbb{E}\{\|H_{l,k}\|_F^2\}/N_l \) given by

\[
\mathbb{E}\{\|H_{l,k}\|_F^2\} = \frac{1}{n_k} \text{tr}(R_{l,k}) \text{tr}(T_{l,k}) + \text{tr}(\bar{H}_{l,k} \bar{H}_{l,k}^H). \tag{3}
\]

Following the standard conventions [14], \( R_{l,k} \), \( T_{l,k} \), and \( \bar{H}_{l,k} \) are normalized such that

\[
\begin{cases}
\text{tr}(R_{l,k}) = \frac{1}{\kappa_{l,k} + 1} g_{l,k} N_l, \\
\text{tr}(T_{l,k}) = n_k, \\
\text{tr}(\bar{H}_{l,k} \bar{H}_{l,k}^H) = \frac{\kappa_{l,k}}{\kappa_{l,k} + 1} g_{l,k} N_l.
\end{cases} \tag{4}
\]

It is noted that \( \kappa_{l,k} \) and \( g_{l,k} \) are independent from the matrix dimensions. Therefore, the normalization is valid for all possible correlation patterns and imposes no restriction on practical applications. Although for convenience purpose we will simply set the same noise level (i.e., \( \sigma^2 \)) at all the receivers, it imposes no restriction since one can adjust \( g_{l,k} \) to get an arbitrary signal-to-noise ratio (SNR) of the \((l,k)\)-th pair. In addition, the setting implies that the LOS components of some link pairs are allowed to be absent.
2.2 Problem Formulation

The sum rate has been a key metric for performance analysis of a MAC. We begin with the sum rate formulation of the large-scale MIMO system and then explain its relation to RMT. For ease of exposition, we define
\[ N \triangleq \sum_{l=1}^{L} N_l, \quad n \triangleq \sum_{k=1}^{K} n_k, \quad H_k \triangleq \begin{bmatrix} H_{1,k}^T & \cdots & H_{L,k}^T \end{bmatrix}^T \in \mathbb{C}^{N \times n_k}, \quad \bar{H}_k \triangleq \begin{bmatrix} \bar{H}_{1,k}^T & \cdots & \bar{H}_{L,k}^T \end{bmatrix}^T \in \mathbb{C}^{N \times n_k}, \]
\[ \mathbf{H} \triangleq [H_1 \cdots H_K] \in \mathbb{C}^{N \times n}, \quad \bar{\mathbf{H}} \triangleq [ar{H}_1 \cdots \bar{H}_K] \in \mathbb{C}^{N \times n}. \]
The channel \( H_k \) represents the joint channel between UE\( k \) and the \( L \) distributed antenna sets interconnected at the BS. Then, the ergodic sum rate of the MIMO MAC can be expressed as \[ V_{BN}(\sigma^2) \equiv \frac{1}{N} \mathbb{E} \left\{ \log \det \left( \mathbf{I}_N + \frac{1}{\sigma^2} \mathbf{B}_N \right) \right\} \] (5)
where \( \sigma^2 \) is the noise variance at the receivers and
\[ \mathbf{B}_N \triangleq \sum_{k=1}^{K} H_k \bar{H}_k^H \in \mathbb{C}^{N \times N}. \] (6)
Specifically, \( V_{BN}(\sigma^2) \) provides a performance metric regarding the total number of nats (or bits if in base 2 of logarithm) per antenna that can be transmitted reliably over the channel matrices \( \{H_k\}_{k=1}^{K} \).

The derivative of \( V_{BN}(\sigma^2) \) with respect to \( \sigma^2 \) is given by
\[ \frac{\partial V_{BN}(\sigma^2)}{\partial \sigma^2} = \frac{1}{N} \mathbb{E} \left\{ \text{tr} \left[ \left( \mathbf{I}_N + \frac{1}{\sigma^2} \mathbf{B}_N \right)^{-1} \right] \right\} - \frac{1}{\sigma^2}. \] (7)
By Fubini’s theorem, we have [12, page 891]
\[ V_{BN}(\sigma^2) = \int_{\sigma^2}^{\infty} \left( \frac{1}{\omega} - \mathbb{E} \{m_{\mathbf{B}_N}(\omega)\} \right) d\omega, \] (8)
where
\[ m_{\mathbf{B}_N}(\omega) \triangleq \frac{1}{N} \text{tr} (\mathbf{B}_N + \omega \mathbf{I}_N)^{-1}. \] (9)
In RMT, \( m_{\mathbf{B}_N} \) is referred to as the Stieltjes transform of \( \mathbf{B}_N \) at point \(-\omega\), which provides a convenient tool to study the behavior of large dimensional random matrices. The relationship by which the mutual information is expressed as a functional of the Stieltjes transform is called the Shannon transform [11, Section 2.2.3].

In this paper, we are interested in understanding the ergodic sum capacity of the MIMO MAC by using large dimensional RMT. In particular, we consider that \( L, K \) are fixed but \( N_1, \ldots, N_L, n_1, \ldots, n_K \) all go
to infinity with ratios \( \{ \beta_{l,k}(N) \equiv \frac{N_l}{n_k} \} \) such that

\[
0 < \min_{l,k} \liminf_{N} \beta_{l,k}(N) < \max_{l,k} \limsup_{N} \beta_{l,k}(N) < \infty. \tag{10}
\]

For convenience, we refer to this large dimensional regime simply as \( \mathcal{N} \to \infty \) in the sequel. To this end, in the next section, we first find a \textit{deterministic} matrix-valued function \( \Psi(\omega) \in \mathbb{C}^{N \times N} \) (to be done later) such that

\[
\mathbb{E} \{ m_{B_N}(\omega) \} - \frac{1}{N} \text{tr}(\Psi(\omega)) \xrightarrow{N \to \infty} 0 \quad \text{for } \omega \in \mathbb{R}^+.
\tag{11}
\]

Following [12] (or [17, Definition 6.1]), we refer to \( \frac{1}{N} \text{tr}(\Psi(\omega)) \) as the deterministic equivalent of \( \mathbb{E} \{ m_{B_N}(\omega) \} \).

To appreciate the contributions of this paper, it is worth emphasizing that \( H_k \), in general, cannot be written in the form (1) using the separable correlation model, because different antenna sets have different spatial correlations, and this is the main obstacle of this class of random matrices – otherwise, there are some existing results [6, 10–16, 18, 19]. Next, using the Shannon transform (8), we will find \( \mathcal{V}_{B_N}(\sigma^2) \) so that

\[
\mathbb{E} \{ \mathcal{V}_{B_N}(\sigma^2) \} - \mathcal{V}_{B_N}(\sigma^2) \to 0 \quad \text{as } \mathcal{N} \to \infty.
\]

Finally, we will use \( \mathcal{V}_{B_N}(\sigma^2) \) to obtain the optimal input covariance matrices that maximize the deterministic approximation of the ergodic sum rate.

3 Deterministic Equivalents and Ergodic Capacity

3.1 Deterministic Equivalents

We first state the assumptions imposed in our system model.

**Assumption 1** Let \( X_{l,k} \equiv \left[ \frac{1}{\sqrt{n_k}} X_{ij}^{(l,k)} \right] \in \mathbb{C}^{N_l \times n_k} \), where \( X_{ij}^{(l,k)} \)'s are i.i.d. complex random variables with independent real and imaginary parts such that

\[
\mathbb{E} \{ X_{11}^{(l,k)} \} = 0, \quad \text{and } \mathbb{E} \{ |X_{11}^{(l,k)}|^2 \} = 1, \tag{12}
\]

and have finite 6-th order moment.

**Assumption 2** The family of deterministic matrices \( \{ T_{l,k}, R_{l,k} \} \) is nonnegative definite. In addition, the spectral norms of \( R_{l,k}, T_{l,k}, \) and \( \bar{H}_{l,k} \bar{H}_{l,k}^H \) are bounded by a constant, i.e.,

\[
\max_{k,l} \max \{ \| R_{l,k} \|, \| T_{l,k} \|, \| \bar{H}_{l,k} \bar{H}_{l,k}^H \| \} \leq C_{\text{max}}. \tag{13}
\]
To facilitate our expressions, we define the notation \( \langle A \rangle_k \) that returns the submatrix of \( A \) obtained by extracting the elements of the rows and columns with indices from \( \sum_{i=1}^{k-1} n_i + 1 \) to \( \sum_{i=1}^{k} n_i \). Similarly, the notation \( \langle \langle A \rangle \rangle_l \) returns the submatrix of \( A \) obtained by extracting the elements of the rows and columns with indices from \( \sum_{j=1}^{l-1} N_j + 1 \) to \( \sum_{j=1}^{l} N_j \). Also, for convenience, in the paper, we often omit \( \omega \) when writing \( m_{B_N}, \Psi, \bar{\Psi}, \Phi, \bar{\Phi}, \Phi_l, \bar{\Phi}_k, e_{l,k}, \bar{e}_{l,k} \), and denote \( \sum_{l,k} \equiv \sum_{l=1}^{L} \sum_{k=1}^{K} \).

**Theorem 1**  
Let \( \beta_{l,k} = \frac{N_l}{n_k} \). Under Assumption 2, the deterministic system of the \( L \times K \) equations

\[
\begin{align*}
e_{l,k} &= \frac{1}{N_l} \text{tr}(R_{l,k} \langle \langle \Psi \rangle \rangle_l), \\
\bar{e}_{l,k} &= \frac{1}{n_k} \text{tr}(T_{l,k} \langle \bar{\Psi} \rangle_k),
\end{align*}
\]

for \( 1 \leq l \leq L \) and \( 1 \leq k \leq K \), where

\[
\begin{align*}
\Psi &= \left( \Phi^{-1} + \omega \bar{\Phi} \bar{\Phi}^H \right)^{-1}, \\
\bar{\Psi} &= \left( \bar{\Phi}^{-1} + \omega \bar{\Phi}^H \bar{\Phi} \right)^{-1}, \\
\Phi &= \text{diag}(\Phi_1, \ldots, \Phi_L), \\
\bar{\Phi} &= \text{diag}(\bar{\Phi}_1, \ldots, \bar{\Phi}_K), \\
\Phi_l &= \left( \omega I_{N_l} + \omega \sum_{k=1}^{K} \bar{e}_{l,k} R_{l,k} \right)^{-1}, \\
\bar{\Phi}_k &= \left( \omega I_{n_k} + \omega \sum_{l=1}^{L} \beta_{l,k} e_{l,k} T_{l,k} \right)^{-1}
\end{align*}
\]

have a unique solution for \( \omega \in \mathbb{R}^+ \).

Under Assumptions 1 and 2, as \( N \to \infty \), we then have

\[
\mathbb{E}\{m_{B_N}\} - \frac{1}{N} \text{tr}(\Psi) = O \left( \frac{1}{\sqrt{N}} \right), \quad \text{for } \omega \in \mathbb{R}^+.
\]

**Proof:** Here, for ease of understanding, we give an outline of the proof. Our strategy is to show that the deterministic equivalent of \( \mathbb{E}\{m_{B_N}\} \) [i.e. \( \frac{1}{N} \text{tr}(\Psi) \)] can be found for the Gaussian random matrices and...
then we prove that the result is also applied for the non-Gaussian distributions.

Let $\mathcal{B}_N$ be an $N \times N$ matrix obtained from $B_N$ in (6) with all $X_{l,k}$'s replaced by $\mathcal{X}_{l,k}$'s, where $\mathcal{X}_{l,k}$'s are matrices with entries being independent standard Gaussian. Using the Gaussian method [27] (the integration by part formula and Poincaré-Nash inequality), we can show that the error term $E \{ m_{B_N} \} - \frac{1}{N} \text{tr}(\Psi)$ is of order $O \left( \frac{1}{N^2} \right)$. The detailed derivation is given in Appendix A.

Next, applying the Lindeberg principle [30, Theorem 2], we prove that $E \{ m_{B_N} \} - E \{ m_{B_N} \} = O \left( \frac{1}{\sqrt{N}} \right)$. The detailed derivation using the Lindeberg principle is provided in Appendix B. Together with the result for the Gaussian case, the proof of (16) can be accomplished by noting that

$$E \{ m_{B_N} \} - \frac{1}{N} \text{tr} (\Psi) = \left( E \{ m_{B_N} \} - E \{ m_{B_N} \} \right) + \left( E \{ m_{B_N} \} - \frac{1}{N} \text{tr} (\Psi) \right).$$

$$= O \left( \frac{1}{\sqrt{N}} \right) = O \left( \frac{1}{N^2} \right).$$

Finally, we consider the existence and uniqueness of the solution to (14) in Appendix C.

**Remark 1** If $X_{ij}^{(l,k)}$'s are Gaussian, the assumption that $X_{ij}^{(l,k)}$'s have finite 6-th order moment is naturally satisfied. When the amplitudes of the channel fading coefficients follow the Nakagami and log-normal distributions, Theorem 1 is applicable since these distributions have finite 6-th order moment. In Appendix B, the proof of $E \{ m_{B_N} \} - \frac{1}{N} \text{tr} (\Psi) = O \left( \frac{1}{\sqrt{N}} \right)$ was given under the assumption that $X_{ij}^{(l,k)}$'s have finite 6-th order moment. In fact, with additional arguments, the more general case can be obtained. Specifically, if $X_{ij}^{(l,k)}$'s have only finite second moment, we can prove that $E \{ m_{B_N} (\omega) \} - \frac{1}{N} \text{tr} (\Psi (\omega)) = O (\varepsilon_n)$, where $\varepsilon_n$ is a positive sequence converging to zero. However, it should be noted that with the finite 6-th order moment assumption, the proof of $E \{ m_{B_N} \} - \frac{1}{N} \text{tr} (\Psi) = O \left( \frac{1}{\sqrt{N}} \right)$ is much simpler than the latter general case. The proof of the general case requires some additional truncation, centralization, and rescaling techniques together with some careful derivations as those in [19]. Since these are beyond the scope of this paper, we do not show the detail proof regarding this general case. Interested readers can refer to [19].

**Remark 2** Theorem 1 is developed under the asymptotic regime where $L, K$ are fixed but $\{ N_l, n_k \}$'s all grow to infinity with fixed ratios. For other applications, we might be interested in the cases with fixed $\{ N_l, n_k \}$'s while $L$ and $K$ grow to infinity. In this case, the entries of $X_{i,k}$’s will be normalized by $\sqrt{n}$ rather than $\sqrt{n_k}$ and a similar deterministic equivalent result as that of Theorem 1 can be obtained.\(^7\)

We then derive a deterministic equivalent of the ergodic sum rate of the large-scale MIMO MAC in the following theorem.

\(^7\)Only different in some scalar adjustment.
**Theorem 2** Assuming that $B_N$ follows the hypotheses of Theorem 1, as $N \to \infty$, the Shannon transform of $B_N$ satisfies

$$\mathbb{E}\{V_{B_N}(\sigma^2)\} - V_N(\sigma^2) = O\left(\frac{1}{\sqrt{N}}\right),$$

where

$$V_N(\sigma^2) = \frac{1}{N} \log \det \left(\frac{\Psi(\sigma^2)^{-1}}{\sigma^2}\right) + \frac{1}{N} \sum_{k=1}^K \log \det \left(\frac{\Phi_k(\sigma^2)^{-1}}{\sigma^2}\right) - \frac{\sigma^2}{N} \sum_{l,k} N_l e_{l,k}(\sigma^2) \tilde{e}_{l,k}(\sigma^2),$$

$$(19a)$$

and

$$= \frac{1}{N} \log \det \left(\frac{\tilde{\Psi}(\sigma^2)^{-1}}{\sigma^2}\right) + \frac{1}{N} \sum_{l=1}^L \log \det \left(\frac{\tilde{\Phi}_l(\sigma^2)^{-1}}{\sigma^2}\right) - \frac{\sigma^2}{N} \sum_{l,k} N_l e_{l,k}(\sigma^2) \tilde{e}_{l,k}(\sigma^2).$$

$$(19b)$$

Furthermore, if $X_{l,k}$'s are Gaussian, we have, as $N \to \infty$,

$$N \left(\mathbb{E}\{V_{B_N}(\sigma^2)\} - V_N(\sigma^2)\right) = O\left(\frac{1}{N}\right).$$

$$(20)$$

**Proof:** By (16) in Theorem 1 together with the dominated convergence theorem, (18) is obtained. Then, we show that $\int_{\sigma^2}^{\infty} \left(\frac{1}{\omega} - \frac{1}{\omega} \text{tr}(\Psi(\omega))\right) d\omega$ can be written more explicitly as (19a). The details of the proof are similar to those in [19, Theorem 3], and thus omitted. Since $\det(I + AB) = \det(I + BA)$, we then have (19b). On the other hand, (20) can be obtained by (17) in Theorem 1. \[\square\]

**Remark 3** With (18), we can get the deterministic equivalent of the ergodic sum rate regarding the number of nats per antenna. However, (20) shows the convergence regarding the total ergodic sum rate and as a consequence has a wider range of applications for the performance evaluation criteria.

Over the last few years, there have been quite many deterministic equivalent results obtained by using large dimensional RMT (e.g., [6, 10–16, 18, 19]). Since our model is fairly general, Theorem 2 may be interpreted as a unified formula that encompasses many such results. For the case with $K = 1$ and $\bar{H} = 0$, $V_N(\sigma^2)$ agrees with that in [18, Theorem 2], in which $\{X_{l,1}\}_{\forall l}$ are assumed to be Gaussian. Theorem 2 thus extends its application to the non-Gaussian scenarios in this sense. Indeed, if $\bar{H} = 0$, (19) was first presented in [33, (23)], where the replica method was used. Also, for the case with $K = 2$, $L = 1$, and $\{R_{l,k} = \mathbf{R}\}_{\forall k}$, Theorem 2 is consistent with the results in [34] by the replica method which is however mathematically incomplete. In contrast, Theorem 2 is not only mathematically rigorous but also more general than the proposition in [33] in the sense that $\bar{H} \neq 0$ and there is no requirement on the Gaussian distribution on the entries of $X_{l,k}$. Finally, if $n_k = 1$ and $N_l = 1$ for all $k, l$, then Theorem 2 degenerates to that in [12] (or [11] without the LOS components). Clearly, in contrast with [11, 12], Theorem 2 allows
the UEs and each antenna set of the BS to be equipped with multiple spatially correlated antennas.

As mentioned before, deterministic equivalent results together with optimization approaches have found numerous applications in system optimization designs [9, 23–25]. For example, based on the deterministic equivalent result of [11], the authors of [9] devised an algorithm to compute the ergodic sum rate subject to a general fairness criterion. Also, based on [11], the authors of [23] derived an analytical expression of a system spectral efficiency when multiple BSs employ joint transmission with linear zero-forcing beamforming. They also developed a downlink scheduling scheme under a fairness criterion. Our deterministic equivalent results provide a promising foundation to these applications while under the more general large-scale MIMO system. In addition, a deterministic equivalent for the SINR at the output of the MMSE receiver can be derived using our deterministic equivalent results. Due to space limitations, such applications through Theorems 1–2 are left out. In the next subsection, our aim is to answer one of the fundamental questions: How should the input covariances be designed so that the ergodic sum rate can be maximized?

### 3.2 Ergodic Capacity

It is well known that the ergodic sum capacity of a MIMO MAC is achieved by selecting proper input covariance matrices so that the ergodic sum rate is maximized [35]. In this subsection, we aim to design the optimal covariance matrices using the deterministic equivalent results. Firstly, we state that the covariance matrices maximizing the deterministic equivalent of the ergodic sum rate yield a result which converges to the ergodic capacity. After that, these optimal covariance matrices will be shown to be structurally equivalent to an iterative waterfilling procedure over a deterministic channel. Finally, we propose an iterative waterfilling algorithm for finding the capacity-achieving input covariance matrices.

Let $Q_k$ be the input covariance matrix of $UE_k$ which satisfies $\text{tr}(Q_k) \leq n_k$. With the input covariance matrices $Q \triangleq \text{diag}(Q_1, \ldots, Q_K)$, we thus write the ergodic sum rate of the large-scale MIMO MAC as

$$E\{V_{BN}(\sigma^2, Q_1, \ldots, Q_K)\} = \frac{1}{N}E\left\{\log \det \left( I_N + \frac{1}{\sigma^2}HQH^H \right) \right\}. \quad (21)$$

Then, the ergodic capacity under the power constraint is given by

$$\max_{Q_k \in \mathbb{Q}, \forall k} E\{V_{BN}(\sigma^2, Q_1, \ldots, Q_K)\}, \quad (22)$$

---

8The power constraint can be replaced by $\text{tr}(Q_k) \leq P_k n_k$ with $P_k$ being any finite positive value independent from the matrix dimension. Note that the current setting $\text{tr}(Q_k) \leq n_k$ is for notational brevity only.
where
\[
Q_k \triangleq \left\{ Q_k \mid \text{tr}(Q_k) \leq n_k \text{ and } Q_k \succeq 0 \right\}
\]
is the feasible set of \( Q_k \). The problem (22) is convex and can be solved using stochastic programming based on convex optimization with Monte-Carlo methods [36]. Specifically, we can apply the method in [37] (called the Vu-Paulraj algorithm), which was developed based on the barrier method [36, Chap. 11] where the related gradient and Hessian are approximated by Monte-Carlo methods. Since \( Q_k \) is a Hermitian matrix of size \( n_k \times n_k \), the optimization involves \( n_k \) real entries on the diagonal and \( n_k(n_k - 1)/2 \) complex entries in the upper triangle. The complexity of such an algorithm is high and requires long execution time.

We thus propose an approximate approach using the deterministic equivalent results in Theorem 2.

In Theorem 2, we have shown that the deterministic equivalent results are invariant to the type of fading distribution. As a result, the asymptotic optimal input covariances, which are designed based on the deterministic equivalent results, are also invariant to the type of fading distribution. To get the deterministic equivalent of \( \mathbb{E} \left\{ V_{B_N}(\sigma^2, Q_1, \ldots, Q_K) \right\} \), the effect of \( Q_k \) has to be included in \( V_{B_N}(\sigma^2) \). With Theorem 2, this can be easily accomplished by the following replacements: for \( 1 \leq l \leq L \),
\[
T_{l,k} := Q_k^{1/2} T_{l,k} Q_k^{1/2}, \quad \text{and} \quad \bar{H}_{l,k} := \bar{H}_{l,k} Q_k^{1/2}.
\]
Now, let \( V_N(\sigma^2, Q_1, \ldots, Q_K) \) be the result obtained from \( V_N(\sigma^2) \) with \( T_{l,k} \) and \( \bar{H}_{l,k} \) based on the above replacements. Then, (19b) becomes
\[
V_N(\sigma^2, Q_1, \ldots, Q_K) = \frac{1}{N} \log \det (I_n + FQ) + \frac{1}{N} \sum_{l=1}^L \log \det \left( \frac{\Phi_l(\sigma^2)^{-1}}{\sigma^2} \right) - \frac{\sigma^2}{N} \sum_{l,k} N_l e_{l,k}(\sigma^2) \bar{e}_{l,k}(\sigma^2),
\]
where
\[
F = \text{diag} \left( \left\{ \sum_{l=1}^L \beta_{l,k} e_{l,k}(\sigma^2) T_{l,k} \right\}_{\forall k} \right) + \bar{H}^H \Phi(\sigma^2) \bar{H}.
\]
Note that \( Q_k \)'s appear in \( \bar{e}_{l,k}(\sigma^2) \)'s, i.e., \( \bar{e}_{l,k}(\omega) = \frac{1}{\sqrt{\nu_k}} \text{tr}(Q_k^{1/2} T_{l,k} Q_k^{1/2} \langle \hat{\Psi}(\omega) \rangle_k) \) and thus are involved in all the three terms of (24). Using the deterministic equivalent result, we have the optimization problem:
\[
\max_{Q_k \in Q_k, \forall k} V_N(\sigma^2, Q_1, \ldots, Q_K).
\]

Before solving the above problem, two important issues must be resolved. One is to establish the concavity of \( V_N(\sigma^2, Q_1, \ldots, Q_K) \) with respect to \( (Q_1, \ldots, Q_K) \), and the other one is to ensure that
E{\mathcal{V}_N(\sigma^2, Q_1^0, \ldots, Q_K^0)} - \mathcal{V}_N(\sigma^2, Q_1^*, \ldots, Q_K^*) goes asymptotically to zero, where, \((Q_1^0, \ldots, Q_K^0)\) and \((Q_1^*, \ldots, Q_K^*)\) are the maximizers of (22) and (26), respectively. The required results are described by the following proposition.

**Proposition 1** We have:

1. The function \((Q_1, \ldots, Q_K) \mapsto \mathcal{V}_N(\sigma^2, Q_1, \ldots, Q_K)\) is strictly concave on \((Q_1, \ldots, Q_K)\).

2. In addition to Assumption 2, suppose that \(Q_k^0\)'s and \(Q_k^*\)'s lay within a set of positive semi-definite matrices with bounded spectral norm. The, we have

\[
E\{\mathcal{V}_N(\sigma^2, Q_1^0, \ldots, Q_K^0)\} - \mathcal{V}_N(\sigma^2, Q_1^*, \ldots, Q_K^*) = O\left(\frac{1}{\sqrt{N}}\right),
\]

Furthermore, if \(X_{l,k}\)'s are Gaussian, then (27) becomes \(O\left(\frac{1}{N^2}\right)\).

**Proof:** The proof is similar to that in [15, Theorem 4 and Proposition 3] and [18, Theorem 3 and Proposition 4], and therefore omitted.

So far, we have stated that \((Q_1^*, \ldots, Q_K^*)\) yield a result which converges to the ergodic capacity. Next, by using tools from convex optimization [36], we will gain a better understanding on the structure of \((Q_1^*, \ldots, Q_K^*)\). In particular, our next proposition will state that the optimal covariance matrices are structurally equivalent to an iterative waterfilling procedure over a deterministic equivalent channel.

To that end, we start with defining the Lagrangians of the optimization problem (26) as

\[
\mathcal{L}(Q, \Upsilon, \mu) = -\mathcal{V}_N(\sigma^2, Q_1, \ldots, Q_K) + \sum_{k=1}^{K} \text{tr}(\Upsilon_k Q_k) + \sum_{k=1}^{K} \mu_k (n_k - \text{tr}(Q_k)),
\]

where \(\Upsilon \triangleq \{\Upsilon_k\}_{\forall k}\) and \(\mu \triangleq \{\mu_k\}_{\forall k}\) are the Lagrange multipliers associated with the problem constraints.

In order to express the partial derivative of \(\mathcal{V}_N(\sigma^2, Q_1, \ldots, Q_K)\) with respect to \(Q_k\), i.e. \(\frac{\partial \mathcal{V}_N}{\partial Q_k}\), we define \(\mathcal{I}(\sigma^2, Q_1, \ldots, Q_K) = \frac{1}{N} \log \det (I_n + FQ)\). From (24), it is noted that the parameters affected by the perturbation of \(Q_k\) are \(\mathcal{I}(\sigma^2, Q_1, \ldots, Q_K), \{e_{l,k}\}_{\forall l,k}\), and \(\{\tilde{e}_{l,k}\}_{\forall l,k}\). As a result, we have

\[
\frac{\partial \mathcal{V}_N}{\partial Q_k} = \frac{\partial \mathcal{V}_N}{\partial \mathcal{I}} \frac{\partial \mathcal{I}}{\partial Q_k} + \sum_{l,k} \frac{\partial \mathcal{V}_N}{\partial e_{l,k}} \frac{\partial e_{l,k}}{\partial Q_k} + \sum_{l,k} \frac{\partial \mathcal{V}_N}{\partial \tilde{e}_{l,k}} \frac{\partial \tilde{e}_{l,k}}{\partial Q_k}.
\]

(29)

It can be checked that \(\frac{\partial \mathcal{V}_N}{\partial e_{l,k}} = 0\) and \(\frac{\partial \mathcal{V}_N}{\partial \tilde{e}_{l,k}} = 0, \forall l, k\). Therefore, the Karush-Kuhn-Tucker (KKT) conditions
of (26) are
\[
\begin{aligned}
- \frac{1}{N} (I_n + FQ)^{-1} F_k + \Upsilon_k - \mu_k I_{n_k} &= 0, \\
\text{tr} (\Upsilon_k Q_k) &= 0, \quad \Upsilon_k \succeq 0, \quad Q_k \succeq 0, \\
\mu_k (n_k - \text{tr}(Q_k)) &= 0, \quad \mu_k \geq 0,
\end{aligned}
\]

(30)

for \(k = 1, \ldots, K\).

Since (26) is a convex optimization problem with constraints satisfying Slater’s condition, the optimal \(Q_k\)'s can be found by solving the KKT conditions [36]. Using Lemma 14, the first line of (30) can be rewritten as
\[
- \frac{1}{N} (I_{n_k} + P_k Q_k)^{-1} P_k + \Upsilon_k - \mu_k I_{n_k} = 0,
\]

(31)

where
\[
P_k \defeq \langle (I_n + FQ)_{\setminus k}^{-1} F \rangle_k,
\]

(32)
\[
Q_{\setminus k} \defeq \text{diag} (Q_1, \ldots, Q_{k-1}, 0, Q_{k+1}, \ldots, Q_K).
\]

(33)

Note that \(P_k\) is a function of \((Q_1, \ldots, Q_K)\) rather than only \(Q_{\setminus k}\), as \(F\) defined in (25) includes the whole \(Q_k\)’s. For brevity, we have omitted its argument when writing \(P_k\). Substituting (31) for the first line of (30), the KKT conditions (30) are now equivalent to those of the following optimization problem:
\[
\max_{Q_k \in Q_k} \frac{1}{N} \log \det (I_{n_k} + P_k Q_k),
\]

(34)

which can be solved by a standard iterative waterfilling procedure. Thus, we get the next proposition.

**Proposition 2** Let \(P_k^*\) be the matrix in (32) by replacing \((Q_1, \ldots, Q_K)\) with \((Q_1, \ldots, Q_k^*, \ldots, Q_K)\) and \(P_k^* = V_{P_k} \Lambda_{P_k} U_{P_k}^H\). The eigenvectors of \(Q_k^*\) coincide with the right singular vectors of matrix \(P_k^*\), i.e.,
\[
Q_k^* = U_{P_k} \Lambda_{Q_k}^* U_{P_k}^H,
\]

(35)

and the eigenvalues are given by
\[
\Lambda_{Q_k}^* = \left( \frac{1}{\mu_k} I_{n_k} - \Lambda_{P_k}^{-1} \right)^+ ,
\]

(36)

where \((a)^+ = \max\{0, a\}\) and \(\mu_k\) is chosen to satisfy the power constraints \(\text{tr}(Q_k^*) = n_k\).

Using Proposition 2, we have the following observations:
• \( H = 0 \) — In this case, \( P_k = \sum_{l=1}^{L} \beta_{l,k} e_{l,k} (\sigma^2) T_{l,k} \). Therefore, the optimal transmit directions align with the eigenvectors of some weighted sum of \( T_{l,k} \)'s. As such, \( \beta_{l,k} e_{l,k} (\sigma^2) \) can be understood as the equivalent channel gain contributed by BS \( l \).

• \( H = \bar{H} \) — This implies that the channels are deterministic. In this case,

\[
P_k = \bar{H}_k^H \left( \sigma^2 I_N + \sum_{j \neq k} \bar{H}_j Q_j \bar{H}_j^H \right)^{-1} \bar{H}_k.
\]  

(37)

It shows that the optimal input covariance matrix of each user follows the water-filling principle that treats the other users as noise. This characteristic agrees with that for finite-size systems [35].

• \( K = 1 \) — In this case, we have

\[
P_1 = \sum_{l=1}^{L} \left( \beta_{l,1} e_{l,1} T_{l,1} + \frac{1}{\sigma^2} \bar{H}_{l,1}^H \left( \bar{e}_{l,1} R_{l,1} + I_{N_l} \right)^{-1} \bar{H}_{l,1} \right).
\]

If \( \{ R_{l,1} = I_{N_l} \} \forall l \), the optimal transmit directions thus align with the eigenvectors of some weighted sum of \( T_{l,1} \)'s and \( \bar{H}_{l,1}^H \bar{H}_{l,1} \)'s. While if \( R_{l,1} \neq I_{N_l} \), the impact of \( R_{l,1} \) on the optimal transmit directions is involved by \( \bar{H}_{l,1} \) via \( \bar{H}_{l,1}^H \left( \bar{e}_{l,1} R_{l,1} + I_{N_l} \right)^{-1} \bar{H}_{l,1} \). It appears that if the link pair does not have LOS, the corresponding correlation pattern at the receiver side does not provide a “direct” impact on the structure of the optimal transmit directions. Nevertheless, this inference is not entirely true, since the optimal transmit directions still can be changed by the correlation pattern at the receiver side through \( \beta_{l,1} e_{l,1} \). We will illustrate this phenomenon by an example in the simulation results.

Through the observations above, Proposition 2 shows its potential in understanding the impact of antenna correlations and LOS components on the structure of the optimal transmit directions. We now introduce an iterative algorithm for optimizing \( V_N(\sigma^2, Q_1, \ldots, Q_K) \) which adapts parameters \( Q \) and \( \{ e_{l,k} \} \forall l,k, \{ \bar{e}_{l,k} \} \forall l,k \) separately.

**Algorithm 1 (Optimization for \( Q \))**

- **Initialization:** \( Q_k^{(0)} = I_{N_k}, e_{l,k}^{(0)} = 1 \) and \( \bar{e}_{l,k}^{(0)} = 1 \) for \( k = 1, \ldots, K \) and \( l = 1, \ldots, L \).

- **Iteration \( t \):**

  - Given that \( Q_k^{(t-1)}, e_{l,k}^{(t-1)} \) and \( \bar{e}_{l,k}^{(t-1)} \) are available, for \( l = 1, \ldots, L \) and \( k = 1, \ldots, K \);
there is no known condition (e.g., spatial correlation pattern) to exclude such behavior of the frozen water-filling algorithm. In [15, 16, 18, 38, 39], named the frozen water-filling. For the frozen water-filling, convergence of Algorithm 1 has been proved in [34]. Note that Algorithm 1 is slightly different from those defined as the unique solutions of (14) at every iteration step.

In this section, computer simulations are conducted to evaluate the accuracy of the approximation $V_N(\sigma^2)$ in Theorem 2, and the effectiveness of the iterative algorithm developed in Algorithm 1. In particular, we are interested in their performances when the numbers of antennas are not so large. The simulation settings

- Calculate $T_{l,k}$ and $\bar{H}_{l,k}$ by the replacements of (23) for $l = 1, \ldots, L$ and $k = 1, \ldots, K$. Then, $\{e_{l,k}^{(t)}\}_{\forall l,k}, \{\tilde{e}_{l,k}^{(t)}\}_{\forall l,k}$ are obtained by

\[
e_{l,k}^{(t)} = \frac{1}{N_l} \text{tr} \left( R_{l,k} \langle \Psi^{(t-1)} \rangle_{l,k} \right),
\]

\[
\tilde{e}_{l,k}^{(t)} = \frac{1}{n_k} \text{tr} \left( T_{l,k} \langle \tilde{\Psi}^{(t-1)} \rangle_{l,k} \right),
\]

where

\[
\Psi^{(t-1)} = \left( \Phi^{(t-1)} \right)^{-1} + \sigma^2 \bar{H} \Phi^{(t-1)} \bar{H}^T \right)^{-1},
\]

\[
\tilde{\Psi}^{(t-1)} = \left( \tilde{\Phi}^{(t-1)} \right)^{-1} + \sigma^2 \tilde{H} \tilde{\Phi}^{(t-1)} \tilde{H}^T \right)^{-1},
\]

\[
\Phi^{(t-1)} = \text{diag} \left( \left\{ \left( \sigma^2 I_{N_l} + \sigma^2 \sum_{k=1}^{K} e_{l,k}^{(t-1)} R_{l,k} \right) \right\}_{\forall l}^{-1} \right),
\]

\[
\tilde{\Phi}^{(t-1)} = \text{diag} \left( \left\{ \left( \sigma^2 I_{n_k} + \sigma^2 \sum_{l=1}^{L} \tilde{\beta}_{l,k} e_{l,k}^{(t-1)} T_{l,k} \right) \right\}_{\forall k}^{-1} \right)
\]

- Calculate $P_k^{(t)}$ based on (32), for $k = 1, \ldots, K$;

- Calculate $Q_k^{(t)}$ based on Proposition 2, for $k = 1, \ldots, K$.

- Update $t := t + 1$ until $|V_N(\sigma^2, Q_1^{(t)}, \ldots, Q_K^{(t)}) - V_N(\sigma^2, Q_1^{(t-1)}, \ldots, Q_K^{(t-1)})|$ is small enough.

A similar iteration procedure was adopted by [34]. For the case with $K = L = 1$ and $\bar{H} = 0$, the convergence of Algorithm 1 has been proved in [34]. Note that Algorithm 1 is slightly different from those in [15, 16, 18, 38, 39], named the frozen water-filling. For the frozen water-filling, $\{e_{l,k}^{(t)}\}_{\forall l,k}, \{\tilde{e}_{l,k}^{(t)}\}_{\forall l,k}$ are defined as the unique solutions of (14) at every iteration step $t$, while in Algorithm 1, $\{e_{l,k}^{(t)}\}_{\forall l,k}, \{\tilde{e}_{l,k}^{(t)}\}_{\forall l,k}$ are obtained by performing a single update. It was pointed out in [34] that the frozen water-filling algorithm does not always converge.\(^9\) The convergence proof of Algorithm 1 is still an open challenge now.

4 Simulation Results

In this section, computer simulations are conducted to evaluate the accuracy of the approximation $V_N(\sigma^2)$ in Theorem 2, and the effectiveness of the iterative algorithm developed in Algorithm 1. In particular, we are interested in their performances when the numbers of antennas are not so large. The simulation settings

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\(^9\)Note that an example of oscillating behavior of the frozen water-filling algorithm is artificially constructed in [34]. However, there is no known condition (e.g., spatial correlation pattern) to exclude such behavior of the frozen water-filling algorithm.
Figure 2: Ergodic sum rate versus SNRs with $N_1 = N_2 = n_1 = n_2 = 2$ and $N_1 = N_2 = n_1 = n_2 = 8$ for a) $\kappa_{l,k} = 0$, $\forall l, k$ and b) $\kappa_{l,k} = 1$, $\forall l, k$. The solid lines plot the deterministic equivalent results, while the markers plot the Monte-Carlo simulation results under different fading distributions.
are based on the propagation model introduced in [40], in which the spatial correlation is generated from a uniform linear array with half wavelength spacing in a wireless scenario where there is one propagation path cluster with Gaussian power azimuthal distribution having mean angle of $\theta_{k,l}$ and root-mean-square spread of $\delta_{k,l}$. Specifically, we take the correlation matrix with elements [10]

$$
[T_{l,k}]_{m,n} \text{ (or } [R_{l,k}]_{m,n}) = \int_{-\frac{\pi}{2}}^{\frac{\pi}{2}} \frac{d\phi}{2\pi\delta_k^2} e^{j\pi(m-n)\sin\left(\frac{\phi}{180}\right) - j\frac{(\phi - \theta_{l,k})^2}{2\delta_k^2}} \tag{38}
$$

with $m, n$ being the indices of antennas. In addition, we use the superscripts $T$ and $R$, respectively, to refer to the corresponding values at the transmit and receive sides. The LOS matrix $\tilde{H}_{l,k}$ is generated according to $\tilde{H}_{l,k} = a_{R,l}(\tilde{\theta}_l^{R}) a_{T,k}(\tilde{\theta}_l^{T})^H$ where

$$
a_{R,l}(\tilde{\theta}_l^{R}) = \begin{bmatrix} e^{-j\pi \sin\left(\frac{\theta_{l,k}}{180}\pi\right)} & \cdots & e^{-j\pi (N_l-1) \sin\left(\frac{\theta_{l,k}}{180}\pi\right)} \end{bmatrix}^T,
$$

$$
a_{T,k}(\tilde{\theta}_l^{T}) = \begin{bmatrix} e^{-j\pi \sin\left(\frac{\theta_{l,k}}{180}\pi\right)} & \cdots & e^{-j\pi (N_k-1) \sin\left(\frac{\theta_{l,k}}{180}\pi\right)} \end{bmatrix}^T.$$

Regarding the fading distribution, we assume that $X_{ij}^{(l,k)}$ is of the form $W_{R,ij}^{(l,k)} \cos(\theta_{R,ij}^{(l,k)}) + jW_{I,ij}^{(l,k)} \sin(\theta_{I,ij}^{(l,k)})$ [41], where $\theta_{R,ij}^{(l,k)}$'s (and $\theta_{I,ij}^{(l,k)}$'s) are the phases modeled as i.i.d. uniform random variables over $[0, 2\pi]$, and those $W_{R,ij}^{(l,k)}$'s (and $W_{I,ij}^{(l,k)}$'s) are the amplitude fading drawn from a distribution with $E\{W_{R,ij}^{(l,k)}\} = 1$. The typical probability distributions of $W_{R,ij}^{(l,k)}$ include the Rayleigh, Nakagami, and log-normal distributions [20, 21]. Throughout this section, all the expected values (e.g., $E\{\mathcal{V}_N(\sigma^2)\}$) are obtained by the Monte-Carlo method in which 10,000 independent realizations of $\mathbf{H}$ are used for averaging.

In Theorem 2, we have shown that in the large-system limit the ergodic sum rate is invariant in distribution and can be well approximated by $\mathcal{V}_N(\sigma^2)$. Therefore, it is important to see how well $\mathcal{V}_N(\sigma^2)$

| $\theta_{R,1,1}^R$ | $\theta_{R,2,1}^R$ | $\theta_{R,1,2}^R$ | $\theta_{R,2,2}^R$ | $\theta_{I,1,1}^I$ | $\theta_{I,2,1}^I$ | $\theta_{I,1,2}^I$ | $\theta_{I,2,2}^I$ |
|-----------------|-----------------|-----------------|-----------------|-----------------|-----------------|-----------------|-----------------|
| $10^\circ$      | $20^\circ$      | $30^\circ$      | $40^\circ$      | $15^\circ$      | $25^\circ$      | $35^\circ$      | $45^\circ$      |
| $\delta_{R,1,1}^I$ | $\delta_{R,2,1}^I$ | $\delta_{R,1,2}^I$ | $\delta_{R,2,2}^I$ | $\delta_{I,1,1}^I$ | $\delta_{I,2,1}^I$ | $\delta_{I,1,2}^I$ | $\delta_{I,2,2}^I$ |
| 0.01            | 0.02            | 0.03            | 0.04            | 0.04            | 0.03            | 0.02            | 0.01            |
| $\theta_{I,1,1}^R$ | $\theta_{I,2,1}^R$ | $\theta_{I,1,2}^R$ | $\theta_{I,2,2}^R$ | $\theta_{I,1,1}^I$ | $\theta_{I,2,1}^I$ | $\theta_{I,1,2}^I$ | $\theta_{I,2,2}^I$ |
| $10^\circ$      | $20^\circ$      | $30^\circ$      | $40^\circ$      | $40^\circ$      | $40^\circ$      | $30^\circ$      | $20^\circ$      | $10^\circ$      |
Table 2: Average execution time in seconds.

|                          | $L = 10, K = 20$ | $L = 20, K = 40$ | $L = 30, K = 60$ |
|--------------------------|------------------|------------------|------------------|
| Monte-Carlo simulation   | 490              | 1941             | 4541             |
| Deterministic approximation | 0.5              | 3.2              | 8.2              |

in (19) approximates to the ergodic sum rate $E\{V_{B_N}(\sigma^2)\}$ when the dimensions of the system are not so large. For this purpose, Figure 2 compares the results of $E\{V_{B_N}(\sigma^2)\}$ with $V_N(\sigma^2)$ for $K = 2$ and $L = 2$ under different fading distributions. Their mean arrival/departure angles and angular spreads are given in Table 1 and their distance-dependent pathlosses are $g_{1,1} = g_{2,2} = 1$ and $g_{1,2} = g_{2,1} = 0.25$. We see that $V_N(\sigma^2)$ produces very good estimates for $E\{V_{B_N}(\sigma^2)\}$ even when only a few antenna elements (e.g., $N_1 = N_2 = n_1 = n_2 = 2$) are located at each UE and antenna set. As expected, when the number of antennas grows large (e.g., $N_1 = N_2 = n_1 = n_2 = 8$) all curves tend to overlap regardless of the distributions. In addition, we notice that for the Nakagami-\(m\) distribution, the difference between the case $m = 0.5$ and $m = 10$ is small even when there are only a few antenna elements.

In the above experiments, we have shown that $V_N(\sigma^2)$ provides a very good approximation for the sum rate of finite-dimensional systems. Before proceeding, it is useful to discuss the computational efficiency of evaluating $E\{V_{B_N}(\sigma^2)\}$ through $V_N(\sigma^2)$. For the considered scenarios in Figure 2 with $K = 2$ and $L = 2$, the execution time for evaluating $E\{V_{B_N}(\sigma^2)\}$ is at the order of decasecond (i.e., $10^1$ seconds). Although the execution time for $V_N(\sigma^2)$ is only at the order of centisecond (i.e., $10^{-2}$ seconds), one may not be convinced to use $V_N(\sigma^2)$ since writing a program to perform $E\{V_{B_N}(\sigma^2)\}$ is much easier than that for $V_N(\sigma^2)$. However, when the numbers of $K$ and $L$ grow, the Monte-Carlo simulations will become very demanding. Table 2 gives the average execution times on a 2.93 GHz Intel CPU with 4 GB of RAM under various system sizes. Here, we set $\{N_l = n_k = 2\}_{l,k}$, and the spatial correlation and LOS are generated from an arbitrary pattern. For typical systems with twenties of distributed antenna sets and forties of users, the simulations become prohibitive, ruling out the possibility for other system optimization designs such as scheduling [9, 23]. Clearly, the proposed deterministic equivalent result is much more efficient in this sense and provides a promising foundation to further applications of system optimization.

Next, we examine if the input covariance design based on the deterministic equivalent results performs well under different fading distributions when the numbers of antennas are not so large. Recall that $\{Q_1^\circ, \ldots, Q_K^\circ\}$ denote the optimal solutions of (22) that maximize the ergodic sum rate; and $\{Q_1^\star, \ldots, Q_K^\star\}$ denote the optimal solutions of (26) that maximize the deterministic equivalent of the ergodic sum rate. Algorithm 1 is used for solving $\{Q_1^\star, \ldots, Q_K^\star\}$, while $\{Q_1^\circ, \ldots, Q_K^\circ\}$ is solved by the Vu-Paulraj algorithm
Figure 3: Achievable rates versus SNRs with $N_1 = N_2 = n_1 = n_2 = 2$ for a) $\{\kappa_{l,k} = 0, \forall l,k\}$ and b) $\{\kappa_{l,k} = 1, \forall l,k\}$. The lines plot the results based on the deterministic equivalent, while the markers on dotted line plot the results for the Vu-Paulraj algorithm.
which is based on the barrier method where the ergodic sum rate and their first and second derivatives are calculated by the Monte-Carlo method. In contrast to \( \{Q^o_1, \ldots, Q^o_K\}, \{Q^*_1, \ldots, Q^*_K\} \) is independent from the true distributions of \( X_{l,k} \)'s. In Figure 3, we depict \( \mathbb{E}\{\mathcal{V}_{B_N}(\sigma^2, Q_1, \ldots, Q_K)\} \) when the input covariance matrices are \( \{Q^o_1, \ldots, Q^o_K\}, \{Q^*_1, \ldots, Q^*_K\} \), and identity matrices, when the amplitude fading distributions are either Rayleigh or log-normal. The reason for considering the two distributions is because, from Figure 2, the values of \( \mathbb{E}\{\mathcal{V}_{B_N}(\sigma^2)\} \) for the two distributions are significantly different and \( \mathcal{V}_{B_N}(\sigma^2) \) does not get very good estimation on \( \mathbb{E}\{\mathcal{V}_{B_N}(\sigma^2)\} \) when the amplitude fading distribution is log-normal. However, regardless of Rayleigh or log-normal distributions, the ergodic sum rate based on \( \{Q^*_1, \ldots, Q^*_K\} \) provides indistinguishable results to that based on \( \{Q^o_1, \ldots, Q^o_K\} \). In addition to its ability of providing good performance, Algorithm 1 is computationally much more efficient than the Vu-Paulraj algorithm.

Finally, we discuss the fact mentioned in Section 3.2 that the optimal transmit directions can be changed by the correlation pattern at the receiver side through \( \beta_{l,1}e_{l,1} \). To understand this better, we consider two scenarios with \( K = 1, L = 2 \) and \( \bar{H} = 0 \). The two scenarios use the same parameters except that the radiation patterns at the receiver of the second antenna set have different beam-widthes. Specifically, the radiation patterns at the receiver of the second antenna set have \( \delta_{2,1}^R = 0.01 \) for scenario 1 and \( \delta_{2,1}^R = 0.1 \) for scenario 2. We find it useful to observe the array patterns by plotting its array factor\footnote{Consider a uniform linear array with half wavelength spacing. Given a vector \( a \in \mathbb{C}^{n \times 1} \), we can get its array factor in direction \( \phi \) by \[ f(\phi) = \sum_{l=1}^{n} a_le^{-j\phi \sin(\frac{\pi l}{\lambda})}. \]}.  

Figure 4: Antenna radiation patterns.
The array patterns of \( T_{1,1} \) and \( T_{2,1} \) are depicted in Figure 4(a), where \( \theta_{T_{1,1}} = 30^\circ, \theta_{T_{2,1}} = 60^\circ, \delta_{T_{1,1}} = 0.04, \) and \( \delta_{T_{2,1}} = 0.03. \) The array patterns of the optimal input covariance \( Q_1 \) for the two scenarios are given in Figure 4(b) and (c), respectively. Figure 4(c) corresponds to the setting with the broader beamwidth of \( R_{2,1}. \) In this case, the optimal covariance is shown to feed the signal largely according to \( T_{2,1}, \) showing that the optimal transmit directions can be changed by the correlation pattern at the receiver side.

5 Conclusion

By using the large dimensional RMT, this paper investigated the deterministic equivalents for the large-scale MIMO MAC. The considered model includes the large-scale MIMO channel such as the general spatial correlation, the LOS components, and the channel entries being non-Gaussian. In particular, we derived the deterministic equivalent of the ergodic sum rate of the large-scale MIMO MAC. In addition, through the deterministic equivalent of the ergodic sum rate, we investigated the capacity-achieving input covariance matrices for the the large-scale MIMO MAC and proposed the iterative waterfilling algorithm for finding them. Finally, computer simulations were conducted to conclude the following three facts: First, the deterministic equivalent of the ergodic sum rate provides a very good approximation even when the numbers of antennas are of practical size. Second, calculating the ergodic sum rate by using the deterministic equivalent result is much more efficient than that by using the Monte-Carlo method when the system sizes are large. Hence, the deterministic equivalent result is of interest to addressing complex system optimization problems. Third, the optimal input covariance matrices predicted by the deterministic equivalent result are indeed remarkably close to those obtained by the corresponding finite-dimensional optimization approach, but in a much more efficient manner.

Investigation of the central limit theorem of the sum rate for the large-scale MIMO MAC by using the mathematical framework in [13], as well as application of the deterministic equivalent results to system-level designs [9, 23–25], are promising topics for future research.
Appendix

A Proof of $E\{m_{B_N}\} - \frac{1}{N} \text{tr} (\Psi) = O\left(\frac{1}{N^2}\right)$ in Theorem 1

We start the proof by reformulating the channel model so that the derivation can be performed systematically. To this end, we denote

$$\mathbf{R}_{l,k} \triangleq \text{diag} (0_{N_1}, \ldots, 0_{N_{l-1}}, \mathbf{R}_{l,k}, 0_{N_{l+1}}, \ldots, 0_{N_L}) .$$

(39)

$$\mathbf{T}_{l,k} \triangleq \text{diag} (0_{n_1}, \ldots, 0_{n_{k-1}}, \mathbf{T}_{l,k}, 0_{n_{k+1}}, \ldots, 0_{n_K}).$$

(40)

Also, let $\bar{\mathbf{H}}_{l,k}$ be the all-zero $N \times n$ matrix except that $\mathbf{H}_{l,k}$ is used for its $(\sum_{i=1}^{l-1} N_i + 1)$ to $(\sum_{i=1}^{l} N_i)$-th row and $(\sum_{j=1}^{k-1} n_j + 1)$ to $(\sum_{j=1}^{k} n_j)$-th column. As a result, $\mathbf{H}$ is statistically equivalent to

$$\mathbf{H} = \sum_{l,k} \mathbf{H}_{l,k} = \sum_{l,k} \left( \mathbf{H}_{l,k} + \bar{\mathbf{H}}_{l,k} \right) ,$$

(41)

where $\mathbf{H}_{l,k} = (\mathbf{R}_{l,k}^{\frac{1}{2}} \mathbf{X}_{l,k} \mathbf{T}_{l,k}^{\frac{1}{2}})^{\frac{1}{2}} \in \mathbb{C}^{N \times n}$, and $\mathbf{X}_{l,k} \equiv \left[ \frac{1}{\sqrt{n_k}} X^{l,(k)}_{ij} \right] \in \mathbb{C}^{N \times n}$ consists of the random components of the channel. Here, $\mathbf{X}_{l,k}$'s are assumed to be mutually independent. From (6), we have

$$\mathbf{B}_N = \left( \sum_{l,k} \left( \mathbf{R}_{l,k}^{\frac{1}{2}} \mathbf{X}_{l,k} \mathbf{T}_{l,k}^{\frac{1}{2}} + \bar{\mathbf{H}}_{l,k} \right) \right)^H \left( \sum_{l,k} \left( \mathbf{R}_{l,k}^{\frac{1}{2}} \mathbf{X}_{l,k} \mathbf{T}_{l,k}^{\frac{1}{2}} + \bar{\mathbf{H}}_{l,k} \right) \right)$$

(42)

and

$$\hat{\mathbf{B}}_N = \left( \sum_{l,k} \left( \mathbf{R}_{l,k}^{\frac{1}{2}} \mathbf{X}_{l,k} \mathbf{T}_{l,k}^{\frac{1}{2}} + \bar{\mathbf{H}}_{l,k} \right) \right)^H \left( \sum_{l,k} \left( \mathbf{R}_{l,k}^{\frac{1}{2}} \mathbf{X}_{l,k} \mathbf{T}_{l,k}^{\frac{1}{2}} + \bar{\mathbf{H}}_{l,k} \right) \right) ,$$

(43)

where $\mathbf{X}_{l,k}$'s and $\hat{\mathbf{X}}_{l,k}$'s are matrices with entries satisfying Assumption 1 but $\hat{\mathbf{X}}_{l,k}$'s are Gaussian.

Let $\mathcal{S}$ and $\hat{\mathcal{S}}$ be the resolvents of matrices $\mathbf{HH}^H$ and $\mathbf{H}^H \mathbf{H}$, respectively, given by

$$\mathcal{S} \triangleq (\mathbf{HH}^H + \omega \mathbf{I}_N)^{-1} ,$$

(44)

$$\hat{\mathcal{S}} \triangleq (\mathbf{H}^H \mathbf{H} + \omega \mathbf{I}_n)^{-1} .$$

(45)

These resolvents clearly satisfy the following useful properties:

$$\mathcal{S} \preceq \frac{1}{\omega} \mathbf{I}_N , \quad \text{and} \quad \hat{\mathcal{S}} \preceq \frac{1}{\omega} \mathbf{I}_n .$$

(46)

To facilitate our notations, we use $\bar{a}$ to denote the zero-mean random variable $a - E\{a\}$, where $a$ is a random variable. To accomplish the proof, the following two lemmas are useful.

Lemma 1 (Integration by Parts Formula for Gaussian Functionals) (see, e.g., [27, Proposition 2.4]) Let $\xi = [\xi_1, \ldots, \xi_M]^T$ be a complex Gaussian random vector such that $E\{\xi\} = 0$ and $E\{\xi \xi^H\} = \Omega$. Denoting
by $\Gamma(\xi)$ a complex function polynomially bounded with its derivatives, we have

$$E\{\xi_p \Gamma(\xi)\} = \sum_{m=1}^{M} \Omega_{pm} E\left\{ \frac{\Gamma(\xi)}{\partial \xi^m} \right\}.$$  (47)

**Lemma 2** (The Poincaré-Nash Inequality) (see, e.g., [27, Proposition 2.5]) Let $\xi = [\xi_1, \ldots, \xi_M]^T$ be a complex Gaussian random vector such that $E\{\xi\} = 0$ and $E\{\xi \xi^H\} = \Omega$. Denoting by $\Gamma(\xi)$ a complex function polynomially bounded with its derivatives, the following inequality holds true:

$$\text{Var}(\Gamma(\xi)) \leq E\left\{ (\nabla_\xi \Gamma(\xi))^T \Omega (\nabla_\xi \Gamma(\xi))^* + E\left\{ (\nabla_\xi^* \Gamma(\xi))^H \Omega (\nabla_\xi^* \Gamma(\xi)) \right\} \right\},$$  (48)

where $\nabla_\xi \Gamma(\xi) = \left[ \frac{\partial \Gamma}{\partial \xi_1}, \ldots, \frac{\partial \Gamma}{\partial \xi_M} \right]^T$ and $\nabla_\xi^* \Gamma(\xi) = \left[ \frac{\partial \Gamma}{\partial \xi^*_1}, \ldots, \frac{\partial \Gamma}{\partial \xi^*_M} \right]^T$.

The rigorous proof of Theorem 1 is rather complex. Although a standard procedure for the MIMO channel without the LOS components [18] is used, several additional manipulations for the LOS components to our present argument are required. To show this, we split the proof into two steps: First, we prove that $\text{tr}(E\{S\} - \Psi) \to 0$; secondly, we refine the convergence rate that $rac{1}{N} \text{tr}(E\{S\} - \Psi) = O \left( \frac{1}{N^2} \right)$. However, it is difficult to prove directly that $\text{tr}(E\{S\} - \Psi) \to 0$. To that end, we employ an intermediate quantity between $E\{S\}$ and $\Psi$ and establish the following two propositions.

**Proposition 3** As $N \to \infty$, we have

$$\text{tr}(E\{S\} - \Xi) \to 0,$$  (49a)
$$\text{tr}(E\{\tilde{S}\} - \tilde{\Xi}) \to 0,$$  (49b)

where

$$\Xi \triangleq \left[ \omega \left( I_N + \text{diag} \left( \sum_{k=1}^{K} \tilde{\alpha}_{l,k} R_{l,k} \right) \right) + \tilde{\Theta} \tilde{H} \tilde{H}^H \right]^{-1},$$  (50a)
$$\tilde{\Xi} \triangleq \left[ \omega \left( I_n + \text{diag} \left( \sum_{l=1}^{L} \alpha_{l,k} T_{l,k} \right) \right) + \tilde{H}^H \Theta H \right]^{-1},$$  (50b)
$$\Theta \triangleq \text{diag}(\Theta_1, \ldots, \Theta_L),$$  (50c)
$$\tilde{\Theta} \triangleq \text{diag}(\tilde{\Theta}_1, \ldots, \tilde{\Theta}_K),$$  (50d)
$$\Theta_l \triangleq \left[ \omega \left( I + \sum_{k=1}^{K} \alpha_{l,k} R_{l,k} \right) \right]^{-1}, \text{ for } l = 1, \ldots, L,$$  (50e)
$$\tilde{\Theta}_k \triangleq \left[ \omega \left( I + \sum_{l=1}^{L} \alpha_{l,k} T_{l,k} \right) \right]^{-1}, \text{ for } k = 1, \ldots, K,$$  (50f)
$$\alpha_{l,k} \triangleq \frac{1}{n_k} \text{tr}(R_{l,k} E\{\langle \mathcal{S} \rangle_{l_1} \}),$$  (50g)
$$\tilde{\alpha}_{l,k} \triangleq \frac{1}{n_k} \text{tr}(T_{l,k} E\{\langle \tilde{\mathcal{S}} \rangle_{k} \}).$$  (50h)
Proof: See Appendix A.1.

Proposition 4 As \( N \to \infty \), we have

\[
\begin{align*}
\text{tr}(\Xi - \Psi) & \to 0, \\
\text{tr}(\tilde{\Xi} - \tilde{\Psi}) & \to 0.
\end{align*}
\]

(51a)

(51b)

Proof: See Appendix A.2.

From (49a) to (51a), the proof of \( \mathbb{E}\{m_{B_N}^N\} - \frac{1}{N} \text{tr}(\Psi) \to 0 \) can be accomplished.

Proposition 5 As \( N \to \infty \), we have

\[
\begin{align*}
\frac{1}{N} \text{tr}(\mathbb{E}\{S\} - \Xi) & = \mathcal{O}\left(\frac{1}{N^2}\right), \\
\frac{1}{N} \text{tr}(\mathbb{E}\{\tilde{S}\} - \tilde{\Xi}) & = \mathcal{O}\left(\frac{1}{N^2}\right), \\
\frac{1}{N} \text{tr}(\Xi - \Psi) & = \mathcal{O}\left(\frac{1}{N^2}\right), \\
\frac{1}{N} \text{tr}(\tilde{\Xi} - \tilde{\Psi}) & = \mathcal{O}\left(\frac{1}{N^2}\right).
\end{align*}
\]

(52a)

(52b)

(52c)

(52d)

Proof: See Appendix A.3.

Consequently, (17) then follows from (52a) and (52c). The proof is complete.

A.1 Proof of Proposition 3

From (41) and (44), we have

\[
\begin{align*}
S & = \frac{1}{\omega} I_N - \frac{1}{\omega} S H H^H, \\
& = \frac{1}{\omega} I_N - \frac{1}{\omega} \sum_{l,k} S_{l,k} H_{l,k}^H,
\end{align*}
\]

(53)

and

\[
\begin{align*}
\mathbb{E}\{S_{pq}\} & = \frac{1}{\omega} \delta_{pq} - \frac{1}{\omega} \mathbb{E}\{[S H H^H]_{pq}\} \\
& = \frac{1}{\omega} \delta_{pq} - \frac{1}{\omega} \sum_{l,k} \mathbb{E}\{[S_{l,k} H_{l,k}^H]_{pq}\} - \frac{1}{\omega} \sum_{l \neq l_1, k \neq k_1} K \mathbb{E}\{[S_{l,k} H_{l,k}^H]_{pq}\}.
\end{align*}
\]

(54)

(55)

We first calculate \( \mathbb{E}\{[S H H^H]_{pq}\} \). Using the integration by parts formula (47), we write

\[
\begin{align*}
\mathbb{E}\{S_{pq} H_{ij}^{(l,k)} H_{i'i'}^{(l,k)*}\} & = \mathbb{E}\{S_{pq} H_{ij}^{(l,k)} H_{i'i'}^{(l,k)*}\} + \mathbb{E}\{S_{pq} H_{i'i'}^{(l,k)*}\} H_{ij}^{(l,k)} + \mathbb{E}\{S_{pq}\} H_{ij}^{(l,k)} H_{i'i'}^{(l,k)*} \\
& = \frac{1}{n_k} \sum_{m,n} H_{ij}^{(l,k)} H_{i'n}^{(l,k)*} + \frac{1}{n_k} \sum_{m,n} H_{ij}^{(l,k)*} H_{i'n}^{(l,k)} + \mathbb{E}\{S_{pq}\} H_{ij}^{(l,k)} H_{i'i'}^{(l,k)*},
\end{align*}
\]

(56)
and similarly, 

\[
E\left\{ S_p H_{ij}^{(l,k)} H_{qr}^{(l_1,k_1)^*} \right\} = \frac{1}{n_k} \sum_{m,n} R_{rm}^{(l,k)} T_{jn}^{(l,k)^*} E\left\{ \frac{\partial S_{pi}}{\partial H_{mn}^{(l,k)^*}} \right\} H_{ij}^{(l_1,k_1)^*} \\
+ \frac{1}{n_{k_1}} \sum_{m,n} R_{qm}^{(l_1,k_1)^*} T_{rn}^{(l,k)_1} E\left\{ \frac{\partial S_{pi}}{\partial H_{mn}^{(l,k)_1}} \right\} H_{ij}^{(l_1,k_1)^*} + E\{S_{pi}\} H_{ij}^{(l,k)} H_{qr}^{(l_1,k_1)^*}. \quad (57)
\]

Now, using the fact that 

\[
E\left\{ \frac{\partial S_{pi}}{\partial H_{mn}^{(l,k)}} \right\} = -E \{ S_{pm} [H^H S]_{ni} \}, \quad (58a)
\]
\[
E\left\{ \frac{\partial S_{pi}}{\partial H_{mn}^{(l,k)^*}} \right\} = -E \{ S_{mi} [SH]_{pm} \}, \quad (58b)
\]
\[
E\left\{ \frac{\partial S_{pi} H_{ij}^{(l,k)^*}}{\partial H_{mn}^{(l,k)}} \right\} = E\left\{ S_{pi} \delta_{qm} \delta_{rn} - S_{mi} [SH]_{pm} H_{ij}^{(l,k)^*} \right\}, \quad (58c)
\]

we have 

\[
E\left\{ S_p H_{ij}^{(l,k)} H_{qr}^{(l,k)^*} \right\} = \frac{1}{n_k} R_{eq}^{(l,k)} T_{jr}^{(l,k)^*} E\{S_p\} - \frac{1}{n_k} E\left\{ [R_{j,k} S]_{ii} [SHT_{j,k}]_{pq} H_{ij}^{(l,k)^*} \right\} \\
- \frac{1}{n_k} E\left\{ [S R_{j,k}]_{pq} [T_{i,k} H^H S]_{ri} \right\} L_{ij}^{(l,k)} + E\{S_{pi}\} \tilde{H}_{ij}^{(l,k)} \tilde{H}_{qr}^{(l,k)^*}, \quad (59)
\]

and 

\[
E\left\{ S_p H_{ij}^{(l,k)} H_{qr}^{(l_1,k_1)^*} \right\} = -\frac{1}{n_k} E\left\{ [R_{j,k} S]_{ii} [SHT_{j,k}]_{pq} H_{ij}^{(l_1,k_1)^*} \right\} \\
- \frac{1}{n_k} E\left\{ [S R_{j,k}]_{pq} [T_{i,k} H^H S]_{ri} \right\} \tilde{H}_{ij}^{(l,k)} + E\{S_{pi}\} \tilde{H}_{ij}^{(l,k)} \tilde{H}_{qr}^{(l_1,k_1)^*}. \quad (60)
\]

Then, summing over \( i \), we have 

\[
E\left\{ [SH]_{j,k} [P] H_{qr}^{(l,k)^*} \right\} = \frac{1}{n_k} T_{jr}^{(l,k)^*} E\{[SR]_{j,k} [P] \} - \frac{1}{n_k} E\left\{ \text{tr}(R_{j,k} S) [SHT_{j,k}]_{pq} H_{ij}^{(l,k)^*} \right\} \\
- \frac{1}{n_k} E\left\{ [SR_{j,k}]_{pq} [T_{i,k} H^H S]_{ri} \right\} \tilde{H}_{ij}^{(l,k)} + E\{[SH]_{j,k} [P] \} \tilde{H}_{qr}^{(l,k)^*}, \quad (61)
\]

and 

\[
E\left\{ [SH]_{j,k} [P] H_{qr}^{(l_1,k_1)^*} \right\} = -\frac{1}{n_k} E\left\{ \text{tr}(R_{j,k} S) [SHT_{j,k}]_{pq} H_{ij}^{(l_1,k_1)^*} \right\} \\
- \frac{1}{n_k} E\left\{ [SR_{j,k}]_{pq} [T_{i,k} H^H S]_{ri} \right\} S_{pq} + E\{[SH]_{j,k} [P] \} \tilde{H}_{qr}^{(l_1,k_1)^*}. \quad (62)
\]
Let \( \alpha_{l,k} \triangleq \frac{1}{n_k} \text{tr} \left( \mathbf{R}_{l,k} E \{ \mathbf{S} \} \right) \) and \( \eta_{l,k} \triangleq \frac{1}{n_k} \text{tr} \left( \mathbf{R}_{l,k} \mathbf{S} \right) - \alpha_{l,k} \). Then, we get

\[
E \left\{ [\mathbf{SH}]_{l,k} [\mathbf{pj}] [\mathbf{Hq}]^* \right\} = \frac{1}{n_k} T_{lj}^{(l,k)^*} E \left\{ [\mathbf{SR}]_{l,k} \mathbf{pq} \right\} - \alpha_{l,k} E \left\{ [\mathbf{SHT}]_{l,k} [\mathbf{pj}] [\mathbf{Hq}]^* \right\} - E \left\{ \eta_{l,k} [\mathbf{SHT}]_{l,k} [\mathbf{pj}] [\mathbf{Hq}]^* \right\}
\]

and

\[
E \left\{ [\mathbf{SH}]_{l,k} [\mathbf{pj}] [\mathbf{Hq}]^* \right\} = -\alpha_{l,k} E \left\{ [\mathbf{SHT}]_{l,k} [\mathbf{pj}] [\mathbf{Hq}]^* \right\} - E \left\{ \eta_{l,k} [\mathbf{SHT}]_{l,k} [\mathbf{pj}] [\mathbf{Hq}]^* \right\}
\]

From (63) and (64), we obtain

\[
E \left\{ [\mathbf{SH}]_{l,k} [\mathbf{pj}] [\mathbf{Hq}]^* \right\} = \frac{1}{n_k} T_{lj}^{(l,k)^*} E \left\{ [\mathbf{SR}]_{l,k} \mathbf{pq} \right\} - \sum_{l,k} \alpha_{l,k} E \left\{ [\mathbf{SHT}]_{l,k} [\mathbf{pj}] [\mathbf{Hq}]^* \right\}
\]

\[
-\sum_{l,k} E \left\{ \eta_{l,k} [\mathbf{SHT}]_{l,k} [\mathbf{pj}] [\mathbf{Hq}]^* \right\} - \sum_{l,k} \frac{1}{n_k} E \left\{ [\mathbf{SR}]_{l,k} \mathbf{pq} [\mathbf{T}]_{l,k} [\mathbf{H}]^* \right\} + E \left\{ [\mathbf{SH}]_{l,k} [\mathbf{pj}] [\mathbf{Hq}]^* \right\}.
\]

Defining

\[
\tilde{\Theta} \triangleq \omega \left( I_n + \sum_{l,k} \alpha_{l,k} [\mathbf{T}]_{l,k} \right)^{-1} = \text{diag} \left( \tilde{\Theta}_1, \ldots, \tilde{\Theta}_K \right),
\]

where \( \tilde{\Theta}_k \) is given by (50f). Multiplying both sides of (65) by \( [\tilde{\Theta}]_{kj} \) and summing over \( j \) and \( r \), we get

\[
E \{ [\mathbf{SHH}]_{pq} \} = \omega \sum_{l,k} \frac{1}{n_k} \text{tr} [\mathbf{T}]_{l,k} \tilde{\Theta} E \{ [\mathbf{SR}]_{l,k} \mathbf{pq} \} - \omega \sum_{l,k} \frac{1}{n_k} E \left\{ \eta_{l,k} [\mathbf{SHT}]_{l,k} \tilde{\Theta} \mathbf{Hq} \right\}
\]

\[
-\omega \sum_{l,k} \frac{1}{n_k} \left\{ \text{tr} [\mathbf{T}]_{l,k} \mathbf{H}^* \mathbf{SH} \tilde{\Theta} \right\} E \{ [\mathbf{SR}]_{l,k} \mathbf{pq} \} + \omega E \{ [\mathbf{SH} \tilde{\Theta} \mathbf{H}]_{pq} \}.
\]
This, together with (54), yields
\[
\mathbb{E}\{S\} = \frac{1}{\omega} I_N - \sum_{l,k} \frac{1}{n_k} \text{tr}(T_{l,k} \tilde{\Theta}) \mathbb{E}\{\mathbf{S}\mathbf{R}_{l,k}\} + \sum_{l,k} \mathbb{E}\left\{\eta_{l,k} SHT_{l,k} \tilde{\Theta} H^H\right\}
\]
\[
+ \sum_{l,k} \frac{1}{n_k} \mathbb{E}\left\{\text{tr}(T_{l,k} \tilde{H}^H S \tilde{\Theta}) \mathbf{S}\mathbf{R}_{l,k}\right\} - \mathbb{E}\left\{SH\tilde{\Theta} H^H\right\}
\]
\[
= \frac{1}{\omega} I_N - \sum_{l,k} \frac{1}{n_k} \text{tr}(T_{l,k} \tilde{\Theta}) \mathbb{E}\{\mathbf{S}\mathbf{R}_{l,k}\} + \sum_{l,k} \mathbb{E}\left\{\eta_{l,k} SHT_{l,k} \tilde{\Theta} H^H\right\}
\]
\[
+ \sum_{l,k} \frac{1}{n_k} \mathbb{E}\left\{\text{tr}(SH\tilde{\Theta} T_{l,k} (\tilde{H} + \tilde{H})^H) \mathbf{S}\mathbf{R}_{l,k}\right\} - \mathbb{E}\left\{SH\tilde{\Theta} H^H\right\}
\]
\[
= \frac{1}{\omega} I_N - \sum_{l,k} \frac{1}{n_k} \text{tr}(T_{l,k} \tilde{\Theta}) \mathbb{E}\{\mathbf{S}\mathbf{R}_{l,k}\} + \sum_{l,k} \mathbb{E}\left\{\eta_{l,k} SHT_{l,k} \tilde{\Theta} H^H\right\}
\]
\[
+ \sum_{l,k} \mathbb{E}\left\{\frac{1}{n_k} \text{tr}(S\tilde{\Theta} T_{l,k} \tilde{H}^H)\right\} \mathbb{E}\{\mathbf{S}\mathbf{R}_{l,k}\} + \sum_{l,k} \mathbb{E}\left\{\rho_{l,k} S\mathbf{R}_{l,k}\right\} - \mathbb{E}\left\{S\tilde{\Theta} H^H\right\},
\]
(68)

where the third equality follows from the following definitions
\[
\rho_{l,k}^{(1)} \triangleq \frac{1}{n_k} \text{tr}\left(S\tilde{\Theta} T_{l,k} \tilde{H}^H\right), \quad \rho_{l,k}^{(2)} \triangleq \frac{1}{n_k} \text{tr}\left(S\tilde{\Theta} T_{l,k} H^H\right).
\]
(69a)

Before proceeding, we establish the following lemma.

**Lemma 3**

\[
\mathbb{E}\left\{\frac{1}{n_k} \text{tr}(S\tilde{\Theta} T_{l,k} \tilde{H}^H)\right\} = -\omega \sum_{l_1,k_1} \alpha_{l_1,k_1} \mathbb{E}\left\{\frac{1}{n_k} \text{tr}(S\tilde{\Theta} T_{l_1,k_1} \tilde{T}_{l_1,k_1} \tilde{H}^H)\right\} - \omega \sum_{l_1,k_1} \mathbb{E}\left\{\eta_{l_1,k_1} \rho_{l,l_1,k_1}\right\},
\]
(70)

where
\[
\rho_{l,k,l_1,k_1}^{(3)} \triangleq \frac{1}{n_k} \text{tr}\left(S\tilde{\Theta} T_{l,k} \tilde{T}_{l_1,k_1} H^H\right).
\]
(71)

**Proof:** Using the integration by parts formula (47), we write
\[
\mathbb{E}\left\{S_{pi}\tilde{H} T_{l,k} \tilde{H}^*_{pr}\right\} = \sum_{l_1,k_1} \mathbb{E}\left\{S_{pi} \tilde{H}^*_{pr(l_1,k_1)}\right\}\tilde{H} T_{l,k} \tilde{H}^*_{pr}\]
\[
= \sum_{l_1,k_1} \frac{1}{n_k} \sum_{m,n} \mathbb{E}\left\{\frac{\partial S_{pi}}{\partial \tilde{H}^{(l_1,k_1)}}\right\}\tilde{T}_{l_1,k_1} \tilde{H}^*_{pr}\]
\[
= -\sum_{l_1,k_1} \frac{1}{n_k} \mathbb{E}\left\{S\mathbf{R}_{l_1,k_1} \tilde{T}_{l_1,k_1} \tilde{H}^*_{pr}\right\}\tilde{H} T_{l,k} \tilde{H}^*_{pr}
\]
(72)
Summing over $p$, we have

$$
E\left\{ [H\tilde{\Theta}T_{l,k}]_{ij} [\tilde{H}^H S]_{rl} \right\} = - \sum_{l_1,k_1} \frac{1}{n_{k_1}} [H\tilde{\Theta}T_{l,k}]_{ij} E\left\{ \text{tr}(SR_{l_1,k_1}) [T_{l_1,k_1} H^H S]_{rl} \right\} 
$$

$$
= - \sum_{l_1,k_1} \alpha_{l_1,k_1} [H\tilde{\Theta}T_{l,k}]_{ij} E\left\{ [R_{l_1,k_1} \tilde{H}^H S]_{rl} \right\} 
- \sum_{l_1,k_1} E\left\{ \eta_{l_1,k_1} [H\tilde{\Theta}T_{l,k}]_{ij} [T_{l_1,k_1} \tilde{H}^H S]_{rl} \right\} 
- \sum_{l_1,k_1} \frac{1}{n_{k_1}} [H\tilde{\Theta}T_{l,k}]_{ij} E\left\{ \text{tr}(SR_{l_1,k_1}) [T_{l_1,k_1} \tilde{H}^H S]_{rl} \right\}.
$$

(73)

After simple algebraic operations and summing over $i,j$ and $l$, we then get

$$
E\left\{ \text{tr} \left( S\tilde{H}\tilde{\Theta}T_{l,k} \tilde{H}^H \right) \right\} = -\omega \sum_{l_1,k_1} E\left\{ \eta_{l_1,k_1} \text{tr} \left( S\tilde{H}\tilde{\Theta}T_{l,k} \tilde{\Theta}T_{l_1,k_1} \tilde{H}^H \right) \right\} 
- \omega \sum_{l_1,k_1} \frac{1}{n_{k_1}} E\left\{ \text{tr} \left( ST_{l_1,k_1} \text{tr} \left( S\tilde{H}\tilde{\Theta}T_{l,k} \tilde{\Theta}T_{l_1,k_1} \tilde{H}^H \right) \right) \right\}.
$$

(74)

Therefore, we have

$$
E\left\{ \frac{1}{n_k} \text{tr} \left( S\tilde{H}\tilde{\Theta}T_{l,k} \tilde{H}^H \right) \right\} = -\omega \sum_{l_1,k_1} \alpha_{l_1,k_1} E\left\{ \frac{1}{n_k} \text{tr} \left( S\tilde{H}\tilde{\Theta}T_{l,k} \tilde{\Theta}T_{l_1,k_1} \tilde{H}^H \right) \right\} 
- \omega \sum_{l_1,k_1} E\left\{ \eta_{l_1,k_1} \rho_{l,k,l_1,k_1}^{(3)} \right\},
$$

(75)

where $\rho_{l,k,l_1,k_1}$ is given by (71). Using the fact that $E\left\{ \eta_{l_1,k_1} \rho_{l,k,l_1,k_1}^{(3)} \right\} = E\left\{ \eta_{l_1,k_1} \rho_{l,k,l_1,k_1}^{(3)} \right\}$, we obtain (70). \(\square\)

Applying this lemma to (68), we get

$$
E\left\{ S \right\} = \frac{1}{\omega} I_N - \sum_{l,k} \frac{1}{n_k} \text{tr}(T_{l,k} \tilde{\Theta}) E\{SR_{l,k}\} + \sum_{l,k} E\left\{ \eta_{l,k} SHT_{l,k} \tilde{\Theta}H^H \right\} 
+ \sum_{l,k} E\left\{ \frac{1}{n_k} \text{tr} \left( S\tilde{H}\tilde{\Theta}T_{l,k} \tilde{H}^H \right) \right\} E\{SR_{l,k}\} 
- \omega \sum_{l_1,k_1} \sum_{k_1} \alpha_{l_1,k_1} E\left\{ \frac{1}{n_k} \text{tr} \left( S\tilde{H}\tilde{\Theta}T_{l,k} \tilde{\Theta}T_{l_1,k_1} \tilde{H}^H \right) \right\} E\{SR_{l,k}\} 
- \omega \sum_{l_1,k_1} \sum_{k_1} \eta_{l_1,k_1} \rho_{l,k,l_1,k_1}^{(3)} E\{SR_{l,k}\} 
+ \sum_{l,k} E\left\{ \rho_{l,k}^{(1)} SR_{l,k} \right\} + \sum_{l,k} E\left\{ \rho_{l,k}^{(2)} SR_{l,k} \right\} - E\left\{ S\tilde{H}\tilde{\Theta}H^H \right\}.
$$

(76)

Define

$$
\Delta \overset{\Delta}{=} \sum_{l,k} E\left\{ \eta_{l,k} SHT_{l,k} \tilde{\Theta}H^H \right\} - \omega \sum_{l_1,k_1} \sum_{k_1} \eta_{l_1,k_1} \rho_{l,k,l_1,k_1}^{(3)} E\{SR_{l,k}\} 
+ \sum_{l,k} E\left\{ \rho_{l,k}^{(1)} SR_{l,k} \right\} + \sum_{l,k} E\left\{ \rho_{l,k}^{(2)} SR_{l,k} \right\}.
$$

(77)
Noting that

\[ \omega \sum_{l,l_1,k,k_1}^K \alpha_{l_1,k_1} E \left\{ \frac{1}{n_k} \text{tr} \left( S \tilde{H} \tilde{\Theta} \tilde{T}_{l,k} \tilde{\Theta} \tilde{T}_{l_1,k_1} \tilde{H}^H \right) \right\} E \{ S R_{l,k} \} \]

\[ = \sum_{l,k} E \left\{ \frac{1}{n_k} \text{tr} \left( S \tilde{H} \tilde{\Theta} \tilde{T}_{l,k} \tilde{H}^H \right) \right\} E \{ S R_{l,k} \} - \omega \sum_{l,k} E \left\{ \frac{1}{n_k} \text{tr} \left( S \tilde{H} \tilde{\Theta} \tilde{T}_{l,k} \tilde{H}^H \right) \right\} E \{ S R_{l,k} \}, \tag{78} \]

we therefore get

\[ E \{ S \} = \frac{1}{\omega} I_N - \sum_{l,k} \frac{1}{n_k} \text{tr}(T_{l,k} \tilde{\Theta}) E \{ S R_{l,k} \} + \omega \sum_{l,k} E \left\{ \frac{1}{n_k} \text{tr} \left( S \tilde{H} \tilde{\Theta} \tilde{T}_{l,k} \tilde{H}^H \right) \right\} E \{ S R_{l,k} \} \]

\[ - E \left\{ S \right\} E \{ S \} + \Delta. \tag{79} \]

Writing

\[ \tilde{\tau}_{l,k} \triangleq \frac{1}{n_k} \text{tr}(T_{l,k} \tilde{\Theta}) - \omega \frac{1}{n_k} \text{tr} \left( E \{ S \} \tilde{H} \tilde{\Theta} T_{l,k} \tilde{H}^H \right) = \frac{1}{n_k} \text{tr}(T_{l,k} \tilde{\Theta} \left( I_n - \omega \tilde{H}^H E \{ S \} \tilde{H} \right)), \tag{80} \]

we have

\[ E \{ S \} \left( I_N + \tilde{H} \tilde{\Theta} \tilde{H}^H \right) = \frac{1}{\omega} I_N - \sum_{l,k} \tilde{\tau}_{l,k} E \{ S \} R_{l,k} + \Delta, \tag{81} \]

and then

\[ E \{ S \} \left( I_N + \sum_{l,k} \tilde{\alpha}_{l,k} R_{l,k} + \tilde{H} \tilde{\Theta} \tilde{H}^H \right) = \frac{1}{\omega} I_N + \sum_{l,k} (\tilde{\alpha}_{l,k} - \tilde{\tau}_{l,k}) E \{ S \} R_{l,k} + \Delta. \tag{82} \]

As a result, we then get

\[ E \{ \tilde{S} \} = \Xi + \omega \sum_{l,k} (\tilde{\alpha}_{l,k} - \tilde{\tau}_{l,k}) E \{ \tilde{S} \} R_{l,k} + \omega \Delta \Xi, \tag{83} \]

where

\[ \Xi = \left( I_N + \sum_{l,k} \tilde{\alpha}_{l,k} R_{l,k} + \tilde{H} \tilde{\Theta} \tilde{H}^H \right)^{-1} \tag{84} \]

and

\[ \tilde{\alpha}_{l,k} \triangleq \frac{1}{n_k} \text{tr} \left( T_{l,k} E \{ \tilde{S} \} \right) = \frac{1}{n_k} \text{tr} \left( T_{l,k} E \{ \langle \tilde{S} \rangle \} \right). \tag{85} \]

To get Proposition 3, it remains to show that \( \tilde{\alpha}_{l,k} - \tilde{\tau}_{l,k} \to 0 \) and \( \text{tr}(\Delta \Xi) \to 0 \). To that end, we have to get a similar expression of \( E \{ \tilde{S} \} \) as that of (83). Following the same derivation of (83) from the beginning, we can get

\[ E \{ \tilde{S} \} = \tilde{\Xi} + \omega \sum_{l,k} (\tilde{\alpha}_{l,k} - \tilde{\tau}_{l,k}) E \{ \tilde{S} \} T_{l,k} + \omega \tilde{\Delta} \Xi, \tag{86} \]

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where

\[
\hat{\Delta} = \sum_{l,k} E \left\{ \hat{\eta}_{l,k} \hat{\mathcal{S}} \hat{H}^H \mathbf{R}_{l,k} \mathbf{H} \right\} - \omega \sum_{l_1} \sum_{k_1} E \left\{ \hat{\eta}_{l_1,k_1} \hat{\mathcal{S}} \hat{H}^H \mathbf{R}_{l_1,k_1} \mathbf{H} \right\} E \left\{ \hat{\mathcal{T}}_{l,k} \right\} \\
+ \sum_{l,k} E \left\{ \hat{\rho}_{l,k} \hat{\mathcal{S}} \hat{\mathcal{T}}_{l,k} \right\} + \sum_{l,k} E \left\{ \hat{\rho}_{l,k} \hat{\mathcal{T}}_{l,k} \hat{\mathcal{T}}_{l,k} \right\}
\]

(87a)

and those \( \hat{\Xi} \) and \( \Theta \) are given by (50b) and (50c) respectively.

From (50h), (80), (83) and (86), write

\[
\tilde{\alpha}_{l,k} = \frac{1}{n_k} \text{tr} \left( \mathbf{T}_{l,k} \hat{\Xi} \right) + \frac{\omega}{n_k} \sum_{i,j} (\alpha_{i,j} - \tilde{\tau}_{i,j}) \text{tr} \left( \mathbf{T}_{l,k} E \{ \hat{\mathcal{S}} \} \mathbf{T}_{i,j} \hat{\Xi} \right) + \frac{\omega}{n_k} \text{tr} \left( \mathbf{T}_{l,k} \hat{\Delta} \hat{\Xi} \right)
\]

(88)

and

\[
\tilde{\eta}_{l,k} = \frac{1}{n_k} \text{tr} \left( \mathbf{T}_{l,k} \hat{\Theta} \left( \mathbf{I}_N - \omega \hat{\mathcal{S}} \mathbf{H}^H \hat{\Theta} \right) \right) - \frac{\omega^2}{n_k} \sum_{i,j} (\tilde{\alpha}_{i,j} - \tilde{\tau}_{i,j}) \text{tr} \left( \mathbf{T}_{l,k} \hat{\Theta} \mathbf{H}^H E \{ \mathbf{S} \} \mathbf{R}_{i,j} \hat{\Xi} \right) \hat{\Theta} \right) \\
- \frac{\omega^2}{n_k} \text{tr} \left( \mathbf{T}_{l,k} \hat{\Theta} \mathbf{H}^H \Delta \hat{\Xi} \right) \\
= \frac{1}{n_k} \text{tr} \left( \mathbf{T}_{l,k} \hat{\Theta} \right) - \frac{\omega^2}{n_k} \sum_{i,j} (\tilde{\alpha}_{i,j} - \tilde{\tau}_{i,j}) \text{tr} \left( \mathbf{T}_{l,k} \hat{\Theta} \mathbf{H}^H E \{ \mathbf{S} \} \mathbf{R}_{i,j} \hat{\Xi} \right) \hat{\Theta} \right) - \frac{\omega^2}{n_k} \text{tr} \left( \mathbf{T}_{l,k} \hat{\Theta} \mathbf{H}^H \Delta \hat{\Xi} \right)
\]

(89)

Similarly,

\[
\alpha_{l,k} = \frac{1}{n_k} \text{tr} \left( \mathbf{R}_{l,k} \hat{\Xi} \right) + \frac{\omega}{n_k} \sum_{i,j} (\hat{\alpha}_{i,j} - \hat{\tau}_{i,j}) \text{tr} \left( \mathbf{R}_{l,k} E \{ \mathbf{S} \} \mathbf{R}_{i,j} \hat{\Xi} \right) + \frac{\omega}{n_k} \text{tr} \left( \mathbf{R}_{l,k} \Delta \hat{\Xi} \right)
\]

(90)
and

\[ \tau_{l,k} = \alpha_{l,k} - \frac{1}{n_k} \sum_{i,j} (\tilde{\alpha}_{i,j} - \tilde{\tau}_{i,j}) \text{tr} \left( R_{l,k} E \{ S \} R_{j,k} \Xi \right) - \frac{\omega}{n_k} \text{tr} \left( R_{l,k} \Delta \Xi \right) \]

\[ - \frac{\omega^2}{n_k} \sum_{i,j} (\alpha_{i,j} - \tau_{i,j}) \text{tr} \left( R_{l,k} \Theta H \{ \hat{S} \} T_{i,j} \hat{H}^H \Theta \right) - \frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \Theta H \Delta \hat{H}^H \Theta \right). \]  \hspace{1cm} (91)

Let \( \eta \triangleq [\text{vec}(A_1)^T, \text{vec}(A_2)^T]^T, \epsilon \triangleq [\text{vec}(C_1)^T, \text{vec}(C_2)^T]^T, \Gamma \triangleq \begin{bmatrix} \Gamma_{11} & \Gamma_{12} \\ \Gamma_{21} & \Gamma_{22} \end{bmatrix}, \) where \( A_1, A_2, C_1, C_2 \in \mathbb{C}^{L \times K}, \Gamma_{11}, \Gamma_{12}, \Gamma_{21}, \Gamma_{22} \in \mathbb{C}^{LK \times LK}, \) with

\[
\begin{align*}
[A_1]_{l,k} &= \tilde{\alpha}_{l,k} - \tilde{\tau}_{l,k}, \quad [A_2]_{l,k} = \alpha_{l,k} - \tau_{l,k}, \hspace{1cm} (92a) \\
[C_1]_{l,k} &= \frac{\omega}{n_k} \text{tr} \left( T_{l,k} \Delta \Xi \right) + \frac{\omega^2}{n_k} \text{tr} \left( T_{l,k} \hat{H}^H \Delta \Xi \hat{H} \Theta \right), \hspace{1cm} (92b) \\
[C_2]_{l,k} &= \frac{\omega}{n_k} \text{tr} \left( R_{l,k} \Delta \Xi \right) + \frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \Theta \hat{H} \Delta \hat{H} \Theta \right), \hspace{1cm} (92c) \\
[\Gamma_{11}]_{l,k,ij} &= \begin{cases} -\frac{\omega^2}{n_k} \text{tr} \left( T_{l,k} \hat{H}^H E \{ \hat{S} \} R_{j,k} \Xi \hat{H} \Theta \right), & (i,j) \neq (l,k); \\
1 - \frac{\omega^2}{n_k} \text{tr} \left( T_{l,k} \hat{H}^H E \{ \hat{S} \} R_{j,k} \Xi \hat{H} \Theta \right), & (i,j) = (l,k), \end{cases} \hspace{1cm} (92d) \\
[\Gamma_{12}]_{l,k,ij} &= -\frac{\omega}{n_k} \text{tr} \left( T_{l,k} E \{ \hat{S} \} T_{i,j} \hat{E} \right), \quad [\Gamma_{21}]_{l,k,ij} = -\frac{\omega}{n_k} \text{tr} \left( R_{l,k} E \{ \hat{S} \} R_{j,k} \Xi \right), \hspace{1cm} (92e) \\
[\Gamma_{22}]_{l,k,ij} &= \begin{cases} -\frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \Theta \hat{H} E \{ \hat{S} \} T_{i,j} \hat{E} \hat{H} \Theta \right), & (i,j) \neq (l,k); \\
1 - \frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \Theta \hat{H} E \{ \hat{S} \} T_{i,j} \hat{E} \hat{H} \Theta \right), & (i,j) = (l,k). \end{cases} \hspace{1cm} (92f)
\end{align*}
\]

From (89) and (91), we get \( \Gamma \eta = \epsilon. \) \hspace{1cm} (93)

If we can show that \( \epsilon \to 0 \) and \( \Gamma \) is invertible, we then get our desired result \( \eta \to 0. \) To show that \( \epsilon \to 0, \) we establish the following lemma.

**Lemma 4** For any uniformly bounded matrices \( Q \) and \( \hat{Q}, \) we have

\[
\begin{align*}
\frac{1}{n_k} \text{tr}(\Delta Q) &= O \left( \frac{1}{N^2} \right), \hspace{1cm} (94a) \\
\frac{1}{n_k} \text{tr}(\hat{\Delta} Q) &= O \left( \frac{1}{N^2} \right). \hspace{1cm} (94b)
\end{align*}
\]
Using Lemma 2 (the Poincaré-Nash inequality), we obtain
\[ \frac{1}{n_k} \text{tr} \left( \Delta \mathbf{Q} \right) = \sum_{l,k} E \left\{ \bar{\eta}_{l,k} \bar{\rho}_{l,k} \right\} - \omega \sum_{l,l_1,k,k_1} L \sum_{i,k} \frac{1}{n_k} \text{tr} \left( E \{ \mathbf{S} \} \mathbf{R}_{i,k} \mathbf{Q} \right) + \sum_{l,k} E \left\{ \bar{\eta}_{l,k} \bar{\rho}_{l,k} \right\} \frac{1}{n_k} \text{tr} \left( \mathbf{S} \mathbf{R}_{l,k} \mathbf{Q} \right) \],

where
\[ \bar{\rho}_{l,k} = \frac{1}{n_k} \text{tr} \left( \mathbf{S} \mathbf{H} \mathbf{T}_{l,k} \mathbf{H}^H \mathbf{Q} \right). \]

We first prove the following facts for any uniformly bounded matrices \( \mathbf{M} \),
\[ \text{Var} \left( \frac{1}{n_k} \text{tr} \left( \mathbf{S} \mathbf{M} \right) \right) = O \left( \frac{1}{N^2} \right). \]

For this, we let \( \Gamma (\mathbf{H}) \triangleq \frac{1}{n_k} \text{tr} \left( \mathbf{S} \mathbf{M} \right) \) which gives
\[ \frac{\partial \Gamma (\mathbf{H})}{\partial \mathbf{H}_{mn}} = \frac{1}{n_k} \sum_{p,q} \mathbf{M}_{pq} \frac{\partial \mathbf{S}_{pq}}{\partial \mathbf{H}_{mn}} = - \frac{1}{n_k} \left[ H^H \mathbf{S} \mathbf{M} \mathbf{S} H \mathbf{M} \right]_{mn}, \]
\[ \frac{\partial \Gamma (\mathbf{H})}{\partial \mathbf{H}_{mn}^*} = \frac{1}{n_k} \sum_{p,q} \mathbf{M}_{pq} \frac{\partial \mathbf{S}_{pq}}{\partial \mathbf{H}_{mn}^*} = - \frac{1}{n_k} \left[ \mathbf{S} \mathbf{M} \mathbf{S} H \mathbf{M} \right]_{mn^*}. \]

Using Lemma 2 (the Poincaré-Nash inequality), we obtain
\[ \text{Var} \left( \frac{1}{n_k} \text{tr} \left( \mathbf{S} \mathbf{M} \right) \right) \leq \sum_{l,k} \frac{1}{n_k} \sum_{m,n,m',n'} R_{l,m,m',n'} \mathbf{L}_{m',n'} \mathbf{E} \left\{ \frac{1}{n_k} \left[ H^H \mathbf{S} \mathbf{M} \mathbf{S} H \mathbf{M} \right]_{mn} \frac{1}{n_k} \left[ H^H \mathbf{S} \mathbf{M} \mathbf{S} H \mathbf{M} \right]_{n'm'} \right\} \]
\[ + \sum_{l,k} \frac{1}{n_k} \sum_{m,n,m',n'} R_{l,m,m',n'} \mathbf{L}_{m',n'} \mathbf{E} \left\{ \frac{1}{n_k} \left[ \mathbf{S} \mathbf{M} \mathbf{S} H \mathbf{M} \right]_{mn} \frac{1}{n_k} \left[ \mathbf{S} \mathbf{M} \mathbf{S} H \mathbf{M} \right]_{n'm'} \right\} \]
\[ = \sum_{l,k} \frac{1}{n_k} \mathbf{E} \left\{ \text{tr} \left( H^H \mathbf{S} \mathbf{M} \mathbf{S} H \mathbf{T}_{l,k} \right) \right\} \]
\[ + \sum_{l,k} \frac{1}{n_k^2} \mathbf{E} \left\{ \text{tr} \left( H^H \mathbf{S} \mathbf{M} \mathbf{S} \mathbf{H} \mathbf{T}_{l,k} \right) \right\}. \]

Noting the fact that (using \( ||\mathbf{S}|| \leq \frac{1}{\omega} \), Lemma 8, and Lemma 11)
\[ \mathbf{E} \left\{ \text{tr} \left( H^H \mathbf{S} \mathbf{M} \mathbf{S} \mathbf{R}_{l,k} \mathbf{S} \mathbf{M}^H \mathbf{S} \mathbf{T}_{l,k} \right) \right\} \leq \frac{2LK}{\omega^4} \| \mathbf{M} \|^2 C_{\max}^2 N, \]
we get
\[ \text{Var} \left( \frac{1}{n_k} \text{tr} \left( \mathbf{S} \mathbf{M} \right) \right) \leq \frac{4L^2K^2}{\omega^4} \| \mathbf{M} \|^2 C_{\max}^2 N = O \left( \frac{1}{N^2} \right). \]
the right-hand side of (95) which is a $O\left(\frac{1}{N^2}\right)$ term. Similar calculations allow to show the second and third terms of the right-hand side of (95) giving the $O\left(\frac{1}{N^2}\right)$ terms. Therefore, we obtain (94a). Similarly, (94b) can be proved and the proof is omitted.

From this lemma, it can be shown that $e = O\left(\frac{1}{N^2}\right)$. In addition, we note that

$$\|\mathcal{S}\|, \|\tilde{\mathcal{S}}\|, \|\mathcal{E}\|, \|\tilde{\mathcal{E}}\|, \|\Theta_l\|, \|\tilde{\Theta}_k\| \leq \frac{1}{\omega}. \quad (102)$$

Using (13), Lemma 8 and Lemma 11, we have

$$[\Gamma_{11}]_{l,k,ij} \geq \begin{cases} \frac{-N_l LKC^3_{\alpha}}{n_k}, & (i,j) \neq (l,k), \\ 1 - \frac{N_l LKC^3_{\alpha}}{n_k}, & (i,j) = (l,k), \end{cases} \quad (103a)$$

$$[\Gamma_{12}]_{l,k,ij} \geq \frac{-C^2_{\alpha}}{\omega}, \quad [\Gamma_{21}]_{l,k,ij} \geq \frac{-N_l C^2_{\alpha}}{n_k}, \quad (103b)$$

$$[\Gamma_{22}]_{l,k,ij} \geq \begin{cases} \frac{-N_l LKC^3_{\alpha}}{n_k}, & (i,j) \neq (l,k), \\ 1 - \frac{N_l LKC^3_{\alpha}}{n_k}, & (i,j) = (l,k). \end{cases} \quad (103c)$$

It is possible to choose $\omega_0$ such that $\omega > \omega_0$ and $\Gamma$ is a strictly diagonally dominant. Thus the eigenvalues of $\Gamma$ are bounded away from 0 [42, Theorem 6.1.10]. It implies that if $\omega > \omega_0$, then $(\alpha_{l,k} - \tau_{l,k})$'s and $(\tilde{\alpha}_{l,k} - \tilde{\tau}_{l,k})$'s are of the same order of magnitude as $O\left(\frac{1}{N^2}\right)$, and therefore converge to 0 when $N \to \infty$.

In the remaining part, we aim to prove that this convergence still holds for $0 < \omega \leq \omega_0$. Firstly, considering $\alpha_{l,k}$ and $\tau_{l,k}$ as functions of the parameter $z = -\omega \in \mathbb{R}^-$, we extend their domain of validity from $\mathbb{R}^-$ to $\mathbb{C} - \mathbb{R}^+$. Similarly to [18, Proposition 11], we have the following lemma.

**Lemma 5** $\alpha_{l,k}$ and $\tau_{l,k}$ are analytic over $\mathbb{C} - \mathbb{R}^+$ and belong to $\mathcal{S}(\mathbb{R}^+)$ with $|\alpha_{l,k}| \leq \frac{\frac{1}{n_k} \text{tr} \mathcal{R}_{l,k}}{d(z,\mathbb{R}^+)} = \frac{N_l P_{l,k}}{n_k(\kappa_{l,k} + 1)d(z,\mathbb{R}^+)}$ and $|\tau_{l,k}| \leq \frac{\frac{1}{n_k} \text{tr} \mathcal{R}_{l,k}}{d(z,\mathbb{R}^+)} \left(1 + \frac{1}{\frac{1}{n_k} \text{tr} \mathcal{R}_{l,k}}\right)$, where $\mathcal{S}(\mathbb{R}^+)$ is the class of all Stieltjes transforms of finite positive measures carried by $\mathbb{R}^+$.

**Proof:** We only prove the results for $\alpha_{l,k}$ since the proof of results of $\tau_{l,k}$ is similar. From the definition of $\mathcal{S}$, $\mathcal{S}$ is invertible for every $z \in \mathbb{C} - \mathbb{R}^+$ and $\mathbb{E}\{\mathcal{S}\}$ is analytic over $\mathbb{C} - \mathbb{R}^+$. Thus $\alpha_{l,k}$ is analytic over $\mathbb{C} - \mathbb{R}^+$. Using the fact that $\mathcal{S} \preceq \frac{1}{d(z,\mathbb{R}^+)} \mathbf{I}_n$ and Lemma 8, we have

$$|\alpha_{l,k}| = \left|\frac{1}{n_k} \text{tr}(\mathcal{R}_{l,k} \mathbb{E}\{\mathcal{S}\})\right| \leq \frac{1}{n_k} \|\mathcal{E}\| \|\mathcal{S}\| \text{tr} \mathcal{R}_{l,k} \leq \frac{\frac{1}{n_k} \text{tr} \mathcal{R}_{l,k}}{d(z,\mathbb{R}^+)} = \frac{N_l P_{l,k}}{n_k(\kappa_{l,k} + 1)d(z,\mathbb{R}^+)},$$

where the last equality is obtained by (4). In order to state $\alpha_{l,k} \in \mathcal{S}(\mathbb{R}^+)$, we only check the following three conditions by [18, Proposition 10]: 1) $\exists \{\alpha_{l,k}(z)\} \geq 0$ if $\exists \{z\} > 0$; 2) $\exists \{z\alpha_{l,k}(z)\} > 0$ if $\exists \{z\} > 0$; 3) $\lim_{y \to \infty} |y\alpha_{l,k}(y)| < \infty$. 

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Let us first compute $\Im \{\alpha_{l,k}(z)\}$: For every $z \in \mathbb{C}^+$,
\[
\Im \{\alpha_{l,k}(z)\} = \Im \left\{ \frac{1}{n_k} \text{tr} \left( \mathbf{R}_{l,k} \mathbf{E} \left\{ \mathbf{S} \left( \mathbf{H}^H \mathbf{H} - z^* \mathbf{I} \mathbf{N} \right) \mathbf{S}^H \right\} \right) \right\} \\
= \Im \left\{ \frac{1}{n_k} \text{tr} \left( \mathbf{R}_{l,k} \mathbf{E} \left\{ \mathbf{S} \mathbf{H} \mathbf{H}^H \mathbf{S}^H \right\} \right) \right\} - \Im \left\{ \frac{1}{n_k} z^* \text{tr} \left( \mathbf{R}_{l,k} \mathbf{E} \left\{ \mathbf{S} \mathbf{S}^H \right\} \right) \right\} \\
= - \frac{1}{n_k} \Im \{z^*\} \text{tr} \left( \mathbf{R}_{l,k} \mathbf{E} \left\{ \mathbf{S} \mathbf{S}^H \right\} \right) > 0.
\]

By similar arguments above, we can prove that $\Im \{z \alpha_{l,k}(z)\} > 0$ if $\Im \{z\} > 0$. Next, we calculate
\[
\lim_{y \to \infty} |jy \alpha_{l,k}(jy)| = \lim_{y \to \infty} \left| \frac{1}{n_k} \text{tr} \left( \mathbf{R}_{l,k} \mathbf{E} \left\{ \left( \frac{1}{jy} \mathbf{H} \mathbf{H}^H - \mathbf{I} \mathbf{N} \right)^{-1} \right\} \right) \right| \\
= \frac{1}{n_k} \text{tr} \mathbf{R}_{l,k} < \infty.
\]

Since the three sufficient conditions have been verified, we have $\alpha_{l,k} \in \mathbb{S}({\mathbb{R}^+})$. \hfill \Box

Using this lemma, $|\alpha_{l,k} - \tau_{l,k}| \leq \frac{N_l P_{l,k}}{n_k (|\alpha_{l,k}| + 1)|d(z, \mathbb{R}^+)|} \left( 2 + |z| L K C_{\text{max}} (d(z, \mathbb{R}^+))^2 \right)$. Moreover, $\{\alpha_{l,k} - \tau_{l,k}\}_{l,k}$ is a family of analytic functions. By Montel’s theorem [43], this convergence still holds for $0 < \omega \leq \omega_0$, and that (49a) and (49b) hold true.

A.2 Proof of Proposition 4

Using the resolvent identity (Lemma 12) $\Xi - \Psi = \Xi (\Psi^{-1} - \Xi^{-1}) \Psi$, we have
\[
\Xi - \Psi = \omega \Xi \text{diag} \left( \left\{ \sum_{k=1}^{K} (\tilde{\alpha}_{l,k} - \tilde{\alpha}_{l,k}) \mathbf{R}_{l,k} \right\}_{\forall l} \right) \Psi + \omega^2 \Xi \text{Hdiag} \left( \left\{ \sum_{l=1}^{L} (\alpha_{l,k} - \beta_{l,k}) \tilde{\Psi}_l \mathbf{T}_{l,k} \mathbf{\Theta}_k \right\}_{\forall k} \right) \mathbf{H}^H \Psi. 
\]

Similarly,
\[
\tilde{\Xi} - \tilde{\Psi} = \omega \tilde{\Xi} \text{diag} \left( \left\{ \sum_{l=1}^{L} (\beta_{l,k} e_{l,k} - \alpha_{l,k}) \mathbf{T}_{l,k} \right\}_{\forall k} \right) \tilde{\Psi} + \omega^2 \tilde{\Xi} \text{Hdiag} \left( \left\{ \sum_{k=1}^{K} (\tilde{\alpha}_{l,k} - \tilde{\alpha}_{l,k}) \tilde{\Xi}_l \mathbf{R}_{l,k} \mathbf{\Theta}_l \right\}_{\forall l} \right) \mathbf{H}^H \tilde{\Psi}. 
\]

Taking the trace, we get
\[
\text{tr} \left( \Xi - \Psi \right) = \omega \sum_{l,k} (\tilde{\alpha}_{l,k} - \tilde{\alpha}_{l,k}) \text{tr} \left( \Xi \mathbf{R}_{l,k} \Psi \right) + \omega^2 \sum_{l,k} (\alpha_{l,k} - \beta_{l,k}) e_{l,k} \text{tr} \left( \Xi \mathbf{H} \mathbf{\Phi} \mathbf{T}_{l,k} \mathbf{\Theta} \mathbf{H}^H \Psi \right), 
\]
\[
\text{tr} \left( \tilde{\Xi} - \tilde{\Psi} \right) = \omega \sum_{l,k} (\beta_{l,k} e_{l,k} - \alpha_{l,k}) \text{tr} \left( \tilde{\Xi} \mathbf{T}_{l,k} \tilde{\Psi} \right) + \omega^2 \sum_{l,k} (\tilde{\alpha}_{l,k} - \tilde{\alpha}_{l,k}) \text{tr} \left( \tilde{\Xi} \mathbf{H} \mathbf{\Phi} \mathbf{R}_{l,k} \mathbf{\Theta} \mathbf{H}^H \tilde{\Psi} \right). 
\]
From Proposition 3, we have

\[
\alpha_{l,k} = \frac{1}{n_k} \text{tr}(R_{l,k} \langle \Xi \rangle_i) + \varepsilon_{l,k}, \quad (107a)
\]

\[
\hat{\alpha}_{l,k} = \frac{1}{n_k} \text{tr}(T_{l,k} \langle \hat{\Xi} \rangle_k) + \tilde{\varepsilon}_{l,k}, \quad (107b)
\]

where \(\varepsilon_{l,k}\) and \(\tilde{\varepsilon}_{l,k}\) converge towards 0. Therefore,

\[
\alpha_{l,k} - \beta_{l,k} e_{l,k} = \frac{1}{n_k} \text{tr}(R_{l,k} \langle \Xi - \Psi \rangle_i) + \varepsilon_{l,k}, \quad (108a)
\]

\[
\hat{\alpha}_{l,k} - \tilde{e}_{l,k} = \frac{1}{n_k} \text{tr}(T_{l,k} \langle \hat{\Xi} - \hat{\Psi} \rangle_k) + \tilde{\varepsilon}_{l,k}. \quad (108b)
\]

Using the same approach as in the proof in Proposition 3, we prove that \((\alpha_{l,k} - \beta_{l,k} e_{l,k})'\)'s and \((\hat{\alpha}_{l,k} - \tilde{e}_{l,k})'\)'s converge towards 0. From (106a) and (106b), we complete the proof of Proposition 4.

### A.3 Proof of Proposition 5

We first establish (52a) and (52b). The equations (89) and (91) can be rewritten as

\[
\hat{\alpha}_{l,k} - \tilde{\tau}_{l,k} = \frac{\omega}{n_k} \sum_{i,j} (\hat{\alpha}_{i,j} - \tilde{\tau}_{i,j}) \text{tr} \left( T_{l,k} E \{ \hat{S} \} T_{i,j} \hat{\Xi} \right) + \frac{\omega^2}{n_k} \text{tr} \left( T_{l,k} \Delta \hat{\Xi} \right) \quad (109a)
\]

\[
\alpha_{l,k} - \tau_{l,k} = \frac{\omega}{n_k} \sum_{i,j} (\alpha_{i,j} - \tau_{i,j}) \text{tr} \left( R_{l,k} E \{ \hat{S} \} R_{i,j} \hat{\Xi} \right) + \frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \Delta \hat{\Xi} \right) \quad (109b)
\]

We can write these two equations in matrix form:

\[
\eta = \Gamma' \eta + \epsilon',
\]
where \( e' \triangleq [\text{vec}(C_1'), \text{vec}(C_2')]^T, \Gamma' \triangleq \begin{bmatrix} \Gamma'_{11} & \Gamma'_{12} \\ \Gamma'_{21} & \Gamma'_{22} \end{bmatrix} \), with \( C_1', C_2' \in \mathbb{C}^{L \times K}, \Gamma'_{11}, \Gamma'_{12}, \Gamma'_{21}, \Gamma'_{22} \in \mathbb{C}^{LK \times LK} \), and

\[
\begin{align*}
[C'_1]_{l,k} &= \frac{\omega^2}{n_k} \text{tr} \left( T_{l,k} \hat{\Xi} \right) + \frac{\omega^2}{n_k} \text{tr} \left( T_{l,k} \Theta H^H \Delta \Xi \Theta \right) \\
&\quad - \frac{1}{n_k} \text{tr} \left( \Gamma_{l,k} \Theta H^H \{S\} R_{l,k} \Xi \Theta \right), \\
[C'_2]_{l,k} &= \frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \Delta \Xi \right) + \frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \Theta H \Delta \Xi H^H \Theta \right) \\
&\quad - \frac{1}{n_k} \text{tr} \left( R_{l,k} \Theta H \{S\} T_{l,k} \Xi H^H \Theta \right),
\end{align*}
\]

\[
\Gamma'_{ij} = \begin{cases} \\
0, & \text{for } (i, j) \neq (l, k); \\
\frac{\omega^2}{n_k} \text{tr} \left( T_{l,k} \Theta H^H \{S\} R_{l,j} \Xi \Theta \right) \\
\frac{1}{n_k} \text{tr} \left( \Gamma_{l,k} \Theta H^H \{S\} R_{l,j} \Xi \Theta \right), & \text{for } (i, j) = (l, k),
\end{cases}
\]

\[
\begin{align*}
[\Gamma'_{11}]_{l,k} &= \frac{\omega^2}{n_k} \text{tr} \left( T_{l,k} \hat{\Xi} \right) \\
&\quad - \frac{1}{n_k} \text{tr} \left( \Gamma_{l,k} \Theta H^H \{S\} T_{l,k} \Xi \Theta \right), \\
[\Gamma'_{12}]_{l,k} &= \frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \hat{\Xi} \right), \\
&\quad - \frac{1}{n_k} \text{tr} \left( R_{l,k} \Theta H \{S\} T_{l,k} \Xi H^H \Theta \right), \\
[\Gamma'_{21}]_{l,k} &= \frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \{S\} \hat{\Xi} \right), \\
&\quad - \frac{1}{n_k} \text{tr} \left( R_{l,k} \Theta H \{S\} T_{l,k} \Xi H^H \Theta \right), \\
[\Gamma'_{22}]_{l,k} &= \frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \Theta H \{S\} T_{l,k} \Xi H^H \Theta \right) \\
&\quad - \frac{1}{n_k} \text{tr} \left( R_{l,k} \Theta H \{S\} T_{l,k} \Xi H^H \Theta \right), \text{ for } (i, j) \neq (l, k);
\end{align*}
\]

For \( (i, j) = (l, k) \), let \( \Gamma'' \) be the matrix by replacing \( E\{S\}, E\{\hat{S}\}, \Xi, \hat{\Xi}, \Theta \) and \( \hat{\Theta} \) in \( \Gamma' \) with \( \Psi, \hat{\Psi}, \Psi, \hat{\Psi}, \Phi \) and \( \hat{\Phi} \), respectively. Using Propositions 3 and 4, we immediately obtain

\[
\Gamma' = \Gamma'' + \delta,
\]

where all entries of \( \delta \) converge to 0 as \( N \to \infty \), and \( \Gamma'' \) is given by

\[
\Gamma'' = \begin{bmatrix} \Gamma''_{11} & \Gamma''_{12} \\ \Gamma''_{21} & \Gamma''_{22} \end{bmatrix},
\]
with $\Gamma''_{11}, \Gamma''_{12}, \Gamma''_{21}, \Gamma''_{22} \in \mathbb{C}^{LK \times LK}$, and

$$
[\Gamma''_{11}]_{l,k,ij} = \begin{cases} 
0, & \text{for } (i, j) \neq (l, k); \\
\frac{u_{lk,ij}^{(2)}}{1 - u_{lk,lk}^{(2)}}, & \text{for } (i, j) = (l, k),
\end{cases}
[\Gamma''_{12}]_{l,k,ij} = \frac{\nu_{lk,ij}^{(1)}}{1 - \nu_{lk,lk}^{(1)}},
$$

(114a)

$$
[\Gamma''_{21}]_{l,k,ij} = \frac{u_{lk,ij}^{(1)}}{1 - u_{lk,lk}^{(2)}}, \quad [\Gamma''_{22}]_{l,k,ij} = \begin{cases} 
0, & \text{for } (i, j) \neq (l, k); \\
\frac{\nu_{lk,ij}^{(2)}}{1 - \nu_{lk,lk}^{(2)}}, & \text{for } (i, j) = (l, k),
\end{cases}
$$

(114b)

$$
u_{lk,ij}^{(1)} = \frac{\omega}{n_k} \text{tr} \left( R_{jk} \Psi \Psi R_{ij} \Psi \right), \quad u_{lk,ij}^{(2)} = \frac{\omega^2}{n_k} \text{tr} \left( T_{pk} \Phi \Phi^H \Psi R_{ij} \Psi \Phi \right),
$$

(114c)

$$v_{lk,ij}^{(1)} = \frac{\omega}{n_k} \text{tr} \left( T_{pk} \Psi \Psi T_{ij} \Psi \right), \quad v_{lk,ij}^{(2)} = \frac{\omega^2}{n_k} \text{tr} \left( R_{jk} \Phi \Phi^H T_{ij} \Psi \Phi \right).
$$

(114d)

**Lemma 6** Let $\Gamma''$ be the matrix defined by (113). Then, we have

$$
\sup_N \left[ \rho \left( \Gamma'' \right) \right] \leq 1 - \frac{\lambda_0 \omega^2}{(\omega + \lambda'_0)^2} < 1,
$$

(115a)

$$
\sup_N \left\| \left( I - \Gamma'' \right)^{-1} \right\|_{\infty} \leq \frac{\left( \omega + \lambda'_0 \right)^2}{\lambda_0 \omega^2},
$$

(115b)

for some constants $\lambda_0, \lambda'_0$.

**Proof:** From (14), a direct calculation yields

$$
\beta_{l,k} e_{l,k} = \frac{1}{n_k} \text{tr} \left( R_{jk} \Psi \Psi^{-1} \Psi \right)
$$

$$
= \frac{1}{n_k} \text{tr} \left( R_{jk} \Psi \left( \omega I_N + \sum_{i,j} \bar{c}_{i,j} R_{ij} + \omega \Phi \Phi^H \Psi \right) \Psi \right)
$$

$$
= \frac{1}{n_k} \text{tr} \left( R_{jk} \Psi \left( \omega I_N + \sum_{i,j} \bar{c}_{i,j} R_{ij} + \omega \Phi \Phi^{-1} \Phi \Psi \Phi^H \right) \Psi \right)
$$

$$
= \frac{\omega}{n_k} \sum_{i,j} \bar{c}_{i,j} \text{tr} \left( R_{jk} \Psi \Psi R_{ij} \Psi \right) + \frac{\omega}{n_k} \text{tr} \left( R_{jk} \Psi \Psi \Phi \Phi^H \Psi \right) + \frac{\omega^2}{n_k} \text{tr} \left( T_{jk} \Psi \Phi \Phi^H T_{ij} \Psi \Phi \right) + \frac{\omega^2}{n_k} \text{tr} \left( T_{jk} \Psi \Phi \Phi^H T_{ij} \Psi \Phi \right),
$$

(116)

where (i) and (ii) are obtained by expanding $\Psi^{-1}$ and $\Phi^{-1}$, respectively. Similarly, we can get

$$
\bar{e}_{l,k} = \frac{\omega}{n_k} \sum_{i,j} \beta_{i,j} e_{i,j} \text{tr} \left( T_{jk} \Psi \Psi R_{ij} \Psi \right) + \frac{\omega}{n_k} \text{tr} \left( T_{jk} \Psi \Psi \Phi \Phi^H \Psi \right) + \frac{\omega^2}{n_k} \text{tr} \left( T_{jk} \Psi \Phi \Phi^H T_{ij} \Psi \Phi \right) + \frac{\omega^2}{n_k} \text{tr} \left( T_{jk} \Psi \Phi \Phi^H T_{ij} \Psi \Phi \right),
$$

(117)
The equations (116) and (117) can be rewritten as

\[
\frac{n_k}{n} \beta_{l,k} e_{l,k} = \sum_{i,j} n_{ij} \tilde{e}_{ij} \omega_{n_j} \operatorname{tr} \left( R_{i,j} \Psi R_{j,k} \Psi \right) + \frac{\omega}{n} \operatorname{tr} \left( R_{i,k} \Psi \Psi \right) \\
+ \sum_{i,j} \frac{n_{ij}}{n} \beta_{ij} e_{ij} \frac{\omega^2}{n_j} \operatorname{tr} \left( T_{ij} \bar{\Psi} H \Psi R_{j,k} \Psi \Phi \right) + \frac{\omega^2}{n} \operatorname{tr} \left( R_{i,k} \Psi \Psi \Phi \Phi H \Psi \right) \\
= \sum_{i,j} n_{ij} \tilde{e}_{ij} u_{ij,lk}^{(1)} + \sum_{i,j} n_{ij} \beta_{ij} e_{ij} u_{ij,lk}^{(2)} + \frac{\omega}{n} \operatorname{tr} \left( R_{i,k} \Psi \Psi \right) + \frac{\omega^2}{n} \operatorname{tr} \left( R_{i,k} \Psi \Psi \Phi \Phi H \Psi \right) 
\]

(118)

and

\[
\frac{n_k}{n} \tilde{e}_{l,k} = \sum_{i,j} n_{ij} \beta_{ij} e_{ij} \frac{\omega}{n_j} \operatorname{tr} \left( T_{ij} \bar{\Psi} T_{j,k} \Psi \right) + \frac{\omega}{n} \operatorname{tr} \left( T_{i,k} \bar{\Psi} \Psi \right) \\
+ \sum_{i,j} \frac{n_{ij}}{n} \tilde{e}_{ij} \frac{\omega^2}{n_j} \operatorname{tr} \left( R_{ij} \Phi \Phi H T_{j,k} \Psi \Phi \Phi \right) + \frac{\omega^2}{n} \operatorname{tr} \left( T_{i,k} \bar{\Psi} H \Phi \Phi \Psi \right) \\
= \sum_{i,j} n_{ij} \beta_{ij} e_{ij} v_{ij,lk}^{(1)} + \sum_{i,j} n_{ij} \tilde{e}_{ij} v_{ij,lk}^{(2)} + \frac{\omega}{n} \operatorname{tr} \left( T_{i,k} \bar{\Psi} \Psi \right) + \frac{\omega^2}{n} \operatorname{tr} \left( T_{i,k} \bar{\Psi} H \Phi \Phi \Psi \right) . 
\]

(119)

Now, let \( \xi \triangleq \left[ \text{vec}(A_3)^T, \text{vec}(A_4)^T \right]^T, \ b \triangleq \left[ \text{vec}(C_3)^T, \text{vec}(C_4)^T \right]^T, \Gamma''' \triangleq \begin{bmatrix} \Gamma_{11}'''' & \Gamma_{12}'''' \\ \Gamma_{21}'''' & \Gamma_{22}'''' \end{bmatrix} \), where \( A_3, A_4, \ C_3, C_4 \in \mathbb{C}^{L \times K}, \Gamma''''_{11}, \Gamma''''_{12}, \Gamma''''_{21}, \Gamma''''_{22} \in \mathbb{C}^{LK \times LK} \) with

\[
[A_3]_{i,k} = \frac{n_k}{n} \beta_{l,k} e_{l,k}, \quad [A_4]_{i,k} = \frac{n_k}{n} \tilde{e}_{l,k}, 
\]

(120a)

\[
[C_3]_{i,k} = \frac{\omega}{n} \operatorname{tr} \left( R_{i,k} \Psi \Psi \right) + \frac{\omega^2}{n} \operatorname{tr} \left( R_{i,k} \Psi \Phi \Phi H \Psi \right), 
1 - u_{ij,ij}^{(2)}
\]

(120b)

\[
[C_4]_{i,k} = \frac{\omega}{n} \operatorname{tr} \left( T_{i,k} \bar{\Psi} \Psi \right) + \frac{\omega^2}{n} \operatorname{tr} \left( T_{i,k} \bar{\Psi} H \Phi \Phi \Psi \right), 
1 - v_{ij,ij}^{(2)}
\]

(120c)

\[
[\Gamma''''_{11}]_{i,k} = \begin{cases} 
0, & (i,j) \neq (l,k); \\
\frac{u_{ij,ij}^{(2)}}{1 - u_{ij,ij}^{(2)}}, & (i,j) = (l,k), \end{cases} \quad [\Gamma''''_{12}]_{i,k} = \frac{u_{ij,ij}^{(1)}}{1 - u_{ij,ij}^{(2)}},
\]

(120d)

\[
[\Gamma''''_{21}]_{i,k} = \frac{v_{ij,ij}^{(1)}}{1 - v_{ij,ij}^{(2)}}, \quad [\Gamma''''_{22}]_{i,k} = \begin{cases} 
0, & (i,j) \neq (l,k); \\
\frac{v_{ij,ij}^{(2)}}{1 - v_{ij,ij}^{(2)}}, & (i,j) = (l,k). \end{cases}
\]

(120e)

Thus, from (118) and (119), we have

\[
\xi = \Gamma'''' \xi + b.
\]

(121)

Using the matrix inversion lemma (Lemma 13), we obtain \( \bar{\Phi} \Phi H \Psi = \bar{\Psi} H \Phi \Phi \). This implies that \( u_{ij,ij}^{(2)} = \)
\(v_{ij,ij}^{(2)}\), for \(\forall i, j\). We immediately get
\[
\Gamma''' = (\Gamma'^n)^T.
\] (122)

Now, define
\[
\Lambda = \text{diag}\left(1 - u_{11,11}^{(2)}, \ldots, 1 - v_{LK,LK}^{(2)}, 1 - v_{11,11}^{(2)}, \ldots, 1 - v_{LK,LK}^{(2)}\right).
\] (123)

Multiplying both sides of (121) by \(\Lambda\) gives
\[
\Lambda \xi = \Lambda \Gamma''' \xi + \Lambda b.
\] (124)

For \(\omega \in \mathbb{R}^+\), the entries of \(\xi, \Lambda \Gamma'''\) and \(\Lambda b\) are positive. Thus, the entries of \(\Lambda \xi\) are positive. Since the entries of \(\xi\) are positive, we conclude that \(1 - u_{ij,ij}^{(2)} > 0\) and \(1 - v_{ij,ij}^{(2)} > 0\), for \(\forall l, k\). From (121), we obtain that the entries of \(\Gamma'''\) and \(b\) are positive, for \(\omega \in \mathbb{R}^+\). Lemma 15 implies \(\rho(\Gamma''') \leq 1 - \frac{\min b}{\max \xi}\).

Using Lemma 8, (13), and the fact that \(\|\Psi\|, \|\tilde{\Psi}\| \leq \frac{1}{\omega}\), we have
\[
\frac{n}{n} \beta_{l,k} e_{l,k} \leq \frac{N_l C_{\text{max}}}{n \omega} \leq \frac{\beta_0 C_{\text{max}}}{\omega}
\] (125)
and
\[
\frac{n}{n} \tilde{e}_{l,k} \leq \frac{n C_{\text{max}}}{n \omega} \leq \frac{C_{\text{max}}}{\omega},
\] (126)
where \(\beta_0 \triangleq \max_{k,l} \{\beta_{l,k}(N)\}\). From (10), we have
\[
\sup_N \max \xi_l \leq \sup_N \max \{1, \beta_0\} \frac{C_{\text{max}}}{\omega} < +\infty.
\] (127)

For \(b_l\), we have
\[
b_{l,k} \geq \frac{\omega}{n} \left(\mathbf{R}_{l,k} \Psi \Psi\right)^{(i)} \frac{\omega}{n} \left(\mathbf{R}_{l,k}\right)^2 (\mathbf{R}_{l,k} \Psi)^{(i)} \geq \frac{\omega}{n} \left(\mathbf{R}_{l,k} \Psi^{-1}\right)^2 \geq \frac{\omega}{n} \left(\mathbf{R}_{l,k} \Psi^{-1}\right)^2 \geq \frac{\omega}{n} \left(\mathbf{R}_{l,k}\right)^2
\] (128)
where \((i)\) and \((ii)\) follow from 1) \(-a)\) of Lemma 8, i.e., \((\mathbf{tr}(\mathbf{A} \mathbf{B}))^2 \leq \mathbf{tr}(\mathbf{A} \mathbf{A}^H) \mathbf{tr}(\mathbf{B} \mathbf{B}^H), (iii)\) is due to 2) of Lemma 8. Similarly,
\[
\tilde{b}_{l,k} \geq \frac{\omega}{n} \left(\mathbf{T}_{l,k}\right)^{\frac{1}{2}} \geq \frac{\omega}{n} \left(\mathbf{T}_{l,k}\right)^{\frac{1}{2}} \geq \frac{\omega}{n} \left(\mathbf{T}_{l,k}\right)^{\frac{1}{2}}.
\] (129)
As a consequence, we have
\[
\inf_N \min b_l \geq \frac{\omega C_5}{(\omega + \sup_N \max \{1, \beta_0\} LKC_{\text{max}}^2 + LKC_{\text{max}}^2)^{\frac{1}{2}}},
\] (130)
where \(C_5 = \inf_N \max \{\frac{1}{n} \mathbf{tr}(\mathbf{T}_{l,k}), \frac{1}{n} \mathbf{tr}(\mathbf{T}_{l,k})\} > 0\).
Combining (127) and (130), we obtain
\[
\sup_N \left[ \rho (\mathbf{Γ}'') \right] \leq 1 - \frac{\lambda_0 \omega^2}{(\omega + \lambda_0')^2} < 1.
\] (131)

According to (122), (115a) holds true. It is easy to get (115b) by \( \rho (\mathbf{Γ}'') < 1 \). A similar proof can be found in [18, 44], and is therefore omitted.

Applying this lemma and (112), there exists \( N_0 \) such that \((\mathbf{I} - \mathbf{Γ}')\) is invertible, for each \( N > N_0 \), and
\[
\sup_{N > N_0} \left[ \left\| (\mathbf{I} - \mathbf{Γ}')^{-1} \right\|_{\infty} \right] \leq \frac{(\omega + \lambda_0')^2}{\lambda_0 \omega^2}. \]

Note that \( \epsilon' = O \left( \frac{1}{N^2} \right) \). Hence, from (110), we obtain \((\alpha_{l,k} - \tau_{l,k})\)’s and \((\tilde{\alpha}_{l,k} - \tilde{\tau}_{l,k})\)’s are of \( O \left( \frac{1}{N^2} \right) \). This establishes (52a) and (52b).

From (52a), (52b), (107a), and (107b), we have \( \epsilon_{l,k} = O \left( \frac{1}{N} \right) \) and \( \tilde{\epsilon}_{l,k} = O \left( \frac{1}{N^2} \right) \). (108a) and (108b) can be rewritten as a matrix form similar to (110). Using the same approach as in the proof of (52a) and (52b), we prove that \((\alpha_{l,k} - \beta_{l,k} \epsilon_{l,k})\)’s and \((\tilde{\alpha}_{l,k} - \tilde{\epsilon}_{l,k})\)’s are of \( O \left( \frac{1}{N} \right) \). This shows that (52c) and (52d) are established and the proof is completed.

\[\text{B Proof of } \mathbb{E}\{m_{\mathcal{B}_N}\} - \mathbb{E}\{m_{\mathcal{B}_N}\} = O \left( \frac{1}{\sqrt{N}} \right) \text{ in Theorem 1}\]

The aim of this appendix is to prove
\[
| \mathbb{E}\{m_{\mathcal{B}_N}(\omega)\} - \mathbb{E}\{m_{\mathcal{B}_N}\} | = O \left( \frac{1}{\sqrt{N}} \right). \] (132)

We mainly make use of the generalized Lindeberg principle given below.

**Lemma 7** (Generalized Lindeberg Principle [30]) Let \( \mathbf{v} = [v_i] \in \mathbb{R}^n \) and \( \mathbf{v} = [\tilde{v}_i] \in \mathbb{R}^n \) be two random vectors with mutually independent components. Define \( \{a_i\}_{1 \leq i \leq n} \) and \( \{b_i\}_{1 \leq i \leq n} \) with
\[
a_i \triangleq | \mathbb{E}\{v_i\} - \mathbb{E}\{\tilde{v}_i\} |, \quad \text{and} \quad b_i \triangleq | \mathbb{E}\{v_i^2\} - \mathbb{E}\{\tilde{v}_i^2\} |. \] (133)

Then, given a twice continuously differentiable function \( f : \mathbb{R}^n \to \mathbb{R} \), we have
\[
| \mathbb{E}\{f(\mathbf{v})\} - \mathbb{E}\{f(\tilde{\mathbf{v}})\} | \leq \sum_{i=1}^{n} a_i \mathbb{E}\{ | \partial f (v_{1}^{i-1}, 0, \mathbf{v}_{i+1}^n) | \} + \frac{1}{2} b_i \mathbb{E}\{ | \partial^2 f (v_{1}^{i-1}, 0, \mathbf{v}_{i+1}^n) | \}
+ \frac{1}{2} \mathbb{E}\left\{ \int_{0}^{v_i} | \partial^2 f (v_{1}^{i-1}, s, \mathbf{v}_{i+1}^n) | (v_i - s)^2 \, ds \right\}
+ \frac{1}{2} \mathbb{E}\left\{ \int_{0}^{\tilde{v}_i} | \partial^2 f (v_{1}^{i-1}, s, \mathbf{v}_{i+1}^n) | (\tilde{v}_i - s)^2 \, ds \right\} , \] (134)

where \( \partial^p_i \) is the \( p \)-fold derivative in the \( i \)-th coordinate, \( v_{1}^{i-1} = (v_1, \ldots, v_{i-1}) \), and \( \tilde{\mathbf{v}}_{i+1} = (\tilde{v}_{i+1}, \ldots, \tilde{v}_n) \).

As \( \mathbf{X}_{l,k} \)’s and \( \mathbf{X}_{l,k} \)’s are matrices with entries satisfying (12), we have \( a_{i,j}^{(l,k)} = b_{i,j}^{(l,k)} = 0 \) for \( i = 1, \ldots, N \) and \( j = 1, \ldots, n \). Therefore, the remaining challenge is to evaluate the third and fourth terms of the
right-hand side of inequality (134). Since the real and imaginary parts of $X_{ij}^{(l,k)}$ are independent, all the results established in the real case can be directly applied for the complex case. Thus, without loss of generality, we only take the derivative with respect to the real part of $X_{ij}^{(l,k)}$ in (134). Before proceeding, we remark that because of the finite 6-th order moment assumption of $X_{ij}^{(l,k)}$’s, the following proof is much simpler than that in [19].

Let

$$f \left( \{A_{l,k}\}_{ql,k} \right) = \frac{1}{N} \tr (G + \omega I_N)^{-1}$$

(135)

where

$$G = \left( \sum_{l,k} \left( R_{l,k}^\frac{1}{2} A_{l,k} T_{l,k}^\frac{1}{2} + \bar{H}_{l,k} \right) \right) \left( \sum_{l,k} \left( R_{l,k}^\frac{1}{2} A_{l,k} T_{l,k}^\frac{1}{2} + \bar{H}_{l,k} \right) \right)^H$$

(136)

for any $A_{l,k} \in \mathbb{R}^{N \times n}$, for $l = 1, \ldots, L$ and $k = 1, \ldots, K$. As such, we have $m_{B_N}(\omega) = f \left( \{X_{l,k}\}_{ql,k} \right)$ and $m_{B_N}(\omega) = f \left( \{A_{l,k}\}_{ql,k} \right)$. To use (134), $\{A_{l,k}\}_{ql,k}$ will take the form $\{A_{l,k} = [A_{ij}^{(l,k)}(l_0, k_0, r, c, s)]\}_{ql,k}$ with

$$A_{l,k}^{(l,k)}(l_0, k_0, r, c, s) = \begin{cases} \frac{X_{ij}^{(l,k)}}{\sqrt{n_k}}, & \text{if } l < l_0, \text{ or } l = l_0, k < k_0, \text{ or } l = l_0, k = k_0, i < r, \\ \text{or } l = l_0, k = k_0, i = r, j < c; \\ s, & \text{if } (l, k) = (l_0, k_0) \text{ and } (i, j) = (r, c); \\ \frac{X_{ij}^{(l,k)}}{\sqrt{n_k}}, & \text{otherwise}. \end{cases}$$

(137)

Taking the third-fold partial derivative of (135) with respect to $A_{ij}^{(l,k)}$, denoted by $\partial_{ij}^{(l,k)}; G$, we have

$$\partial_{ij}^{(l,k)} f = -\frac{6}{N} \tr \left( (\partial_{ij}^{(l,k)} G) (G + \omega I_N)^{-1} (\partial_{ij}^{(l,k)} G) (G + \omega I_N)^{-1} (\partial_{ij}^{(l,k)} G) (G + \omega I_N)^{-1} \right)$$

$$+ \frac{3}{N} \tr \left( (\partial_{ij}^{(l,k)^2} G) (G + \omega I_N)^{-1} (\partial_{ij}^{(l,k)} G) (G + \omega I_N)^{-1} \right)$$

$$+ \frac{3}{N} \tr \left( (\partial_{ij}^{(l,k)} G) (G + \omega I_N)^{-1} (\partial_{ij}^{(l,k)^2} G) (G + \omega I_N)^{-1} \right)$$

(138)

where

$$\partial_{ij}^{(l,k)} G = \left( R_{ij}^\frac{1}{2} E_{ij} R_{ij}^\frac{1}{2} \right) \sum_{l_1,k_1} \left( R_{ij}^\frac{1}{2} A_{l_1,k_1} T_{l_1,k_1}^\frac{1}{2} + \bar{H}_{l_1,k_1} \right)^H$$

$$+ \sum_{l_1,k_1} \left( R_{ij}^\frac{1}{2} A_{l_1,k_1} T_{l_1,k_1}^\frac{1}{2} + \bar{H}_{l_1,k_1} \right) \left( T_{l,k}^\frac{1}{2} E_{ij} R_{ij}^\frac{1}{2} \right),$$

(139a)

$$\partial_{ij}^{(l,k)^2} G = 2 \sum_{l_1,k_1} R_{ij}^\frac{1}{2} E_{ij} R_{ij}^\frac{1}{2}.$$
Using Lemma 8, the first term of \( \partial_{ij}^{(l,k)} f \) can be bounded by
\[
\left| \text{tr} \left( \left( \partial_{ij}^{(l,k)} G \right) \left( G + \omega I_N \right)^{-1} \left( \partial_{ij}^{(l,k)} G \right) \left( G + \omega I_N \right)^{-2} \right) \right| \leq \frac{1}{\omega^4} \left\| \left( \partial_{ij}^{(l,k)} G \right) \right\|_F^3, \tag{140}
\]
and the second and third terms of \( \partial_{ij}^{(l,k)} f \) can be bounded by
\[
\left| \text{tr} \left( \left( \partial_{ij}^{(l,k)2} G \right) \left( G + \omega I_N \right)^{-1} \left( \partial_{ij}^{(l,k)} G \right) \left( G + \omega I_N \right)^{-2} \right) \right| \leq \frac{1}{\omega^4} \left\| \left( \partial_{ij}^{(l,k)2} G \right) \right\|_F \left\| \left( \partial_{ij}^{(l,k)} G \right) \right\|_F. \tag{141}
\]

From (139a) and (139b), using Lemma 8 and (13), we obtain
\[
\left\| \left( \partial_{ij}^{(l,k)} G \right) \right\|_F \leq \sum_{l_1,k_1} \left( \left\| R_{l_1,k_1} \right\|_F \left\| E_{l_1,k_1} T_{l_1,k_1} H_{l_1,k_1} \right\|_F \right) \tag{i}
\]
and
\[
\left\| \left( \partial_{ij}^{(l,k)2} G \right) \right\|_F = 2 \sum_{l_1,k_1} \left( \left\| R_{l_1,k_1} \right\|_F \left\| E_{l_1,k_1} T_{l_1,k_1} H_{l_1,k_1} \right\|_F \right) \tag{ii}
\]
where (i) is obtained by the triangle inequality of the Frobenius norm, (ii) follows from 1(b) of Lemma 8 and (13), and (iii) follows from 2 of Lemma 8 and (13). Combining everything together, we get
\[
\mathbb{E} \left\{ \partial_{rc}^{(l_0,k_0)} f \right\} \leq \frac{C_1}{N} \mathbb{E} \left\{ \left( \sum_{l,k} \left( \sum_{i=1}^N \left( A_{ic}^{(l,k)} \right)^2 \right)^{\frac{1}{2}} \right)^3 + C_2 \right\} \tag{142}
\]
and
\[
\left\| \left( \partial_{ij}^{(l,k)2} G \right) \right\|_F = 2 \sum_{l_1,k_1} \left( \left\| R_{l_1,k_1} \right\|_F \left\| E_{l_1,k_1} T_{l_1,k_1} H_{l_1,k_1} \right\|_F \right) \leq 2 \left\| T_{l,k} \right\|_F \left\| R_{l,k} \right\|_F \leq 2 C_{max}, \tag{143}
\]
where (i) is obtained by the triangle inequality of the Frobenius norm, (ii) follows from 1(b) of Lemma 8 and (13), and (iii) follows from 2 of Lemma 8 and (13). Combining everything together, we get
\[
\mathbb{E} \left\{ \partial_{rc}^{(l_0,k_0)} f \right\} \leq \frac{C_1}{N} \mathbb{E} \left\{ \left( \sum_{l,k} \left( \sum_{i=1}^N \left( A_{ic}^{(l,k)} \right)^2 \right)^{\frac{1}{2}} + C_2 \right)^3 \right\} \tag{144}
\]

\[
\leq \frac{C_1(LK + 1)^2}{N} \mathbb{E} \left\{ \left( \sum_{l,k} \left( \sum_{i=1}^N \left( A_{ic}^{(l,k)} \right)^2 \right)^{\frac{1}{2}} + C_2 \right)^3 \right\} \tag{i}
\]
and
\[
\mathbb{E} \left\{ \partial_{rc}^{(l_0,k_0)} f \right\} \leq \frac{C_3}{N} \left( \left| s \right|^3 + \mathbb{E} \left\{ \left( \sum_{i \neq r} \left( A_{ic}^{(l_0,k_0)} \right)^2 \right)^{\frac{1}{2}} + \mathbb{E} \left\{ \sum_{l \neq l_0} \sum_{k \neq k_0} \left( \sum_{i=1}^N \left( A_{ic}^{(l,k)} \right)^2 \right)^{\frac{1}{2}} \right\} + C_2 \right) \right\} \tag{ii}
\]
and
\[
\leq \frac{C_3}{N} \left( \left| s \right|^3 + C_4 \right), \tag{iii}
\]

\[44\]
where $C_1, C_2, C_3, C_4$ denote constants, (i) is obtained by Lemma 9, (ii) follows from the definition of $A_{ij}^{(l,k)}$ (137), and (iii) is due to the fact that $X_{ij}^{(l,k)}$ and $X_{ij}^{(l,k)}$ have finite 6-th order moment, thus giving the second and third terms of third line of (144) as $O(1)$.

Finally, using (134) and (144), we obtain

$$|E\{\Im\{m_{B_N}(\omega)\}\} - E\{\Im\{m_{B_N}(\omega)\}\}| \leq C_3^2 \sum_{l,k} \sum_{r=1}^N \sum_{c=1}^n \left(E\left\{\int_0^{\frac{|X_{rc}^{(l,k)}|/\sqrt{nk}}}{|s|^3 + C_4} \left(\frac{X_{rc}^{(l,k)}}{\sqrt{nk}} - s\right)^2 ds\right\}ight)
+ E\left\{\int_0^{\frac{|X_{rc}^{(l,k)}|/\sqrt{nk}}}{|s|^3 + C_4} \left(\frac{X_{rc}^{(l,k)}}{\sqrt{nk}} - s\right)^2 ds\right\}
\leq C_3 \frac{1}{2N} \sum_{l,k} \sum_{r=1}^N \sum_{c=1}^n \left(\frac{|X_{rc}^{(l,k)}|}{\sqrt{nk}}\right)^6 + C_4 E\left\{\left(\frac{|X_{rc}^{(l,k)}|}{\sqrt{nk}}\right)^3\right\}
+ \frac{1}{6} E\left\{\left(\frac{|X_{rc}^{(l,k)}|}{\sqrt{nk}}\right)^6\right\} + \frac{C_4}{3} E\left\{\left(\frac{|X_{rc}^{(l,k)}|}{\sqrt{nk}}\right)^3\right\}
= O\left(\frac{1}{\sqrt{N}}\right). \quad (145)$$

The quantity $|E\{\Im\{m_{B_N}(\omega)\}\} - E\{\Im\{m_{B_N}(\omega)\}\}|$ also admits the same upper bound. Thus, (132) is true.

C Existence and Uniqueness

C.1 Existence

Following [16] and using Proposition 3 the existence of $(e_{l,k}, \tilde{e}_{l,k})_{\forall l,k}$ can be shown.

C.2 Uniqueness

Let $(e_{l,k}, \tilde{e}_{l,k})$ and $(e_{l,k}^0, \tilde{e}_{l,k}^0)$ be two solutions satisfying (14), and $\Psi^0, \tilde{\Psi}^0, \Phi^0, \tilde{\Phi}^0$ be the matrices obtained by replacing $e_{l,k}(\omega)$’s and $\tilde{e}_{l,k}(\omega)$’s in $\Psi, \tilde{\Psi}, \Phi, \tilde{\Phi}$ with $e_{l,k}^0(\omega)$’s and $\tilde{e}_{l,k}^0(\omega)$’s respectively. To prove the uniqueness, we need to show that $e_{l,k} - e_{l,k}^0 = 0$ and $\tilde{e}_{l,k} - \tilde{e}_{l,k}^0 = 0$, for any $l$ and $k$. Our proof is inspired by [18].
A standard calculation involving Lemma 11 yields

\[
\beta_{l,k}(e_{l,k} - e_{l,k}^0) = -\frac{\omega}{n_k} \sum_{i,j} (\tilde{e}_{i,j} - \tilde{e}_{i,j}^0) \text{tr} \left( R_{l,k} \Phi H \Psi\tilde{R}_{l,k} \Phi^0 H^0 \Psi^0 \right) \\
+ \frac{\omega^2}{n_k} \sum_{i,j} \beta_{i,j}(e_{i,j} - e_{i,j}^0) \text{tr} \left( R_{l,k} \Phi H \Psi\tilde{T}_{l,k} \Phi^0 H^0 \Psi^0 \right),
\]

(146a)

\[
\tilde{e}_{l,k} - \tilde{e}_{l,k}^0 = -\frac{\omega}{n_k} \sum_{i,j} \beta_{i,j}(e_{i,j} - e_{i,j}^0) \text{tr} \left( T_{l,k} \Phi H \Psi\tilde{T}_{l,k} \Phi^0 H^0 \Psi^0 \right) \\
+ \frac{\omega^2}{n_k} \sum_{i,j} (\tilde{e}_{i,j} - \tilde{e}_{i,j}^0) \text{tr} \left( T_{l,k} \Phi H \Psi\tilde{T}_{l,k} \Phi^0 H^0 \Psi^0 \right).
\]

(146b)

Now, let \( \zeta \triangleq [\text{vec}(A_5)^T, \text{vec}(A_6)^T]^T \), \( \Pi \triangleq \begin{bmatrix} \Pi_{11} & \Pi_{12} \\ \Pi_{21} & \Pi_{22} \end{bmatrix} \), where \( A_5, A_6 \in \mathbb{C}^{L \times K} \), \( \Pi_{11}, \Pi_{12}, \Pi_{21}, \Pi_{22} \in \mathbb{C}^{L \times K} \) with

\[
[A_5]_{l,k} = \beta_{l,k}(e_{l,k} - e_{l,k}^0), \quad [A_6]_{l,k} = \tilde{e}_{l,k} - \tilde{e}_{l,k}^0,
\]

(147a)

\[
[I_{ij}]_{l,k} = \begin{cases} 
0, & \text{for } (i, j) \neq (l, k); \\
\frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \Phi H \Psi\tilde{T}_{l,k} \Phi^0 H^0 \Psi^0 \right) \\
1 - \frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \Phi H \Psi\tilde{T}_{l,k} \Phi^0 H^0 \Psi^0 \right), & \text{for } (i, j) = (l, k),
\end{cases}
\]

(147b)

\[
[I_{ij}]_{l,k} = \begin{cases} 
\frac{\omega}{n_k} \text{tr} \left( R_{l,k} \Phi H \Psi\tilde{R}_{l,k} \Phi^0 H^0 \Psi^0 \right) \\
1 - \frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \Phi H \Psi\tilde{R}_{l,k} \Phi^0 H^0 \Psi^0 \right),
\end{cases}
\]

(147d)

\[
[I_{ij}]_{l,k} = \begin{cases} 
0, & \text{for } (i, j) \neq (l, k); \\
\frac{\omega^2}{n_k} \text{tr} \left( T_{l,k} \Phi H \Psi\tilde{T}_{l,k} \Phi^0 H^0 \Psi^0 \right) \\
1 - \frac{\omega^2}{n_k} \text{tr} \left( T_{l,k} \Phi H \Psi\tilde{T}_{l,k} \Phi^0 H^0 \Psi^0 \right), & \text{for } (i, j) = (l, k).
\end{cases}
\]

(147e)

Thus, (146a) and (146b) can be written together as

\[
\zeta = \Pi \zeta.
\]

(148)

To complete the proof, it remains to prove that \( \rho(\Pi) < 1 \). To do so, we first write (116) and (117) in matrix form as follows:

\[
\xi' = K \xi' + b',
\]

(149)

where \( \xi' = [\text{vec}(A_7)^T, \text{vec}(A_8)^T]^T \), \( b' = [\text{vec}(C_7)^T, \text{vec}(C_8)^T]^T \), \( K = \begin{bmatrix} K_{11} & K_{12} \\ K_{21} & K_{22} \end{bmatrix} \), and \( A_7, A_8, C_7, C_8 \in \mathbb{C}^{L \times K} \).
\( \mathcal{C}^{L \times K}, K_{11}, K_{12}, K_{21}, K_{22} \in \mathcal{C}^{LK \times LK} \) with

\[
[A_7]_{l,k} = \beta l, k e_{l,k}, \quad [A_8]_{l,k} = \beta l, k, \tag{150a}
\]

\[
[C_7]_{l,k} = \frac{\omega}{n_k} \text{tr} \left( R_{k,k} \Psi \Phi \right) + \frac{\omega^2}{n_k} \text{tr} \left( R_{k,k} \Psi \Phi \Phi \Psi \Phi \Phi \Psi \right), \quad 1 - u_{l,k,lk}^{(2)}, \tag{150b}
\]

\[
[C_8]_{l,k} = \frac{\omega}{n_k} \text{tr} \left( T_{l,k} \Phi \Phi \right) + \frac{\omega^2}{n_k} \text{tr} \left( T_{l,k} \Phi \Phi \Phi \Phi \Psi \Phi \Psi \right), \quad 1 - v_{l,k,lk}^{(2)}. \tag{150c}
\]

\[
[K_{11}]_{l,k,i,j} = \begin{cases} 
0, & \text{for } (i, j) \neq (l, k); \\
\frac{u_{l,k,lk}^{(2)}}{1 - u_{l,k,lk}^{(2)}}, & \text{for } (i, j) = (l, k), \quad [K_{12}]_{l,k,i,j} = \frac{u_{l,k,lk}^{(2)}}{1 - u_{l,k,lk}^{(2)}}, \tag{150d}
\end{cases}
\]

\[
[K_{21}]_{l,k,i,j} = \frac{v_{l,k,lk}^{(2)}}{1 - v_{l,k,lk}^{(2)}}, \quad [K_{22}]_{l,k,i,j} = \begin{cases} 
0, & \text{for } (i, j) \neq (l, k); \\
\frac{v_{l,k,lk}^{(2)}}{1 - v_{l,k,lk}^{(2)}}, & \text{for } (i, j) = (l, k), \tag{150e}
\end{cases}
\]

\[
u_{l,k,lk}^{(1)} = \frac{\omega}{n_k} \text{tr} \left( R_{k,k} \Psi \rho R_{l,l} \right), \quad \nu_{l,k,lk}^{(2)} = \frac{\omega^2}{n_k} \text{tr} \left( R_{k,k} \Psi \Phi \tilde{R}_{l,l} \Phi \Psi \phi \Psi \phi \Psi \right), \tag{150f}
\]

\[
\nu_{l,k,lk}^{(1)} = \frac{\omega}{n_k} \text{tr} \left( T_{l,k} \Phi \phi \right), \quad \nu_{l,k,lk}^{(2)} = \frac{\omega^2}{n_k} \text{tr} \left( T_{l,k} \Phi \phi \phi \Psi \phi \phi \phi \phi \Psi \right). \tag{150g}
\]

Using a similar approach of (121), we get that \( 1 - u_{l,k,lk}^{(2)} > 0, 1 - v_{l,k,lk}^{(2)} > 0, \forall l, k, \) and the entries of \( \xi', K \) and \( b' \) are positive, for \( \omega \in \mathbb{R}^+ \). Therefore, from (149) and Lemma 16, we have \( \rho(K) < 1 \). Similarly, we also have \( \rho(K^0) < 1 \), where \( K^0 \) as well as \( K^0_{11}, K^0_{12}, K^0_{21}, \) and \( K^0_{22} \) are the matrices by replacing \( \Psi, \Psi, \Phi, \) and \( \Phi \) with \( \Psi^0, \Psi^0, \Phi^0, \) and \( \Phi^0 \), respectively.

For the denominator of \( [\Pi_{11}]_{l,k,i,j} \), applying Lemma 17 with \( A = \omega \sqrt{\text{tr} R_{l,l} \Psi \Phi \tilde{R}_{l,l} \Phi \Psi} \) and \( B = \omega \sqrt{\text{tr} T_{l,k} \Phi \phi \phi \Psi \phi \phi \phi \phi \Psi} \) satisfying \( \text{tr}(AA^H) = u_{l,k,lk}^{(2)} < 1 \) and \( \text{tr}(BB^H) = v_{l,k,lk}^{(2)} < 1 \), we have

\[
1 - \frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \Psi \Phi \tilde{R}_{l,k} \Phi \phi \phi \Psi \phi \phi \phi \phi \Psi \right) \geq \left( 1 - \frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \Psi \Phi \tilde{R}_{l,k} \Phi \phi \phi \Psi \phi \phi \phi \phi \Psi \right) \right)^\frac{1}{2} \left( 1 - \frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \Psi \Phi \tilde{R}_{l,k} \Phi \phi \phi \Psi \phi \phi \phi \phi \Psi \right) \right)^\frac{1}{2}. \tag{151}
\]

Applying the Cauchy-Schwarz inequality to the numerator of \( [\Pi_{11}]_{l,k,i,j} \) and from (151), we obtain

\[
[\Pi_{11}]_{l,k,i,j} \leq \left( \frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \Psi \Phi \tilde{R}_{l,k} \Phi \phi \phi \Psi \phi \phi \phi \phi \Psi \right) \right)^\frac{1}{2} \left( 1 - \frac{\omega^2}{n_k} \text{tr} \left( R_{l,k} \Psi \Phi \tilde{R}_{l,k} \Phi \phi \phi \Psi \phi \phi \phi \phi \Psi \right) \right)^\frac{1}{2} \left( [K_{11}]_{l,k,i,j} \right)^\frac{1}{2} \left( [K_{11}]_{l,k,i,j} \right)^\frac{1}{2}, \tag{152}
\]
Likewise, we have

\[ \| \Pi_{12} \|_{k,ij} \leq \| K_{12} \|_{k,ij}^{\frac{1}{2}} \| K_{12}^\circ \|_{k,ij}^{\frac{1}{2}}, \tag{153a} \]

\[ \| \Pi_{21} \|_{k,ij} \leq \| K_{21} \|_{k,ij}^{\frac{1}{2}} \| K_{21}^\circ \|_{k,ij}^{\frac{1}{2}}, \tag{153b} \]

\[ \| \Pi_{22} \|_{k,ij} \leq \| K_{22} \|_{k,ij}^{\frac{1}{2}} \| K_{22}^\circ \|_{k,ij}^{\frac{1}{2}}. \tag{153c} \]

Using Lemma 18 and Lemma 19, we obtain

\[ \rho(\Pi) \leq \rho(\Pi) \leq \rho(K)^{\frac{1}{2}} \rho(K^\circ)^{\frac{1}{2}} < 1. \tag{154} \]

This contradicts to the statement that \( \Pi \) has an eigenvalue equal to 1. Therefore, we have \( e_{\omega,l,k} - e_{\omega,l,k}^\circ = 0 \) and \( \tilde{e}_{\omega,l,k} - e_{\omega,l,k}^\circ = 0 \), for any \( l, k \) and \( \omega \in \mathbb{R}^+ \).

\section*{D Mathematical Tools}

In this appendix, we provide some mathematical tools used in the proof of the appendices.

\textbf{Lemma 8} \cite{45}

1. Let \( A = [A_{ij}] \) and \( B \) be any matrices such that the product is a square matrix. Then,

   \( (a) \quad | \text{tr}(AB) | \leq \| A \|_F \| B \|_F, \)

   \( (b) \quad \| AB \|_F \leq \| A \|_F \| B \|_F, \)

   \( (c) \quad \| AB \|_F \leq \| A \|_F \| B \|_F, \)

   \( (d) \quad | A_{ij} | \leq \| A \|. \)

2. If \( A \) is nonnegative definite, we have \( | \text{tr}(AB) | \leq \| B \| \text{tr}(A). \)

3. Let \( A \) be any matrix such that the product \( AB \) exists. Then, \( \| AB \| \leq \| A \| \| B \|. \)

\textbf{Lemma 9} For any \( p \geq 1 \) and real numbers \( a_i \)'s, we have

\[ \left| \sum_{i=1}^{n} a_i \right|^p \leq n^{p-1} \sum_{i=1}^{n} |a_i|^p. \tag{155} \]

\textbf{Lemma 10} \cite[Theorem 4.3.1]{42} Let \( A \) and \( B \) be Hermitian matrix and let the eigenvalues \( \lambda_i(A), \lambda_i(B), \) and \( \lambda_i(A + B) \) be arranged in decreasing order. For each \( k = 1, 2, \ldots, n \), we have

\[ \lambda_k(A) + \lambda_n(B) \leq \lambda_k(A + B) \leq \lambda_k(A) + \lambda_1(B). \tag{156} \]

\textbf{Lemma 11} Let matrix \( A_{l,k} \in \mathbb{C}^{N_l \times n_k} \) for \( l = 1, \ldots, L, k = 1, \ldots, K \), and let \( A_k = \left[ A_{l,k}^T \cdots A_{L,k}^T \right]^T \in \mathbb{C}^{N \times n}, A = [A_1, \ldots, A_K] \in \mathbb{C}^{N \times n} \), with \( N = \sum_{l=1}^{L} N_l \) and \( n = \sum_{k=1}^{K} n_k \). If \( \| A_{l,k} A_{l,k}^H \| \leq C \), then we have \( \| AA^H \| \leq LKC. \)
Proof: Notice that $AA^H$ and $A_{l,k}A_{l,k}^H$ are Hermitian matrices. Therefore, a standard computation involving Lemma 10 yields

$$
\|AA^H\| = \lambda_1 (AA^H) = \lambda_1 \left( \sum_{k=1}^K A_k A_k^H \right)
\leq \sum_{k=1}^K \lambda_1 (A_k A_k^H) = K \sum_{k=1}^K \lambda_1 (A_k^H A_k)
= K \sum_{k=1}^L \left( \sum_{l=1}^L A_{l,k} A_{l,k}^H \right) \leq \sum_{l,k} \lambda_1 (A_{l,k}^H A_{l,k})
= \sum_{l,k} \|A_{l,k}^H A_{l,k}\| \leq LKC.
$$

(157)

Lemma 12 (Resolvent Identity) For invertible $A$ and $B$ matrices, we have the identity

$$
A^{-1} - B^{-1} = A^{-1}(B - A)B^{-1}.
$$

(158)

Lemma 13 (Matrix Inversion) For invertible $A, B$ and $R$ matrices, suppose that $B = A + XRY$, then

$$
B^{-1} = A^{-1} - A^{-1}X(R^{-1} + YA^{-1}X)^{-1}YA^{-1}.
$$

Lemma 14 Assume that $A$ is a positive semi-definite $M \times M$ matrix and $B = \text{diag}(B_1, \ldots, B_K)$ is a block-diagonal matrix, where $B_k$ is a positive semi-definite $M_k \times M_k$ matrix and $M = \sum_{k=1}^K M_k$. Let $C_k = \langle (I + AB)_{k}^{-1}A \rangle_k$, where $B_{k} = \text{diag}(B_1, \ldots, B_{k-1}, 0, B_{k+1}, \ldots, B_K)$. Then, we have

$$
\langle (I + AB)^{-1} A \rangle_k = (I + C_k B_k)^{-1} C_k.
$$

(159)

Proof: Letting $B_k = \text{diag}(0, \ldots, 0, B_k, 0, \ldots, 0)$, we have

$$(I + AB)^{-1} A = (I + AB_{\setminus k} + AB_k)^{-1} A$$

\[\begin{align*}
\overset{(i)}{=} & C - C(I + B_k C)^{-1} B_k C = C(I - (B_k C)^{-1} + I)^{-1} \\
\overset{(ii)}{=} & C (B_k C)^{-1} ((B_k C)^{-1} + I)^{-1} = (I + CB_k)^{-1} C
\end{align*}
\]

(160)

where $(i)$ follows from Lemma 13 and defining $C = (I + AB_{\setminus k})^{-1} A$, $(ii)$ is due to Lemma 12. Substituting (160) into (159), we obtain

$$
\langle (I + AB)^{-1} A \rangle_k = \langle (I + CB_k)^{-1} C \rangle_k = (I + C_k B_k)^{-1} C_k,
$$

where the last step is obtained by calculating the inverse of $(I + CB_k)^{-1}$.  \[\square\]
Lemma 15 [42, Corollary 8.1.29] Let $A \in \mathbb{R}^{n \times n}$, $x \in \mathbb{R}^n$, for $A \succeq 0$ and $x > 0$. If $\alpha, \beta \geq 0$ are such that $\alpha x \leq Ax \leq \beta x$, then $\alpha \leq \rho(A) \leq \beta$. If $\alpha x < Ax$, then $\alpha < \rho(A)$. If $Ax < \beta x$, then $\rho(A) < \beta$.

Lemma 16 [16, Lemma 9] If the components of $C, x$, and $b$ are all positive, then $x = Cx + b$ implies $\rho(C) < 1$.

Lemma 17 [19, Lemma 16] Let $A$ and $B$ be any matrices such that $AB^H$ exists and is a squared matrix. If $\text{tr}(AA^H) \leq 1$ and $\text{tr}(BB^H) \leq 1$, then

$$|1 - \text{tr}(AB^H)| \geq (1 - \text{tr}(AA^H))^{\frac{1}{2}} (1 - \text{tr}(BB^H))^{\frac{1}{2}}. \quad (161)$$

Lemma 18 [42, Theorem 8.1.18] Let $A = [A_{ij}]$ and $B = [B_{ij}]$ be square matrices. If $|A_{ij}| \leq B_{ij}, \forall i, j$, then $\rho(A) \leq \rho(|A|) \leq \rho(B)$.

Lemma 19 [45, Lemma 5.7.9] Let $A = [A_{ij}]$ and $B = [B_{ij}]$ be matrices with nonnegative elements. Then $\rho([\frac{1}{2} A_{ij}], B_{ij}]) \leq \rho(A)^{\frac{1}{2}} \rho(B)^{\frac{1}{2}}$.

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