A Central Limit Theorem for the SINR at the LMMSE Estimator Output for Large-Dimensional Signals

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Abstract—This paper is devoted to the performance study of the linear minimum mean squared error (LMMSE) estimator for multidimensional signals in the large-dimension regime. Such an estimator is frequently encountered in wireless communications and in array processing, and the signal-to-interference-plus-noise ratio (SINR) at its output is a popular performance index. The SINR can be modeled as a random quadratic form which can be studied with the help of large random matrix theory, if one assumes that the dimension of the received and transmitted signals go to infinity at the same pace. This paper considers the asymptotic behavior of the SINR for a wide class of multidimensional signal models that includes general multiple-antenna as well as spread-spectrum transmission models.

The expression of the deterministic approximation of the SINR in the large-dimension regime is recalled and the SINR fluctuations around this deterministic approximation are studied. These fluctuations are shown to converge in distribution to the Gaussian law in the large-dimension regime, and their variance is shown to decrease as the inverse of the signal dimension.

Index Terms—Antenna arrays, central limit theorem, code-division multiple access (CDMA), linear minimum mean squared error (LMMSE), martingales, multiple-carrier (MC)-CDMA, multiple-input multiple-output (MIMO), random matrix theory.

I. INTRODUCTION

LARGE Random Matrix Theory (LRMT) is a powerful mathematical tool used to study the performance of multiuser and multiple-access communication systems such as multiple-input multiple-output (MIMO) digital wireless systems, antenna arrays for source detection and localization, spread spectrum communication systems as code-division multiple access (CDMA) and multiple-carrier CDMA (MC-CDMA) systems. In most of these communication systems, the N-dimensional received random vector $\boldsymbol{r} \in \mathbb{C}^N$ is described by the model

$$\boldsymbol{r} = \boldsymbol{\Sigma}\boldsymbol{s} + \boldsymbol{n} \tag{1}$$

where $\boldsymbol{s} = [s_0, s_1, \ldots, s_K]^T$ is the unknown random vector of transmitted symbols with size $K+1$ satisfying $\mathbb{E}\boldsymbol{s}\boldsymbol{s}^* = \boldsymbol{I}_{K+1}$, the noise $\boldsymbol{n}$ is an independent additive white Gaussian noise

(AWGN) with covariance matrix $\mathbb{E}\boldsymbol{n}\boldsymbol{n}^* = \rho\boldsymbol{I}_N$ whose variance $\rho > 0$ is known, and matrix $\boldsymbol{\Sigma}$ represents the known "channel" in the wide sense whose structure depends on the particular system under study. One typical problem addressed by LRMT concerns the estimation performance by the receiver of a given transmitted symbol, say s_0.

In this paper, we focus on one of the most popular estimators, namely, the linear Wiener estimator, also called the linear minimum mean squared error LMMSE for linear minimum mean squared error estimator: the LMMSE estimate $\hat{s}_0 = \boldsymbol{g}^*\boldsymbol{r}$ of signal s_0 is the one for which the $N \times 1$ vector $\boldsymbol{g}$ minimizes $\mathbb{E}|\hat{s}_0 - s_0|^2$. If we partition the channel matrix as $\boldsymbol{\Sigma} = [\boldsymbol{y}\,\boldsymbol{Y}]$, where $\boldsymbol{y}$ is the first column of $\boldsymbol{\Sigma}$ and where matrix $\boldsymbol{Y}$ has dimensions $N \times K$, then it is well known that vector $\boldsymbol{g}$ is given by $\boldsymbol{g} = (\boldsymbol{\Sigma}\boldsymbol{\Sigma}^* + \rho\boldsymbol{I}_N)^{-1}\boldsymbol{y}$. Usually, the performance of this estimator is evaluated in terms of the signal-to-interference-plus-noise ratio (SINR) at its output. Writing the received vector $\boldsymbol{r}$ as $\boldsymbol{r} = s_0\boldsymbol{y} + \boldsymbol{r}_{\text{in}}$ where $s_0\boldsymbol{y}$ is the relevant term and $\boldsymbol{r}_{\text{in}}$ represents the so-called interference-plus-noise term, the SINR is given by $\beta_K = |\boldsymbol{g}^*\boldsymbol{y}|^2/\mathbb{E}|\boldsymbol{g}^*\boldsymbol{r}_{\text{in}}|^2$. Plugging the expression of $\boldsymbol{g}$ given above into this expression, one can prove that the SINR β_K is given by the well-known expression

$$\beta_K = \boldsymbol{y}^* \left(\boldsymbol{Y}\boldsymbol{Y}^* + \rho\boldsymbol{I}_N\right)^{-1} \boldsymbol{y}. \tag{2}$$

In general, this expression does not provide a clear insight into the impact of the channel model parameters (such as the load factor KN^{-1}, the power distribution of the transmission data streams, or the correlation structure of the channel paths in the context of multiple-antenna transmissions) on the performance of the LMMSE estimator.

An alternative approach, justified by the fluctuating nature of the channel paths in the context of MIMO communications and by the pseudorandom nature of the spreading sequences in spread-spectrum applications consists to model matrix $\boldsymbol{\Sigma}$ as a random matrix (in this case, β_K becomes a random SINR). The simplest random matrix model for $\boldsymbol{\Sigma}$, corresponding to the most canonical MIMO or CDMA transmission channels, corresponds to independent and identically distributed (i.i.d.) entries with mean zero and variance N^{-1}. In that case, LRMT shows that when $K \to \infty$ and the load factor KN^{-1} converges to a limiting load factor $\alpha > 0$, the SINR β_K converges almost surely (a.s.) to an explicit deterministic quantity $\overline{\beta}(\alpha, \rho)$ which simply depends on the limiting load factor α and on the noise variance ρ. As a result, the impact of these two parameters on the LMMSE performance can be easily evaluated [1], [2].

The LMMSE SINR large-dimensional behavior for more sophisticated random matrix models has also been thoroughly studied (cf. [1], [3]–[9]) and it has been proved that there

exists a deterministic sequence $(\overline{\beta}_K)$, generally defined as the solution of an implicit equation, such that $\beta_K - \overline{\beta}_K \to 0$ a.s. as $K \to \infty$ and $\frac{K}{N}$ remains bounded away from zero and from infinity.

Beyond the convergence $\beta_K - \overline{\beta}_K \to 0$, a natural question arises concerning the accuracy of $\overline{\beta}_K$ for finite values of K. A first answer to this question consists in evaluating the Mean Squared Error (MSE) of the SINR $\mathbb{E}|\beta_K - \overline{\beta}_K|^2$ for large K. A further problem is the computation of outage probability, that is the probability for $\beta_K - \overline{\beta}_K$ to be below a certain level. Both problems can be addressed by establishing a Central Limit Theorem (CLT) for $\beta_K - \overline{\beta}_K$. In this paper, we establish such a CLT (Theorem 3 below) for a large class of random matrices Σ. We prove that there exists a sequence $\Theta_K^2 = \mathcal{O}(1)$ such that $\frac{\sqrt{K}}{\Theta_K}(\beta_K - \overline{\beta}_K)$ converges in distribution to the standard normal law $\mathcal{N}(0,1)$ in the asymptotic regime. One can therefore infer that the mean square error (MSE) asymptotically behaves like $\frac{\Theta_K^2}{K}$ and that the outage probability can be simply approximated by a Gaussian tail function.

The class of random matrices Σ we consider in this paper is described by the following statistical model: Assume that

$$\Sigma = \left(\Sigma_{nk}\right)_{n=1,\,k=0}^{N,K} = \left(\frac{\sigma_{nk}}{\sqrt{K}} W_{nk}\right)_{n=1,\,k=0}^{N,K} \quad (3)$$

where the complex random variables W_{nk} are i.i.d. with $\mathbb{E}W_{nk} = 0$, $\mathbb{E}W_{nk}^2 = 0$, and $\mathbb{E}|W_{nk}|^2 = 1$, and where $(\sigma_{nk}^2;\ 1 \le n \le N;\ 0 \le k \le K)$ is an array of real numbers. Due to the fact that $\mathbb{E}|\Sigma_{nk}|^2 = \frac{\sigma_{nk}^2}{K}$, the array (σ_{nk}^2) is referred to as a variance profile. An important particular case is when σ_{nk}^2 is *separable* , that is, writes

$$\sigma_{nk}^2 = d_n \tilde{d}_k \quad (4)$$

where $(d_1, \ldots, d_N)$ and $(\tilde{d}_0, \ldots, \tilde{d}_K)$ are two vectors of real positive numbers.

Applicative Contexts:

Among the applicative contexts where the channel is described appropriately by model (3) or by its particular case (4), let us mention the following.

- Multiple-antenna transmissions with $K + 1$ distant sources sending their signals toward an array of N antennas. The corresponding transmission model is $r = \Xi s + n$ where $\Xi = \frac{1}{\sqrt{K}} H P^{1/2}$, matrix H is an $N \times (K + 1)$ random matrix with complex Gaussian elements representing the radio channel, $P = \text{diag}(p_0, \ldots, p_K)$ is the (deterministic) matrix of the powers given to the different sources, and n is the usual AWGN satisfying $\mathbb{E}nn^* = \rho I_N$. Write $H = [h_0 \cdots h_K]$, and assume that the columns h_k are independent, which is realistic when the sources are distant one from another. Let C_k be the covariance matrix $C_k = \mathbb{E}h_k h_k^*$ and let $C_k = U_k \Lambda_k U_k$ be a spectral decomposition of C_k where $\Lambda_k = \text{diag}(\lambda_{nk};\ 1 \le n \le N)$ is the matrix of eigenvalues. Assume now that the eigenvector matrices $U_0, \ldots, U_K$ are all equal (to some matrix U, for instance), a case considered in, e.g., [10] (note that sometimes they are all identified with the Fourier $N \times N$ matrix

[11]). Let $\Sigma = U^* \Xi$. Then matrix Σ is described by the statistical model (3) where the W_{nk} are standard Gaussian i.i.d., and $\sigma_{nk}^2 = \lambda_{nk} p_k$. If we partition Ξ as $\Xi = [x\ X]$ similarly to the partition $\Sigma = [y\ Y]$ above, then the SINR β at the output of the LMMSE estimator for the first element of vector s in the transmission model $r = \Xi s + n$ is

$$\beta = x^* \left(XX^* + \rho I_N\right)^{-1} x = y^* \left(YY^* + \rho I_N\right)^{-1} y$$

due to the fact that U is a unitary matrix. Therefore, the problem of LMMSE SINR convergence for this MIMO model is a particular case of the general problem of convergence of the right-hand member of (2) for model (3).

It is also worth to say a few words about the particular case (4) in this context. If we assume that $\Lambda_0 = \cdots = \Lambda_K$ and these matrices are equal to $\Lambda = \text{diag}(\lambda_1, \ldots, \lambda_N)$, then the model for H is the well-known Kronecker model with correlations at reception [12]. In this case

$$\Sigma = U^* \Xi = \frac{1}{\sqrt{K}} U^* H P^{1/2} = \frac{1}{\sqrt{K}} \Lambda^{1/2} W P^{1/2} \quad (5)$$

where W is a random matrix with i.i.d. standard Gaussian elements. This model coincides with the separable variance profile model (4) with $d_n = \lambda_n$ and $\tilde{d}_k = p_k$.

- CDMA transmissions on flat-fading channels. Here N is the spreading factor, $K + 1$ is the number of users, and

$$\Sigma = V P^{1/2} \quad (6)$$

where V is the $N \times (K + 1)$ signature matrix assumed here to have random i.i.d. elements with mean zero and variance N^{-1}, and where $P = \text{diag}(p_0, \ldots, p_K)$ is the user's power matrix. In this case, the variance profile is separable with $d_n = 1$ and $\tilde{d}_k = \frac{K}{N} p_k$. Note that elements of V are not Gaussian in general.

- Cellular MC-CDMA transmissions on frequency-selective channels. In the uplink direction, the matrix Σ is written as

$$\Sigma = [H_0 v_0 \cdots H_{K+1} v_{K+1}] \quad (7)$$

where $H_k = \text{diag}(h_k(\exp(2\imath\pi(n-1)/N);\ 1 \le n \le N)$ is the radio channel matrix of user k ($\imath = \sqrt{-1}$) in the discrete Fourier domain (here N is the number of frequency bins) and $V = [v_0, \ldots, v_K]$ is the $N \times (K + 1)$ signature matrix with i.i.d. elements as in the CDMA case above. Modeling this time the channel transfer functions as deterministic functions, we have

$$\sigma_{nk}^2 = \frac{K}{N} |h_k(\exp(2\imath\pi(n-1)/N))|^2.$$

In the downlink direction, we have

$$\Sigma = H V P^{1/2} \quad (8)$$

where $H = \text{diag}(h(\exp(2\imath\pi(n-1)/N);\ 1 \le n \le N)$ is the radio channel matrix in the discrete Fourier domain, the $N \times (K + 1)$ signature matrix V is as above, and $P = \text{diag}(p_0, \ldots, p_K)$ is the matrix of the powers given to

the different users. Model (8) coincides with the separable variance profile model(4) with

$$d_n = \frac{K}{N}|h(\exp(2\imath\pi(n-1)/N))|^2$$

and $d_k = p_k$.

About the Literature:

The asymptotic approximation $\overline{\beta}_K$ (first-order result) is connected with the asymptotic eigenvalue distribution of Gram matrices YY^* where elements of Y are described by the model (3), and can be found in the mathematical LRMT literature in the work of Girko [13] (see also [14] and [15]). Applications in the field of wireless communications can be found in, e.g., [6] in the separable case and in [8] in the general variance profile case.

Concerning the CLT for $\beta_K - \overline{\beta}_K$ (second-order result), only some particular cases of the general model (3) have been considered in the literature among which the i.i.d. case ($\sigma_{nk}^2 = 1$) is studied in [16] (and based on a result of [17] pertaining to the asymptotic behavior of the eigenvectors of YY^*). The more general CDMA model (6) has been considered in [18], using a result of [19]. The model used in this paper includes the models of [16] and [18] as particular cases.

Approximations of the distribution of the SINR at the LMMSE output have been studied in [20]. The authors of [20] propose to approximate the SINR distribution with Gamma and generalized Gamma distributions by adjusting the asymptotic moments of this SINR.

In another line of thought, performance of nonlinear detectors has been studied in the large-dimensional regime by using statistical mechanics techniques, as in [21], [22].

Fluctuations of other performance indexes such as Shannon's mutual information $\mathbb{E}\log\det\left(\frac{\Sigma\Sigma^*}{\rho} + I_N\right)$ have also been studied at length. Let us cite [23] where the CLT is established in the separable case and [24] for a CLT in the general variance profile case. Similar results concerning the mutual information are found in [25] and [26].

Limiting Expressions Versus K-Dependent Expressions:

As one may check in Theorems 2 and 3 below, we deliberately chose to provide deterministic expressions $\overline{\beta}_K$ and Θ_K^2 which remain bounded but do not necessarily converge as $K \to \infty$. For instance, Theorem 2 only states that $\beta_K - \overline{\beta}_K \to 0$ a.s. No conditions which would guarantee the convergence of β_K are added. This approach has two advantages: 1) such expressions for $\overline{\beta}_K$ and Θ_K^2 exist for very general variance profiles (σ_{nk}^2) while limiting expressions may not, and 2) they provide a natural discretization which can easily be implemented.

Statements about these deterministic approximations are valid within the following asymptotic regime:

$$K \to \infty, \quad \liminf \frac{K}{N} > 0 \quad \text{and} \quad \limsup \frac{K}{N} < \infty. \quad (9)$$

Note that $\frac{K}{N}$ is not required to converge. In the remainder of the paper, the notation "$K \to \infty$" will refer to (9).

We note that in the particular case where $\frac{K}{N} \to \alpha > 0$ and the variance profile is obtained by a regular sampling of a continuous function f i.e., $\sigma_{nk}^2 = f\left(\frac{n}{N}, \frac{k}{K+1}\right)$, it is possible to prove that $\overline{\beta}_K$ and Θ_K^2 converge towards limits that can be characterized by integral equations.

Principle of the Approach:

The approach used here is simple and powerful. It is based on the approximation of β_K by the sum of a martingale difference sequence and on the use of the CLT for martingales [27]. We note that apart from the LRMT context, such a technique has been used recently in [28] to establish a CLT on general quadratic forms of the type $z^* A z$ where A is a deterministic matrix and z is a random vector with i.i.d. elements.

Paper Organization:

In Section II, first-order results, whose presentation and understanding are necessary to state the CLT, are recalled. The CLT, which is the main contribution of this paper, is provided in Section III. In Section IV, simulations and numerical illustrations are provided. The proof of the main theorem (Theorem 3) in given in Section V while the Appendix gathers proofs of intermediate results.

Notations:

Given a complex $N \times N$ matrix $X = [x_{ij}]_{i,j=1}^N$, denote by $\|X\|$ its spectral norm, and by $\||X|\|_\infty$ its maximum row sum norm, i.e., $\||X|\|_\infty = \max_{1 \le i \le N} \sum_{j=1}^N |x_{ij}|$. Denote by $\|\cdot\|$ the Euclidean norm of a vector and by $\|\cdot\|_\infty$ its max (or ℓ_∞) norm.

II. FIRST-ORDER RESULTS: THE SINR DETERMINISTIC APPROXIMATION

In the sequel, we shall often show explicitly the dependence on K in the notations. Consider the quadratic form (2)

$$\beta_K = y^* (YY^* + \rho I_N)^{-1} y$$

where the sequence of matrices $\Sigma(K) = [y(K)\, Y(K)]$ is given by

$$\Sigma(K) = (\Sigma_{nk}(K))_{n=1,\,k=0}^{N,K} = \left(\frac{\sigma_{nk}(K)}{\sqrt{K}} W_{nk}\right)_{n=1,\,k=0}^{N,K}.$$

Let us state the main assumptions:

A1: The complex random variables $(W_{nk};\ n \ge 1,\ k \ge 0)$ are i.i.d. with $\mathbb{E}W_{10} = 0$, $\mathbb{E}W_{10}^2 = 0$, $\mathbb{E}|W_{10}|^2 = 1$, and $\mathbb{E}|W_{10}|^8 < \infty$.

A2: There exists a real number $\sigma_{\max} < \infty$ such that

$$\sup_{K \ge 1} \max_{\substack{1 \le n \le N \\ 0 \le k \le K}} |\sigma_{nk}(K)| \le \sigma_{\max}.$$

Let $(a_m; 1 \le m \le M)$ be complex numbers, then $\operatorname{diag}(a_m; 1 \le m \le M)$ refers to the $M \times M$ diagonal matrix whose diagonal elements are the a_m's. If $A = (a_{ij})$ is a square matrix, then $\operatorname{diag}(A)$ refers to the matrix $\operatorname{diag}(a_{ii})$. Consider

the following diagonal matrices based on the variance profile along the columns and the rows of $\boldsymbol{\Sigma}$:

$$\boldsymbol{D}_k(K) = \operatorname{diag}(\sigma_{1k}^2(K), \ldots, \sigma_{Nk}^2(K)), \quad 0 \le k \le K$$
$$\widetilde{\boldsymbol{D}}_n(K) = \operatorname{diag}(\sigma_{n1}^2(K), \ldots, \sigma_{nK}^2(K)), \quad 1 \le n \le N. \quad (10)$$

A3: The variance profile satisfies

$$\liminf_{K \ge 1} \min_{0 \le k \le K} \frac{1}{K} \operatorname{tr} \boldsymbol{D}_k(K) > 0.$$

Since $\mathbb{E}|W_{10}|^2 = 1$, one has $\mathbb{E}|W_{10}|^4 \ge 1$. The following assumption is needed.

A4: At least one of the following conditions is satisfied:

$$\mathbb{E}|W_{10}|^4 > 1 \quad \text{or} \quad \liminf_K \frac{1}{K^2} \operatorname{tr} \left(\boldsymbol{D}_0(K) \sum_{k=1}^K \boldsymbol{D}_k(K) \right) > 0.$$

Remark 1: If needed, one can attenuate the assumption on the eighth moment in **A1**. For instance, one can adapt without difficulty the proofs in this paper to the case where $\mathbb{E}|W_{10}|^{4+\epsilon} < \infty$ for $\varepsilon > 0$. We assumed $\mathbb{E}|W_{10}|^8 < \infty$ because at some places we rely on results of [24] which are stated with the assumption on the eighth moment.

Assumption **A3** is technical. It has already appeared in [29].

Assumption **A4** is necessary to get a nonvanishing variance Θ_K^2 in Theorem 3.

The following definitions will be of help in the sequel. A complex function $t(z)$ belongs to class $\mathcal{S}$ if $t(z)$ is analytical in the upper half plane $\mathbb{C}_+ = \{z \in \mathbb{C} \, ; \, \operatorname{im}(z) > 0\}$, if $t(z) \in \mathbb{C}_+$ for all $z \in \mathbb{C}_+$ and if $\operatorname{im}(z)|t(z)|$ is bounded over the upper half plane $\mathbb{C}_+$.

The *Stieltjes transform* (ST) f_μ of a probability measure μ is the complex function

$$f_\mu(z) = \int \frac{1}{t-z} \mu(dt), \quad z \in \mathbb{C}_+.$$

One can check that the ST of a probability measure belongs to the class $\mathcal{S}$.

Let $\{\lambda_{1,K}, \ldots, \lambda_{N,K}\}$ be the set of eigenvalues of the Gram matrix $\boldsymbol{Y}(K)\boldsymbol{Y}(K)^*$. The *spectral measure* of this matrix is the random probability measure μ_K defined as

$$\mu_K = \frac{1}{N} \sum_{n=1}^N \boldsymbol{\delta}_{\lambda_{n,K}}$$

where $\boldsymbol{\delta}_x$ is the Dirac measure at x. Denote by $\boldsymbol{Q}_K(z)$ and $\widetilde{\boldsymbol{Q}}_K(z)$ the *resolvents* of $\boldsymbol{Y}(K)\boldsymbol{Y}(K)^*$ and $\boldsymbol{Y}(K)^*\boldsymbol{Y}(K)$, respectively, that is, the $N \times N$ and $K \times K$ matrices defined by

$$\boldsymbol{Q}_K(z) = (\boldsymbol{Y}(K)\boldsymbol{Y}(K)^* - z\boldsymbol{I}_N)^{-1} \quad \text{and}$$
$$\widetilde{\boldsymbol{Q}}_K(z) = (\boldsymbol{Y}(K)^*\boldsymbol{Y}(K) - z\boldsymbol{I}_K)^{-1}.$$

We have

$$\frac{1}{N} \operatorname{tr} \boldsymbol{Q}_K(z) = \frac{1}{N} \sum_{n=1}^N \frac{1}{\lambda_{n,K} - z} = \int \frac{1}{t-z} \mu_K(dt) = f_{\mu_K}(z)$$

in other words, the normalized trace of the resolvent is the ST of the spectral measure of $\boldsymbol{Y}(K)\boldsymbol{Y}(K)^*$.

A. The SINR Deterministic Approximation

It is known [13], [29] that there exists a deterministic diagonal $N \times N$ matrix function $\boldsymbol{T}(z)$ with the following properties: First, the normalized trace $\frac{1}{N} \operatorname{tr} \boldsymbol{T}(z)$ is the ST of a deterministic probability measure π_K. Second, $\boldsymbol{T}(z)$ approximates the resolvent $\boldsymbol{Q}(z)$ in the following sense: Given any diagonal deterministic matrix $\boldsymbol{S}$ with bounded spectral norm, the quantity $\frac{1}{K} \operatorname{tr} \boldsymbol{S}(\boldsymbol{Q}(z) - \boldsymbol{T}(z))$ converges a.s. to zero as $K \to \infty$. By this means (take $\boldsymbol{S} = \boldsymbol{I}$ and consider the STs), one can show that π_K approximates the spectral measure μ_K of $\boldsymbol{Y}(K)\boldsymbol{Y}(K)^*$.

It is also known that the approximation $\overline{\beta}_K$ of the SINR β_K is simply related to $\boldsymbol{T}(z)$ (cf. Theorem 2). As we shall see, matrix $\boldsymbol{T}(z)$ also plays a fundamental role in the second-order result (Theorem 3). In the following theorem, we recall the definition and some of the main properties of $\boldsymbol{T}(z)$.

Theorem 1: The following statements hold true.
1) [29, Theorem 2.4] Let $(\sigma_{nk}^2(K); \, 1 \le n \le N; \, 1 \le k \le K)$ be a sequence of arrays of real numbers and consider the matrices $\boldsymbol{D}_k(K)$ and $\widetilde{\boldsymbol{D}}_n(K)$ defined in (10). The system of $N + K$ functional equations

$$\begin{cases} t_{n,K}(z) = \dfrac{-1}{z \left(1 + \frac{1}{K} \operatorname{tr}(\widetilde{\boldsymbol{D}}_n(K)\widetilde{\boldsymbol{T}}_K(z))\right)}, & 1 \le n \le N \\ \tilde{t}_{k,K}(z) = \dfrac{-1}{z \left(1 + \frac{1}{K} \operatorname{tr}(\boldsymbol{D}_k(K)\boldsymbol{T}_K(z))\right)}, & 1 \le k \le K \end{cases} \quad (11)$$

where

$$\boldsymbol{T}_K(z) = \operatorname{diag}(t_{1,K}(z), \ldots, t_{N,K}(z))$$
$$\widetilde{\boldsymbol{T}}_K(z) = \operatorname{diag}(\tilde{t}_{1,K}(z), \ldots, \tilde{t}_{K,K}(z))$$

admits a unique solution $(\boldsymbol{T}, \widetilde{\boldsymbol{T}})$ among the diagonal matrices for which the $t_{n,K}$'s and the $\tilde{t}_{k,K}$'s belong to class $\mathcal{S}$. Moreover, functions $t_{n,K}(z)$ and $\tilde{t}_{k,K}(z)$ admit analytical continuations over $\mathbb{C} - [0, \infty)$ which are real and positive for $z \in (-\infty, 0)$.
2) [29, Theorem 2.4] There exist probability measures π_K and $\tilde{\pi}_K$ with support in $\mathbb{R}_+ = [0, \infty)$ such that

$$f_{\pi_K}(z) = \frac{1}{N} \operatorname{tr} \boldsymbol{T}(z) \quad \text{and} \quad f_{\tilde{\pi}_K}(z) = \frac{1}{K} \operatorname{tr} \widetilde{\boldsymbol{T}}(z).$$

3) [29, Theorem 2.5] Assume from now on that Assumptions **A1** and **A2** hold true. Consider the sequence of random matrices $\boldsymbol{Y}(K)\boldsymbol{Y}(K)^*$ where $\boldsymbol{Y}$ has dimensions $N \times K$ and whose entries are given by $Y_{nk} = \frac{\sigma_{nk}}{\sqrt{K}} W_{nk}$. For every sequence $\boldsymbol{S}_K$ of $N \times N$ diagonal matrices and every sequence $\widetilde{\boldsymbol{S}}_K$ of $K \times K$ diagonal matrices with

$$\sup_K \max \left(\|\boldsymbol{S}_K\|, \|\widetilde{\boldsymbol{S}}_K\| \right) < \infty$$

the following limits hold true a.s.:

$$\lim_{K \to \infty} \frac{1}{K} \operatorname{tr} \boldsymbol{S}_K (\boldsymbol{Q}_K(z) - \boldsymbol{T}_K(z)) = 0,$$
$$\forall z \in \mathbb{C} - \mathbb{R}_+$$
$$\lim_{K \to \infty} \frac{1}{K} \operatorname{tr} \widetilde{\boldsymbol{S}}_K \left(\widetilde{\boldsymbol{Q}}_K(z) - \widetilde{\boldsymbol{T}}_K(z) \right) = 0,$$
$$\forall z \in \mathbb{C} - \mathbb{R}_+.$$

4) [29, Corollary 2.7] Denote by μ_K and $\tilde{\mu}_K$ the spectral measures of $\boldsymbol{Y}(K)\boldsymbol{Y}(K)^*$ and $\boldsymbol{Y}(K)^*\boldsymbol{Y}(K)$, respec-

tively. Then for any bounded and continuous function $g : \mathbb{R}_+ \to \mathbb{R}$

$$\int g(t)\mu_K(dt) - \int g(t)\pi_K(dt) \xrightarrow[K \to \infty]{} 0 \quad \text{and}$$

$$\int g(t)\tilde{\mu}_K(dt) - \int g(t)\tilde{\pi}_K(dt) \xrightarrow[K \to \infty]{} 0$$

almost surely.

The following lemma which reproduces [30, Lemma 2.7] will be used throughout the paper. It characterizes the asymptotic behavior of an important class of quadratic forms.

Lemma 1: Let $\boldsymbol{x} = [X_1, \ldots, X_N]^T$ be an $N \times 1$ vector where the X_n are centered i.i.d. complex random variables with unit variance. Let $\boldsymbol{A}$ be a deterministic $N \times N$ complex matrix. Then, for any $p \geq 2$, there exists a constant C_p depending on p only such that

$$\mathbb{E}\left| \frac{1}{N}\boldsymbol{x}^*\boldsymbol{A}\boldsymbol{x} - \frac{1}{N}\text{tr}(\boldsymbol{A}) \right|^p$$
$$\leq \frac{C_p}{N^p}\left(\left(\mathbb{E}|X_1|^4\text{tr}(\boldsymbol{A}\boldsymbol{A}^*) \right)^{p/2} + \mathbb{E}|X_1|^{2p}\text{tr}\left((\boldsymbol{A}\boldsymbol{A}^*)^{p/2} \right) \right). \tag{12}$$

Noticing that $\text{tr}(\boldsymbol{A}\boldsymbol{A}^*) \leq N\|\boldsymbol{A}\|^2$ and that $\text{tr}\left((\boldsymbol{A}\boldsymbol{A}^*)^{p/2} \right) \leq N\|\boldsymbol{A}\|^p$, we obtain the simpler inequality

$$\mathbb{E}\left| \frac{1}{N}\boldsymbol{x}^*\boldsymbol{A}\boldsymbol{x} - \frac{1}{N}\text{tr}(\boldsymbol{A}) \right|^p$$
$$\leq \frac{C_p}{N^{p/2}}\|\boldsymbol{A}\|^p \left(\left(\mathbb{E}|X_1|^4\right)^{p/2} + \mathbb{E}|X_1|^{2p} \right) \tag{13}$$

which is useful in case one has bounds on $\|\boldsymbol{A}\|$.

Using Theorem 1 and Lemma 1, we are in position to characterize the asymptotic behavior of the quadratic form β_K given by (2). We begin by rewriting β_K as

$$\beta_K = \frac{1}{K}\boldsymbol{w}_0^*\boldsymbol{D}_0^{1/2}\left(\boldsymbol{Y}\boldsymbol{Y}^* + \rho\boldsymbol{I}_N\right)^{-1}\boldsymbol{D}_0^{1/2}\boldsymbol{w}_0$$
$$= \frac{1}{K}\boldsymbol{w}_0^*\boldsymbol{D}_0^{1/2}\boldsymbol{Q}(-\rho)\boldsymbol{D}_0^{1/2}\boldsymbol{w}_0 \tag{14}$$

where the $N \times 1$ vector $\boldsymbol{w}_0$ is given by $\boldsymbol{w}_0 = [W_{10}, \ldots, W_{N0}]^T$ and the diagonal matrix $\boldsymbol{D}_0$ is given by (10). Recall that $\boldsymbol{w}_0$ and $\boldsymbol{Q}$ are independent and that $\|\boldsymbol{D}_0\| \leq \sigma_{\max}^2$ by **A2**. Furthermore, one can easily notice that $\|\boldsymbol{Q}(-\rho)\| = \|(\boldsymbol{Y}\boldsymbol{Y}^* + \rho\boldsymbol{I})^{-1}\| \leq 1/\rho$.

Denote by $\mathbb{E}_{\boldsymbol{Q}}$ the conditional expectation with respect to $\boldsymbol{Q}$, i.e., $\mathbb{E}_{\boldsymbol{Q}} = \mathbb{E}(\cdot \| \boldsymbol{Q})$. From inequality (13), there exists a constant $C > 0$ for which

$$\mathbb{E}\mathbb{E}_{\boldsymbol{Q}}\left| \beta_K - \frac{1}{K}\text{tr}\boldsymbol{D}_0\boldsymbol{Q}(-\rho) \right|^4$$
$$\leq \frac{C}{K^2}\left(\frac{N}{K} \right)^2 \mathbb{E}\|\boldsymbol{D}_0\boldsymbol{Q}\|^4 \left((\mathbb{E}|W_{10}|^4)^2 + \mathbb{E}|W_{10}|^8 \right)$$
$$\leq \frac{C}{K^2}\left(\frac{N}{K} \right)^2 \left(\frac{\sigma_{\max}^2}{\rho} \right)^4 \left((\mathbb{E}|W_{10}|^4)^2 + \mathbb{E}|W_{10}|^8 \right)$$
$$= \mathcal{O}\left(\frac{1}{K^2} \right).$$

By the Borel–Cantelli lemma, we therefore have

$$\beta_K - \frac{1}{K}\text{tr}(\boldsymbol{D}_0\boldsymbol{Q}(-\rho)) \xrightarrow[K \to \infty]{} 0 \quad \text{a.s.}$$

Using this result, simply apply Theorem 1–part 3 with $\boldsymbol{S} = \boldsymbol{D}_0$ (recall that $\|\boldsymbol{D}_0\| \leq \sigma_{\max}^2$) to obtain the following.

Theorem 2: Let $\overline{\beta}_K = \frac{1}{K}\text{tr}(\boldsymbol{D}_0(K)\boldsymbol{T}_K(-\rho))$ where $\boldsymbol{T}_K$ is given by Theorem 1–part 1. Assume **A1** and **A2**. Then

$$\beta_K - \overline{\beta}_K \xrightarrow[K \to \infty]{} 0 \quad \text{a.s.}$$

B. The Deterministic Approximation in the Separable Case

In the separable case $\sigma_{nk}(K) = d_n(K)\tilde{d}_k(K)$, matrices $\boldsymbol{D}_k(K)$ and $\widetilde{\boldsymbol{D}}_n(K)$ are written as $\boldsymbol{D}_k(K) = \tilde{d}_k(K)\boldsymbol{D}(K)$ and $\widetilde{\boldsymbol{D}}_n(K) = d_n(K)\widetilde{\boldsymbol{D}}(K)$ where $\boldsymbol{D}(K)$ and $\widetilde{\boldsymbol{D}}(K)$ are the diagonal matrices

$$\boldsymbol{D}(K) = \text{diag}(d_1(K), \ldots, d_N(K))$$
$$\widetilde{\boldsymbol{D}}(K) = \text{diag}(\tilde{d}_1(K), \ldots, \tilde{d}_K(K)) \tag{15}$$

and one can check that the system of $N + K$ equations leading to $\boldsymbol{T}_K$ and $\widetilde{\boldsymbol{T}}_K$ simplifies into a system of two equations, and Theorem 1 takes the following form.

Proposition 1 [29, Sec. 3.2]:
1) Assume $\sigma_{nk}^2(K) = d_n(K)\tilde{d}_k(K)$. Given $\rho > 0$, the system of two equations

$$\begin{cases} \delta_K(\rho) = \frac{1}{K}\text{tr}\left(\boldsymbol{D}\left(\rho(\boldsymbol{I}_N + \tilde{\delta}_K(\rho)\boldsymbol{D}) \right)^{-1} \right) \\ \tilde{\delta}_K(\rho) = \frac{1}{K}\text{tr}\left(\widetilde{\boldsymbol{D}}\left(\rho(\boldsymbol{I}_K + \delta_K(\rho)\widetilde{\boldsymbol{D}}) \right)^{-1} \right) \end{cases} \tag{16}$$

where $\boldsymbol{D}$ and $\widetilde{\boldsymbol{D}}$ are given by (15) admits a unique solution $(\delta_K(\rho), \tilde{\delta}_K(\rho))$. Moreover, in this case, matrices $\boldsymbol{T}(-\rho)$ and $\widetilde{\boldsymbol{T}}(-\rho)$ provided by Theorem 1–part 1 coincide with

$$\boldsymbol{T}(-\rho) = \frac{1}{\rho}(\boldsymbol{I} + \tilde{\delta}(\rho)\boldsymbol{D})^{-1} \quad \text{and}$$
$$\widetilde{\boldsymbol{T}}(-\rho) = \frac{1}{\rho}(\boldsymbol{I} + \delta(\rho)\widetilde{\boldsymbol{D}})^{-1}. \tag{17}$$

2) Assume that **A1** and **A2** hold true. Let matrices $\boldsymbol{S}_K$ and $\widetilde{\boldsymbol{S}}_K$ be as in Theorem 1–part 3. Then, almost surely $\frac{1}{K}\text{tr}\left(\boldsymbol{S}_K\left(\boldsymbol{Q}_K(-\rho) - \boldsymbol{T}_K(-\rho) \right) \right) \to 0$ and $\frac{1}{K}\text{tr}\left(\widetilde{\boldsymbol{S}}_K\left(\widetilde{\boldsymbol{Q}}_K(-\rho) - \widetilde{\boldsymbol{T}}_K(-\rho) \right) \right) \to 0$ as $K \to \infty$.

With these equations we can adapt the result of Theorem 2 to the separable case. Notice that $\boldsymbol{D}_0 = \tilde{d}_0\boldsymbol{D}$ and that $\delta(\rho)$ given by the system (16) coincides with $\frac{1}{K}\text{tr}(\boldsymbol{D}\boldsymbol{T})$, hence we have the following.

Proposition 2: Assume that $\sigma_{nk}^2(K) = d_n(K)\tilde{d}_k(K)$, and that **A1** and **A2** hold true. Then

$$\frac{\beta_K}{\tilde{d}_0} - \delta_K(\rho) \xrightarrow[K \to \infty]{} 0 \quad \text{a.s.}$$

where $\delta_K(\rho)$ is given by Proposition 1–part 1.

Let us provide a more explicit expression of δ_K which will be used in Section IV to illustrate the SINR behavior for the MIMO model (5) and for MC-CDMA downlink model (8). By

combining the two equations in system (16), it turns out that $\delta = \delta_K(\rho)$ is the unique solution of the implicit equation

$$\delta = \frac{1}{K} \sum_{n=1}^{N} \frac{d_n}{\rho + \frac{1}{K} d_n \sum_{k=1}^{K} \frac{p_k}{1+p_k\delta}}. \quad (18)$$

Recall that in the case of the MIMO model (5), $d_n = \lambda_n$ and $\tilde{d}_k = p_k$, while in the case of the MC-CDMA downlink model (8), $d_n = \frac{K}{N}|h(\exp(2\iota\pi(n-1)/N)|^2$ and $\tilde{d}_k = p_k$ again. Here $\tilde{d}_0 = p_0$ is the power of the user of interest (user 0), and therefore $\beta_K/\tilde{d}_0$ is the normalized SINR of this user. Notice that $\delta_K(\rho)$ is almost the same for all users, hence the normalized SINRs for all users are close to each other for large K. Their common deterministic approximation is given by (18) which is the discrete analogue of the integral (16) in [6].

This example will be continued in Section III.

III. SECOND-ORDER RESULTS: THE CENTRAL LIMIT THEOREM

The following theorem is the main result of this paper. Its proof is postponed to Section V.

Theorem 3:

1) Assume that **A2**, **A3**, and **A4** hold true. Let $\boldsymbol{A}_K$ and $\boldsymbol{\Delta}_K$ be the $K \times K$ matrices

$$\boldsymbol{A}_K = \left[\frac{1}{K} \frac{\frac{1}{K}\mathrm{tr}\boldsymbol{D}_\ell\boldsymbol{D}_m\boldsymbol{T}(-\rho)^2}{\left(1+\frac{1}{K}\mathrm{tr}\boldsymbol{D}_\ell\boldsymbol{T}(-\rho)\right)^2} \right]_{\ell,m=1}^{K} \quad \text{and}$$

$$\boldsymbol{\Delta}_K = \mathrm{diag}\left(\left(1+\frac{1}{K}\mathrm{tr}\boldsymbol{D}_\ell\boldsymbol{T}(-\rho)\right)^2 ; 1 \leq \ell \leq K \right) \quad (19)$$

where $\boldsymbol{T}$ is defined in Theorem 1–part 1. Let $\boldsymbol{g}_K$ be the $K \times 1$ vector

$$\boldsymbol{g}_K = \left[\frac{1}{K}\mathrm{tr}\boldsymbol{D}_0\boldsymbol{D}_1\boldsymbol{T}(-\rho)^2, \ldots, \frac{1}{K}\mathrm{tr}\boldsymbol{D}_0\boldsymbol{D}_K\boldsymbol{T}(-\rho)^2 \right]^T.$$

Then the sequence of real numbers

$$\Theta_K^2 = \frac{1}{K}\boldsymbol{g}_K^T(\boldsymbol{I}_K - \boldsymbol{A}_K)^{-1}\boldsymbol{\Delta}_K^{-1}\boldsymbol{g}_K$$
$$+ (\mathbb{E}|W_{10}|^4 - 1)\frac{1}{K}\mathrm{tr}\boldsymbol{D}_0(K)^2\boldsymbol{T}_K(-\rho)^2 \quad (20)$$

is well defined and furthermore

$$0 < \liminf_K \Theta_K^2 \leq \limsup_K \Theta_K^2 < \infty.$$

2) Assume in addition **A1**. Then the sequence $\beta_K = \boldsymbol{y}^*(\boldsymbol{YY}^* + \rho\boldsymbol{I})^{-1}\boldsymbol{y}$ satisfies

$$\frac{\sqrt{K}}{\Theta_K}\left(\beta_K - \overline{\beta}_K\right) \underset{K \to \infty}{\longrightarrow} \mathcal{N}(0,1)$$

in distribution where $\overline{\beta}_K = \frac{1}{K}\mathrm{tr}\boldsymbol{D}_0\boldsymbol{T}_K$ is defined in the statement of Theorem 2.

Remark 2: (Comparison With Other Performance Indexes): It is interesting to compare the "mean squared error" (MSE)

related to the SINR β_K: $\mathrm{MSE}(\beta_K) = \mathbb{E}(\beta_K - \overline{\beta}_K)^2$, with the MSE related to Shannon's mutual information per transmit dimension $I = \frac{1}{K}\log\det(\rho\boldsymbol{\Sigma\Sigma}^* + \boldsymbol{I})$ (studied in [24], [25] for instance)

$$\mathrm{MSE}(\beta_K) \propto \mathcal{O}\left(\frac{1}{K}\right) \quad \text{while} \quad \mathrm{MSE}(I) \propto \mathcal{O}\left(\frac{1}{K^2}\right).$$

Remark 3 (On the Achievability of the Minimum of the Variance): Recall that the variance writes

$$\Theta_K^2 = \frac{1}{K}\boldsymbol{g}^T(\boldsymbol{I}_K - \boldsymbol{A})^{-1}\boldsymbol{\Delta}^{-1}\boldsymbol{g} + (\mathbb{E}|W_{10}|^4 - 1)\frac{1}{K}\mathrm{tr}\boldsymbol{D}_0^2\boldsymbol{T}^2.$$

As $\mathbb{E}|W_{10}|^2 = 1$, one clearly has $\mathbb{E}|W_{10}|^4 - 1 \geq 0$ with equality if and only if $|W_{10}| = 1$ with probability one. Moreover, we shall prove in the sequel (Section V-B) that $\liminf_K \frac{1}{K}\boldsymbol{D}_0(K)\boldsymbol{T}_K^2 > 0$. Therefore, $(\mathbb{E}|W_{10}|^4 - 1)\frac{1}{K}\mathrm{tr}\boldsymbol{D}_0^2\boldsymbol{T}^2$ is nonnegative, and is zero if and only if $|W_{10}| = 1$ with probability one. As a consequence, Θ_K^2 is minimum with respect to the distribution of the W_{nk} if and only if these random variables have their values on the unit circle. In the context of CDMA and MC-CDMA, this is the case when the signature matrix elements are elements of a phase-shift keying (PSK) constellation. In multiple-antenna systems, the W_{nk}'s are frequently considered as Gaussian which induces a penalty on the SINR asymptotic MSE with respect to the unit norm case.

In the separable case, $\Theta_K^2 = \tilde{d}_0^2\Omega_K^2$ where Ω_K^2 is given by the following corollary.

Corollary 1: Assume that **A2** is satisfied and that $\sigma_{nk}^2 = d_n\tilde{d}_k$. Assume moreover that

$$\min\left(\liminf_K \frac{1}{K}\mathrm{tr}(\boldsymbol{D}(K)), \liminf_K \frac{1}{K}\mathrm{tr}(\widetilde{\boldsymbol{D}}(K))\right) > 0 \quad (21)$$

where $\boldsymbol{D}$ and $\widetilde{\boldsymbol{D}}$ are given by (15). Let $\gamma = \frac{1}{K}\mathrm{tr}\boldsymbol{D}^2\boldsymbol{T}^2$ and $\tilde{\gamma} = \frac{1}{K}\mathrm{tr}\widetilde{\boldsymbol{D}}^2\widetilde{\boldsymbol{T}}^2$. Then the sequence

$$\Omega_K^2 = \gamma\left(\frac{\rho^2\gamma\tilde{\gamma}}{1 - \rho^2\gamma\tilde{\gamma}} + \left(\mathbb{E}|W_{10}|^4 - 1\right)\right) \quad (22)$$

satisfies $0 < \liminf_K \Omega_K^2 \leq \limsup_K \Omega_K^2 < \infty$. If, in addition, **A1** holds true, then

$$\frac{\sqrt{K}}{\Omega_K}\left(\frac{\beta_K}{\tilde{d}_0} - \delta_K\right) \underset{K \to \infty}{\longrightarrow} \mathcal{N}(0,1)$$

in distribution.

Remark 4: Condition (21) is the counterpart of Assumption **A3** in the case of a separable variance profile and suffices to establish $0 < \liminf_K(1 - \rho^2\gamma\tilde{\gamma}) \leq \limsup_K(1 - \rho^2\gamma\tilde{\gamma}) < 1$ (see for instance [23]), hence, the fact that $0 < \liminf_K \Omega_K^2 \leq \limsup_K \Omega_K^2 < \infty$. The remainder of the proof of Corollary 1 is postponed to Appendix B.

Remark 5: As a direct application of Corollary 1 (to be used in Section IV below), let us provide the expressions of γ and $\tilde{\gamma}$ for the MIMO model (5) or MC-CDMA downlink model (8). From (15)–(17), we get

$$\gamma = \frac{1}{K} \sum_{n=1}^{N} \left(\frac{d_n}{\rho + \rho d_n\tilde{\delta}}\right)^2$$

$$= \frac{1}{K} \sum_{n=1}^{N} \left(\frac{d_n}{\rho + \frac{1}{K} d_n \sum_{k=1}^{K} \frac{p_k}{1 + p_k \delta}} \right)^2$$

$$\tilde{\gamma} = \frac{1}{K} \sum_{k=1}^{K} \left(\frac{p_k}{\rho + \rho p_k \delta} \right)^2$$

where we recall that $d_n = \lambda_n$ for model (5), $d_n = \frac{K}{N} |h(\exp(2\imath\pi(n-1)/N)|^2$ for model (8), and δ is the solution of (18).

In the context of Corollary 1, if we further assume that $D = I$, then we recover the results of [18] and [16] (the latter being specific to the case where $\tilde{D} = \tilde{I}$ in addition).

Corollary 2: Assume the setting of Corollary 1 with $D = I$. Then

$$\Omega_K^2 = \frac{N}{K} \int \frac{1}{(t+\rho)^2} \pi_K(dt)$$
$$+ \frac{N}{K} \left(\mathbb{E}|W_{10}|^4 - 2 \right) \left(\int \frac{1}{t+\rho} \pi_K(dt) \right)^2 \quad (23)$$

where π_K is the probability measure which ST is $\frac{1}{N} \mathrm{tr} T(z)$ as shown in the statement of Theorem 1.

This corollary will be proven in Appendix C.

IV. SIMULATIONS

A. The General (Not Necessarily Separable) Case

In this section, the accuracy of the Gaussian approximation is verified by simulation. In order to validate the results of Theorems 2 and 3 for practical values of K, we consider the example of a MC-CDMA transmission in the uplink direction. We recall that K is the number of interfering users in this context. In the simulation, the discrete time channel impulse response of user k is represented by the vector with $L = 5$ coefficients $\boldsymbol{g}_k = [g_{k,0}, \ldots, g_{k,L-1}]^T$. In the simulations, these vectors are generated pseudo-randomly according to the complex multivariate Gaussian law $\mathcal{CN}(0, 1/L\boldsymbol{I}_L)$. Setting the number of frequency bins to N, the channel matrix $\boldsymbol{H}_k$ for user k in the frequency domain (see (7)) is

$$\boldsymbol{H}_k = \mathrm{diag}(h_k(\exp(2\imath\pi(n-1)/N); 1 \leq n \leq N)$$

where

$$h_k(z) = \frac{\sqrt{P_k}}{\|\boldsymbol{g}_k\|} \sum_{l=0}^{L-1} g_{k,l} z^{-l}$$

the norm $\|\boldsymbol{g}_k\|$ is the Euclidean norm of $\boldsymbol{g}_k$, and P_k is the power received from user k. Concerning the distribution of the user powers P_k, we assume that these are arranged into five power classes with powers $P, 2P, 4P, 8P$, and $16P$ with relative frequencies given by Table I. The user of interest (user 0) is assumed to belong to Class 1. Finally, we assume that the number K of interfering users is set to $K = N/2$.

In Table II, the signal to noise ratio (SNR) P/ρ for the user of interest is fixed to 10 dB. The evolution of $K\mathbb{E}(\beta_K - \bar{\beta}_K)^2/\Theta_K^2$

TABLE I
POWER CLASSES AND RELATIVE FREQUENCIES

Class	1	2	3	4	5
Power	P	$2P$	$4P$	$8P$	$16P$
Relative frequency	$1/8$	$1/4$	$1/4$	$1/8$	$1/4$

TABLE II
SINR NORMALIZED MSE VERSUS K (SNR = 10 dB)

K	8	16	32	64	128	256
$K\mathbb{E}(\beta_K - \bar{\beta}_K)^2/\Theta_K^2$	0.9761	0.9845	1.0464	1.0187	1.0127	0.9919

TABLE III
SINR NORMALIZED MSE VERSUS SNR ($K = 64$)

SNR	0	5	10	15	20	25	30
$K\mathbb{E}(\beta_K - \bar{\beta}_K)^2/\Theta_K^2$	1.0283	1.0294	1.0373	1.0358	1.0347	1.0348	1.0350

for this user (where $\mathbb{E}(\beta_K - \bar{\beta}_K)^2$ is measured numerically) is shown with respect to K. We note that this quantity is close to one for values of K as small as $K = 8$.

In Table III, K is set to $K = 64$, and the SINR normalized MSE $K\mathbb{E}(\beta_K - \bar{\beta}_K)^2/\Theta_K^2$ is plotted with respect to the input SNR P/ρ. This figure also confirms the fact that the MSE asymptotic approximation is highly accurate.

Fig. 1 shows the histogram of $\sqrt{K}(\beta_K - \bar{\beta}_K)/\Theta_K$ for $N = 16$ and $N = 64$. This figure gives an idea of the similarity between the distribution of $\sqrt{K}(\beta_K - \bar{\beta}_K)/\Theta_K$ and $\mathcal{N}(0, 1)$.

More precisely, Fig. 2 quantifies this similarity through a quantile–quantile (Q-Q) plot.

B. The Separable Case

In order to test the results of Proposition 2 and Corollary 1, we consider the following multiple-antenna (MIMO) model with exponentially decaying correlation at reception:

$$\boldsymbol{\Sigma} = \frac{1}{\sqrt{K}} \boldsymbol{\Psi}^{1/2} \boldsymbol{W} \boldsymbol{P}^{1/2}$$

where $\boldsymbol{\Psi} = [a^{m-n}]_{m,n=0}^{N-1}$ with $0 < a < 1$ is the covariance matrix that accounts for the correlations at the receiver side, $P = \mathrm{diag}(p_0, \ldots, p_K)$ is the matrix of the powers given to the different sources, and W is an $N \times (K+1)$ matrix with Gaussian standard i.i.d. elements. Let $\boldsymbol{P_u}$ denote the vector containing the powers of the interfering sources. We set $\boldsymbol{P_u}$ (up to a permutation of its elements) as shown in the equation at the bottom of the following page. For $K = 2^p$ with $3 \leq p \leq 7$, we assume that the powers of the interfering sources are arranged into five classes as in Table I. We set the SNR P/ρ to 10 dB and a to 0.1. We investigate in this section the accuracy of the Gaussian approximation in terms of the outage probability. In Fig. 3, we compare the empirical 1% outage SINR with the one predicted by the CLT. We note that the Gaussian approximation tends to underestimate the 1% outage SINR. We also note that it has a good accuracy for small values of α and for enough large values of N ($N \geq 64$).

Observe that all these simulations confirm a fact announced in Remark 3 above: compared with functionals of the channel singular values such as Shannon's mutual information, larger signal dimensions are needed to attain the asymptotic regime for quadratic forms such as the SINR (see for instance outage
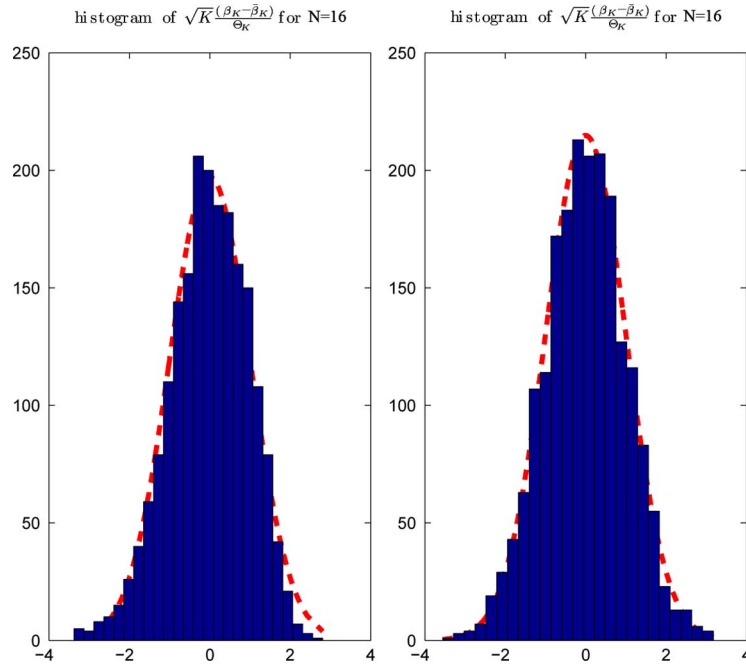
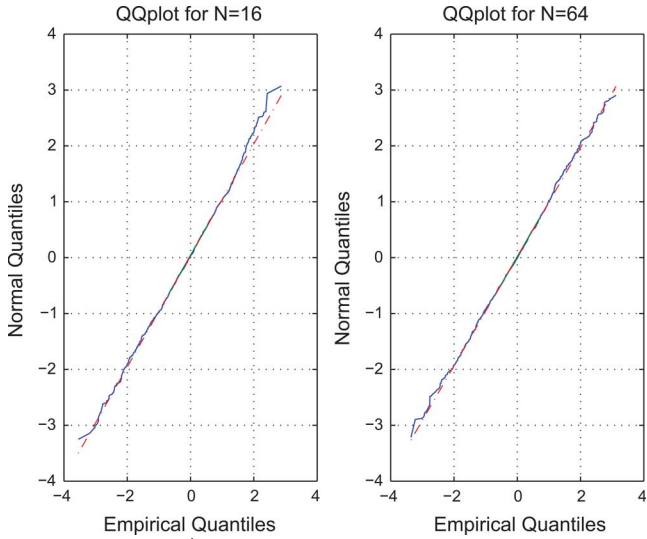Fig. 1. Histogram of $\sqrt{K}(\beta_K - \bar\beta_K)$ for $N = 16$ and $N = 64$.

Fig. 2. Q-Q plot for $\sqrt{K}(\beta_K - \bar\beta_K)$, $N = 16$ and $N = 64$; dash-dotted line is the 45° line.

probability approximations for mutual information in [25] and in [26]). This observation holds for first-order as well as second-order results.

V. PROOF OF THEOREM 3

This section is devoted to the proof of Theorem 3. We begin with mathematical preliminaries.

A. Preliminaries

The following lemma gathers useful matrix results, whose proofs can be found in [31].

Lemma 2: Assume $\boldsymbol{X} = [x_{ij}]_{i,j=1}^N$ and $\boldsymbol{Y}$ are complex $N \times N$ matrices. Then
1) For every $i, j \leq N$, $|x_{ij}| \leq \|\boldsymbol{X}\|$. In particular, $\|\mathrm{diag}(\boldsymbol{X})\| \leq \|\bar{\boldsymbol{X}}\|$.
2) $\|\boldsymbol{XY}\| \leq \|\boldsymbol{X}\| \|\boldsymbol{Y}\|$
3) For $\rho > 0$, the resolvent $(XX^* + \rho\boldsymbol{I})^{-1}$ satisfies $\|(XX^* + \rho\boldsymbol{I})^{-1}\| \leq \rho^{-1}$.
4) If $\boldsymbol{Y}$ is Hermitian nonnegative, then $|\mathrm{tr}(\boldsymbol{XY})| \leq \|\boldsymbol{X}\|\mathrm{tr}(\boldsymbol{Y})$.

Let $\boldsymbol{X} = \boldsymbol{U\Lambda V}^*$ be a spectral decomposition of $\boldsymbol{X}$ where $\boldsymbol{\Lambda} = \mathrm{diag}(\lambda_1, \ldots, \lambda_n)$ is the matrix of singular values of $\boldsymbol{X}$. For a real $p \geq 1$, the Schatten ℓ_p-norm of $\boldsymbol{X}$ is defined as $\|\boldsymbol{X}\|_p = \left(\sum \lambda_i^p\right)^{1/p}$. The following bound over the Schatten ℓ_p-norm of a triangular matrix will be of help (for a proof, see [28], [32, p. 278]).

Lemma 3: Let $\boldsymbol{X} = [x_{ij}]_{i,j=1}^N$ be an $N \times N$ complex matrix and let $\widetilde{\boldsymbol{X}} = [x_{ij}\mathbf{1}_{i>j}]_{i,j=1}^N$ be the strictly lower triangular matrix extracted from $\boldsymbol{X}$. Then for every $p \geq 1$, there exists a constant C_p depending on p only such that

$$\|\widetilde{\boldsymbol{X}}\|_p \leq C_p\|\boldsymbol{X}\|_p.$$

$$\boldsymbol{P_u} = \begin{cases} [4P \quad 5P], & \text{if } K = 2 \\ [PZ \quad P \quad 2P \quad 4P], & \text{if } K = 4 \\ [P \quad P \quad 2P \quad 2P \quad 2P \quad 4P \quad 4P \quad 4P \quad 8P \quad 16P \quad 16P \quad 16P], & \text{if } K = 12. \end{cases}$$
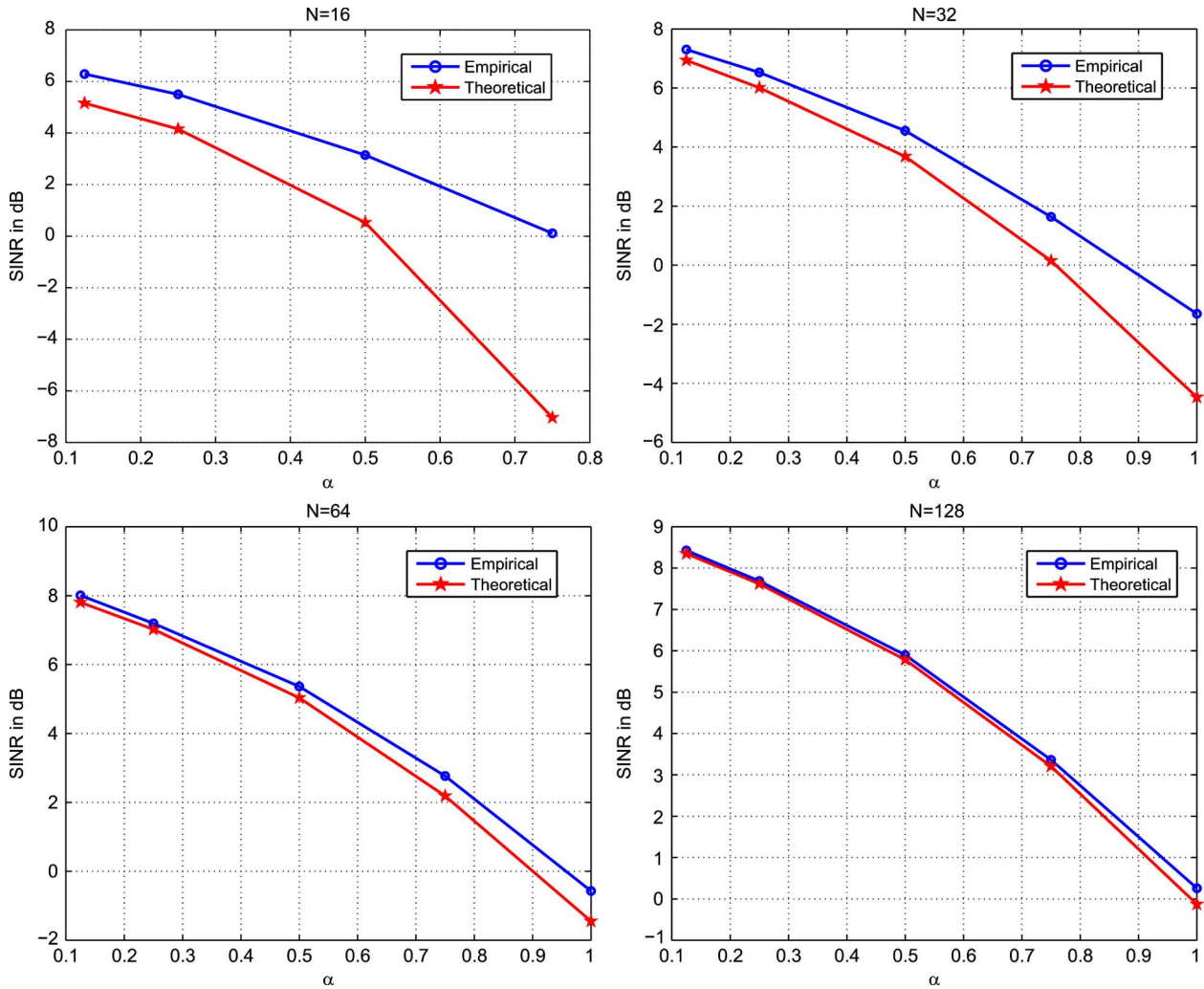
Fig. 3.   Theoretical and empirical 1% outage SINR.

The following lemma lists some properties of the resolvent $\boldsymbol{Q}$ and the deterministic approximation matrix $\boldsymbol{T}$. Its proof is postponed to Appendix A.

Lemma 4: The following facts hold true.
1) Assume **A2**. Consider matrices

$$\boldsymbol{T}_K(-\rho) = \operatorname{diag}(t_1(-\rho), \ldots, t_N(-\rho))$$

defined by Theorem 1–part 1. Then for every $1 \le n \le N$

$$\frac{1}{\rho + \sigma_{\max}^2} \le t_n(-\rho) \le \frac{1}{\rho}. \qquad (24)$$

2) Assume in addition **A1** and **A3**. Let $\boldsymbol{Q}_K(-\rho) = (\boldsymbol{Y}\boldsymbol{Y}^* + \rho\boldsymbol{I})^{-1}$ and let matrices $\boldsymbol{S}_K$ be as in the statement of Theorem 1–part 3. Then

$$\sup_K \mathbb{E} \left| \operatorname{tr} \boldsymbol{S}_K(\boldsymbol{Q}_K - \boldsymbol{T}_K) \right|^2 < \infty. \qquad (25)$$

B. Proof of Theorem 3–Part 1

We introduce the following notations. Assume that $\boldsymbol{X}$ is a real matrix, by $\boldsymbol{X} \succcurlyeq \boldsymbol{0}$ we mean $X_{ij} \ge 0$ for every element X_{ij}. For a vector $\boldsymbol{x}$, $\boldsymbol{x} \succcurlyeq \boldsymbol{0}$ is defined similarly. In the remainder of

the paper, $C = C(\rho, \sigma_{\max}^2, \liminf \frac{N}{K}, \sup \frac{N}{K}) < \infty$ denotes a positive constant whose value may change from line to line.

The following lemma, which directly follows from [24, Lemma 5.2 and Proposition 5.5], states some important properties of the matrices $\boldsymbol{A}_K$ defined in the statement of Theorem 3.

Lemma 5: Assume **A2** and **A3**. Consider matrices $\boldsymbol{A}_K$ defined by (19). Then the following facts hold true.
1) Matrix $\boldsymbol{I}_K - \boldsymbol{A}_K$ is invertible, and $(\boldsymbol{I}_K - \boldsymbol{A}_K)^{-1} \succcurlyeq \boldsymbol{0}$.
2) Element (k, k) of the inverse satisfies $\left[(\boldsymbol{I}_K - \boldsymbol{A}_K)^{-1} \right]_{k,k} \ge 1$ for every $1 \le k \le K$.
3) The maximum row sum norm of the inverse satisfies $\limsup_K \left\| \left\| (\boldsymbol{I}_K - \boldsymbol{A}_K)^{-1} \right\| \right\|_\infty < \infty$.

Due to Lemma 5–part 1, Θ_K^2 is well defined. Let us prove that $\limsup_K \Theta_K^2 < \infty$. The first term of the right-hand side of (20) satisfies

$$\frac{1}{K} \boldsymbol{g}^T (\boldsymbol{I}_K - \boldsymbol{A}_K)^{-1} \boldsymbol{\Delta}^{-1} \boldsymbol{g}$$

$$\le \|\boldsymbol{g}\|_\infty \| (\boldsymbol{I}_K - \boldsymbol{A}_K)^{-1} \boldsymbol{\Delta}^{-1} \boldsymbol{g} \|_\infty$$

$$\le \|\boldsymbol{g}\|_\infty \left\| \left\| (\boldsymbol{I}_K - \boldsymbol{A}_K)^{-1} \right\| \right\|_\infty \| \boldsymbol{\Delta}^{-1} \boldsymbol{g} \|_\infty$$

$$\le \|\boldsymbol{g}\|_\infty^2 \left\| \left\| (\boldsymbol{I}_K - \boldsymbol{A}_K)^{-1} \right\| \right\|_\infty \qquad (26)$$

due to $\||\boldsymbol{\Delta}^{-1}\||_{\infty} \leq 1$. Recall that $\|\boldsymbol{T}\| \leq \rho^{-1}$ by Lemma 4–part 1. Therefore, any element of $\boldsymbol{g}$ satisfies

$$\frac{1}{K}\mathrm{tr}\boldsymbol{D}_0\boldsymbol{D}_k\boldsymbol{T}^2 \leq \frac{N}{K}\|\boldsymbol{D}_0\|\|\boldsymbol{D}_k\|\|\boldsymbol{T}\|^2 \leq \frac{N}{K}\frac{\sigma_{\max}^4}{\rho^2} \quad (27)$$

by **A2**, hence $\sup_K \|\boldsymbol{g}\| \leq C$. From Lemma 5–part 3 and (26), we then obtain

$$\limsup_K \frac{1}{K}\boldsymbol{g}^T(\boldsymbol{I}_K - \boldsymbol{A}_K)^{-1}\boldsymbol{\Delta}^{-1}\boldsymbol{g} \leq C. \quad (28)$$

We can prove similarly that the second term in the right-hand side of (20) satisfies $\sup_K((\mathbb{E}|W_{10}|^4 - 1)\frac{1}{K}\mathrm{tr}\boldsymbol{D}_0^2\boldsymbol{T}(-\rho)^2) \leq C$. Hence, $\limsup_K \Theta_K^2 < \infty$.

Let us prove that $\liminf_K \Theta_K^2 > 0$. We have

$$\frac{1}{K}\boldsymbol{g}^T(\boldsymbol{I}_K - \boldsymbol{A}_K)^{-1}\boldsymbol{\Delta}^{-1}\boldsymbol{g}$$

$$\overset{(a)}{\geq} \frac{1}{K}\boldsymbol{g}^T\mathrm{diag}\left((\boldsymbol{I}_K - \boldsymbol{A}_K)^{-1}\right)\boldsymbol{\Delta}^{-1}\boldsymbol{g}$$

$$\overset{(b)}{\geq} \frac{1}{\left(1 + \frac{N}{K}\frac{\sigma_{\max}^2}{\rho}\right)^2}\frac{1}{K}\sum_{k=1}^{K}\left(\frac{1}{K}\mathrm{tr}\boldsymbol{D}_0\boldsymbol{D}_k\boldsymbol{T}^2\right)^2$$

$$\overset{(c)}{\geq} \frac{1}{\left(1 + \frac{N}{K}\frac{\sigma_{\max}^2}{\rho}\right)^2}\left(\frac{1}{K^2}\mathrm{tr}\,\boldsymbol{D}_0\left(\sum_{k=1}^{K}\boldsymbol{D}_k\right)\boldsymbol{T}^2\right)^2$$

$$\overset{(d)}{\geq} \frac{1}{\left(1 + \frac{N}{K}\frac{\sigma_{\max}^2}{\rho}\right)^2(\rho + \sigma_{\max}^2)^4}\left(\frac{1}{K^2}\mathrm{tr}\,\boldsymbol{D}_0\sum_{k=1}^{K}\boldsymbol{D}_k\right)^2$$

$$\geq C\left(\frac{1}{K^2}\mathrm{tr}\,\boldsymbol{D}_0\sum_{k=1}^{K}\boldsymbol{D}_k\right)^2$$

where (a) follows from the fact that $(\boldsymbol{I}_K - \boldsymbol{A}_K)^{-1} \succcurlyeq \boldsymbol{0}$ (Lemma 5–part 1, and the straightforward inequalities $\boldsymbol{\Delta}^{-1} \succcurlyeq \boldsymbol{0}$ and $\boldsymbol{g} \succcurlyeq \boldsymbol{0}$), (b) follows from Lemma 5–part 2 and $\|\boldsymbol{\Delta}\| \leq (1 + \frac{N}{K}\frac{\sigma_{\max}^2}{\rho})^2$, (c) follows from the elementary inequality $n^{-1}\sum x_i^2 \geq (n^{-1}\sum x_i)^2$, and (d) is due to Lemma 4–part 1 and Lemma 2-part 4. Similar derivations yield

$$(\mathbb{E}|W_{10}|^4 - 1)\frac{1}{K}\mathrm{tr}\boldsymbol{D}_0^2\boldsymbol{T}$$

$$\geq \frac{\mathbb{E}|W_{10}|^4 - 1}{(\rho + \sigma_{\max}^2)^2}\left(\frac{1}{K}\mathrm{tr}\boldsymbol{D}_0\right)^2 \geq C(\mathbb{E}|W_{10}|^4 - 1)$$

by **A3**. Therefore, if **A4** holds true, then $\liminf_K \Theta_K^2 > 0$ and Theorem 3–part 1 is proved.

C. Proof of Theorem 3–Part 2

Recall that the SINR β_K is given by (14). The random variable $\frac{\sqrt{K}}{\Theta_K}(\beta_K - \overline{\beta}_K)$ can therefore be decomposed as

$$\frac{\sqrt{K}}{\Theta_K}(\beta_K - \overline{\beta}_K) = \frac{1}{\sqrt{K}\Theta_K}\left(\boldsymbol{w}_0^*\boldsymbol{D}_0^{1/2}\boldsymbol{Q}\boldsymbol{D}_0^{1/2}\boldsymbol{w}_0 - \mathrm{tr}(\boldsymbol{D}_0\boldsymbol{Q})\right)$$

$$+ \frac{1}{\sqrt{K}\Theta_K}\left(\mathrm{tr}(\boldsymbol{D}_0(\boldsymbol{Q} - \boldsymbol{T}))\right)$$

$$= U_{1,K} + U_{2,K}. \quad (29)$$

Thanks to Lemma 4–part 2 and to the fact that $\liminf_K \Theta_K^2 > 0$, we have $\mathbb{E}U_{K,2}^2 < CK^{-1}$ which implies that $U_{K,2} \to 0$ in probability as $K \to \infty$. Hence, in order to conclude that

$$\frac{\sqrt{K}}{\Theta_K}(\beta_K - \overline{\beta}_K) \underset{K \to \infty}{\longrightarrow} \mathcal{N}(0, 1) \quad \text{in distribution}$$

it is sufficient by Slutsky's theorem to prove that $U_{1,K} \to \mathcal{N}(0, 1)$ in distribution. The remainder of the section is devoted to this point.

Remark 6: Decomposition (29) and the convergence to zero (in probability) of $U_{2,K}$ yield the following interpretation: The fluctuations of $\sqrt{K}(\beta_K - \overline{\beta}_K)$ are mainly due to the fluctuations of vector $\boldsymbol{w}_0$. Indeed, the contribution of the fluctuations[1] of $\frac{1}{K}\mathrm{tr}\boldsymbol{D}_0\boldsymbol{Q}$, due to the random nature of $\boldsymbol{Y}$, is negligible.

Denote by $\mathbb{E}_n$ the conditional expectation $\mathbb{E}_n[\cdot] = \mathbb{E}[\cdot \| W_{n,0}, W_{n+1,0}, \ldots, W_{N,0}, \boldsymbol{Y}]$. Put $\mathbb{E}_{N+1}[\cdot] = \mathbb{E}[\cdot \| \boldsymbol{Y}]$ and note that $\mathbb{E}_{N+1}(\boldsymbol{w}_0^*\boldsymbol{D}_0^{1/2}\boldsymbol{Q}\boldsymbol{D}_0^{1/2}\boldsymbol{w}_0) = \mathrm{tr}\boldsymbol{D}_0\boldsymbol{Q}$. With these notations at hand, we have

$$U_{1,K} = \frac{1}{\Theta_K}\sum_{n=1}^{N}(\mathbb{E}_n - \mathbb{E}_{n+1})\frac{\boldsymbol{w}_0^*\boldsymbol{D}_0^{1/2}\boldsymbol{Q}\boldsymbol{D}_0^{1/2}\boldsymbol{w}_0}{\sqrt{K}}$$

$$\triangleq \frac{1}{\Theta_K}\sum_{n=1}^{N}Z_{n,K}. \quad (30)$$

Consider the increasing sequence of σ-fields

$$\mathcal{F}_{N,K} = \sigma(W_{N,0}, \boldsymbol{Y}), \ldots, \mathcal{F}_{1,K}$$

$$= \sigma(W_{1,0}, \ldots, W_{N,0}, \boldsymbol{Y}).$$

Then the random variable $Z_{n,K}$ is integrable and measurable with respect to $\mathcal{F}_{n,K}$; moreover, it readily satisfies $\mathbb{E}_{n+1}Z_{n,K} = 0$. In particular, the sequence $(Z_{N,K}, \ldots, Z_{1,K})$ is a martingale difference sequence with respect to $(\mathcal{F}_{N,K}, \ldots, \mathcal{F}_{1,K})$. The following CLT for martingales is the key tool to study the asymptotic behavior of $U_{1,K}$:

Theorem 4: Let $X_{N,K}, X_{N-1,K}, \ldots, X_{1,K}$ be a martingale difference sequence with respect to the increasing filtration $\mathcal{G}_{N,K}, \ldots, \mathcal{G}_{1,K}$. Assume that there exists a sequence of real positive numbers s_K^2 such that

$$\frac{1}{s_K^2}\sum_{n=1}^{N}\mathbb{E}\left[X_{n,K}^2\|\mathcal{G}_{n+1,K}\right] \underset{K \to \infty}{\longrightarrow} 1$$

in probability. Assume further that the Lyapunov condition holds

$$\exists \alpha > 0, \quad \frac{1}{s_K^{2+\alpha}}\sum_{n=1}^{N}\mathbb{E}|X_{n,K}|^{2+\alpha} \underset{K \to \infty}{\longrightarrow} 0.$$

Then $s_K^{-1}\sum_{n=1}^{N}X_{n,K}$ converges in distribution to $\mathcal{N}(0, 1)$ as $K \to \infty$.

Remark 7: This theorem is proved in [27], gathering Theorem 35.12 (which is expressed under the weaker Lindeberg

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[1]In fact, one may prove that the fluctuation of $\frac{1}{K}\mathrm{tr}\boldsymbol{D}_0(\boldsymbol{Q} - \boldsymbol{T})$ are of order K, i.e., $\mathrm{tr}\boldsymbol{D}_0(\boldsymbol{Q} - \boldsymbol{T})$ asymptotically behaves as a Gaussian random variable. Such a speed of fluctuations already appears in [24], when studying the fluctuations of the mutual information.

condition) together with the arguments of Section 27 (where it is proved that Lyapunov's condition implies Lindeberg's condition).

In order to prove that

$$U_{1,K} = \frac{1}{\Theta_K} \sum_{n=1}^{N} Z_{n,K} \underset{K \to \infty}{\longrightarrow} \mathcal{N}(0,1) \quad \text{in distribution} \quad (31)$$

we shall apply Theorem 4 to the sum $\frac{1}{\Theta_K} \sum_{n=1}^{N} Z_{n,K}$ and the filtration $(\mathcal{F}_{n,K})$. The proof is carried out into four steps:

Step 1: We first establish Lyapunov's condition. Due to the fact that $\liminf_K \Theta_K^2 > 0$, we only need to show that

$$\exists\, \alpha > 0, \quad \sum_{n=1}^{N} \mathbb{E}|Z_{n,K}|^{2+\alpha} \underset{K \to \infty}{\longrightarrow} 0. \quad (32)$$

Step 2: We prove that $V_K = \sum_{n=1}^{N} \mathbb{E}_{n+1} Z_{n,K}^2$ satisfies

$$V_K - \left(\frac{(\mathbb{E}|W_{10}|^4 - 2)}{K} \text{tr}\big(\boldsymbol{D}_0^2 (\text{diag}(\boldsymbol{Q}))^2\big) \right.$$
$$\left. + \frac{1}{K} \text{tr}(\boldsymbol{D}_0 \boldsymbol{Q} \boldsymbol{D}_0 \boldsymbol{Q}) \right) \underset{K \to \infty}{\longrightarrow} 0 \quad \text{in probability.} \quad (33)$$

Step 3: We first show that

$$\frac{1}{K} \text{tr} \boldsymbol{D}_0^2 (\text{diag}(\boldsymbol{Q}))^2 - \frac{1}{K} \text{tr} \boldsymbol{D}_0^2 \boldsymbol{T}^2 \underset{K \to \infty}{\longrightarrow} 0 \quad \text{in probability.} \quad (34)$$

In order to study the asymptotic behavior of $\frac{1}{K} \text{tr}(\boldsymbol{D}_0 \boldsymbol{Q} \boldsymbol{D}_0 \boldsymbol{Q})$, we introduce the random variables $U_\ell = \frac{1}{K} \text{tr}(\boldsymbol{D}_0 \boldsymbol{Q} \boldsymbol{D}_\ell \boldsymbol{Q})$ for $0 \le \ell \le K$ (the one of interest being U_0). We then prove that the U_ℓ's satisfy the following system of equations:

$$U_\ell = \sum_{k=1}^{K} c_{\ell k} U_k + \frac{1}{K} \text{tr} \boldsymbol{D}_0 \boldsymbol{D}_\ell \boldsymbol{T}^2 + \epsilon_\ell, \quad 0 \le \ell \le K \quad (35)$$

where

$$c_{\ell k} = \frac{1}{K} \frac{\frac{1}{K} \text{tr} \boldsymbol{D}_\ell \boldsymbol{D}_k \boldsymbol{T}(-\rho)^2}{\left(1 + \frac{1}{K} \text{tr} \boldsymbol{D}_k \boldsymbol{T}(-\rho)\right)^2}, \quad 0 \le \ell \le K, \, 1 \le k \le K \quad (36)$$

and the perturbations ϵ_ℓ satisfy $\mathbb{E}|\epsilon_\ell| \le CK^{-\frac{1}{2}}$ where we recall that C is independent of ℓ.

Step 4: We prove that $U_0 = \frac{1}{K} \text{tr} \boldsymbol{D}_0 \boldsymbol{Q} \boldsymbol{D}_0 \boldsymbol{Q}$ satisfies

$$U_0 = \frac{1}{K} \text{tr} \boldsymbol{D}_0^2 \boldsymbol{T}^2 + \frac{1}{K} \boldsymbol{g}^T (\boldsymbol{I} - \boldsymbol{A})^{-1} \boldsymbol{\Delta}^{-1} \boldsymbol{g} + \epsilon \quad (37)$$

with $\mathbb{E}|\epsilon| \le CK^{-\frac{1}{2}}$. This equation combined with (33) and (34) yields $\sum_n \mathbb{E}_{n+1} Z_{n,K}^2 - \Theta_K^2 \to 0$ in probability. As $\liminf_K \Theta_K^2 > 0$, this implies $\frac{1}{\Theta_K} \sum_n \mathbb{E}_{n+1} Z_{n,K}^2 \to 1$ in probability, which proves (31) and thus ends the proof of Theorem 3.

Write $\boldsymbol{B} = [b_{ij}]_{i,j=1}^{N} = \boldsymbol{D}_0^{1/2} \boldsymbol{Q} \boldsymbol{D}_0^{1/2}$ and recall from (30) that $Z_{n,K} = \frac{1}{\sqrt{K}} (\mathbb{E}_n - \mathbb{E}_{n+1}) \boldsymbol{w}_0^* \boldsymbol{B} \boldsymbol{w}_0$. We have

$$\mathbb{E}_n \boldsymbol{w}_0^* \boldsymbol{B} \boldsymbol{w}_0 = \sum_{\ell=1}^{n-1} b_{\ell\ell} + \sum_{\ell_1,\ell_2=n}^{N} W_{\ell_1 0}^* W_{\ell_2 0} b_{\ell_1 \ell_2}.$$

Hence

$$Z_{n,K} = \frac{1}{\sqrt{K}} \bigg((|W_{n0}|^2 - 1) b_{nn} $$
$$+ W_{n0}^* \sum_{\ell=n+1}^{N} W_{\ell 0} b_{n\ell} + W_{n0} \sum_{\ell=n+1}^{N} W_{\ell 0}^* b_{\ell n} \bigg). \quad (38)$$

Step 1: Validation of the Lyapunov Condition: The following inequality will be of help to check Lyapunov's condition.

Lemma 6 (Burkholder's Inequality): Let X_k be a complex martingale difference sequence with respect to the increasing sequence of σ–fields $\mathcal{F}_k$. Then for $p \ge 2$, there exists a constant C_p for which

$$\mathbb{E} \left| \sum_k X_k \right|^p$$
$$\le C_p \left(\mathbb{E} \left(\sum_k \mathbb{E}\left[|X_k|^2 | \mathcal{F}_{k-1}\right] \right)^{p/2} + \mathbb{E} \sum_k |X_k|^p \right).$$

Recall Assumption **A1**. Equation (38) yields

$$|Z_{n,K}|^4 \le \frac{1}{K^2} \left(\frac{|W_{n0}|^2 + 1}{\rho \sigma_{\max}^2} + 2 \left| W_{n0} \sum_{\ell=n+1}^{N} W_{\ell 0} b_{n\ell} \right| \right)^4$$
$$\le \frac{2^3}{K^2} \left(\left(\frac{|W_{n0}|^2 + 1}{\rho \sigma_{\max}^2} \right)^4 + 2^4 \left| W_{n0} \sum_{\ell=n+1}^{N} W_{\ell 0} b_{n\ell} \right|^4 \right) \quad (39)$$

where we use the fact that $|b_{nn}| \le (\rho \sigma_{\max}^2)^{-1}$ (cf. Lemma 2–part 1) and the convexity of $x \mapsto x^4$. Due to Assumption **A1**, we have

$$\mathbb{E}\left(|W_{n0}|^2 + 1\right)^4 \le 2^3 \left(\mathbb{E}|W_{n0}|^8 + 1\right) < \infty. \quad (40)$$

Considering the second term at the right-hand side of (39), we write the equation at the top of the following page, where (a) follows from Lemma 6 (Burkholder's inequality), the filtration being $\mathcal{F}_{N,K}, \dots, \mathcal{F}_{n+1,K}$, and (b) follows from the bound $|b_{n\ell}|^4 \le |b_{n\ell}|^2 \max |b_{n\ell}|^2 \le |b_{n\ell}|^2 (\sigma_{\max}^2 \rho^{-1})^2$ (cf. Lemma 2–part 1). Now, notice that

$$\sum_{\ell=n+1}^{N} |b_{n\ell}|^2 < \sum_{\ell=1}^{N} |b_{n\ell}|^2$$
$$= \left[\boldsymbol{D}_0^{1/2} \boldsymbol{Q} \boldsymbol{D}_0 \boldsymbol{Q} \boldsymbol{D}_0^{1/2} \right]_{nn}$$
$$\le \| \boldsymbol{D}_0^{1/2} \boldsymbol{Q} \boldsymbol{D}_0 \boldsymbol{Q} \boldsymbol{D}_0^{1/2} \| \le \frac{\sigma_{\max}^4}{\rho^2}.$$

This yields $\mathbb{E}|W_{n0} \sum_{\ell=n+1}^{N} W_{\ell 0} b_{n\ell}|^4 \le C$. Gathering this result with (40), getting back to (39), taking the expectation and summing up finally yields

$$\sum_{n=1}^{N} \mathbb{E}|Z_{n,K}|^4 \le \frac{C}{K} \underset{K \to \infty}{\longrightarrow} 0$$

which establishes Lyapunov's condition (32) with $\alpha = 2$.

$$\mathbb{E}\left|W_{n0}\sum_{\ell=n+1}^{N}W_{\ell 0}b_{n\ell}\right|^4 = \mathbb{E}|W_{n0}|^4\,\mathbb{E}\left|\sum_{\ell=n+1}^{N}W_{\ell 0}b_{n\ell}\right|^4$$

$$\overset{(a)}{\leq} C\left(\mathbb{E}\left(\sum_{\ell=n+1}^{N}(\mathbb{E}|W_{\ell 0}|^2)|b_{n\ell}|^2\right)^2 + \sum_{\ell=n+1}^{N}(\mathbb{E}|W_{\ell 0}|^4)(\mathbb{E}|b_{n\ell}|^4)\right)$$

$$\overset{(b)}{\leq} C\left(\mathbb{E}\left(\sum_{\ell=n+1}^{N}|b_{n\ell}|^2\right)^2 + \sum_{\ell=n+1}^{N}\mathbb{E}|b_{n\ell}|^2\right)$$

Step 2: Proof of (33): Equation (38) yields

$$\mathbb{E}_{n+1}Z_{n,K}^2 = \frac{1}{K}\left((\mathbb{E}|W_{10}|^4 - 1)\,b_{nn}^2 + \mathbb{E}_{n+1}\right.$$

$$\left(W_{n0}^*\sum_{\ell=n+1}^{N}W_{\ell 0}b_{n\ell} + W_{n0}\sum_{\ell=n+1}^{N}W_{\ell 0}^*b_{\ell n}\right)^2$$

$$+2b_{nn}\left(\mathbb{E}\,W_{10}^*|W_{10}|^2\right)\sum_{\ell=n+1}^{N}W_{\ell 0}b_{n\ell}$$

$$\left.+2b_{nn}\left(\mathbb{E}\,W_{10}|W_{10}|^2\right)\sum_{\ell=n+1}^{N}W_{\ell 0}^*b_{\ell n}\right).$$

Note that the second term of the right-hand side writes

$$\mathbb{E}_{n+1}\left(W_{n0}^*\sum_{\ell=n+1}^{N}W_{\ell 0}b_{n\ell} + W_{n0}\sum_{\ell=n+1}^{N}W_{\ell 0}^*b_{\ell n}\right)^2$$

$$= 2\sum_{\ell_1,\ell_2=n+1}^{N}W_{\ell_1 0}W_{\ell_2 0}^*b_{n\ell_1}b_{\ell_2 n}.$$

Therefore, $V_K = \sum_{n=1}^{N}\mathbb{E}_{n+1}Z_{n,K}^2$ writes

$$V_K = \frac{(\mathbb{E}|W_{10}|^4 - 1)}{K}\sum_{n=1}^{N}b_{nn}^2$$

$$+\frac{2}{K}\sum_{n=1}^{N}\sum_{\ell_1,\ell_2=n+1}^{N}W_{\ell_1 0}W_{\ell_2 0}^*b_{n\ell_1}b_{\ell_2 n}$$

$$+\frac{2}{K}\Re\left((\mathbb{E}\,W_{10}^*|W_{10}|^2)\sum_{n=1}^{N}b_{nn}\sum_{\ell=n+1}^{N}W_{\ell 0}b_{n\ell}\right)$$

where $\Re$ denotes the real part of a complex number. We introduce the following notations:

$$\boldsymbol{R} = (r_{ij})_{i,j=1}^{N} \triangleq (b_{ij}\mathbf{1}_{i>j})_{i,j=1}^{N} \quad \text{and}$$

$$\Gamma_K = \frac{1}{K}\sum_{n=1}^{N}b_{nn}\sum_{\ell=n+1}^{N}W_{\ell 0}b_{n\ell}.$$

Note in particular that $\boldsymbol{R}$ is the strictly lower triangular matrix extracted from $\boldsymbol{D}_0^{1/2}\boldsymbol{Q}\boldsymbol{D}_0^{1/2}$. We can now rewrite V_K as

$$V_K = \frac{(\mathbb{E}|W_{10}|^4 - 1)}{K}\mathrm{tr}\left(\boldsymbol{D}_0^2(\mathrm{diag}(\boldsymbol{Q}))^2\right)$$

$$+\frac{2}{K}\boldsymbol{w}_0^*\boldsymbol{R}\boldsymbol{R}^*\boldsymbol{w}_0 + 2\Re\left(\Gamma_K\mathbb{E}W_{10}^*|W_{10}|^2\right). \quad (41)$$

We now prove that the third term of the right-hand side vanishes, and find an asymptotic equivalent for the second one. Using Lemma 2, we have

$$\mathbb{E}_{N+1}|\Gamma_K|^2 = \frac{1}{K^2}\sum_{n,m=1}^{N}b_{nn}b_{mm}\sum_{\ell=1}^{N}b_{n\ell}b_{m\ell}^*\mathbf{1}_{\ell>n}\mathbf{1}_{\ell>m}$$

$$= \frac{1}{K^2}\mathrm{tr}\left(\mathrm{diag}(\boldsymbol{B})\boldsymbol{R}^*\boldsymbol{R}\,\mathrm{diag}(\boldsymbol{B})\right)$$

$$= \frac{1}{K^2}\mathrm{tr}\left(\boldsymbol{D}_0^{1/2}\mathrm{diag}(\boldsymbol{Q})\boldsymbol{D}_0^{1/2}\right.$$

$$\left.\times\boldsymbol{R}^*\boldsymbol{R}\boldsymbol{D}_0^{1/2}\mathrm{diag}(\boldsymbol{Q})\boldsymbol{D}_0^{1/2}\right)$$

$$\leq \frac{1}{K^2}\|\boldsymbol{D}_0\|^2\|\boldsymbol{Q}\|^2\mathrm{tr}(\boldsymbol{R}^*\boldsymbol{R})$$

$$\leq \frac{1}{K^2}\|\boldsymbol{D}_0\|^2\|\boldsymbol{Q}\|^2\mathrm{tr}(\boldsymbol{B}^2)$$

$$\leq \frac{1}{K^2}\|\boldsymbol{D}_0\|^4\|\boldsymbol{Q}\|^2\mathrm{tr}(\boldsymbol{Q}^2)$$

$$\leq \frac{1}{K}\|\boldsymbol{D}_0\|^2\|\boldsymbol{Q}\|^4$$

$$\leq \frac{1}{K}\frac{\sigma_{\max}^4}{\rho^4}\underset{K\to\infty}{\longrightarrow} 0.$$

In particular, $\mathbb{E}|\Gamma_K|^2 \to 0$ and

$$\Re\left((\mathbb{E}W_{10}^*|W_{10}|^2)\,\Gamma_K\right)\underset{K\to\infty}{\longrightarrow} 0 \quad \text{in probability.} \quad (42)$$

Consider now the second term of the right-hand side of (41). We prove that

$$\frac{1}{K}\boldsymbol{w}_0^*\boldsymbol{R}\boldsymbol{R}^*\boldsymbol{w}_0 - \frac{1}{K}\mathrm{tr}(\boldsymbol{R}\boldsymbol{R}^*)\underset{K\to\infty}{\longrightarrow} 0 \quad \text{in probability.} \quad (43)$$

By Lemma 1 (12), we have

$$\mathbb{E}\left(\frac{1}{K}\boldsymbol{w}_0^*\boldsymbol{R}\boldsymbol{R}^*\boldsymbol{w}_0 - \frac{1}{K}\mathrm{tr}(\boldsymbol{R}\boldsymbol{R}^*)\right)^2$$

$$\leq \frac{C}{K^2}(\mathbb{E}|W_{10}|^4)\mathrm{tr}(\boldsymbol{R}\boldsymbol{R}^*\boldsymbol{R}\boldsymbol{R}^*).$$

Notice that $\mathrm{tr}(\boldsymbol{R}\boldsymbol{R}^*\boldsymbol{R}\boldsymbol{R}^*) = \|\boldsymbol{R}\|_4^4$ where $\|\boldsymbol{R}\|_4$ is the Schatten ℓ_4-norm of $\boldsymbol{R}$. Using Lemma 3, we have

$$\|\boldsymbol{R}\|_4^4 \leq C\|\boldsymbol{D}_0^{1/2}\boldsymbol{Q}\boldsymbol{D}_0^{1/2}\|_4^4 \leq NC\|\boldsymbol{D}_0^{1/2}\boldsymbol{Q}\boldsymbol{D}_0^{1/2}\|^4 \leq N\frac{C\sigma_{\max}^8}{\rho^4}.$$

Therefore

$$\mathbb{E}\left(\frac{1}{K}\boldsymbol{w}_0^*\boldsymbol{R}\boldsymbol{R}^*\boldsymbol{w}_0 - \frac{1}{K}\mathrm{tr}(\boldsymbol{R}\boldsymbol{R}^*)\right)^2 \le C\frac{N}{K^2}\xrightarrow[K\to\infty]{} 0$$

which implies (43). Now, due to the fact that $\boldsymbol{B} = \boldsymbol{B}^*$, we have

$$\begin{aligned}\frac{2}{K}\mathrm{tr}\,\boldsymbol{R}\boldsymbol{R}^* &= \frac{2}{K}\sum_{n=1}^N\sum_{\ell=n+1}^N |b_{n\ell}|^2\\&= \frac{1}{K}\sum_{n,\ell=1}^N |b_{n\ell}|^2 - \frac{1}{K}\sum_{n=1}^N |b_{nn}|^2\\&= \frac{1}{K}\mathrm{tr}\,\boldsymbol{D}_0\boldsymbol{Q}\boldsymbol{D}_0\boldsymbol{Q} - \frac{1}{K}\mathrm{tr}\,\boldsymbol{D}_0^2(\mathrm{diag}(\boldsymbol{Q}))^2. \quad (44)\end{aligned}$$

Gathering (41)–(44), we obtain (33). Step 2 is proved.

Step 3: Proof of (34) and (35): We begin with some identities. Write $\boldsymbol{Q}(z) = [q_{ij}(z)]_{i,j=1}^N$ and $\widetilde{\boldsymbol{Q}}(z) = [\tilde{q}_{ij}(z)]_{i,j=1}^K$. Denote by $\boldsymbol{y}_k$ the column number k of $\boldsymbol{Y}$ and by $\boldsymbol{\xi}_n$ the row number n of $\boldsymbol{Y}$. Denote by $\boldsymbol{Y}^k$ the matrix that remains after deleting column k from $\boldsymbol{Y}$ and by $\boldsymbol{Y}_n$ the matrix that remains after deleting row n from $\boldsymbol{Y}$. Finally, write $\boldsymbol{Q}_k(z) = (\boldsymbol{Y}^k\boldsymbol{Y}^{k*} - z\boldsymbol{I})^{-1}$ and $\widetilde{\boldsymbol{Q}}_n(z) = (\boldsymbol{Y}_n^*\boldsymbol{Y}_n - z\boldsymbol{I})^{-1}$. The following formulas can be established easily (see for instance [31, Secs. 0.7.3 and 0.7.4])

$$q_{nn}(-\rho) = \frac{1}{\rho(1 + \boldsymbol{\xi}_n\widetilde{\boldsymbol{Q}}_n(-\rho)\boldsymbol{\xi}_n^*)}$$

$$\tilde{q}_{kk}(-\rho) = \frac{1}{\rho(1 + \boldsymbol{y}_k^*\boldsymbol{Q}_k(-\rho)\boldsymbol{y}_k)} \quad (45)$$

$$\boldsymbol{Q} = \boldsymbol{Q}_k - \frac{\boldsymbol{Q}_k\boldsymbol{y}_k\boldsymbol{y}_k^*\boldsymbol{Q}_k}{1 + \boldsymbol{y}_k^*\boldsymbol{Q}_k\boldsymbol{y}_k}. \quad (46)$$

Lemma 7: The following hold true.
1) (*Rank one perturbation inequality*) The resolvent $\boldsymbol{Q}_k(-\rho)$ satisfies $|\mathrm{tr}\,\boldsymbol{A}(\boldsymbol{Q} - \boldsymbol{Q}_k)| \le \|\boldsymbol{A}\|/\rho$ for any $N \times N$ matrix $\boldsymbol{A}$.
2) Let Assumptions **A1**–**A3** hold. Then
$$\max_{1\le n\le N}\mathbb{E}(q_{nn}(-\rho) - t_n(-\rho))^2 \le \frac{C}{K}. \quad (47)$$
The same conclusion holds true if q_{nn} and t_n are replaced with $\tilde{q}_{kk}$ and $\tilde{t}_k$, respectively.

We are now in position to prove (34). First, notice that

$$\begin{aligned}\mathbb{E}\left|q_{nn}^2 - t_n^2\right| &= \mathbb{E}|q_{nn} - t_n|(q_{nn} + t_n)\\&\le \sqrt{\mathbb{E}(q_{nn} - t_n)^2}\sqrt{\mathbb{E}(q_{nn} + t_n)^2}\\&\le \frac{2}{\rho}\sqrt{\mathbb{E}(q_{nn} - t_n)^2}. \quad (48)\end{aligned}$$

Now

$$\begin{aligned}&\frac{1}{K}\mathbb{E}\left|\mathrm{tr}\,\boldsymbol{D}_0^2(\mathrm{diag}(\boldsymbol{Q})^2 - \boldsymbol{T}^2)\right|\\&\le \frac{1}{K}\sum_{n=1}^N \sigma_{0,n}^4\mathbb{E}\left|q_{nn}^2 - t_n^2\right|\\&\le \frac{\sigma_{\max}^4 N}{K}\max_{1\le n\le N}\mathbb{E}\left|q_{nn}^2 - t_n^2\right|\\&\le \frac{2\sigma_{\max}^4 N}{\rho K}\sqrt{\max_{1\le n\le N}\mathbb{E}(q_{nn} - t_n)^2}\xrightarrow[K\to\infty]{} 0\end{aligned}$$

where the last inequality follows from (48) together with Lemma 7–part 2. Convergence (34) is established.

We now establish the system of (35). Our starting point is the identity

$$\begin{aligned}\boldsymbol{Q} &= \boldsymbol{T} + \boldsymbol{T}(\boldsymbol{T}^{-1} - \boldsymbol{Q}^{-1})\boldsymbol{Q}\\&= \boldsymbol{T} + \frac{\rho}{K}\boldsymbol{T}\,\mathrm{diag}(\mathrm{tr}\,\widetilde{\boldsymbol{D}}_1\widetilde{\boldsymbol{T}},\ldots,\mathrm{tr}\,\widetilde{\boldsymbol{D}}_N\widetilde{\boldsymbol{T}})\boldsymbol{Q} - \boldsymbol{T}\boldsymbol{Y}\boldsymbol{Y}^*\boldsymbol{Q}.\end{aligned}$$

Using this identity, we develop $U_\ell = \frac{1}{K}\mathrm{tr}\,\boldsymbol{D}_0\boldsymbol{Q}\boldsymbol{D}_\ell\boldsymbol{Q}$ as

$$\begin{aligned}U_\ell &= \frac{1}{K}\mathrm{tr}\,\boldsymbol{D}_0\boldsymbol{Q}\boldsymbol{D}_\ell\boldsymbol{T}\\&\quad + \frac{\rho}{K^2}\mathrm{tr}\,\boldsymbol{D}_0\boldsymbol{Q}\boldsymbol{D}_\ell\boldsymbol{T}\,\mathrm{diag}(\mathrm{tr}\,\widetilde{\boldsymbol{D}}_1\widetilde{\boldsymbol{T}},\ldots,\mathrm{tr}\,\widetilde{\boldsymbol{D}}_N\widetilde{\boldsymbol{T}})\boldsymbol{Q}\\&\quad - \frac{1}{K}\mathrm{tr}\,\boldsymbol{D}_0\boldsymbol{Q}\boldsymbol{D}_\ell\boldsymbol{T}\boldsymbol{Y}\boldsymbol{Y}^*\boldsymbol{Q}\\&\triangleq X_1 + X_2 - X_3. \quad (49)\end{aligned}$$

Lemma 4–part 2 with $\boldsymbol{S} = \boldsymbol{D}_0\boldsymbol{D}_\ell\boldsymbol{T}$ yields

$$X_1 = \frac{1}{K}\mathrm{tr}\,\boldsymbol{D}_0\boldsymbol{D}_\ell\boldsymbol{T}^2 + \epsilon_1 \quad (50)$$

where $\mathbb{E}|\epsilon_1| \le \sqrt{\mathbb{E}\epsilon_1^2} \le C/K$. Consider now the term

$$X_3 = \frac{1}{K}\sum_{k=1}^K \mathrm{tr}\,\boldsymbol{D}_0\boldsymbol{Q}\boldsymbol{D}_\ell\boldsymbol{T}\boldsymbol{y}_k\boldsymbol{y}_k^*\boldsymbol{Q}.$$

Using (45) and (46), we have

$$\boldsymbol{y}_k^*\boldsymbol{Q} = \left(1 - \frac{\boldsymbol{y}_k^*\boldsymbol{Q}\boldsymbol{y}_k}{1 + \boldsymbol{y}_k^*\boldsymbol{Q}\boldsymbol{y}_k}\right)\boldsymbol{y}_k^*\boldsymbol{Q}_k = \rho\tilde{q}_{kk}\boldsymbol{y}_k^*\boldsymbol{Q}_k.$$

Hence

$$\begin{aligned}X_3 &= \frac{\rho}{K}\sum_{k=1}^K \tilde{q}_{kk}\boldsymbol{y}_k^*\boldsymbol{Q}_k\boldsymbol{D}_0\boldsymbol{Q}\boldsymbol{D}_\ell\boldsymbol{T}\boldsymbol{y}_k\\&= \frac{\rho}{K}\sum_{k=1}^K \tilde{t}_k\boldsymbol{y}_k^*\boldsymbol{Q}_k\boldsymbol{D}_0\boldsymbol{Q}\boldsymbol{D}_\ell\boldsymbol{T}\boldsymbol{y}_k\\&\quad + \frac{\rho}{K}\sum_{k=1}^K (\tilde{q}_{kk} - \tilde{t}_k)\boldsymbol{y}_k^*\boldsymbol{Q}_k\boldsymbol{D}_0\boldsymbol{Q}\boldsymbol{D}_\ell\boldsymbol{T}\boldsymbol{y}_k\\&\triangleq X_3' + \epsilon_2. \quad (51)\end{aligned}$$

By Cauchy–Schwartz inequality

$$\mathbb{E}|\epsilon_2| \le \frac{\rho}{K}\sum_{k=1}^K \sqrt{\mathbb{E}(\tilde{q}_{kk} - \tilde{t}_k)^2}\sqrt{\mathbb{E}(\boldsymbol{y}_k^*\boldsymbol{Q}_k\boldsymbol{D}_0\boldsymbol{Q}\boldsymbol{D}_\ell\boldsymbol{T}\boldsymbol{y}_k)^2}.$$

We have $\mathbb{E}(\boldsymbol{y}_k^*\boldsymbol{Q}_k\boldsymbol{D}_0\boldsymbol{Q}\boldsymbol{D}_\ell\boldsymbol{T}\boldsymbol{y}_k)^2 \le \sigma_{\max}^8\rho^{-6}\mathbb{E}\|\boldsymbol{y}_k\|^4 \le C$. Using in addition Lemma 7–part 2, we obtain

$$\mathbb{E}|\epsilon_2| \le \frac{C}{\sqrt{K}}.$$

Consider X_3'. From (45) and (46), we have $\boldsymbol{Q} = \boldsymbol{Q}_k - \rho\tilde{q}_{kk}\boldsymbol{Q}_k\boldsymbol{y}_k\boldsymbol{y}_k^*\boldsymbol{Q}_k$. Hence, we can develop X_3' as

$$\begin{aligned}X_3' &= \frac{\rho}{K}\sum_{k=1}^K \tilde{t}_k\boldsymbol{y}_k^*\boldsymbol{Q}_k\boldsymbol{D}_0\boldsymbol{Q}_k\boldsymbol{D}_\ell\boldsymbol{T}\boldsymbol{y}_k\\&\quad - \frac{\rho^2}{K}\sum_{k=1}^K \tilde{t}_k\tilde{q}_{kk}\boldsymbol{y}_k^*\boldsymbol{Q}_k\boldsymbol{D}_0\boldsymbol{Q}_k\boldsymbol{y}_k\boldsymbol{y}_k^*\boldsymbol{Q}_k\boldsymbol{D}_\ell\boldsymbol{T}\boldsymbol{y}_k\\&\triangleq X_4 + X_5. \quad (52)\end{aligned}$$

Consider X_4. Notice that $\boldsymbol{y}_k$ and $\boldsymbol{Q}_k$ are independent. Therefore, by Lemma 1, we obtain

$$
\begin{aligned}
&\boldsymbol{y}_k^* \boldsymbol{Q}_k \boldsymbol{D}_0 \boldsymbol{Q}_k \boldsymbol{D}_\ell \boldsymbol{T} \boldsymbol{y}_k \\
&= \frac{1}{K} \operatorname{tr} \boldsymbol{D}_k \boldsymbol{Q}_k \boldsymbol{D}_0 \boldsymbol{Q}_k \boldsymbol{D}_\ell \boldsymbol{T} + \epsilon_3 \\
&= \frac{1}{K} \operatorname{tr} \boldsymbol{D}_k \boldsymbol{Q} \boldsymbol{D}_0 \boldsymbol{Q} \boldsymbol{D}_\ell \boldsymbol{T} + \epsilon_3 + \epsilon_4
\end{aligned}
$$

where $\mathbb{E}\epsilon_3^2 < CK^{-1}$ by (13). Applying twice Lemma 7–part 1 to

$$
\epsilon_4 = \frac{1}{K} (\operatorname{tr} \boldsymbol{D}_k \boldsymbol{Q}_k \boldsymbol{D}_0 \boldsymbol{Q}_k \boldsymbol{D}_\ell \boldsymbol{T} - \operatorname{tr} \boldsymbol{D}_k \boldsymbol{Q} \boldsymbol{D}_0 \boldsymbol{Q} \boldsymbol{D}_\ell \boldsymbol{T})
$$

yields $|\epsilon_4| < CK^{-1}$. Note in addition that $\sum \tilde{t}_k \boldsymbol{D}_k = \operatorname{diag}(\operatorname{tr} \widetilde{\boldsymbol{D}}_1 \widetilde{\boldsymbol{T}}, \ldots, \operatorname{tr} \widetilde{\boldsymbol{D}}_N \widetilde{\boldsymbol{T}})$. Thus, we obtain

$$
\begin{aligned}
X_4 &= \frac{\rho}{K^2} \operatorname{tr} \left(\sum_{k=1}^K \tilde{t}_k \boldsymbol{D}_k \right) \boldsymbol{Q} \boldsymbol{D}_0 \boldsymbol{Q} \boldsymbol{D}_\ell \boldsymbol{T} + \epsilon_5 \\
&= X_2 + \epsilon_5,
\end{aligned}
\tag{53}
$$

where $\epsilon_5 = \epsilon_3 + \epsilon_4$, which yields $\mathbb{E}|\epsilon_5| \le CK^{-\frac{1}{2}}$.

We now turn to X_5. First introduce the following random variable:

$$
\begin{aligned}
\epsilon_6 =\,& \tilde{t}_k \tilde{q}_{kk} \boldsymbol{y}_k^* \boldsymbol{Q}_k \boldsymbol{D}_0 \boldsymbol{Q}_k \boldsymbol{y}_k \boldsymbol{y}_k^* \boldsymbol{Q}_k \boldsymbol{D}_\ell \boldsymbol{T} \boldsymbol{y}_k \\
&- \tilde{t}_k \tilde{q}_{kk} \left(\frac{1}{K} \operatorname{tr} \boldsymbol{D}_k \boldsymbol{Q}_k \boldsymbol{D}_0 \boldsymbol{Q}_k \right) \left(\frac{1}{K} \operatorname{tr} \boldsymbol{D}_k \boldsymbol{Q}_k \boldsymbol{D}_\ell \boldsymbol{T} \right).
\end{aligned}
$$

Then

$$
\begin{aligned}
|\epsilon_6| \le\,& \frac{1}{\rho^2} \boldsymbol{y}_k^* \boldsymbol{Q}_k \boldsymbol{D}_0 \boldsymbol{Q}_k \boldsymbol{y}_k \left| \boldsymbol{y}_k^* \boldsymbol{Q}_k \boldsymbol{D}_\ell \boldsymbol{T} \boldsymbol{y}_k - \frac{1}{K} \operatorname{tr} \boldsymbol{D}_k \boldsymbol{Q}_k \boldsymbol{D}_\ell \boldsymbol{T} \right| \\
&+ \frac{1}{\rho^2} \left| \boldsymbol{y}_k^* \boldsymbol{Q}_k \boldsymbol{D}_0 \boldsymbol{Q}_k \boldsymbol{y}_k^* - \frac{1}{K} \operatorname{tr} \boldsymbol{D}_k \boldsymbol{Q}_k \boldsymbol{D}_0 \boldsymbol{Q}_k \right| \\
&\times \frac{1}{K} \operatorname{tr} \boldsymbol{D}_k \boldsymbol{Q}_k \boldsymbol{D}_\ell \boldsymbol{T}
\end{aligned}
$$

and one can prove that $\mathbb{E}|\epsilon_6| < CK^{-\frac{1}{2}}$ with help of Lemma 1, together with Cauchy–Schwarz inequality. In addition, we can prove with the help of Lemma 7 that

$$
\begin{aligned}
&\tilde{t}_k \tilde{q}_{kk} \left(\frac{1}{K} \operatorname{tr} \boldsymbol{D}_k \boldsymbol{Q}_k \boldsymbol{D}_0 \boldsymbol{Q}_k \right) \left(\frac{1}{K} \operatorname{tr} \boldsymbol{D}_k \boldsymbol{Q}_k \boldsymbol{D}_\ell \boldsymbol{T} \right) \\
&= \tilde{t}_k^2 \left(\frac{1}{K} \operatorname{tr} \boldsymbol{D}_k \boldsymbol{Q} \boldsymbol{D}_0 \boldsymbol{Q} \right) \left(\frac{1}{K} \operatorname{tr} \boldsymbol{D}_k \boldsymbol{Q} \boldsymbol{D}_\ell \boldsymbol{T} \right) + \epsilon_7 \\
&= \tilde{t}_k^2 \left(\frac{1}{K} \operatorname{tr} \boldsymbol{D}_k \boldsymbol{Q} \boldsymbol{D}_0 \boldsymbol{Q} \right) \left(\frac{1}{K} \operatorname{tr} \boldsymbol{D}_k \boldsymbol{D}_\ell \boldsymbol{T}^2 \right) + \epsilon_7 + \epsilon_8
\end{aligned}
$$

where ϵ_7 and ϵ_8 are random variables satisfying $\mathbb{E}|\epsilon_7| < CK^{-\frac{1}{2}}$ by Lemma 7, and $\max_{k,\ell} \mathbb{E}|\epsilon_8| \le \max_{k,\ell} \sqrt{\mathbb{E}|\epsilon_8|^2} \le CK^{-\frac{1}{2}}$ by Lemma 4–part 2. Using the fact that $\rho^2 \tilde{t}_k^2 = (1 + \frac{1}{K} \operatorname{tr} \boldsymbol{D}_k \boldsymbol{T})^{-2}$, we end up with

$$
\begin{aligned}
X_5 &= -\frac{\rho^2}{K} \sum_{k=1}^K \tilde{t}_k^2 \left(\frac{1}{K} \operatorname{tr} \boldsymbol{D}_k \boldsymbol{Q} \boldsymbol{D}_0 \boldsymbol{Q} \right) \left(\frac{1}{K} \operatorname{tr} \boldsymbol{D}_k \boldsymbol{D}_\ell \boldsymbol{T}^2 \right) + \epsilon_9 \\
&= -\sum_{k=1}^K c_{\ell k} U_k + \epsilon_9
\end{aligned}
\tag{54}
$$

where $c_{\ell k}$ is given by (36), and where $\mathbb{E}|\epsilon_9| < CK^{-\frac{1}{2}}$.

Plugging (50)–(54) into (49), we end up with $U_\ell = \sum_{k=1}^K c_{\ell k} U_k + \frac{1}{K} \operatorname{tr} \boldsymbol{D}_0 \boldsymbol{D}_\ell \boldsymbol{T}^2 + \epsilon$ with $\mathbb{E}|\epsilon| < CK^{-\frac{1}{2}}$. Step 3 is established.

Step 4. Proof of (37): We rely on results of Section V-B, in particular on Lemma 5. Define the following $(K+1) \times 1$ vectors:

$$
\boldsymbol{u} = [U_k]_{k=0}^K, \quad \boldsymbol{d} = \left[\frac{1}{K} \operatorname{tr} \boldsymbol{D}_0 \boldsymbol{D}_k \boldsymbol{T}^2 \right]_{k=0}^K, \quad \boldsymbol{\epsilon} = [\epsilon_k]_{k=0}^K
$$

where the U_k's and ϵ_k's are defined in (35). Recall the definition of the $c_{\ell k}$'s for $0 \le \ell \le K$ and $1 \le k \le K$, define $c_{\ell 0} = 0$ for $0 \le \ell \le K$ and consider the $(K+1) \times (K+1)$ matrix $\boldsymbol{C} = [c_{\ell k}]_{\ell,k=0}^K$.

With these notations, system (35) writes

$$
(\boldsymbol{I}_{K+1} - \boldsymbol{C}) \boldsymbol{u} = \boldsymbol{d} + \boldsymbol{\epsilon}.
\tag{55}
$$

Let $\alpha = \frac{1}{K} \operatorname{tr} \boldsymbol{D}_0^2 \boldsymbol{T}^2$. We have in particular

$$
\boldsymbol{d} = \begin{bmatrix} \alpha \\ \boldsymbol{g} \end{bmatrix}, \quad \boldsymbol{C} = \begin{bmatrix} 0 & \frac{1}{K} \boldsymbol{g}^T \boldsymbol{\Delta}^{-1} \\ \boldsymbol{0} & \boldsymbol{A}^T \end{bmatrix}
$$

(recall that $\boldsymbol{A}$, $\boldsymbol{\Delta}$, and $\boldsymbol{g}$ are defined in the statement of Theorem 3).

Consider a square matrix $\boldsymbol{X}$ which first column is equal to $[1, 0, \ldots, 0]^T$, and partition $\boldsymbol{X}$ as $\boldsymbol{X} = \begin{bmatrix} 1 & \boldsymbol{x}_{01}^T \\ \boldsymbol{0} & \boldsymbol{X}_{11} \end{bmatrix}$. Recall that the inverse of $\boldsymbol{X}$ exists if and only if $\boldsymbol{X}_{11}^{-1}$ exists, and in this case the first row $[\boldsymbol{X}^{-1}]_0$ of $\boldsymbol{X}^{-1}$ is given by

$$
\left[\boldsymbol{X}^{-1} \right]_0 = \begin{bmatrix} 1 & -\boldsymbol{x}_{01}^T \boldsymbol{X}_{11}^{-1} \end{bmatrix}
$$

(see for instance [31]). We now apply these results to the system (55). Due to (55), U_0 can be expressed as

$$
U_0 = [(\boldsymbol{I} - \boldsymbol{C})^{-1}]_0 (\boldsymbol{d} + \boldsymbol{\epsilon}).
$$

By Lemma 5–part 1, $(\boldsymbol{I}_K - \boldsymbol{A}^T)^{-1}$ exists hence $(\boldsymbol{I} - \boldsymbol{C})^{-1}$ exists

$$
\left[(\boldsymbol{I}_{K+1} - \boldsymbol{C})^{-1} \right]_0 = \begin{bmatrix} 1 & \frac{1}{K} \boldsymbol{g}^T \boldsymbol{\Delta}^{-1} (\boldsymbol{I}_K - \boldsymbol{A}^T)^{-1} \end{bmatrix},
$$

and

$$
\begin{aligned}
U_0 =\,& \alpha + \frac{1}{K} \boldsymbol{g}^T \boldsymbol{\Delta}^{-1} \left(\boldsymbol{I} - \boldsymbol{A}^T \right)^{-1} \boldsymbol{g} \\
&+ \epsilon_0 + \frac{1}{K} \boldsymbol{g}^T \boldsymbol{\Delta}^{-1} \left(\boldsymbol{I} - \boldsymbol{A}^T \right)^{-1} \boldsymbol{\epsilon}'
\end{aligned}
$$

with $\boldsymbol{\epsilon}' = [\epsilon_1, \ldots, \epsilon_K]^T$. Gathering the estimates of Section V-B together with the fact that $\|\mathbb{E}\boldsymbol{\epsilon}\|_\infty \le CK^{-\frac{1}{2}}$, we get (37). Step 4 is established, so is Theorem 3.

APPENDIX

A. Proof of Lemma 4

Let us establish (24). The lower bound immediately follows from the representation

$$
t_n = \frac{1}{\rho + \frac{1}{K} \sum_{k=1}^K \frac{\sigma_{nk}^2}{1 + \frac{1}{K} \sum_{\ell=1}^N \sigma_{\ell k}^2 t_\ell}} \overset{(a)}{\ge} \frac{1}{\rho + \sigma_{\max}^2}
$$

where (a) follows from **A2** and $t_\ell(-\rho) \geq 0$. The upper bound requires an extra argument: As proved in [29, Theorem 2.4], the t_n's are STs of probability measures supported by $\mathbb{R}_+$, i.e., there exists a probability measure μ_n over $\mathbb{R}_+$ such that $t_n(z) = \int \frac{\mu_n(dt)}{t-z}$. Thus

$$t_n(-\rho) = \int_0^\infty \frac{\mu_n(dt)}{t+\rho} \leq \frac{1}{\rho}$$

and (24) is proved.

We now briefly justify (25). We have

$$\mathbb{E} |\mathrm{tr}\, \boldsymbol{S}(\boldsymbol{Q}-\boldsymbol{T})|^2 = \mathbb{E} |\mathrm{tr}\boldsymbol{S}(\boldsymbol{Q}-\mathbb{E}\boldsymbol{Q})|^2 + |\mathrm{tr}\boldsymbol{S}(\mathbb{E}\boldsymbol{Q}-\boldsymbol{T})|^2 .$$

In [24, Lemma 6.3] it is stated that $\sup_K \mathbb{E} |\mathrm{tr}\boldsymbol{S}(\boldsymbol{Q}-\mathbb{E}\boldsymbol{Q})|^2 < \infty$. Furthermore, in the proof of [24, Theorem 3.3] it is shown that $\sup_K K\|\mathbb{E}\boldsymbol{Q}-\boldsymbol{T}\| < \infty$, hence

$$|\mathrm{tr}\boldsymbol{S}(\mathbb{E}\boldsymbol{Q}-\boldsymbol{T})| \leq K\|\boldsymbol{S}(\mathbb{E}\boldsymbol{Q}-\boldsymbol{T})\| \leq K\|\mathbb{E}\boldsymbol{Q}-\boldsymbol{T}\|\|\boldsymbol{S}\| < \infty$$

by Lemma 2–part 2. The result follows.

B. Proof of Corollary 1

Recall that in the separable case, $\boldsymbol{D}_k = \tilde{d}_k \boldsymbol{D}$ and $\widetilde{\boldsymbol{D}}_n = d_n \widetilde{\boldsymbol{D}}$. Let $\tilde{\boldsymbol{d}}$ be the $K \times 1$ vector $\tilde{\boldsymbol{d}} = [\tilde{d}_k]_{k=1}^K$. In the separable case, (20) is written as

$$\frac{\Theta^2}{\tilde{d}_0^2} = \frac{1}{K\tilde{d}_0^2}\boldsymbol{g}^T(\boldsymbol{I}-\boldsymbol{A})^{-1}\boldsymbol{\Delta}^{-1}\boldsymbol{g} + \gamma(\mathbb{E}|W_{10}|^4 - 1) \quad (56)$$

where γ is defined in the statement of the corollary. Here, vector $\boldsymbol{g}$ and matrix $\boldsymbol{A}$ are given by

$$\boldsymbol{g} = \gamma \tilde{d}_0 \tilde{\boldsymbol{d}} \quad \text{and}$$
$$\boldsymbol{A} = \left[\frac{1}{K}\frac{\frac{1}{K}\mathrm{tr}\boldsymbol{D}_\ell \boldsymbol{D}_m \boldsymbol{T}^2}{(1+\frac{1}{K}\mathrm{tr}\boldsymbol{D}_\ell \boldsymbol{T})^2}\right]_{\ell,m=1}^K = \frac{\gamma}{K}\boldsymbol{\Delta}^{-1}\tilde{\boldsymbol{d}}\tilde{\boldsymbol{d}}^T .$$

By the matrix inversion lemma [31], we have

$$\frac{1}{K\tilde{d}_0^2}\boldsymbol{g}^T(\boldsymbol{I}-\boldsymbol{A})^{-1}\boldsymbol{\Delta}^{-1}\boldsymbol{g}$$
$$= \frac{\gamma^2}{K}\tilde{\boldsymbol{d}}^T\left(\boldsymbol{\Delta} - \frac{\gamma}{K}\tilde{\boldsymbol{d}}\tilde{\boldsymbol{d}}^T\right)^{-1}\tilde{\boldsymbol{d}}$$
$$= \frac{\gamma^2}{K}\tilde{\boldsymbol{d}}^T\left(\boldsymbol{\Delta}^{-1} + \frac{\gamma}{K}\frac{1}{1-\frac{\gamma}{K}\tilde{\boldsymbol{d}}^T\boldsymbol{\Delta}^{-1}\tilde{\boldsymbol{d}}}\boldsymbol{\Delta}^{-1}\tilde{\boldsymbol{d}}\tilde{\boldsymbol{d}}^T\boldsymbol{\Delta}^{-1}\right)\tilde{\boldsymbol{d}} .$$

Noticing that

$$\frac{1}{K}\tilde{\boldsymbol{d}}^T\boldsymbol{\Delta}^{-1}\tilde{\boldsymbol{d}} = \frac{1}{K}\sum_{k=1}^K \frac{\tilde{d}_k^2}{\left(1+\frac{1}{K}\mathrm{tr}\boldsymbol{D}_k\boldsymbol{T}\right)^2}$$
$$= \frac{\rho^2}{K}\sum_{k=1}^K \tilde{d}_k^2 \tilde{t}_k^2 = \rho^2\tilde{\gamma}$$

we obtain

$$\frac{1}{K\tilde{d}_0^2}\boldsymbol{g}^T(\boldsymbol{I}-\boldsymbol{A})^{-1}\boldsymbol{\Delta}^{-1}\boldsymbol{g} = \gamma\frac{\rho^2\gamma\tilde{\gamma}}{1-\rho^2\gamma\tilde{\gamma}} .$$

Plugging this equation into (56), we obtain (22).

C. Proof of Corollary 2

Expression (22) can be rewritten as

$$\Omega_K^2 = \left(\frac{\gamma}{1-\rho^2\gamma\tilde{\gamma}} + \gamma\left(\mathbb{E}|W_{10}|^4 - 2\right)\right) .$$

As $\boldsymbol{D} = \boldsymbol{I}$, we have

$$\delta(\rho) = \frac{1}{K}\mathrm{tr}\boldsymbol{T}(-\rho) = \frac{N}{K}\int \frac{1}{t+\rho}\pi_K(dt) \quad (57)$$

and from (17) we have $\boldsymbol{T}(-\rho) = \left(\rho(1+\tilde{\delta}(\rho))\right)^{-1}\boldsymbol{I}$. It results in

$$\gamma = \frac{1}{K}\mathrm{tr}\boldsymbol{T}^2 = \frac{K}{N}\left(\frac{1}{K}\mathrm{tr}\boldsymbol{T}\right)^2 = \frac{K}{N}\delta^2$$

hence, $\gamma\left(\mathbb{E}|W_{10}|^4 - 2\right)$ coincides with the second term at the right-hand side of (23). We now turn to the term $\gamma/(1-\rho^2\gamma\tilde{\gamma})$. By (17) we have

$$\tilde{\delta}'(\rho) = \frac{d\tilde{\delta}(\rho)}{d\rho} = \frac{1}{K}\sum_{k=1}^K\left(\frac{\tilde{d}_k}{\rho(1+\tilde{d}_k\delta(\rho))}\right)'$$
$$= -\frac{1}{K}\sum_{k=1}^K\frac{\tilde{d}_k}{\rho^2(1+\tilde{d}_k\delta)} - \frac{\delta'}{K}\sum_{k=1}^K\frac{\tilde{d}_k^2}{\rho(1+\tilde{d}_k\delta)^2}$$
$$= -\frac{1}{\rho}\tilde{\delta} - \rho\delta'\tilde{\gamma}. \quad (58)$$

Symmetrically, we have $\delta' = -\frac{1}{\rho}\delta - \rho\tilde{\delta}'\gamma$. Replacing $\tilde{\delta}'$ by its value in (58), we obtain

$$(1-\rho^2\gamma\tilde{\gamma})\delta' = -\frac{1}{\rho}\delta + \gamma\tilde{\delta}. \quad (59)$$

Now we have $\tilde{\delta}\boldsymbol{T} = \left(\tilde{\delta}/(\rho(1+\tilde{\delta}))\right)\boldsymbol{I} = \rho^{-1}\boldsymbol{I} - \boldsymbol{T}$. Hence

$$\gamma\tilde{\delta} = \frac{1}{K}\mathrm{tr}\left(\boldsymbol{T}(\tilde{\delta}\boldsymbol{T})\right) = \frac{1}{K}\mathrm{tr}\left(\boldsymbol{T}\left(\frac{1}{\rho}\boldsymbol{I}-\boldsymbol{T}\right)\right) = \frac{1}{\rho}\delta - \gamma.$$

Plugging into (59) and using (57), we end up with

$$\frac{\gamma}{1-\rho^2\gamma\tilde{\gamma}} = -\delta' = \frac{N}{K}\int\frac{1}{(t+\rho)^2}\pi_K(dt)$$

where the second equality is due to (57). Equation (23) is proven.

D. Proof of Lemma 7

The proof of Part 1 can be found in [24, Proof of Lemma 6.3] (see also [14, Lemma 2.6]). Let us prove Part 2. We have from (11) and (45)

$$|q_{nn}(-\rho) - t_n(-\rho)|$$
$$= \frac{1}{\rho(1+\frac{1}{K}\mathrm{tr}\widetilde{\boldsymbol{D}}_n\widetilde{\boldsymbol{T}})(1+\boldsymbol{\xi}_n\widetilde{\boldsymbol{Q}}_n\boldsymbol{\xi}_n^*)}\left|\boldsymbol{\xi}_n\widetilde{\boldsymbol{Q}}_n\boldsymbol{\xi}_n^* - \frac{1}{K}\mathrm{tr}\widetilde{\boldsymbol{D}}_n\widetilde{\boldsymbol{T}}\right|$$
$$\leq \frac{1}{\rho}\left|\boldsymbol{\xi}_n\widetilde{\boldsymbol{Q}}_n\boldsymbol{\xi}_n^* - \frac{1}{K}\mathrm{tr}\widetilde{\boldsymbol{D}}_n\widetilde{\boldsymbol{T}}\right|.$$

Hence

$$\mathbb{E}(q_{nn}-t_n)^2 \leq \frac{2}{\rho}\mathbb{E}\left(\boldsymbol{\xi}_n\widetilde{\boldsymbol{Q}}_n\boldsymbol{\xi}_n^* - \frac{1}{K}\mathrm{tr}\widetilde{\boldsymbol{D}}_n\widetilde{\boldsymbol{Q}}\right)^2$$
$$+ \frac{2}{\rho K^2}\mathbb{E}\left(\mathrm{tr}\widetilde{\boldsymbol{D}}_n(\widetilde{\boldsymbol{Q}}-\widetilde{\boldsymbol{T}})\right)^2 \leq \frac{C}{K}$$

by Lemma 1 and Lemma 4–part 2, which proves (47).

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