Positive versus Negative Transfer in Multi-Task Learning in the High-Dimensional Regime

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1 Introduction

Multi-task learning has recently become a powerful paradigm to solve complex prediction tasks in computer vision [1, 2], natural language processing [3, 4] and numerous other areas [5]. By combining multiple information sources, multi-task learning allows new information to be shared across different sources in a model [6]. At the same time, it is well-known that the performance of multi-task learning methods depends on the relationship between the information sources. *Negative transfer*, the phenomenon where for a particular prediction task, multi-task learning performs worse than single-task learning, is prevalent when the information sources are heterogeneous [7, 8]. While numerous studies have sought to alleviate negative transfer when it occurs [9], a rigorous understanding to the contributing causes of negative transfer has remained elusive in the literature [10]. In this work, we consider learning multiple high-dimensional linear regression tasks to better understand when and why negative transfer happens. We provide theoretical and practical insights to show how task data affects the transfer of information.

Identifying when negative transfer occurs requires comparing the performance of multi-task learning directly to single-task learning. The technical challenge is to develop generalization bounds that are able to scale tightly with the qualities of task data, such as the number of datapoints. In classical Rademacher or VC based theory of multi-task learning [11, 12, 13], the generalization bounds are usually presented in a way so that the error goes down as more labeled data is added. On the other hand, we have observed that adding more labeled data does not always improve performance in multi-task learning. More recent work has shown the benefit of learning multi-task representations for certain half-spaces [14] and multipl sparse regressions [15, 16].

In this work, we consider multiple high-dimensional linear regression tasks as input and focus on predicting a particular task that only has limited amount of labeled data. Following [19, 20], we assume that for every task $1 \le i \le t$, its features are random vectors $x = \Sigma_i^{1/2} z$, where $z \in \mathbb{R}^p$ consists of i.i.d. entries with mean zero and unit variance, and $\Sigma_i \in \mathbb{R}^{p \times p}$ is a positive semidefinite matrix. The label of x is given by $y_i = x^\top \beta_i + \varepsilon$, where $\beta_i \in \mathbb{R}^p$ denotes the ground truth parameters for task i and ε denotes random noise with mean zero and variance σ^2 . Importantly, we assume that for the target task t, its training data size n_t is a small constant ρ_t times p. We use a hard parameter sharing architecture that contains a shared body $B \in \mathbb{R}^{p \times r}$ for all tasks and a separate prediction head $\{W_i \in \mathbb{R}^r\}_{i=1}^t$ for each task [10, 17]. This corresponds to minimizing the following objective.

$$f(B; W_1, \dots, W_t) = \sum_{i=1}^t ||X_i B W_i - Y_i||^2.$$
(1.1)

Let $\hat{\beta}_t^{\text{MTL}} = BW_t$ denote the optimal predictor obtained from solving equation (1.1) for task t. Our goal is to compare the test error of $\hat{\beta}_t^{\text{MTL}}$, denoted as $te(\hat{\beta}_t^{\text{MTL}})$, to $te(\hat{\beta}_t^{\text{STL}})$, where $\hat{\beta}_t^{\text{STL}}$ is likewise obtained from equation (1.1) but with task t in isolation.

Main results. We begin by observing that $te(\hat{\beta}_t^{\text{MTL}})$ can be decomposed into two parts, a variance part that is reduced from $te(\hat{\beta}_t^{\text{STL}})$, and a bias part that captures the difference between β_t and the rest of β 's. We term the bias part as *model shift bias*. Intuitively, whether $te(\hat{\beta}_t^{\text{MTL}})$ is lower than $te(\hat{\beta}_t^{\text{STL}})$ depends on the trade-off between the amount of variance reduced and the model shift bias part. For the high-dimensional regime when p goes to infinity, we derive the asymptotic limit of $te(\hat{\beta}_t^{\text{MTL}}) - te(\hat{\beta}_t^{\text{STL}})$ as a function of the number of datapoints $\{n_1, n_2\}$, the covariance matrices $\{\Sigma_1, \Sigma_2\}$, the ground truth parameters $\{\beta_1, \beta_2\}$, and a certain fixed value derived from solving equation (1.1) (see Theorem 3.1 for the statement). Then, we show a similar result for any number of tasks that have the same covariates, i.e. the X_i 's are equal to each other. This setting is prevalent in applications of multi-task learning to image classification, where there are multiple prediction labels/tasks for every image [1, 18]. Based on the technical result, we provide theoretical insights over when multi-task learning performs better than single-task learning.

Our first theoretical insight is to explain a sharp transition from positive to negative transfer as a parameter of task model dissimilarity $\|\beta_1 - \beta_2\|^2$. In Proposition 3.3, we provide the trade-off between $\|\beta_1 - \beta_2\|^2$ and a certain function $\Phi(\rho_1, \rho_2)$ to determine the type of transfer.

Our second theoretical insight is to show that performing multi-task learning reduces the need for labeled data, which is a key finding of Taskonomy [2]. We define the *data efficiency ratio* as the smallest $\alpha \in (0,1)$ such that if we only use an α fraction of labeled data, then the test error of $\hat{\beta}_t^{\text{MTL}}$ matches that of the $\hat{\beta}_t^{\text{STL}}$ on the entire set. In Proposition 3.6, we show that the data entire ratio of an illustrative example is at most $\frac{1}{2\rho_t}$.



Our third theoretical insight is that as the data size becomes more imbalanced between the source and target task, having mis-aligned task covariance matrices can result in suboptimal performance. In Proposition 3.8, we show that as n_1/n_2 becomes large, the best performing source task has have the same covariance matrix as the target task. On the other hand, when n_1/n_2 is small, there are counter examples where having the same covariance matrix is not necessarily the optimal choice.

Finally, while focus is multi-task learning, our setup and technical tools are also applicable to transfer learning settings. In Section 4, we provide a case study of the transfer function used in Taskonomy [2]. We observe a similar trade-off between model-shift bias and plane for this setting.

A crucial technical tool that we develop for showing Theorem 3.1 is the asymptotic limit of the trace of $(X_1^\top X_1 + X_2^\top X_2)^{-1}$ (for the setting of two tasks), which may be of independent interest. The result extends a well-known result on the trace of $(X_1^\top X_1)^{-1}$ [22] Furthermore, the error bound for the asymptotic limit is nearly optimal.

Experimental results. We show practical implications of our theoretical insights on text and image classification tasks. First, we provide a metric to determine the test accuracies of single-task models. Second, and the test accuracies of the tasks becomes more imbalanced, aligning the covariances of the tasks becomes more beneficial.

2 Preliminaries



Recall that we have t labeled tasks available, denoted by $(X_1,Y_1),(X_2,Y_2),\ldots,(X_t,Y_t)$. Following [19, 20], we assume that for each task $i=1,2,\ldots,t$, every feature vector is generated as $x=\sum_i^{1/2}z$, where $z\in\mathbb{R}^p$ is a random vector with i.i.d. entries of mean zero and unit variance and $\Sigma_i\in\mathbb{R}^{p\times p}$ is a positive semidefinite matrix. Without loss of generality, let the t-th task denote the target task. For an estimator $\hat{\beta}\in\mathbb{R}^p$, the test error for target task is defined as

$$te_t(\hat{\beta}) := \underset{z}{\mathbb{E}} \left[\underset{\varepsilon_t}{\mathbb{E}} \left[\left(\left(\Sigma_t^{1/2} z \right)^\top \hat{\beta} - \left(\Sigma_t^{1/2} \right)^\top \beta_t \right)^2 \right] \right] = \underset{\varepsilon_t}{\mathbb{E}} \left[\left(\hat{\beta} - \beta_t \right)^\top \Sigma_t (\hat{\beta} - \beta_t) \right].$$

Let $\hat{\beta}_t^{\text{STL}} = (X_t^{\top} X_t)^{-1} X_t^{\top} Y_t$ denote the linear regression estimator, which can be also be obtained from equation (1.1) by using the target task alone. The bias-variance trade-off [21] says that

$$te_t(\hat{\beta}) = \left\| \underset{\varepsilon_t}{\mathbb{E}} \left[\hat{\beta} \right] - \beta_t \right\|^2 + \underset{\varepsilon_t}{\mathbb{E}} \left[\hat{\beta} - \underset{\varepsilon_t}{\mathbb{E}} \left[\hat{\beta} \right] \right].$$

Recall that we assume $n_i > p$ for every $1 \leqslant i \leqslant t$. Hence the bias of $te_t(\hat{\beta}_t^{\text{STL}})$ is zero and its test error is equal to variance, given by $\sigma^2 \cdot \text{Tr}[(X_i^\top X_i)^{-1}\Sigma_i]$. We focus on a setting where the data size n_t of the target task is equal to a small constant $\rho_t > 1$ times p. This setting captures the need for adding more labeled data to reduce the test error of the target task.

To illustrate our intuition, we begin by considering the multi-task learning setting with one source task. Recall that $\hat{\beta}_t^{\text{MTL}}$ is defined as BW_t after solving equation (1.1). We decompose the test error of $\hat{\beta}_t^{\text{MTL}}$ on the targe task into two parts (to be derived in Appendix A.1)

$$te_t(\hat{\beta}_t^{\text{MTL}}) = \hat{v}^2 \left\| \Sigma_2^{1/2} (\hat{v}^2 X_1^\top X_1 + X_2^\top X_2)^{-1} X_1^\top X_1 (\beta_1 - \hat{v}\beta_2) \right\|^2$$
 (2.1)

$$+ \sigma^{2} \cdot \text{Tr} \left[(\hat{v}^{2} X_{1}^{\top} X_{1} + X_{2}^{\top} X_{2})^{-1} \Sigma_{2} \right], \tag{2.2}$$

where $\hat{v} = W_2/W_1$ denotes the ratio of the output layer weights. By comparing the bias and variance of $te_t(\hat{\beta}_t^{\text{MTL}})$ to $te_t(\hat{\beta}_t^{\text{STL}})$, we observe the following two effects, respectively.

Model-shift bias. We term equation (2.1) as *model-shift bias*, which captures how similar β_1 and β_2 are. This part introduces a negative effect to multi-task learning. The scaling term \hat{v} corresponds to the intuition that the shared module B learns a subspace for both tasks. The output layer of each task scales the direction of B suitably to fit the task.

Variance reduction. Next, it is not hard to verify that equation (2.2), the variance of $\hat{\beta}_t^{\text{MTL}}$, is always smaller than the variance of $\hat{\beta}_t^{\text{STL}}$. This part introduces a positive variance reduction effect to performing multi-task learning. Hence, whether $te(\hat{\beta}_t^{\text{MTL}})$ is lower than $te(\hat{\beta}_t^{\text{STL}})$ is determined by the trade-off between two effects: (i) the positive effect from variance reduction; (ii) the negative effect from model shift bias.

Notations. We refer to the last (t-th) task as the target task and the first task as the source task when there are only two tasks. When there is no ambiguity, we drop the subscript t from $te_t(\hat{\beta}_t^{\text{MTL}})$ to $te(\hat{\beta}_t^{\text{MTL}})$ for simplicity. We call $M = \Sigma_1^{1/2} \Sigma_2^{-1/2}$ the covariate shift matrix.

3 Revisiting the Bias-Variance Trade-off for Multi-Tas earning

We revisit the classical bias-variance trade-off in the context of multi-task learning for linear regressions, in particular the setup of equation (1.1). For the high-dimensional setting where the data size of the target task is a small constant times dimension p, we quantify the trade-off as a function of the task models $\{\beta_i\}_{i=1}^t$, data size $\{n_i\}_{i=1}^t$, and covariance matrices $\{\Sigma_i\}_{i=1}^t$. After describing this technical result, we provide theoretical insights over when multi-task learning performs better than single-task learning. Our analysis reveals three contributing causes of negative transfer, including task model dissimilarity, unbalanced data size, and covariate shift, which we describe in detail for the rest of the section.

3.1 Determining Positive versus Negative Transfer

In order to study the trade-off between model-shift bias and variance reduction, we need tight concentration bounds to quantify both effects. For this purpose, we consider the high-dimensional setting where p goes to infinity. A well-known result for this setting states that the asymptotic limit of $te(\hat{\beta}_t^{\rm STL})$ is $\frac{\sigma^2}{\rho_t-1}$ (e.g. Chapter 6 of [22]), which scales with the data size and noise level of the target task. However, this result only applies to the single-task setting. Therefore, we focus on deriving the asymptotic limit of $te(\hat{\beta}_t^{\rm MTL})$ when p goes to infinity.

Our first result applies to the setting of two tasks where their covariance matrices may be arbitrarily different. We derive a precise limit of $\operatorname{Tr}\left[(X_1^{\top}X_1 + \bigcup_{i=1}^{T}X_i)^{-1}\Sigma_i\right]$, which is a deterministic function

that only depends on Σ_1, Σ_2 and $n_1/p, n_2/p$ (see Lemma A.2 in Appendix A.1 for the result). Based on the result, we show how to determine positive versus negative transfer as follows.

Theorem 3.1 (Informal). Let $X_i \in \mathbb{R}^{n_i \times p}$ and $Y_i = X_i \beta_i + \varepsilon_i$, for i = 1, 2. Suppose that $n_1 = \rho_1 p$ and $n_2 = \rho_2 p$, where $\rho_1 > 1$ and $\rho_2 > 1$ are fixed constants. There exists two deterministic functions Δ_{β} and Δ_{var} and a small deterministic error δ that only depend on $\{\hat{v}, \Sigma_1, n_1, n_2, \beta_1, \beta_2\}$ such that

- If $\Delta_{var} \Delta_{\beta} \geqslant \delta$, then who $te(\hat{\beta}_{t}^{MTL}) < te(\hat{\beta}_{t}^{STL})$.
- If $\Delta_{var} \Delta_{\beta} \leq \delta$, then who $te(\hat{\beta}_{t}^{MTL}) > te(\hat{\beta}_{t}^{STL})$.

Theorem 3.1 shows nearly tight bounds on the trade-off between model-shift bias and variance reduction. The bounds get tighter and tighter as ρ_1 increases. A formal version of Theorem 3.1 is presented in Theorem A.4 and its proof is presented in Appendix A.1.1. While the general form of Δ_{var} and Δ_{β} can be quite complex, we show that they provide nice interpretation for several simple settings in the rest of the section.

Our second result applies to the setting of any number of tasks with the same covariates. Since the tasks all have the same number of datapoints and covariance matrix, the trade-off between model shift bias and variance will be captured by their task models $\{\beta_i\}_{i=1}^k$. Let $B^* = [\beta_1, \beta_2, \dots, \beta_k] \in \mathbb{R}^{p \times k}$ denote the underlying task model parameters. We can also derive a tight trade-off between model shift bias and variance reduction for this setting. The formal statement is stated in Theorem A.7 and its proof can be found in Appendix A.1.2.

3.2 Task Model Similarity

assumptions as shown in the supplementary material.

A folklore intuition sing the seminal work of Caruana [6] is that task relatedness determines how well multi-task learning forms. We formalize the intuition using our setup and technical tools. For a simple example of two tasks, we show that whether $\hat{\beta}_t^{\text{MTL}} \leq \hat{\beta}_t^{\text{STL}}$ is determined by $\|\beta_1 - \beta_2\|^2$, the distance of the task models. We derive a sharp threshold when positive transfer transitions to negative transfer, as a function of model distance $\|\beta_1 - \beta_2\|^2$, data size n_1 abd n_2 and the noise level σ . Example 3.2. Consider a setting where $\Sigma_1 = \Sigma_2 = \operatorname{Id}$, in other words there is no covariate shift between the two tasks. For the noises, we assume that ε_i has i.i.d. entries with mean zero and variances σ_i^2 , i = 1, 2. For the task models, suppose that β_2 has i.i.d. entries with mean zero avariance κ^2 and $\beta_1 - \beta_2$ also has i.i.d. entries with mean 0 and variance d^2 . We have $n_i = c_i$ data points from each task, for i = 1, 2. For simplicity, we assume that all the random variables have

We illustrate the example with a simulation. We consider a setting where p = 200, $n_1 = 90p$, $n_2 = 30p$. We fix the target task and vary the source task, in particular the parameter d which determines $\|\beta_1 - \beta_2\|$. Figure 1a shows the result. Based on Theorem A.4, we derive the transition threshold in the following proposition. Define the following function

subexponential decay, while keeping in mind that our results can be applied under weaker moments

$$\Phi(\rho_1, \rho_2) = \frac{(\rho_1 + \rho_2 - 1)^2}{\rho_1(\rho_1 + \rho_2)(\rho_2 - 1)}$$

Proposition 3.3. In the setting of Example 3.2 with $\sigma_1 = \sigma_2 = \sigma$, assume that $\rho_1 > 40$ is a fixed constant. Whether $te(\hat{\beta}_t^{MTL})$ is lower than $te(\hat{\beta}_t^{STL})$ is determined by the ratio between the model distance and the noise level:

- If $d^2 < \frac{2}{p} \Phi(\rho_1, \rho_2)$, then whp we have that $te(\hat{\beta}_t^{MTL}) < te(\hat{\beta}_t^{STL})$. If $d^2 \geqslant \frac{2\sigma^2}{p} \cdot \Phi(\rho_1, \rho_2)$, then whp we have that $te(\hat{\beta}_t^{MTL}) \geqslant te(\hat{\beta}_t^{STL})$.

We remark that the constants 1/2 and 2 in the statement above are chosen for the simplicity of presentation. In general we can close the gap between the two constants by increasing ρ_1 . While the statement of Theorem 3.1 only applies to the asymptotic regime of p going to infinity, Proposition 3.3 explains most of the observations in Figure 1a when p = 200.

The proof of Proposition 3.3 involves two parts. First, in equation (2.2), the positive variance reduction effect scales with $n_1 = \rho_1 p$, the number of source task data points. Second, we show that

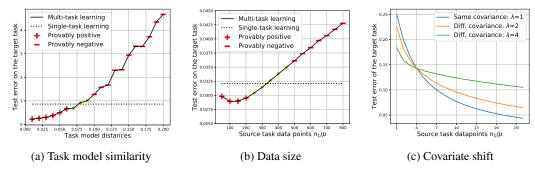


Figure 1: Comparing the test error of multi-task learning to single-task learning: (a) positive to negative transfer as the task model dissimilarity increases (cf. Proposition 3.3) (b) positive to negative transfer as the source task data size increases (cf. Proposition 3.5); (c)

the negative effect of model-shift bias scales with pd^2 , which is the expectation of $\|\beta_1 - \beta_2\|^2$. The proof, which is based on Theorem 3.1, can be found in Appendix A.2.

Implications. We furthe that multi-task learning is particular powerful when the labeled data of the source task is less misy compared to the target task. Consider a more general setting of Proposition 3.3 where the noise level σ_1 of task 1 differs from the noise level σ_2 of task 2. We derive a sharp transition from positive to negative transfer as a parameter of σ_1^2 .

Proposition 3.4. In the setting of Example 3.2 with d being fixed but σ_1 varies, assume that $\rho_1 > 40$ is a fixed constant and $d^2 < \frac{2\sigma_2^2}{3p} \cdot \Phi(\rho_1, \rho_2)$. We derive the following transition as a parameter of σ_1^2 :

$$\bullet \ \ \textit{If} \ \sigma_1^2 \leqslant pd^2 \cdot \rho_1 + \left(1 + \tfrac{1}{2} \cdot \Phi(\rho_1, \rho_2)\right) \cdot \sigma_2^2, \ \textit{then whp } te(\hat{\beta}_t^{\textit{MTL}}) < te(\hat{\beta}_t^{\textit{STL}}).$$

• If
$$\sigma_1^2 > pd^2 \cdot \rho_1 + (1 + 2\Phi(\rho_1, \rho_2)) \cdot \sigma_2^2$$
, then who $te(\hat{\beta}_t^{MTL}) > te(\hat{\beta}_t^{STL})$.

As a corollary, if $\sigma_1^2 \leqslant \sigma_2^2$, then we always get positive transfer. The proof of Proposition 3.4 is similar to Proposition 3.3. The details can be found in Appendix A.2.

3.3 Data Size

In classical Rademacher or VC based theory of multi-task learning, adding more labeled data improves the generalization performance of a model. On the other hand, we have observed that adding more labeled data does not always improve performance in multi-task learning. Using Example 3.2, we analyze the effect of varying the source task data size.

Proposition 3.5. In the setting of Example 3.2, assume that $\rho_1 > 40$ and $\rho_2 > 500$ are fixed constants. We have the following conditions to determine whether $te(\hat{\beta}_t^{MTL})$ is lower than $te(\hat{\beta}_t^{STL})$:

• If
$$d^2 \leqslant \frac{\sigma^2}{2p(\rho_2-1)}$$
, then whp $te(\hat{\beta}_t^{MTL}) < te(\hat{\beta}_t^{STL})$.

• If $d^2 > \frac{2\sigma^2}{n(\alpha_2-1)}$, then we have the following transition depending on ρ_1 :

$$\begin{split} & - \textit{ If } \rho_1 < \frac{(\rho_2 - 2)\sigma^2}{(1 + \rho_1^{-0.5})^4(\rho_2 - 1)pd^2 - \sigma^2}, \textit{ then whp } te(\hat{\beta}_t^{\textit{MTL}}) < te(\hat{\beta}_t^{\textit{STL}}). \\ & - \textit{ If } \rho_1 > \frac{(\rho_2 - 2)\sigma^2}{(1 - \rho_1^{-0.5})^4(\rho_2 - 3)pd^2 - \sigma^2}, \textit{ then whp } te(\hat{\beta}_t^{\textit{MTL}}) > te(\hat{\beta}_t^{\textit{STL}}). \end{split}$$

- If
$$\rho_1 > \frac{(\rho_2 - 2)\sigma^2}{(1 - \rho_1^{-0.5})^4(\rho_2 - 3)pd^2 - \sigma^2}$$
, then who $te(\hat{\beta}_t^{MTL}) > te(\hat{\beta}_t^{STL})$.

The proof of Proposition 3.5 is similar to Proposition 3.3. We compare the model shift bias and the amount of reduced variance of $\hat{\beta}_t^{\text{MTL}}$. An intuitive interpretation of Proposition 3.5 is that: i) If the two task models are sufficiently similar (as specified under the first bullet), adding the source task always provides positive transfer; ii) Otherwise, as we increase the number of source task data points, the transfer is positive initially, but transitions to negative eventually. We leave the proof of Proposition 3.5 to Appendix A.3.

Implications. We use our tools to explain a key result of task my [2], which shows that by learning from multiple related tasks, one can reduce the amount of led data from each task. This

is formalized by a metric called the data efficiency ratio as follows. Given several tasks, let α^* be the largest factors such that the total number of labeled datapoints needed for solving all the tasks can be reduced by an α^* factor (compared to training independently) while keeping the performance nearly the same. More precisely, suppose we have n_i datapoints for each task, for i=1,2. If we only use αn_i datapoints from every task to train the multi-task learning estimator $\hat{\beta}(\alpha)$, then $\alpha \in (0,1)$ will be the smallest number such that

$$\alpha^{\star} := \underset{\alpha \in (0,1)}{\operatorname{arg \, min}} \ te_1(\hat{\beta}(\alpha)) + te_2(\hat{\beta}(\alpha)) \leqslant te_1(\hat{\beta}_t^{\text{STL}}) + te_2(\hat{\beta}_t^{\text{STL}}).$$

We quantify the data efficiency ratio of $\hat{\beta}_t^{\text{MTL}}$ for Example 3.2 as follows.

Proposition 3.6. In the setting of Example 3.2, assume that $\rho_1 = \rho_2 = \rho \geqslant 200$ and $d^2 < 8\sigma^2/(3p\rho)$. Then the data efficiency ratio is at most $\frac{1}{2\rho} + \frac{\sigma^2}{2\sigma^2 - 3pd^2\rho/4}$.

Note that we have stated the result assuming that $\rho_1 = \rho_2$. Similar results can also be obtained when they are different. We omit the details. The proof of Proposition 3.6 can be found in Appendix A.3.

3.4 Covariate Shift

So far we have considered settings where $\Sigma_1 = \Sigma_2$. This setting is relevant for multi-class image classification settings, where different tasks share the same input features. In general, the covariance matrices of the two tasks may be different, e.g. in text classification. In this part, we use our tools to provide a case study on the effect of applying multi-task learning for two tasks when $\Sigma_1 \neq \Sigma_2$.

For this setting, covariate shift is captured by the matrix $M = \Sigma_1^{1/2} \Sigma_2^{-1/2}$. We ask: is it better to have M as being close to identity, or should M involve varying levels of singular values? Understanding this question has implications for applying normalization methods in multi-task learning [23, 24, 9]. Our result shows that if n_1 is much larger than n_2 , then the optimal M matrix should be proportional to identity, under certain assumptions on its range of singular values (to be formulated in Proposition 3.8). On the other hand, if n_1 is comparable or even smaller than n_2 , we show an example where having "complementary" covariance matrices is better performing than having the same covariance matrices.

Example 3.7. To compare different choices of M on the performance of $\hat{\beta}_t^{\text{MTL}}$, we assume an upper bound on the scale of M. Consider the following family of matrices

$$S_{\mu} := \{ M \mid \det(M^{\top}M) \leqslant \mu^{p}, \lambda(M) \in [\mu_{\min}, \mu_{\max}] \},$$

where $\mu, \mu_{\min}, \mu_{\max}$ are fixed values that do not grow with p. For the task models, we assume that β_2 has i.i.d. entries with mean zero and variance κ^2 and $\beta_1 - \beta_2$ has i.i.d. entries with mean 0 and variance d^2 . The following proposition shows that when n_1 is large enough compared to n_2 , $te(\hat{\beta}^{\text{MTL}})$ is minimized approximately within the family of \mathcal{S}_{μ} when $M = \mu \operatorname{Id}$.

Proposition 3.8. In the setting of Example 3.2, assume that $\rho_1 > 3$ and $\rho_2 > 3$, and $\|\Sigma_1\| \leqslant C_1$ for some constant $C_1 > 0$. The test error $te(\hat{\beta}_t^{MTL})$, as a function of $M \in \mathcal{S}_{\mu}$, is minimized at the identity matrix μ Id approximately by $1 + O\left(\frac{\rho_2}{\rho_1} + \frac{1}{\sqrt{\rho_1}}\right)$.

Proposition 3.8 shows that when $n_1\gg n_2$, $te(\hat{\beta}_t^{\rm MTL})$ is minimized when Σ_1 and Σ_2 are proportional to each other. In other words, there is no covariate shift between the source task data and target task data. This provides evidence that covariate shift is unfavorable when there are many source task datapoints, To complement the result, we show an example when the statement is not true if $n_1\leqslant n_2$. Example 3.9. Within the setting of Example 3.7, we compare two cases: (i) when $M=\mathrm{Id}$; (ii) when M has p/2 singular values that are equal to $1/\lambda$. For simplicity, we assume that d=0. Hence the two tasks have the same model parameters.

In Figure 2, we plot the test error of the target task for $n_2=4p$ and n_1 ranging from p to 20p. Second, we observe the following two phases as we increase n_1/p . When $n_1\leqslant n_2$, having complementary covariance matrices leads to lower test error compared to the case when $\Sigma_1=\Sigma_2$. When $n_1>n_2$, having complementary covariance matrices leads to higher test error compared to the case when $\Sigma_1=\Sigma_2$. A theoretical justification can be found in Appendix A.4.

Implications.

4 Extending the Bias-Variance Trade-off to Transfer Learning

While our focus is on multi-task learning, we show that our results extend to transfer learning settings naturally. We provide an analysis of the transfer function of Taskonomy [2] using our setup. Specifically, the source task encoder consists of the representations learnt from one or more source tasks. The transfer function then tries to fit the target task data to the source task encoder. For more details, we refer the reader to Figure 4 in Taskonomy.

We map the procedure to our setup as follows. First, we obtain the single-task estimator $\hat{\beta}_i$ from the source tasks, for $1 \leqslant i \leqslant t-1$. This forms the shared representation $B_{t-1}^{\star} = [\hat{\beta}_1, \hat{\beta}_2, \dots, \hat{\beta}_{t-1}]$. Then, we learn the output layer W_t on the target task by minimizing the following objective

$$g(W_t) = \|X_t B W_t - Y_t\|^2. (4.1)$$

After solving W_t , we use $\hat{\beta}_t^{\text{TL-FT}} = BW_t$ as the estimator for the target task. By comparing $te(\hat{\beta}_t^{\text{TL-FT}})$ to $te(\hat{\beta}_t^{\text{STL}})$, we observe a similar trade-off between model-shift bias and variance reduction for this setting. The formal statement is presented in Theorem B.1 and its proof in Appendix B.

5 Experiments

5.1 Experimental Setup



Datasets and models. We describe the datasets and models we use in the experiments.

Sentiment Analysis: This dataset includes six tasks: movie review sentiment (MR), sentence subjectivity (SUBJ), customer reviews polarity (CR), question type (TREC), opinion polarity (MPQA), and the Stanford sentiment treebank (SST) tasks.

For each task, the goal is to categorize sentiment opinions expressed in the text. We use an embedding layer (with GloVe embeddings¹) followed by an LSTM layer proposed by [25].

ChestX-ray14: This dataset contains 112,120 frontal-view X-ray images and each image has up to 14 diseases. This is a 14-task multi-label image classification problem.

For all models, we share the main module across all tasks and assign a separate regression or classification layer on top of the shared module for each tasks.

5.2 Results and Analysis

A metric to determine positive versus negative transfer. We propose a simple metric to determine whether multi-task learning performs better than single-task learning. The metric is as follows. First, we train 2 single-task models on 2 separate tasks. We take the prediction accuracy of the two tasks. Let τ be a fixed threshold. If the accuracy of the source task is higher than the accuracy of the target task by τ , then we predict positive transfer. On the other hand, if the accuracy of the source task is lower than the accuracy of the target task by τ , then we predict negative transfer.

Table 1 shows the results on a sentiment analysis and an image classification task.

Multi-task learning provides labeled data efficiency. We measure the data efficiency ratio on text classification tasks. We find that by performing multi-task learning, only 40% of the data is needed to achieve comparable performance to single-task learning over all six tasks.

Covariance alignment provides more improvement when data size is imbalanced. We validate that the performance of aligning task covariances depends on the size of source task data. Recall that the covariance alignment procedure in [26] adds an additional module between the word embedding representation and the shared module.

In Figure 2, we observe that the benefit from aligning task covariances becomes more significant as we increase the size of the source task dataset.

¹http://nlp.stanford.edu/data/wordvecs/glove.6B.zip

Threshold	Text classification		ChestX-ray14	
	Precision	Recall	Precision	Recall
0.0	0.596	1.000	0.593	1.000
0.1	0.756	0.388	0.738	0.462
0.2	0.919	0.065	0.875	0.044
0.3	1.000	0.004	-	-

Table 1: Ablation study on when should use MTL via different source/target task accuracy. Note: For text classification tasks, the source task training data size ranges from 500 to 1,500 and target task training data size is 1000; For ChestX-ray14, the training data size is 10,000.

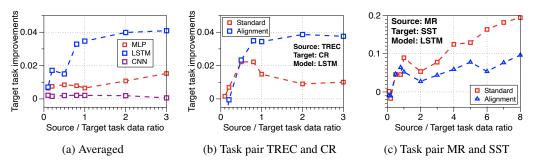


Figure 2: The performance of aligning task covariances depends on data size. As the ratio between source task data size and target task data size increases, the performance improvement from aligning task covariances increases.

6 Related Work

Adding a regularization over B, e.g. [15, 16]. Moreover, [27] observed that controlling the capacity can outperform the implicit capacity control of adding regularization over B.

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A Supplementary Materials for Section 3

From [26], we know that we need to explicitly restrict the capacity r of B so that there is transfer between the two tasks. for the rest of the section, we shall consider the case when r=1 we are considering the case of two tasks. Here, equation (1.1) simplifies to the following

$$f(B; w_1, w_2) = \|X_1 B w_1 - Y_1\|^2 + \|X_2 B w_2 - Y_2\|^2,$$
(A.1)

where $B \in \mathbb{R}^p$ and w_1, w_2 are both real numbers. To solve the above, suppose that w_1, w_2 are fixed, by local optimality, we solve B as

$$\begin{split} &\hat{B}(w_1, w_2) \\ &= &(w_1^2 X_1^\top X_1 + w_2^2 X_2^\top X_2)^{-1} (w_1 X_1^\top Y_1 + w_2 X_2^\top Y_2) \\ &= & \frac{1}{w_2} ((\frac{w_1}{w_2})^2 X_1^\top X_1 + X_2^\top X_2)^{-1} (\frac{w_1}{w_2} X_1^\top Y_1 + X_2^\top Y_2) \\ &= & \frac{1}{w_2} \left(\beta_t + ((\frac{w_1}{w_2})^2 X_1^\top X_1 + X_2^\top X_2)^{-1} \left(X_1^\top X_1 (\frac{w_1}{w_2} \beta_s - w^2 \beta_t) + (\frac{w_1}{w_2} X_1^\top \varepsilon_1 + X_2^\top \varepsilon_2) \right) \right). \end{split}$$

As a remark, when $w_1 = w_2 = 1$, we recover the linear regression estimator. The advantage of using $f(B; w_1, w_2)$ is that if θ_1 is a scaling of θ_2 , then this case can be solved optimally using equation (A.1) [27].

Defining the multi-task learning estimator. Using a validation set that is sub-sampled from the original training dataset, we get a validation loss as follows

$$val(\hat{B}; w_1, w_2) = n_1 \cdot \left\| \Sigma_1^{1/2} \left(\left(\frac{w_1}{w_2} \right)^2 X_1^\top X_1 + X_2^\top X_2 \right)^{-1} X_2^\top X_2 (\beta_s - \frac{w_1}{w_2} \beta_t) \right\|^2$$

$$+ n_1 \sigma^2 \cdot \operatorname{Tr} \left[\left(\left(\frac{w_1}{w_2} \right)^2 X_1^\top X_1 + X_2^\top X_2 \right)^{-1} \Sigma_1 \right]$$

$$+ n_2 \cdot w^2 \left\| \Sigma_2^{1/2} \left(\left(\frac{w_1}{w_2} \right)^2 X_1^\top X_1 + X_2^\top X_2 \right)^{-1} X_1^\top X_1 (\beta_s - \frac{w_1}{w_2} \beta_t) \right\|^2$$

$$+ n_2 \cdot \sigma^2 \cdot \operatorname{Tr} \left[\left(\left(\frac{w_1}{w_2} \right)^2 X_1^\top X_1 + X_2^\top X_2 \right)^{-1} \Sigma_2 \right].$$
(A.2)

Let \hat{w}_1/\hat{w}_2 be the global minimizer of $val(\hat{B}; w_1, w_2)$. We will define the multi-task learning estimator for the target task as

$$\hat{\beta}_t^{\text{MTL}} = \hat{w}_2 \hat{B}(\hat{w}_1, \hat{w}_2).$$

The intuition for deriving $\hat{\beta}_t^{\text{MTL}}$ is akin to performing multi-task training in practice. Let $\hat{v} = \hat{w_1}/\hat{w_2}$ for the simplicity of notation. The test loss of using $\hat{\beta}_t^{\text{MTL}}$ for the target task is

$$te(\hat{\beta}_t^{\text{MTL}}) = \hat{v}^2 \left\| \Sigma_2^{1/2} (\hat{v}^2 X_1^\top X_1 + X_2^\top X_2)^{-1} X_1^\top X_1 (\beta_s - \hat{v}\beta_t) \right\|^2 + \sigma^2 \cdot \text{Tr} \left[(\hat{v}^2 X_1^\top X_1 + X_2^\top X_2)^{-1} \Sigma_2 \right]. \tag{A.3}$$

Our goal is to study under model and covariate shifts, whether multi-task learning helps learn the target task better than single-task learning. The baseline where we solve the target task with its own data is

$$te(\hat{\beta}_t^{\rm STL}) = \sigma^2 \cdot {\rm Tr} \left[(X_1^\top X_1 + X_2^\top X_2)^{-1} \right], \text{ where } \hat{\beta}_t^{\rm STL} = (X_2^\top X_2)^{-1} X_2^\top Y_2.$$

We state several helper lemmas to get a bound on the variance of $\hat{\beta}_t^{\text{STL}}$ and $\hat{\beta}_t^{\text{MTL}}$. The first lemma, which is a folklore result in random matrix theory, helps determine the asymptotic limit of $te(\hat{\beta}_t^{\text{STL}})$, as p goes to infinity. When the entries of X are multivariate Gaussian, this lemma recovers the classical result for the mean of inverse Wishart distribution [?]. For general non-Gaussian random matrices, it can be obtained from Stieltjes transform method; see e.g., Lemma 3.11 of [?]. Here we shall state a result that is a consequence of Theorem 2.4 in [29], which gives an almost sharp error bound.

Lemma A.1. Let $X = \Sigma^{1/2}Z \in \mathbb{R}^{n \times p}$ be a random matrix that contains i.i.d. row vectors with mean 0 and covariance $\Sigma \in \mathbb{R}^{p \times p}$. Suppose the entries of Z have finite moments up to any order. Let A be a $p \times p$ matrix that is independent of X and satisfies ||A|| = O(1). In the setting when $n = \rho p$ for some constant $\rho > 1$, we have that with high probability,

$$\operatorname{Tr}\left[(X^{\top}X)^{-1}A\right] = \frac{1}{\rho - 1} \frac{1}{p} \operatorname{Tr}(A\Sigma^{-1}) + \operatorname{O}\left(p^{-1/2 + \varepsilon}\right) \tag{A.4}$$

for any constant $\varepsilon > 0$.

In this section, we will always consider random variables with finite moments up to any order. On the other hand, if the entries of Z have finite a-th moment for some constant a>4, then (A.4) still holds except that the error bound becomes $O(p^{-\frac{1}{2}+\frac{2}{a}+\varepsilon})$; see our Corollary C.8 below.

The second lemma, which deals the inverse of the sum of two random matrices, can be viewed as a special case of Theorem A.4.

Lemma A.2. Let $X_i \in \mathbb{R}^{n_i \times p}$ be a random matrix that contains i.i.d. row vectors with mean 0 and variance $\Sigma_i \in \mathbb{R}^{p \times p}$, for i = 1, 2. Denote by $M = \Sigma_1^{1/2} \Sigma_2^{-1/2}$ and let $\lambda_1, \lambda_2, \ldots, \lambda_p$ be the singular values of $M^\top M$ in decreasing order. When $n_1 = c_1 p$ and $n_2 = c_2 p$, we have that with high probability over the randomness of X_1 and X_2 , the following equation holds

$$\operatorname{Tr}\left[(X_1^{\top} X_1 + X_2^{\top} X_2)^{-1} \Sigma_2 \right] = \frac{1}{c_1 + c_2} \cdot \operatorname{Tr}\left[\frac{1}{p} \cdot (a_1 M^{\top} M + a_2)^{-1} \right] + \operatorname{O}\left(n^{-1/2 + \varepsilon} \right), \quad (A.5)$$

for any constant $\varepsilon > 0$, where a_1, a_2 are solutions to the following deterministic equations:

$$a_1 + a_2 = 1 - \frac{1}{c_1 + c_2}, \ a_1 + \frac{1}{p} \cdot \sum_{i=1}^{p} \frac{a_1}{(c_1 + c_2)(a_1 + a_2/\lambda_i^2)} = \frac{c_1}{c_1 + c_2}.$$
 (A.6)

We will give the proof of Lemma A.2 in Section C.

Finally, we can get a bound tighter than Theorem A.4 as follows. [Todo: restate this]

Lemma A.3. In the setting of Theorem A.4, assume that $\Sigma_1 = \operatorname{Id}$, β_t is i.i.d. with mean 0 and variance κ^2 and $\beta_s - \beta_t$ is i.i.d. with mean 0 and variance d^2 . We set $\Delta_\beta = \left((1-\hat{w})^2\kappa^2 + d^2\right)\right)\operatorname{Tr}[Z]$ and we have

$$te(\hat{\beta}_t^{MTL}) \leqslant te(\hat{\beta}_t^{STL}) \text{ when: } \Delta_{var} \geqslant \left(1 + \sqrt{\frac{p}{n_1}}\right)^4 \Delta_{\beta},$$

$$te(\hat{\beta}_t^{MTL}) \geqslant te(\hat{\beta}_t^{STL}) \text{ when: } \Delta_{var} \leqslant \left(1 - \sqrt{\frac{p}{n_1}}\right)^4 \Delta_{\beta}.$$

A.1 Proofs for Section 3.1

Let $M = \hat{w} \Sigma_1^{1/2} \Sigma_2^{-1/2}$ denote the weighted covariate shift matrix. Denote by $\lambda_1 \geqslant \lambda_2 \geqslant \ldots \geqslant \lambda_p$ the singular values of $M^T M$. Let (a_1, a_2) be the solutions to the following system of equations

$$a_1 + a_2 = 1 - \frac{p}{n_1 + n_2}, \ a_1 + \sum_{i=1}^{p} \frac{a_1}{(n_1 + n_2)(a_1 + a_2/\lambda_i^2)} = \frac{n_1}{n_1 + n_2}.$$
 (A.7)

After obtaining (a_1, a_2) , we can solve the following linear equations to get (a_3, a_4) :

$$\left(\frac{n_2}{a_2^2} - \sum_{i=1}^p \frac{1}{(a_2 + \lambda_i^2 a_1)^2}\right) a_3 - \left(\sum_{i=1}^p \frac{\lambda_i^2}{(a_2 + \lambda_i^2 a_1)^2}\right) a_4 = \sum_{i=1}^p \frac{1}{(a_2 + \lambda_i^2 a_1)^2}, \tag{A.8}$$

$$\left(\frac{n_1}{a_1^2} - \sum_{i=1}^p \frac{\lambda_i^4}{(a_2 + \lambda_i^2 a_1)^2}\right) a_4 - \left(\sum_{i=1}^p \frac{\lambda_i^2}{(a_2 + \lambda_i^2 a_1)^2}\right) a_3 = \sum_{i=1}^p \frac{\lambda_i^2}{(a_2 + \lambda_i^2 a_1)^2}.$$
(A.9)

Then we introduce the following matrix

$$Z = \frac{n_1^2}{(n_1 + n_2)^2} \cdot M \frac{(1 + a_3)\operatorname{Id} + a_4 M^{\top} M}{(a_2 + a_1 M^{\top} M)^2} M^{\top},$$

which can be regarded as the asymptotic limit of $\hat{w}\Sigma_2^{1/2}(\hat{w}^2X_1^{\top}X_1 + X_2^{\top}X_2)^{-1}\Sigma_1^{1/2}$. Finally we introduce

$$\delta := \left[\frac{n_1 \lambda_1}{(\sqrt{n_1} - \sqrt{p})^2 \lambda_p + (\sqrt{n_2} - \sqrt{p})^2} \right]^2 \cdot \|\Sigma_1^{1/2} (\beta_s - \hat{w}\beta_t)\|^2.$$

may be able to get a better bound, but the statement will be long

We now state our main result for two tasks with both covariate and model shift in the following theorem.

Theorem A.4. Let n_1, n_2 be the number of data points for the source, target task, respectively. Let \hat{w} denote the optimal solution for the ratio w_1/w_2 in equation (A.2). The information transfer is solely determined by two deterministic quantities Δ_{β} and Δ_{var} , which show the change of model shift bias and variace, respectively. With high probability we have

$$te(\hat{\beta}_t^{MTL}) \leqslant te(\hat{\beta}_t^{STL}) \text{ when: } \Delta_{var} - \Delta_{\beta} \geqslant \left(\left(1 + \sqrt{\frac{p}{n_1}}\right)^4 - 1\right)\delta$$
 (A.10)

$$te(\hat{\beta}_t^{\textit{MTL}}) \geqslant te(\hat{\beta}_t^{\textit{STL}}) \; \textit{when:} \; \Delta_{\textit{var}} - \Delta_{\beta} \leqslant -2 \left(2 \sqrt{\frac{p}{n_1}} + \frac{p}{n_1}\right) \delta, \tag{A.11}$$

where

$$\Delta_{var} := \sigma^2 \left(\frac{p}{n_2 - p} - \frac{1}{n_1 + n_2} \operatorname{Tr} \left[(a_1 M^\top M + a_2 \operatorname{Id})^{-1} \right] \right)$$
$$\Delta_{\beta} := (\beta_s - \hat{w}\beta_t)^\top \Sigma_1^{1/2} Z \Sigma_1^{1/2} (\beta_s - \hat{w}\beta_t).$$

A.1.1 Proof of Theorem A.4

[Todo: A proof outline; including the following key lemma.] To prove Theorem ??, we study the spectrum of the random matrix model:

$$Q = \Sigma_1^{1/2} Z_1^{\mathsf{T}} Z_1 \Sigma_1^{1/2} + \Sigma_2^{1/2} Z_2^{\mathsf{T}} Z_2 \Sigma_2^{1/2},$$

where $\Sigma_{1,2}$ are $p \times p$ deterministic covariance matrices, and $X_1 = (x_{ij})_{1 \leqslant i \leqslant n_1, 1 \leqslant j \leqslant p}$ and $X_2 = (x_{ij})_{n_1+1 \leqslant i \leqslant n_1+n_2, 1 \leqslant j \leqslant p}$ are $n_1 \times p$ and $n_2 \times p$ random matrices, respectively, where the entries $x_{ij}, 1 \leqslant i \leqslant n_1+n_2 \equiv n, 1 \leqslant j \leqslant p$, are real independent random variables satisfying

$$\mathbb{E}z_{ij} = 0, \qquad \mathbb{E}|z_{ij}|^2 = 1. \tag{A.12}$$

The proof of Theorem A.4 involves two parts.

Part I: Bounding the bias from model shift. We relate the first term in equation (2.1) to Δ_{β} .

Proposition A.5. In the setting of Theorem A.4, denote by $K = (\hat{w}^2 X_1^\top X_1 + X_2^\top X_1)^{-1}$, and

$$\begin{split} \delta_1 &= \hat{w}^2 \left\| \Sigma_2^{1/2} K X_1^\top X_1 (\beta_s - \hat{w} \beta_t) \right\|^2, \\ \delta_2 &= n_1^2 \cdot \hat{w}^2 \left\| \Sigma_2^{1/2} K \Sigma_1 (\beta_s - \hat{w} \beta_t) \right\|, \\ \delta_3 &= n_1^2 \cdot \hat{w}^2 \left\| \Sigma_1^{1/2} K \Sigma_2 K \Sigma_1^{1/2} \right\| \cdot \left\| \Sigma_1^{1/2} (\beta_s - \hat{w} \beta_t) \right\|^2. \end{split}$$

We have that

$$-2n_1^2 \left(2\sqrt{\frac{p}{n_1}} + \frac{p}{n_1}\right) \delta_3 \leqslant \delta_1 - \delta_2 \leqslant n_1^2 \left(2\sqrt{\frac{p}{n_1}} + \frac{p}{n_1}\right) \left(2 + 2\sqrt{\frac{p}{n_1}} + \frac{p}{n_1}\right) \delta_3.$$

For the special case when $\Sigma_1 = \operatorname{Id}$ and $\beta_s - \beta_t$ is i.i.d. with mean 0 and variance d^2 , we further have

$$\left(1 - \sqrt{\frac{p}{n_1}}\right)^4 \Delta_{\beta} \leqslant \left\| \Sigma_2^{1/2} (X_1^{\top} X_1 + X_2^{\top} X_2)^{-1} X_1^{\top} X_1 (\beta_s - \beta_t) \right\|^2.$$

Proof. The proof follows by applying equation (A.14). Recall that $X_1^{\top}X_1 = \Sigma_1^{1/2}Z_1^{\top}Z_1\Sigma_1^{1/2}$. Denote by $\mathcal{E} = Z_1^{\top}Z_1 - n_1$ Id. Let We have

$$\delta_{1} = \delta_{2} + 2\hat{w}^{2}n_{1}(\beta_{s} - \hat{w}\beta_{t})^{\top} \Sigma_{1}^{1/2} \mathcal{E} \Sigma_{1}^{1/2} K \Sigma_{2} K \Sigma_{1}(\beta_{s} - \hat{w}\beta_{t}) + \hat{w}^{2} \left\| \Sigma_{2}^{1/2} K \Sigma_{1}^{1/2} \mathcal{E} \Sigma_{1}^{1/2} (\beta_{s} - \hat{w}\beta_{t}) \right\|^{2}$$
(A.13)

Here we use the following on the second term in equation (A.13)

$$\begin{split} & \left| (\beta_{s} - \hat{w}\beta_{t})^{\top} \Sigma_{1}^{1/2} \mathcal{E} \Sigma_{1}^{1/2} K \Sigma_{2} K \Sigma_{1} (\beta_{s} - \hat{w}\beta_{t}) \right| \\ = & \left| \operatorname{Tr} \left[\mathcal{E} \Sigma_{1}^{1/2} K \Sigma_{2} K \Sigma_{1} (\beta_{s} - \hat{w}\beta_{t}) (\beta_{s} - \hat{w}\beta_{t})^{\top} \Sigma_{1}^{1/2} \right] \right| \\ \leqslant & \| \mathcal{E} \| \cdot \left\| \Sigma_{1}^{1/2} K \Sigma_{2} K \Sigma_{1} (\beta_{s} - \hat{w}\beta_{t}) (\beta_{s} - \hat{w}\beta_{t})^{\top} \Sigma_{1}^{1/2} \right\|_{\star} \\ \leqslant & n_{1} \left(2 \sqrt{\frac{p}{n_{1}}} + \frac{p}{n_{1}} \right) \cdot \left\| \Sigma_{1}^{1/2} K \Sigma_{2} K \Sigma_{1} (\beta_{s} - \hat{w}\beta_{t}) (\beta_{s} - \hat{w}\beta_{t})^{\top} \Sigma_{1}^{1/2} \right\|_{\star} \quad \text{(by equation (A.14))} \\ \leqslant & n_{1} \left(2 \sqrt{\frac{p}{n_{1}}} + \frac{p}{n_{1}} \right) \left\| \Sigma_{1}^{1/2} K \Sigma_{2} K \Sigma_{1}^{1/2} \right\| \cdot \left\| \Sigma_{1}^{1/2} (\beta_{s} - \hat{w}\beta_{t}) \right\|^{2} \end{split}$$

(since the matrix inside is rank 1)

The third term in equation (A.13) can be bounded with

$$\left\| \Sigma_2^{1/2} K \Sigma_1^{1/2} \mathcal{E} \Sigma_1^{1/2} (\beta_s - \hat{w} \beta_t) \right\|^2 \leqslant n_1^2 \left(2 \sqrt{\frac{p}{n_1}} + \frac{p}{n_1} \right)^2 \left\| \Sigma_1^{1/2} K \Sigma_2 K \Sigma_1^{1/2} \right\| \cdot \left\| \Sigma_1^{1/2} (\beta_s - \hat{w} \beta_t) \right\|^2.$$

Combined together we have shown the right direction for $\delta_1 - \delta_2$. For the left direction, we simply note that the third term in equation (A.13) is positive. And the second term is bigger than $-2n_1^2(2\sqrt{\frac{p}{n_1}} + \frac{p}{n_1})\alpha$ using equation (A.14).

Part II: The limit of $\left\|\Sigma_2^{1/2}(\hat{w}^2X_1^\top X_1 + X_2^\top X_2)^{-1}\Sigma_1(\beta_s - \hat{w}\beta_t)\right\|^2$ using random matrix theory. We consider the same setting as in previous subsection:

$$X_1^\top X_1 := \Sigma_1^{1/2} Z_1^T Z_1 \Sigma_1^{1/2}, \quad X_2^\top X_2 = \Sigma_2^{1/2} Z_2^T Z_2 \Sigma_2^{1/2},$$

where z_{ij} , $1 \le i \le n_1 + n_2 \equiv n$, $1 \le j \le p$, are real independent random variables satisfying (A.12). For now, we assume that the random variables z_{ij} are i.i.d. Gaussian, but we know that universality holds for generally distributed entries. Assume that p/n_1 is a small number such that $Z_1^T Z_1$ is roughly an isometry, that is, under (A.12), If we assume the variances of the entries of Z_1 are 1, then we have

$$-n_1 \left(2\sqrt{\frac{p}{n_1}} - \frac{p}{n_1} \right) \leqslant Z_1^T Z_1 - n_1 \operatorname{Id} \leqslant n_1 \left(2\sqrt{\frac{p}{n_1}} + \frac{p}{n_1} \right). \tag{A.14}$$

Lemma A.6. In the setting of Theorem A.4, we have with high probability 1 - o(1),

$$\widehat{w}^{2}(n_{1}+n_{2})^{2} \left\| \Sigma_{2}^{1/2} (\widehat{w}X_{1}^{\top}X_{1} + X_{2}^{\top}X_{2})^{-1} \Sigma_{1} (\beta_{s} - w\beta_{t}) \right\|^{2}$$

$$= (\beta_{s} - w\beta_{t})^{\top} \Sigma_{1}^{1/2} M \frac{(1+a_{3}) \operatorname{Id} + a_{4}M^{\top}M}{(a_{2} + a_{1}M^{\top}M)^{2}} M^{\top} \Sigma_{1}^{1/2} (\beta_{s} - \widehat{w}\beta_{t}) + \operatorname{O}(n^{-1/2+\varepsilon}),$$
(A.15)

for any constant $\varepsilon > 0$.

We will give the proof of this lemma in Section C.

Proof of Proposition A.3. the proof for tighter bound

add some arguments with ε -net.

A.1.2 Proof of Theorem A.7

Theorem A.7. Let $n = c \cdot p$. Let $X \in \mathbb{R}^{n \times p}$ and $Y_i = X\beta_i + \varepsilon_i$, for i = 1, ..., k. Let $U_rU_r^{\top}$ denote the best rank-r approximation subspace of $B^*\Sigma B^*$, where $U_r \in \mathbb{R}^{k \times r}$. Let $U_r(i)$ denote the i-th row vector of U_r . We have the following

• If
$$(1 - \|U_r(i)\|^2) \cdot \frac{\sigma^2}{c-1} \geqslant \|\Sigma(B^*U_rU_r(i) - \beta_i)\|^2$$
, then who $te(\hat{\beta}_t^{MTL}) < te(\hat{\beta}_t^{STL})$.

• If
$$(1 - \|U_r(i)\|^2) \cdot \frac{\sigma^2}{c-1} < \|\Sigma(B^*U_rU_r(i) - \beta_i)\|^2$$
, then who $te(\hat{\beta}_t^{MTL}) > te(\hat{\beta}_t^{STL})$.

As a remark, since the spectral norm of U_r is less than 1, we have that $||U_r(i)|| < 1$, for any $1 \le i \le k$. Compared to Theorem 3.1, we can get a simple expression for the two functions Δ_{vari} and Δ_{β} . The proof of Theorem A.7 can be found in Appendix A.1.2.

For this setting, the problem reduces to the following.

$$f(B; W_1, \dots, W_k) = \sum_{i=1}^k \|XBW_i - Y_i\|^2.$$
 (A.16)

In order to prove Theorem A.7, we will derive a closed form solution for equation (A.16).

Proof of Theorem A.7. By fixing W_1, W_2, \dots, W_k , we can derive a closed form solution for B as

$$\hat{B}(W_1, \dots, W_k) = (X^{\top} X)^{-1} X^{\top} \left(\sum_{i=1}^k Y_i W_i^{\top} \right) (Z Z^{\top})^{-1}$$

$$= \sum_{i=1}^k \left(\beta_i W_i^{\top} \right) (Z Z^{\top})^{-1} + (X^{\top} X)^{-1} X^{\top} \left(\sum_{i=1}^k \varepsilon_i W_i^{\top} \right) (Z Z^{\top})^{-1}$$

where we denote $Z \in \mathbb{R}^{r \times k}$ as the k vectors W_1, W_2, \dots, W_k stacked together. Similar to Section 2, we consider minimizing the validation loss over W_1, W_2, \dots, W_k provided with \hat{B} .

Denote by $\varepsilon(W) = \sum_{i=1}^k \varepsilon_i W_i^{\top}$. We shall decompose the validation loss $val(\hat{B}; W_1, \dots, W_k)$ into two parts. The first part is the model shift bias, which is equal to

$$\sum_{j=1}^{k} \left(\left\| \Sigma^{1/2} \left(\sum_{i=1}^{k} (\beta_i W_i^{\top}) (ZZ^{\top})^{-1} W_j - \beta_j \right) \right\|^2 \right)$$

The second part is the variance, which is equal to

$$\sum_{j=1}^{k} \underset{\varepsilon_{i}, \forall i}{\mathbb{E}} \left[\left(\left(\sum_{i=1}^{k} \varepsilon_{i} W_{i}^{\top} \right) (ZZ^{\top})^{-1} W_{j} \right)^{2} \right] \cdot \operatorname{Tr} \left[\Sigma (X^{\top} X)^{-1} \right]$$

$$= \sigma^{2} \cdot \operatorname{Tr} \left[\Sigma (X^{\top} X)^{-1} \right].$$

Therefore we shall focus on the minimizer for the model shift bias since the variance part does not depend the weights. Let us denote $Q = Z^\top (ZZ^\top)^{-1}Z \in \mathbb{R}^{k \times k}$ where the (i,j)-th entry is equal to $W_i^\top (ZZ^\top)^{-1}W_j$, for any $1 \leqslant i,j \leqslant k$. Let $B^\star = [\beta_1,\beta_2,\ldots,\beta_k] \in \mathbb{R}^{p \times k}$ denote the true model parameters. We can now write the validation loss succinctly as follows.

$$val(\hat{B}; W_1, \dots, W_k) = \left\| \Sigma^{1/2} \left(B^* Q - B^* \right) \right\|_F^2 + \sigma^2 \cdot \text{Tr} \left[\Sigma (X^\top X)^{-1} \right]$$

From the above we can solve for Q optimally as $U_rU_r^{\top}$. Furthermore, we can solve $\hat{\beta}_i^{\text{MTL}}$ as $B^*U_rU_r(i)$. Now we get that

$$\begin{split} te(\hat{\beta}_t^{\text{MTL}}) &= \left\| \Sigma^{1/2} \left(\sum_{i=1}^k W_i^\top (ZZ^\top)^{-1} W_j \beta_i - \beta_j \right) \right\|^2 + \sigma^2 W_j^\top (ZZ^\top)^{-1} W_j \cdot \text{Tr} \left[\Sigma (X^\top X)^{-1} \right] \\ &= \left\| \Sigma^{1/2} \left(B^\star U_r U_r(i) \right) \right\|^2 + \sigma^2 \|U_r(i)\|^2 \cdot \text{Tr} \left[\Sigma (X^\top X)^{-1} \right]. \end{split}$$

By using Lemma A.1, we conclude the proof.

A.2 Proofs for Section 3.2

The proof will consist of two main steps.

- First, we show that \hat{v} is close to 1.
- Second, we plug \hat{v} back into $te(\hat{\beta}_t^{\text{MTL}})$ to show the result.

We denote

$$val(w) = n_1 \left[d^2 + (w - 1)^2 \kappa^2 \right] \cdot \text{Tr} \left[(w^2 X_1^\top X_1 + X_2^\top X_2)^{-2} (X_2^\top X_2)^2 \right]$$

$$+ n_2 w^2 \left[d^2 + (w - 1)^2 \kappa^2 \right] \cdot \text{Tr} \left[(w^2 X_1^\top X_1 + X_2^\top X_2)^{-2} (X_1^\top X_1)^2 \right]$$

$$+ (n_1 + n_2) \sigma^2 \cdot \text{Tr} \left[(w^2 X_1^\top X_1 + X_2^\top X_2)^{-1} \right].$$

Under the setting of Proposition A.3, using concentration of random vector with i.i.d. Gaussian entries, Lemma D.6, we have that

$$val(\hat{B}; w_1, w_2) = val(w) \left(1 + O(p^{-1/2})\right)$$
 with probability $1 - o(1)$.

Thus it suffices to study the behavior of val(w). For the minimizer \hat{w} of val(w), we have a similar result as in Proposition ??.

Lemma A.8. Suppose the assumptions of Proposition A.3 hold. Assume that $\kappa^2 \sim pd^2 \sim \sigma^2$ are of the same order. Then we have that the optimal ratio for val(w) satisfies

$$|\hat{w} - 1| = \mathcal{O}(p^{-1}). \tag{A.17}$$

Proof. The proof is also similar to the one for Proposition ??. First it is easy to observe that $val(w) \le val(-w)$ for $w \ge 0$. Hence it suffices to consider the $w \ge 0$ case.

We first consider the case $w \ge 1$. We write

$$val(w) = n_1 \left[\frac{d^2}{w^4} + \frac{(w-1)^2}{w^4} \kappa^2 \right] \cdot \text{Tr} \left[(X_1^\top X_1 + w^{-2} X_2^\top X_2)^{-2} (X_2^\top X_2)^2 \right]$$
$$+ n_2 \left[\frac{d^2}{w^2} + \frac{(w-1)^2}{w^2} \kappa^2 \right] \cdot \text{Tr} \left[(X_1^\top X_1 + w^{-2} X_2^\top X_2)^{-2} (X_1^\top X_1)^2 \right]$$
$$+ \sigma^2 (n_1 + n_2) \cdot \text{Tr} \left[(w^2 X_1^\top X_1 + X_2^\top X_2)^{-1} \right].$$

Taking derivative of val(w) with respect to w, we obtain that

$$\begin{aligned} val'(w) &\geqslant n_1 \left[\frac{2(w-1)(2-w)}{w^5} \kappa^2 - \frac{4d^2}{w^5} \right] \operatorname{Tr} \left[(X_1^\top X_1 + w^{-2} X_2^\top X_2)^{-2} (X_2^\top X_2)^2 \right] \\ &+ n_2 \left[\frac{2(w-1)}{w^3} \kappa^2 - \frac{2d^2}{w^3} \right] \cdot \operatorname{Tr} \left[(X_1^\top X_1 + w^{-2} X_2^\top X_2)^{-2} (X_1^\top X_1)^2 \right] \\ &- 2(n_1 + n_2) \frac{\sigma^2}{w^3} \cdot \operatorname{Tr} \left[(X_1^\top X_1 + w^{-2} X_2^\top X_2)^{-2} X_1^\top X_1 \right] = n_1 \operatorname{Tr} \left[(X_1^\top X_1 + w^{-2} X_2^\top X_2)^{-2} \mathcal{A} \right], \end{aligned}$$

where the matrix A is

$$\mathcal{A} := \left\lceil \frac{2(w-1)(2-w)}{w^5} \kappa^2 - \frac{4d^2}{w^5} \right\rceil (X_2^\top X_2)^2 + \frac{n_2}{n_1} \left\lceil \frac{2(w-1)}{w^3} \kappa^2 - \frac{2d^2}{w^3} \right\rceil (X_1^\top X_1)^2. \\ - 2\frac{n_1 + n_2}{n_1} \frac{\sigma^2}{w^3} X_1^\top X_1 + \frac{n_2}{m_1} \frac{\sigma^2}{w^3} X_1^\top X_2 + \frac{n_2}{m_1} \frac{\sigma^2}{w^3} X_1^\top X_1 + \frac{n_2}{m_1} \frac{\sigma^2}{w^3} X_1^\top X_2 + \frac{n_2}{m_1} \frac{\sigma^2}{w^3} X_1 + \frac{n_2}{m_1} \frac{\sigma^2}{w^3} X$$

Using the estimate (A.14), we get that A is lower bounded as

$$\mathcal{A} \succeq -\frac{4d^2}{w^5}(\sqrt{n_2} + \sqrt{p})^4 + \frac{n_2}{n_1} \left[\frac{2(w-1)}{w^3} \kappa^2 - \frac{2d^2}{w^3} \right] (\sqrt{n_1} - \sqrt{p})^4 - 2\frac{n_1 + n_2}{n_1} \frac{\sigma^2}{w^3} (\sqrt{n_1} + \sqrt{p})^2 \succ 0,$$

as long as

$$w > w_1 := 1 + \frac{d^2}{\kappa^2} + \frac{\sigma^2}{\kappa^2} \frac{(n_1 + n_2)(\sqrt{n_1} + \sqrt{p})^2}{n_2(\sqrt{n_1} - \sqrt{p})^4} + \frac{2d^2}{\kappa^2} \frac{n_1(\sqrt{n_2} + \sqrt{p})^4}{n_2(\sqrt{n_1} - \sqrt{p})^4}.$$

Hence val'(w) > 0 on (w_1, ∞) , i.e. val(w) is strictly increasing for $w > w_1$. Hence we must have $\hat{w} \leq w_1$. Note that under our assumptions, we have $w_1 = 1 + \mathrm{O}(p^{-1})$.

Then we consider the case $w \leq 1$. Taking derivative of val(w) with respect to w, we obtain that

$$val'(w) \leqslant n_1 \left[2 (w - 1) \kappa^2 \right] \cdot \text{Tr} \left[(w^2 X_1^\top X_1 + X_2^\top X_2)^{-2} (X_2^\top X_2)^2 \right]$$

+ $n_2 \left[2d^2 w + 2w (w - 1) (2w - 1) \kappa^2 \right] \cdot \text{Tr} \left[(w^2 X_1^\top X_1 + X_2^\top X_2)^{-2} (X_1^\top X_1)^2 \right]$
= $n_1 \text{Tr} \left[(w^2 X_1^\top X_1 + X_2^\top X_2)^{-1} \mathcal{B} \right],$

where the matrix \mathcal{B} is

$$\mathcal{B} = 2(w-1)\kappa^{2}(X_{2}^{\top}X_{2})^{2} + \frac{n_{2}}{n_{1}} \left[2d^{2}w + 2w(w-1)(2w-1)\kappa^{2}\right](X_{1}^{\top}X_{1})^{2}.$$

Using the estimate (A.14), we get that \mathcal{B} is upper bounded as

$$\mathcal{B} \leq -2(1-w)\kappa^2(\sqrt{n_2}-\sqrt{p})^4 + 2d^2w\frac{n_2}{n_1}(\sqrt{n_1}+\sqrt{p})^4 < 0,$$

as long as

$$w < w_2 := 1 - \frac{d^2}{\kappa^2} \frac{n_2(\sqrt{n_1} + \sqrt{p})^4}{n_1(\sqrt{n_2} - \sqrt{p})^4}.$$

Hence val'(w) < 0 on $[0, w_2)$, i.e. val(w) is strictly increasing for $w < w_2$. Hence we must have $\hat{w} \le w_2$. Note that under our assumptions, we have $w_2 = 1 - \mathrm{O}(p^{-1})$.

Based on Lemma A.8, we can prove Proposition 3.3.

Proof of Proposition 3.3. Since $\Sigma_1 = \Sigma_2$, we know that $M = \Sigma_1^{1/2} \Sigma_2^{-1/2} = \hat{v} \operatorname{Id}$. Using Lemma A.1 and A.2, we can track the reduction of variance from $\hat{\beta}_t^{\text{MTL}}$ to $\hat{\beta}_t^{\text{STL}}$ as whp

$$\Delta_{\text{var}} := \sigma^2 \left(\frac{1}{c_2 - 1} - \frac{1}{c_1 + c_2} \cdot \frac{1}{a_1 \hat{v}^2 + a_2} + \mathcal{O}\left(p^{-1/2 + \varepsilon}\right) \right)$$

$$= \sigma^2 \left(\frac{c_1}{(c_2 - 1)(c_1 + c_2 - 1)} + \mathcal{O}\left(p^{-1/2 + \varepsilon}\right) \right), \tag{A.18}$$

where we use equation (A.6) and Lemma A.8. Next we consider the model shift bias

$$\Delta_{\beta} := \hat{v}^2 \left\| (\hat{v}^2 X_1^\top X_1 + X_2^\top X_2)^{-1} X_1^\top X_1 (\beta_s - \hat{v}\beta_t) \right\|^2$$

$$= d^2 \cdot \left\| (\hat{v}^2 X_1^\top X_1 + X_2^\top X_2)^{-1} X_1^\top X_1 \right\|_F^2 + \mathcal{O}\left(p^{-1/2 + \varepsilon} d^2 \right)$$

Using Lemma A.3, we get an upper and lower bound on Δ_{β} as

$$\left(1 - \sqrt{\frac{1}{c_1}}\right)^4 \leqslant \Delta_{\beta} / \left(p \cdot d^2 \cdot \frac{c_1^2}{(c_1 + c_2)^2} \cdot \frac{1 + a_3 + a_4}{(a_1 + a_2)^2} + \mathcal{O}\left(p^{1/2 + \varepsilon} d^2\right)\right) \leqslant \left(1 + \sqrt{\frac{1}{c_1}}\right)^4. \tag{A.19}$$

By solving equations (A.6), (A.8) and (A.9), we get

$$a_1 = \frac{c_1(c_1 + c_2 - 1)}{(c_1 + c_2)^2} + O(p^{-1/2 + \varepsilon}), a_2 = \frac{c_2(c_1 + c_2 - 1)}{(c_1 + c_2)^2} + O(p^{-1/2 + \varepsilon}),$$

$$a_3 = \frac{c_2}{(c_1 + c_2)(c_1 + c_2 - 1)} + O(p^{-1/2 + \varepsilon}), a_4 = \frac{c_1}{(c_1 + c_2)(c_1 + c_2 - 1)} + O(p^{-1/2 + \varepsilon}).$$

Hence we obtain that

$$\frac{1+a_3+a_4}{(a_1+a_2)^2} = \frac{(c_1+c_2)^3}{(c_1+c_2-1)^3} + O\left(p^{-1/2+\varepsilon}\right).$$

To sum up, we have shown that when

$$\frac{p \cdot d^2}{\sigma^2} < \frac{(c_1 + c_2 - 1)^2}{c_1(c_1 + c_2)(c_2 - 1)} \cdot \left(1 + \sqrt{\frac{1}{c_1}}\right)^{-4},$$

we have that $\Delta_{\text{var}} > \Delta_{\beta}$, which implies that $te(\hat{\beta}_t^{\text{MTL}})$ is lower than $te(\hat{\beta}_t^{\text{STL}})$. When

$$\frac{p \cdot d^2}{\sigma^2} > \frac{(c_1 + c_2 - 1)^2}{c_1(c_1 + c_2)(c_2 - 1)} \cdot \left(1 - \sqrt{\frac{1}{c_1}}\right)^{-4},$$

we have that $\Delta_{\text{var}} < \Delta_{\beta}$, which implies that $te(\hat{\beta}_t^{\text{MTL}})$ is higher than $te(\hat{\beta}_t^{\text{STL}})$. The result follows by taking $c_1 = c_2 = c$ for the above equation.

Proof of Proposition 3.4. Under the setting of Proposition 3.4, we have

$$\begin{split} te(\hat{\beta}_t^{\mathsf{MTL}}) &= \; \hat{v}^2 \left\| (\hat{v}^2 X_1^\top X_1 + X_2^\top X_2)^{-1} X_1^\top X_1 (\beta_s - \hat{v}\beta_t) \right\|^2 \\ &+ \sigma_2^2 \cdot \mathrm{Tr} \left[(\hat{v}^2 X_1^\top X_1 + X_2^\top X_2)^{-1} \right] + (\sigma_1^2 - \sigma_2^2) \hat{v}^2 \cdot \mathrm{Tr} \left[(\hat{v}^2 X_1^\top X_1 + X_2^\top X_2)^{-2} X_1^\top X_1 \right]. \end{split}$$

Using concentration of random vector with i.i.d. Gaussian entries, Lemma D.6, we obtain that

$$\begin{split} te(\hat{\beta}_t^{\text{MTL}}) &= \; \hat{v}^2 \left[d^2 + (w-1)^2 \, \kappa^2 \right] \text{Tr} \left[(\hat{v}^2 X_1^\top X_1 + X_2^\top X_2)^{-2} (X_1^\top X_1)^2 \right] \cdot \left(1 + \mathcal{O}(p^{-1/2}) \right) \\ &+ (\sigma_1^2 - \sigma_2^2) \hat{v}^2 \cdot \text{Tr} \left[(\hat{v}^2 X_1^\top X_1 + X_2^\top X_2)^{-2} X_1^\top X_1 \right] + \sigma_2^2 \cdot \text{Tr} \left[(\hat{v}^2 X_1^\top X_1 + X_2^\top X_2)^{-1} \right]. \end{split}$$

In the setting of Proposition 3.4, we can also obtain that

$$|\hat{w} - 1| = O(p^{-1}).$$

as in (A.17). We omit the details of the proof, since it is almost the same as the one in the proof of Lemma A.8. Thus $te(\hat{\beta}_t^{\text{MTL}})$ can be reduced to

$$te(\hat{\beta}_t^{\text{MTL}}) = \left[\Delta_\beta + \Delta_{diff} + \sigma_2^2 \cdot \text{Tr}(X_1^\top X_1 + X_2^\top X_2)^{-1}\right] \cdot \left(1 + \mathcal{O}(p^{-1/2})\right),$$

where $\Delta_{\beta} := \sigma_2^2 \cdot \operatorname{Tr} \left[(\hat{v}^2 X_1^{\top} X_1 + X_2^{\top} X_2)^{-1} \right]$ has been estimated in (A.19) and Δ_{diff} is

$$\Delta_{diff} := (\sigma_1^2 - \sigma_2^2) \cdot \text{Tr} \left[(X_1^\top X_1 + X_2^\top X_2)^{-2} X_1^\top X_1 \right].$$

Using (A.14) and Lemma A.6, a similar bound holds for Δ_{diff} as in (A.19):

$$\left(1 - \sqrt{\frac{1}{c_1}}\right)^2 \leqslant \frac{\Delta_{diff}}{\sigma_1^2 - \sigma_2^2} \cdot \left(\frac{c_1(c_1 + c_2)}{(c_1 + c_2 - 1)^3} + \mathcal{O}(p^{-1/2 + \varepsilon})\right)^{-1} \leqslant \left(1 + \sqrt{\frac{1}{c_1}}\right)^2.$$

Next we compare Δ_{Var} with $\Delta_{\beta} + \Delta_{diff}$. Under the condition $d^2 < \frac{2\sigma^2}{3p} \frac{(c_1 + c_2 - 1)^2}{c_1(c_1 + c_2)(c_2 - 1)}$, if $\sigma_1^2 \leqslant \sigma_2^2$ then we always have that $\Delta_{\mathrm{var}} > \Delta_{\beta}$ and hence $te(\hat{\beta}_t^{\mathrm{MTL}})$ is smaller than $te(\hat{\beta}_t^{\mathrm{STL}})$. It remains to consider the case $\sigma_1^2 \geqslant \sigma_2^2$. If

$$\sigma_{2}^{2} \frac{c_{1}}{(c_{2}-1)(c_{1}+c_{2}-1)} \left(1+O\left(p^{-1/2+\varepsilon}\right)\right)$$

$$> \left[\frac{c_{1}^{2}(c_{1}+c_{2})}{(c_{1}+c_{2}-1)^{3}}+O\left(p^{-1/2+\varepsilon}\right)\right] \left[pd^{2}\left(1+\sqrt{\frac{1}{c_{1}}}\right)^{4}+\frac{\sigma_{1}^{2}-\sigma_{2}^{2}}{c_{1}}\left(1+\sqrt{\frac{1}{c_{1}}}\right)^{2}\right],$$
(A.20)

then we have $\Delta_{\text{var}} > \Delta_{\beta}$ and we have $te(\hat{\beta}_t^{\text{MTL}})$ is smaller than $te(\hat{\beta}_t^{\text{STL}})$. For $c_1 \geqslant 100$, we see that (A.20) is implied by the following condition

$$\sigma_1^2 \le pd^2 \cdot c_1 + \left[1 + \frac{2}{3} \frac{(c_1 + c_2 - 1)^2}{(c_1 + c_2)(c_2 - 1)}\right] \sigma_2^2$$

as long as p is large enough. On the other hand, if

$$\sigma_{2}^{2} \frac{c_{1}}{(c_{2}-1)(c_{1}+c_{2}-1)} \left(1+O\left(p^{-1/2+\varepsilon}\right)\right)$$

$$< \left[\frac{c_{1}^{2}(c_{1}+c_{2})}{(c_{1}+c_{2}-1)^{3}}+O\left(p^{-1/2+\varepsilon}\right)\right] \left[pd^{2}\left(1-\sqrt{\frac{1}{c_{1}}}\right)^{4}+\frac{\sigma_{1}^{2}-\sigma_{2}^{2}}{c_{1}}\left(1-\sqrt{\frac{1}{c_{1}}}\right)^{2}\right],$$
(A.21)

then we have $\Delta_{\text{var}} < \Delta_{\beta}$ and we have $te(\hat{\beta}_t^{\text{MTL}})$ is larger than $te(\hat{\beta}_t^{\text{STL}})$. For $c_1 \geqslant 100$, we see that (A.21) is implied by the following condition

$$\sigma_1^2 \geqslant pd^2 \cdot c_1 + \left[1 + \frac{3}{2} \frac{(c_1 + c_2 - 1)^2}{(c_1 + c_2)(c_2 - 1)}\right] \sigma_2^2.$$

A.3 Proofs for Section 3.3

Proof of Proposition 3.5. Recall that under the setting of Example 3.2, (A.18) and (A.19) hold. Then we get that $te(\hat{\beta}_t^{\text{MTL}}) < te(\hat{\beta}_t^{\text{STL}})$ whp if

$$\left(pd^{2}\frac{c_{1}(c_{1}+c_{2})(c_{2}-1)}{(c_{1}+c_{2}-1)^{2}}+O\left(p^{1/2+\varepsilon}d^{2}\right)\right)\left(1+\sqrt{\frac{1}{c_{1}}}\right)^{4}<\sigma^{2}\left(1+O\left(p^{-1/2+\varepsilon}\right)\right); (A.22)$$

otherwise if

$$\left(pd^{2}\frac{c_{1}(c_{1}+c_{2})(c_{2}-1)}{(c_{1}+c_{2}-1)^{2}}+O\left(p^{1/2+\varepsilon}d^{2}\right)\right)\left(1-\sqrt{\frac{1}{c_{1}}}\right)^{4}>\sigma^{2}\left(1+O\left(p^{-1/2+\varepsilon}\right)\right), (A.23)$$

then we have $te(\hat{\beta}_t^{\text{MTL}}) > te(\hat{\beta}_t^{\text{STL}})$ whp.

Now we prove the first statement of Proposition 3.5. Notice that the function

$$\frac{c_1(c_1+c_2)(c_2-1)}{(c_1+c_2-1)^2} = (c_2-1)\left(1+\frac{c_2-2}{c_1}+\frac{1}{c_1(c_1+c_2)}\right)^{-1}$$

is strictly increasing with respect to c_1 as long as $c_2 \geqslant 3$. In particular, by taking $c_1 \to \infty$, we get the bound:

$$\frac{c_1(c_1+c_2)(c_2-1)}{(c_1+c_2-1)^2}\cdot \left(1+\sqrt{\frac{1}{c_1}}\right)^4 < (c_2-1)\left(1+\sqrt{\frac{1}{c_1}}\right)^4.$$

Hence we conclude that $te(\hat{\beta}_t^{\text{MTL}}) < te(\hat{\beta}_t^{\text{STL}})$ whp as long as

$$pd^2 + o(1) \le \left(1 + \sqrt{\frac{1}{c_1}}\right)^{-4} \frac{\sigma^2}{c_2 - 1}.$$

This gives the first statement by taking $c_1 > a$.

The second statement can be proved in a similar way. Suppose $c_1 > a$ and $pd^2 > (1-a^{-1/2})^{-4} \frac{\sigma^2}{c_2-1}$. If $c_1 > \frac{(c_2-2)\sigma^2}{(1-a^{-1/2})^4(1-(a+c_2-2)^{-2})(c_2-1)pd^2-\sigma^2}$, we have

$$pd^{2} \cdot \frac{c_{1}(c_{1}+c_{2})(c_{2}-1)}{(c_{1}+c_{2}-1)^{2}} \cdot \left(1-\sqrt{\frac{1}{c_{1}}}\right)^{4}$$

$$> \left(1-\sqrt{\frac{1}{c_{1}}}\right)^{4} \left(1-\frac{1}{(c_{1}+c_{2}-2)^{2}}\right) \cdot \frac{pd^{2}(c_{2}-1)c_{1}}{c_{1}+(c_{2}-2)} > \sigma^{2} \cdot (1+o(1)),$$

where in the first step we used that

$$\frac{(c_1+c_2)(c_1+c_2-2)}{(c_1+c_2-1)^2} > 1 - \frac{1}{(c_1+c_2-2)^2}.$$

This shows that (A.23) holds as long as p is large enough, and hence $te(\hat{\beta}_t^{\text{MTL}}) > te(\hat{\beta}_t^{\text{STL}})$ holds. On the other hand, if $c_1 < \frac{(c_2-2)\sigma^2}{(1+a^{-1/2})^4(c_2-1)pd^2-\sigma^2}$, then we have

$$pd^{2} \cdot \frac{c_{1}(c_{1}+c_{2})(c_{2}-1)}{(c_{1}+c_{2}-1)^{2}} \cdot \left(1+\sqrt{\frac{1}{c_{1}}}\right)^{4} < \left(1+\sqrt{\frac{1}{c_{1}}}\right)^{4} \cdot \frac{pd^{2}(c_{2}-1)c_{1}}{c_{1}+(c_{2}-2)} < \sigma^{2} \cdot (1-o(1)).$$

This shows that (A.22) holds as long as p is large enough, and hence $te(\hat{\beta}_t^{\text{MTL}}) < te(\hat{\beta}_t^{\text{STL}})$ holds. \square

Remark. Furthermore, as a function of c_1 over the range $[??,\infty]$, the maximum of $te(\hat{\beta}_t^{\text{STL}}) - te(\hat{\beta}_t^{\text{MTL}})$ is attained when $c_1 = c_2\sigma^2/\max(2(c_2-1)pd^2-\sigma^2,0)$. (cannot get this because we only have some bounds. If we let $c_1 \to \infty$, then the curve of $te(\hat{\beta}_t^{\text{STL}}) - te(\hat{\beta}_t^{\text{MTL}})$ already becomes flat, and it is meaningless to discuss the minimum of this function at this point?)

Proof of Proposition 3.6. Suppose we have reduced number of datapoints— αn_1 for task 1 and αn_2 for task 2 with $n_1 = n_2$. In the setting of Proposition 3.6, we still have (A.17), and by symmetry it suffices to focus on one of tasks, say task 2. Using Lemmas A.2 and A.3, we have that whp,

$$te_2(\hat{\beta}(\alpha)) = \sigma^2 \left(\frac{1}{2\alpha c - 1} + O\left(p^{-1/2 + \varepsilon}\right) \right) + \Delta_{\beta}^{(2)}, \tag{A.24}$$

where $\Delta_{\beta}^{(2)}$ satisfies

$$\left(1 - \sqrt{\frac{1}{\alpha c}}\right)^4 \leqslant \Delta_{\beta}^{(2)} / \left(pd^2 \cdot \frac{2\alpha^3 c^3}{(2\alpha c - 1)^3} + \mathcal{O}\left(p^{1/2 + \varepsilon} d^2\right)\right) \leqslant \left(1 + \sqrt{\frac{1}{\alpha c}}\right)^4.$$

On the other hand, using Lemma A.1, we have whp

$$te_2(\hat{\beta}_t^{\text{STL}}) = \frac{\sigma^2}{c-1} \left(1 + \mathcal{O}\left(p^{-1/2+\varepsilon}\right) \right).$$
 (A.25)

Comparing (A.24) and (A.25), we observe that $te_2(\hat{\beta}(\alpha)) \ge te_2(\hat{\beta}_t^{\text{STL}})$ for $\alpha \le 1/2 - o(1)$, which gives $\alpha^* \ge 1/2 - o(1)$.

Now for $c \ge 200$ and $\alpha \ge 1/2$, we have whp

$$te_2(\hat{\beta}(\alpha)) < \sigma^2 \frac{1}{2\alpha c - 1} + \frac{3}{8}pd^2$$
 (A.26)

for large enough p. Moreover, if

$$\alpha \geqslant \alpha_0 := \frac{1}{2c} + \left(2 - \frac{3}{4}c \cdot \frac{pd^2}{\sigma^2}\right)^{-1},$$

we can check from (A.26) that $te_2(\hat{\beta}(\alpha)) < te_2(\hat{\beta}_t^{STL})$. By symmetry, we also have $te_1(\hat{\beta}(\alpha)) < te_1(\hat{\beta}_t^{STL})$. Thus we conclude that

$$te_1(\hat{\beta}(\alpha)) + te_2(\hat{\beta}(\alpha)) \leq te_1(\hat{\beta}_t^{STL}) + te_2(\hat{\beta}_t^{STL})$$

for $\alpha \geqslant \alpha_0$, which shows that $\alpha^* \leqslant \alpha_0$.

A.4 Proofs for Section 3.4

Proof of Proposition 3.8. Let

$$M_0 := \underset{M \in \mathcal{S}_u}{\operatorname{arg\,min}} (te(\hat{\beta}_t^{\text{MTL}}))(M).$$

We now calculate $(te(\hat{\beta}_t^{\text{MTL}}))(M_0)$. In this case, as in Lemma A.8 we also have that

$$|\hat{v} - 1| = \mathcal{O}(p^{-1}) \tag{A.27}$$

in the setting Proposition 3.8. In fact, Lemma A.8 was proved assuming that $M=\mathrm{Id}$, but its proof can be easily extended to the case with general $M\in\mathcal{S}_{\mu}$ by using that $\lambda(M)\in[\mu_{\min},\mu_{\max}]$. We omit the details here.

Now using (A.27), Lemma A.2 and Lemma A.3, we get that whp,

$$(te(\hat{\beta}_t^{\text{MTL}}))(M_0) = \frac{\sigma^2}{c_1 + c_2} \cdot \frac{1}{p} \operatorname{Tr} \left(\frac{1}{a_1 M_0^{\top} M_0 + a_2} \right) + \Delta_{\beta}(M_0),$$

where $\Delta_{\beta}(M_0)$ satisfies

$$\left| \Delta_{\beta}(M_0) - \frac{d^2 \cdot c_1^2}{(c_1 + c_2)^2} \operatorname{Tr} \left[M_0 \frac{(1 + a_3) \operatorname{Id} + a_4 M_0^{\top} M_0}{(a_2 + a_1 M_0^{\top} M_0)^2} M_0^{\top} \right] \right|$$

$$\leq \left(\left(1 + \sqrt{\frac{1}{c_1}} \right)^4 - 1 \right) \left(\frac{c_1 \mu_{\max}}{(\sqrt{c_1} - 1)^2 \mu_{\min} + (\sqrt{c_2} - 1)^2} \right)^2 \cdot d^2 \operatorname{Tr} (\Sigma_1).$$

From equation (A.7), we get

$$a_1 \geqslant \frac{c_1 - 1}{c_1 + c_2}, \quad a_2 \leqslant \frac{c_2}{c_1 + c_2}.$$

Then solving equations (A.8) and (A.9), we get that

$$0 \leqslant a_3 = \frac{\frac{1}{n_2} \sum_{i=1}^p \frac{a_2^2}{(a_2 + \lambda_i^2 a_1)^2} \left(1 - \frac{1}{n_1} \sum_{i=1}^p \frac{\lambda_i^4 a_1^2}{(a_2 + \lambda_i^2 a_1)^2}\right) + \left(\frac{1}{n_1} \sum_{i=1}^p \frac{\lambda_i^2 a_1^2}{(a_2 + \lambda_i^2 a_1)^2}\right) \left(\frac{1}{n_2} \sum_{i=1}^p \frac{\lambda_i^2 a_2^2}{(a_2 + \lambda_i^2 a_1)^2}\right)}{\left(1 - \frac{1}{n_2} \sum_{i=1}^p \frac{a_2^2}{(a_2 + \lambda_i^2 a_1)^2}\right) \left(1 - \frac{1}{n_1} \sum_{i=1}^p \frac{\lambda_i^4 a_1^2}{(a_2 + \lambda_i^2 a_1)^2}\right) - \left(\frac{1}{n_2} \sum_{i=1}^p \frac{\lambda_i^2 a_2^2}{(a_2 + \lambda_i^2 a_1)^2}\right) \left(\frac{1}{n_1} \sum_{i=1}^p \frac{\lambda_i^2 a_1^2}{(a_2 + \lambda_i^2 a_1)^2}\right)} \leqslant \frac{c_2^{-1}}{1 - c_1^{-1} - c_2^{-1}}$$

and

$$0 \leqslant a_4 = \frac{\frac{1}{n_1} \sum_{i=1}^p \frac{\lambda_i^2 a_1^2}{(a_2 + \lambda_i^2 a_1)^2}}{\left(1 - \frac{1}{n_2} \sum_{i=1}^p \frac{a_2^2}{(a_2 + \lambda_i^2 a_1)^2}\right) \left(1 - \frac{1}{n_1} \sum_{i=1}^p \frac{\lambda_i^4 a_1^2}{(a_2 + \lambda_i^2 a_1)^2}\right) - \left(\frac{1}{n_2} \sum_{i=1}^p \frac{\lambda_i^2 a_2^2}{(a_2 + \lambda_i^2 a_1)^2}\right) \left(\frac{1}{n_1} \sum_{i=1}^p \frac{\lambda_i^2 a_1^2}{(a_2 + \lambda_i^2 a_1)^2}\right)}{\left(1 - \frac{1}{n_2} \sum_{i=1}^p \frac{\lambda_i^2 a_2^2}{(a_2 + \lambda_i^2 a_1)^2}\right) \left(\frac{1}{n_1} \sum_{i=1}^p \frac{\lambda_i^2 a_1^2}{(a_2 + \lambda_i^2 a_1)^2}\right)} \leqslant \frac{c_1^{-1} \cdot \mu_{\min}^2}{1 - c_1^{-1} - c_2^{-1}},$$

where we also used that

$$\left(\sum_{i=1}^{p} \frac{\lambda_i^2}{(a_2 + \lambda_i^2 a_1)^2}\right)^2 \leqslant \sum_{i=1}^{p} \frac{\lambda_i^4}{(a_2 + \lambda_i^2 a_1)^2} \cdot \sum_{i=1}^{p} \frac{1}{(a_2 + \lambda_i^2 a_1)^2}$$

by Cauchy-Schwarz inequality.

Combining the above estimates, we get that

$$(te(\hat{\beta}_t^{\text{MTL}}))(M_0) = te(M_0) + \mathcal{E}, \tag{A.28}$$

where

$$te(M_0) := \frac{\sigma^2}{c_1 + c_2} \cdot \frac{1}{p} \operatorname{Tr} \left(\frac{1}{a_1 M_0^\top M_0} \right) + \frac{d^2 \cdot c_1^2}{(c_1 + c_2)^2} \operatorname{Tr} \left(\frac{1 + a_3}{a_1^2 M_0^\top M_0} \right),$$

and the error satisfies

$$|\mathcal{E}| \leqslant C\left(\frac{c_2}{c_1} + c_1^{-1/2}\right) te(M_0),$$

Here the constant C > 0 depends only on μ_{max} , μ_{min} and $\|\Sigma_1\|$, but otherwise does not depend on c_1 and c_2 .

Finally using AM-GM inequality, we observe that

$$\operatorname{Tr}\left(\frac{1}{M^{\top}M}\right) = \sum_{i=1}^{p} \frac{1}{\lambda_i}$$

is minimized when $\lambda_1 = \cdots = \lambda_p = \mu$ under the restriction $\prod_{i=1}^p \lambda_i \leq \mu^p$. Hence we get that

$$te(M_0) \leqslant te(\mu \operatorname{Id}).$$

Together with (A.28), we conclude (??).

Theoretical justification of Example ??. In the setting of Example ??, equations in (A.7) become

$$a_1 + a_2 = 1 - \frac{p}{n_1 + n_2}, \ a_1 + \frac{p}{2(n_1 + n_2)} \cdot \left(\frac{a_1}{a_1 + \lambda^2 a_2} + \frac{a_1}{a_1 + \frac{a_2}{2}}\right) = \frac{n_1}{n_1 + n_2}.$$
 (A.29)

It's not hard to verify that there is only one valid solution (a_1, a_2) to (A.29). After solving these, we get the test error for the target task as follows.

$$te(\lambda) = \frac{p}{2(n_1 + n_2)} \cdot \left(\frac{1}{\frac{a_1}{\lambda^2} + a_2} + \frac{1}{a_1\lambda^2 + a_2}\right).$$
 (A.30)

First we notice that the curves in Figure ?? all cross at the point $n_1 = n_2$. In fact, if $n_1 = n_2$, then it is easy to observe that $a_1 = a_2 = (1 - \gamma)/2$ is the solution to equation (A.29), where we denote $\gamma = p/(n_1 + n_2)$. Then for any λ , the test error in (A.30) takes the value

$$te(\lambda) = \frac{\gamma}{2} \frac{1}{(1-\gamma)/2} = \frac{p}{n_1 + n_2 - p}.$$

This phenomenon can be also explained using our theory. With (A.29), we can write

$$te(\lambda) = \frac{\gamma}{2} \cdot \left(\frac{1}{\frac{a_1}{\lambda^2} + (1 - \gamma - a_1)} + \frac{1}{a_1 \lambda^2 + (1 - \gamma - a_1)} \right).$$

We can compute that

$$te(\lambda) - te(1) = \frac{\gamma}{2(1-\gamma)} (\lambda^2 - 1) a_1 \cdot \left(\frac{1}{-a_1(\lambda^2 - 1) + (1-\gamma)\lambda^2} - \frac{1}{a_1(\lambda^2 - 1) + (1-\gamma)} \right)$$
$$= \frac{\gamma}{2(1-\gamma)} (\lambda^2 - 1)^2 a_1 \cdot \frac{2a_1 - (1-\gamma)}{[-a_1(\lambda^2 - 1) + (1-\gamma)\lambda^2][a_1(\lambda^2 - 1) + (1-\gamma)]}.$$

If $n_1 > n_2$, we have $a_1 > (1-\gamma)/2$ (because $a_1 > a_2$ as observed from the equation (A.29)), and hence $te(\lambda) > te(1)$. Otherwise if $n_1 < n_2$, we have $a_1 < (1-\gamma)/2$, and hence $te(\lambda) < te(1)$. \square

B Supplementary Materials for Section 4

Proposition B.1. Let $X_i \in \mathbb{R}^{n_i \times p}$ and $Y_i = X_i \beta_i + \varepsilon_i$, for $i = 1, \dots, t$, where ε_i are random vectors with i.i.d. entries of mean zero and variance σ^2 . Let $\rho_i > 1$, $1 \le i \le t$ be fixed constants. Assume that $\max_{1 \le i \le t} \|\beta_i\| = \mathrm{O}(1)$, $\sigma^2 = \mathrm{O}(1)$ and

$$\|[(B_{t-1}^{\star})^{\top}B_{t-1}^{\star}]^{-1}\| = O(1), \quad B_{t-1}^{\star} := [\beta_1, \beta_2, \dots, \beta_{t-1}].$$
 (B.1)

Suppose that all random variables have subexponential decay. Then we have with high probability,

$$te(BW_t) = \|\varepsilon_{\beta} + \sigma^2 \varepsilon_{var}^{(1)}\|^2 + \sigma^2 \|\varepsilon_{var}^{(2)}\|^2 + \mathcal{O}(p^{-1/2+\varepsilon}), \tag{B.2}$$

where $B_{t-1}^{\star} := [\beta_1, \beta_2, \dots, \beta_{t-1}]$, and

$$\begin{split} \varepsilon_{\beta} &:= \Sigma_{k}^{1/2} \left\{ 1 - B_{t-1}^{\star} \left[(B_{t-1}^{\star})^{\top} B_{t-1}^{\star} \right]^{-1} (B_{t-1}^{\star})^{\top} \right\} \beta_{t}, \\ \varepsilon_{var}^{(1)} &:= \Sigma_{k}^{1/2} B_{t-1}^{\star} \left[(B_{t-1}^{\star})^{\top} B_{t-1}^{\star} \right]^{-1} \mathcal{M}_{1} \left[(B_{t-1}^{\star})^{\top} B_{t-1}^{\star} + \sigma^{2} \mathcal{M}_{1} \right]^{-1} (B_{t-1}^{\star})^{\top} \beta_{t}, \\ \varepsilon_{var}^{(2)} &:= \mathcal{M}_{2}^{1/2} \left[(B_{t-1}^{\star})^{\top} B_{t-1}^{\star} + \sigma^{2} \mathcal{M}_{1} \right]^{-1} (B_{t-1}^{\star})^{\top} \beta_{t}. \end{split}$$

Here \mathcal{M}_1 and \mathcal{M}_2 are $k \times k$ diagonal matrices with $(\mathcal{M}_1)_{ii} = \frac{1}{\rho_i - 1} \cdot \frac{1}{p} \operatorname{Tr} \left(\Sigma_i^{-1} \right)$ and $(\mathcal{M}_2)_{ii} = \frac{1}{\rho_i - 1} \cdot \frac{1}{p} \operatorname{Tr} \left(\Sigma_i^{-1} \Sigma_k \right)$.

Proof of Proposition B.1. For each task $i, 1 \le i \le t-1$, we can find that

$$\hat{\beta}_i = \beta_i + (X_i^\top X_i)^{-1} X_i^\top \varepsilon_i.$$

Then we calculate that

$$\|\hat{\beta}_i\|^2 = \|\beta_i\|^2 + 2\beta_i^\top \cdot (X_i^\top X_i)^{-1} X_i^\top \varepsilon_i + \|(X_i^\top X_i)^{-1} X_i^\top \varepsilon_i\|^2.$$

We can use concentration of random vector with independent entries, Lemma D.6, to get that for any deterministic vector β with $\|\beta\| = O(1)$,

$$\beta^{\top} (X_i^{\top} X_i)^{-1} X_i^{\top} \varepsilon_i \leqslant p^{\varepsilon/2} \cdot \sigma \left[\beta^{\top} (X_i^{\top} X_i)^{-1} \beta \right]^{1/2} \leqslant \sigma p^{-1/2 + \varepsilon}, \tag{B.3}$$

with high probability for any constant $\varepsilon > 0$, where we used Lemma D.2 to bound

$$\|(X_i^\top X_i)^{-1}\| \lesssim n_i^{-1} \quad \text{whp.} \tag{B.4} \label{eq:B.4}$$

Similarly, using Lemma D.6 we get with high probability,

$$\left| \| (X_i^{\top} X_i)^{-1} X_i^{\top} \varepsilon_i \|^2 - \sigma^2 \operatorname{Tr} \left((X_i^{\top} X_i)^{-1} \right) \right| \leqslant p^{\varepsilon/2} \cdot \sigma^2 \left\{ \operatorname{Tr} \left[X_i (X_i^{\top} X_i)^{-2} X_i^{\top} \right]^2 \right\}^{1/2}$$

$$= p^{\varepsilon/2} \cdot \sigma^2 \left[\operatorname{Tr} (X_i^{\top} X_i)^{-2} \right]^{1/2} \leqslant \sigma^2 p^{-1/2 + \varepsilon},$$
(B.5)

and for $i \neq j$,

$$\left| \varepsilon_j^\top X_j (X_j^\top X_j)^{-1} (X_i^\top X_i)^{-1} X_i^\top \varepsilon_i \right| \leq p^{\varepsilon/2} \cdot \sigma^2 \left\{ \operatorname{Tr} \left[(X_i^\top X_i)^{-1} (X_j^\top X_j)^{-1} \right] \right\}^{1/2}$$

$$\leq \sigma^2 p^{-1/2 + \varepsilon},$$
(B.6)

With (B.3)-(B.6), we get that with high probability,

$$\|\hat{\beta}_{i}\|^{2} = \|\beta_{i}\|^{2} + \sigma^{2} \cdot \operatorname{Tr}\left((X_{i}^{\top}X_{i})^{-1}\right) + \operatorname{O}(\sigma p^{-1/2 + \varepsilon})$$

$$= \|\beta_{i}\|^{2} + \frac{\sigma^{2}}{\rho_{i} - 1} \cdot \frac{1}{n} \operatorname{Tr}\left(\Sigma_{i}^{-1}\right) + \operatorname{O}(\sigma p^{-1/2 + \varepsilon}), \tag{B.7}$$

and

$$\hat{\beta}_i^{\top} \hat{\beta}_j = \beta_i^{\top} \beta_j + \mathcal{O}(\sigma p^{-1/2 + \varepsilon}), \quad i \neq j.$$
 (B.8)

With (B.7) and (B.8), we obtain that with high probability,

$$B^{\top}B = (B_{t-1}^{\star})^{\top}B_{t-1}^{\star} + \sigma^{2}\mathcal{M}_{1} + \mathcal{O}(\sigma p^{-1/2+\varepsilon}), \tag{B.9}$$

where $O(\sigma p^{-1/2+\varepsilon})$ means a $(t-1)\times (t-1)$ matrix, say \mathcal{E} , satisfying $\|\mathcal{E}\|\leqslant C\sigma p^{-1/2+\varepsilon}$. Notice that by (B.1), we also have with high probability,

$$\|(B^{\top}B)^{-1}\| \lesssim \|[(B_{t-1}^{\star})^{\top}B_{t-1}^{\star}]^{-1}\| = O(1).$$

Moreover, using (B.3) we get that

$$B^{\top} \beta = (B_{t-1}^{\star})^{\top} \beta + \mathcal{O}(\sigma p^{-1/2 + \varepsilon}), \tag{B.10}$$

for any deterministic vector β with $\|\beta\| = O(1)$.

Now we are ready to calculate $te(BW_t)$ using the above concentrations results. For the t-th target model, by optimizing over W_t we get

$$\hat{W}_t = (B^\top X_t^\top X_t B)^{-1} B^\top X_t^\top Y_k = (B^\top X_t^\top X_t B)^{-1} B^\top X_t^\top (X_t \beta_t + \varepsilon_k).$$

We then calculate that

$$\mathbb{E}_{\varepsilon_k} \left[\left\| \Sigma_k^{1/2} (B \hat{W}_t - \beta_t) \right\|^2 \right] = \left\| \Sigma_k^{1/2} \left[\operatorname{Id} - B (B^\top X_t^\top X_t B)^{-1} B^\top X_t^\top X_t \right] \beta_t \right\|^2 + \sigma^2 \cdot \operatorname{Tr} \left[\Sigma_k (B^\top X_t^\top X_t B)^{-1} \right].$$
(B.11)

By the concentration of random vector with independent entries or the restricted isometry property for X_tB , we have with high probability,

$$B^{\top} X_t^{\top} X_t B = n_k \left(B^{\top} B + \mathcal{O}(p^{-1/2 + \varepsilon}) \right), \quad B^{\top} X_t^{\top} X_t \beta_t = n_k \left(B^{\top} \beta_t + \mathcal{O}(p^{-1/2 + \varepsilon}) \right).$$

Together with (B.4), (B.9) and (B.10), we can simplify the right-hand side of (B.11) as

$$\mathbb{E}_{\varepsilon_{k}} \left[\left\| \Sigma_{k}^{1/2} (B \hat{W}_{k} - \beta_{t}) \right\|^{2} \right] = \left\| \Sigma_{k}^{1/2} \beta_{t} - \Sigma_{k}^{1/2} B (B^{\top} B)^{-1} B^{\top} \beta_{t} \right\|^{2} + \mathcal{O}(p^{-1/2 + \varepsilon})$$

$$= \left\| \Sigma_{k}^{1/2} \beta_{t} - \Sigma_{k}^{1/2} B \left[(B_{t-1}^{\star})^{\top} B_{t-1}^{\star} + \sigma^{2} \mathcal{M}_{1})^{-1} \right] (B_{t-1}^{\star})^{\top} \beta_{t} \right\|^{2} + \mathcal{O}(p^{-1/2 + \varepsilon}). \tag{B.12}$$

As in (B.9), we can show by concentration that

$$B^{\top} \Sigma_k B = (B_{t-1}^{\star})^{\top} \Sigma_k B_{t-1}^{\star} + \sigma^2 \mathcal{M}_2 + \mathcal{O}(\sigma p^{-1/2 + \varepsilon}).$$
 (B.13)

Together with (B.9) and (B.10), we can further simplify (B.12) as

$$\begin{split} & \underset{\varepsilon_{i},\forall 1 \leqslant i \leqslant k}{\mathbb{E}} \left[\left\| \Sigma_{k}^{1/2} (B \hat{W}_{k} - \beta_{t}) \right\|^{2} \right] \\ & = \mathrm{O}(p^{-1/2 + \varepsilon}) + \left\| \Sigma_{k}^{1/2} \beta_{t} \right\|^{2} - 2\beta_{t}^{\top} \Sigma_{k} B_{t-1}^{\star} \left[(B_{t-1}^{\star})^{\top} B_{t-1}^{\star} + \sigma^{2} \mathcal{M}_{1} \right]^{-1} (B_{t-1}^{\star})^{\top} \beta_{t} \\ & + \beta_{t}^{\top} B_{t-1}^{\star} \left[(B_{t-1}^{\star})^{\top} B_{t-1}^{\star} + \sigma^{2} \mathcal{M}_{1} \right]^{-1} \left((B_{t-1}^{\star})^{\top} \Sigma_{k} B_{t-1}^{\star} + \sigma^{2} \mathcal{M}_{2} \right) \left[(B_{t-1}^{\star})^{\top} B_{t-1}^{\star} + \sigma^{2} \mathcal{M}_{1} \right]^{-1} (B_{t-1}^{\star})^{\top} \beta_{t} \\ & = \| \varepsilon_{\beta} + \sigma^{2} \varepsilon_{\text{var}}^{(1)} \|^{2} + \sigma^{2} \| \varepsilon_{\text{var}}^{(2)} \|^{2} + \mathrm{O}(p^{-1/2 + \varepsilon}). \end{split}$$
 This concludes the proof.

C Proof of Lemma A.2 and Lemma A.6

We consider two $p \times p$ random sample covariance matrices $\mathcal{Q}_1 := \Sigma_1^{1/2} Z_1^\top Z_1 \Sigma_1^{1/2}$ and $\mathcal{Q}_2 := \Sigma_2^{1/2} Z_2^\top Z_2 \Sigma_2^{1/2}$, where Σ_1 and Σ_2 are $p \times p$ deterministic non-negative definite (real) symmetric matrices. We assume that $Z_1 = (z_{ij}^{(1)})$ and $Z_2 = (z_{ij}^{(2)})$ are $n_1 \times p$ and $n_2 \times p$ random matrix with (real) i.i.d. entries satisfying

$$\mathbb{E}z_{ij}^{(\alpha)} = 0, \quad \mathbb{E}|z_{ij}^{(\alpha)}|^2 = n^{-1},$$
 (C.1)

where we denote $n := n_1 + n_2$. Here we have chosen the scaling that is more standard in the random matrix theory literature—under this $n^{-1/2}$ scaling, the eigenvalues of Q_1 and Q_2 are all of order 1. Moreover, we assume that the fourth moment exists:

$$\mathbb{E}|\sqrt{n}z_{ij}^{(\alpha)}|^4 \leqslant C \tag{C.2}$$

for some constant C>0. We assume that the aspect ratios $d_1:=p/n_1$ and $d_2:=p/n_2$ satisfy that

$$0 \le d_1 \le \tau^{-1}, \quad 1 + \tau \le d_2 \le \tau^{-1},$$
 (C.3)

for some small constant $0 < \tau < 1$. Here the lower bound $1 + \tau \le d_2$ is to ensure that the covariance matrix Q_2 for the target task is non-singular with high probability; see Lemma D.2 below.

We assume that Σ_1 and Σ_2 have eigendecompositions

$$\Sigma_{1} = O_{1}\Lambda_{1}O_{1}^{\top}, \quad \Sigma_{2} = O_{2}\Lambda_{2}O_{2}^{\top}, \quad \Lambda_{1} = \operatorname{diag}(\sigma_{1}^{(1)}, \dots, \sigma_{n}^{(1)}), \quad \widetilde{\Sigma} = \operatorname{diag}(\sigma_{1}^{(2)}, \dots, \sigma_{N}^{(2)}), \tag{C.4}$$

where the eigenvalues satisfy that

$$\tau^{-1} \geqslant \sigma_1^{(1)} \geqslant \sigma_2^{(1)} \geqslant \ldots \geqslant \sigma_p^{(1)} \geqslant 0, \quad \tau^{-1} \geqslant \sigma_1^{(2)} \geqslant \sigma_2^{(2)} \geqslant \ldots \geqslant \sigma_p^{(2)} \geqslant \tau, \tag{C.5}$$

for some small constant $0<\tau<1.$ We assume that $M:=\Sigma_1^{1/2}\Sigma_2^{-1/2}$ has singular value decomposition

$$M = U\Lambda V^{\top}, \quad \Lambda = \operatorname{diag}(\sigma_1, \dots, \sigma_n),$$
 (C.6)

where the singular values satisfy that

$$\tau \leqslant \sigma_p \leqslant \sigma_1 \leqslant \tau^{-1} \tag{C.7}$$

for some small constant $0 < \tau < 1$.

We summarize our basic assumptions here for future reference.

Assumption C.1. We assume that Z_1 and Z_2 are independent $n_1 \times p$ and $n_2 \times p$ random matrices with real i.i.d. entries satisfying (C.1) and (C.2), Σ_1 and Σ_2 are deterministic non-negative definite symmetric matrices satisfying (C.4)-(C.7), and $d_{1,2}$ satisfy (C.3).

Before giving the main proof, we first introduce some notations and tools.

C.1 Notations

We will use the following notion of stochastic domination, which was first introduced in [28] and subsequently used in many works on random matrix theory, such as [29, 30, 31, 32, 33, 34]. It simplifies the presentation of the results and their proofs by systematizing statements of the form " ξ is bounded by ζ with high probability up to a small power of n".

Definition C.2 (Stochastic domination). (i) Let

$$\xi = \left(\xi^{(n)}(u) : n \in \mathbb{N}, u \in U^{(n)}\right), \quad \zeta = \left(\zeta^{(n)}(u) : n \in \mathbb{N}, u \in U^{(n)}\right)$$

be two families of nonnegative random variables, where $U^{(n)}$ is a possibly N-dependent parameter set. We say ξ is stochastically dominated by ζ , uniformly in u, if for any fixed (small) $\varepsilon > 0$ and (large) D > 0,

$$\sup_{u \in U^{(n)}} \mathbb{P}\left[\xi^{(n)}(u) > N^{\varepsilon} \zeta^{(n)}(u)\right] \leqslant N^{-D}$$

for large enough $n \ge n_0(\varepsilon, D)$, and we shall use the notation $\xi \prec \zeta$. Throughout this paper, the stochastic domination will always be uniform in all parameters that are not explicitly fixed (such as matrix indices, and z that takes values in some compact set). If for some complex family ξ we have $|\xi| \prec \zeta$, then we will also write $\xi \prec \zeta$ or $\xi = O_{\prec}(\zeta)$.

(ii) We say an event Ξ holds with high probability if for any constant D > 0, $\mathbb{P}(\Xi) \geqslant 1 - n^{-D}$ for large enough n. We say Ξ holds with high probability on an event Ω if for any constant D > 0, $\mathbb{P}(\Omega \setminus \Xi) \leqslant n^{-D}$ for large enough n

The following lemma collects basic properties of stochastic domination \prec , which will be used tacitly in the proof.

Lemma C.3 (Lemma 3.2 in [29]). Let ξ and ζ be families of nonnegative random variables.

- (i) Suppose that $\xi(u,v) \prec \zeta(u,v)$ uniformly in $u \in U$ and $v \in V$. If $|V| \leqslant n^C$ for some constant C, then $\sum_{v \in V} \xi(u,v) \prec \sum_{v \in V} \zeta(u,v)$ uniformly in u.
- (ii) If $\xi_1(u) \prec \zeta_1(u)$ and $\xi_2(u) \prec \zeta_2(u)$ uniformly in $u \in U$, then $\xi_1(u)\xi_2(u) \prec \zeta_1(u)\zeta_2(u)$ uniformly in u.
- (iii) Suppose that $\Psi(u) \geqslant n^{-C}$ is deterministic and $\xi(u)$ satisfies $\mathbb{E}\xi(u)^2 \leqslant n^C$ for all u. Then if $\xi(u) \prec \Psi(u)$ uniformly in u, we have $\mathbb{E}\xi(u) \prec \Psi(u)$ uniformly in u.

Definition C.4 (Bounded support condition). We say a random matrix Z satisfies the bounded support condition with q, if

$$\max_{i,j} |x_{ij}| \prec q. \tag{C.8}$$

Here $q \equiv q(N)$ is a deterministic parameter and usually satisfies $n^{-1/2} \leqslant q \leqslant n^{-\phi}$ for some (small) constant $\phi > 0$. Whenever (C.8) holds, we say that X has support q.

Our main goal is to study the following matrix inverse

$$(Q_1 + Q_2)^{-1} = \left(\Sigma_1^{1/2} Z_1^{\top} Z_1 \Sigma_1^{1/2} + \Sigma_2^{1/2} Z_2^{\top} Z_2 \Sigma_2^{1/2}\right)^{-1}.$$

Using (C.6), we can rewrite it as

$$\Sigma_{2}^{-1/2}V\left(\Lambda U^{\top}Z_{1}^{\top}Z_{1}U\Lambda + V^{\top}Z_{2}^{\top}Z_{2}V\right)^{-1}V^{\top}\Sigma_{2}^{-1/2}.\tag{C.9}$$

For this purpose, we shall study the following matrix for $z \in \mathbb{C}_+$,

$$\mathcal{G}(z) := \left(\Lambda U^{\top} Z_1^{\top} Z_1 U \Lambda + V^{\top} Z_2^{\top} Z_2 V - z\right)^{-1}, \quad z \in \mathbb{C}_+, \tag{C.10}$$

which we shall refer to as resolvent (or Green's function).

Now we introduce a convenient self-adjoint linearization trick. This idea dates back at least to Girko, see e.g., the works [35, 36, 37] and references therein. It has been proved to be useful in studying the local laws of random matrices of the Gram type [38, 39, 34, 40]. We define the following $(p+n)\times(p+n)$ self-adjoint block matrix, which is a linear function of X:

$$H \equiv H(Z_1, Z_2) := \begin{pmatrix} 0 & \Lambda U^{\top} Z_1^{\top} & V^{\top} Z_2^{\top} \\ Z_1 U \Lambda & 0 & 0 \\ Z_2 V & 0 & 0 \end{pmatrix}.$$
 (C.11)

Then we define its resolvent (Green's function) as

$$G \equiv G(Z_1, Z_2, z) := \begin{bmatrix} H(Z_1, Z_2) - \begin{pmatrix} zI_{p \times p} & 0 & 0 \\ 0 & I_{n_1 \times n_1} & 0 \\ 0 & 0 & I_{n_2 \times n_2} \end{pmatrix} \end{bmatrix}^{-1}, \quad z \in \mathbb{C}_+. \quad (C.12)$$

For simplicity of notations, we define the index sets

$$\mathcal{I}_1 := [1, p], \qquad \mathcal{I}_2 := [p+1, p+n_1], \quad \mathcal{I}_3 := [p+n_1+1, p+n_1+n_2], \quad \mathcal{I} := \mathcal{I}_1 \cup \mathcal{I}_2 \cup \mathcal{I}_3.$$

We will consistently use the latin letters $i, j \in \mathcal{I}_1$ and greek letters $\mu, \nu \in \mathcal{I}_2 \cup \mathcal{I}_3$. Moreover, we shall use the notations $\mathfrak{a}, \mathfrak{b} \in \mathcal{I} := \bigcup_{i=1}^3 \mathcal{I}_i$. We label the indices of the matrices according to

$$Z_1 = (z_{\mu i} : i \in \mathcal{I}_1, \mu \in \mathcal{I}_2), \quad Z_2 = (z_{\nu i} : i \in \mathcal{I}_1, \nu \in \mathcal{I}_3).$$

Then we denote the $\mathcal{I}_1 \times \mathcal{I}_1$ block of G(z) by $\mathcal{G}_L(z)$, the $\mathcal{I}_1 \times (\mathcal{I}_2 \cup \mathcal{I}_3)$ by \mathcal{G}_{LR} , the $(\mathcal{I}_2 \cup \mathcal{I}_3) \times \mathcal{I}_1$ block by \mathcal{G}_{RL} , and the $(\mathcal{I}_2 \cup \mathcal{I}_3) \times (\mathcal{I}_2 \cup \mathcal{I}_3)$ block by \mathcal{G}_R . For simplicity, we abbreviate $Y_1 := Z_1 U \Lambda$, $Y_2 := Z_2 V$ and $W := (Y_1^\top, Y_2^\top)$. By Schur complement formula, one can find that (recall (C.10))

$$\mathcal{G}_{11} = \begin{pmatrix} WW^{\top} - z \end{pmatrix}^{-1} = \mathcal{G}, \quad \mathcal{G}_{LR} = \mathcal{G}_{RL}^{\top} = \mathcal{G}W, \quad \mathcal{G}_{R} := \begin{pmatrix} \mathcal{G}_{22} & \mathcal{G}_{23} \\ \mathcal{G}_{32} & \mathcal{G}_{33} \end{pmatrix} = z \begin{pmatrix} W^{\top}W - z \end{pmatrix}^{-1}.$$
(C.13)

Thus a control of G yields directly a control of the resolvent G. We also introduce the following random quantities (some partial traces and weighted partial traces):

$$m(z) := \frac{1}{p} \sum_{i \in \mathcal{I}_1} G_{ii}(z), \quad m_1(z) := \frac{1}{p} \sum_{i \in \mathcal{I}_1} \sigma_i^2 G_{ii}(z),$$

$$m_2(z) := \frac{1}{n_1} \sum_{\mu \in \mathcal{I}_2} G_{\mu\mu}(z), \quad m_3(z) := \frac{1}{n_2} \sum_{\mu \in \mathcal{I}_3} G_{\mu\mu}(z).$$
(C.14)

Next we introduce the spectral decomposition of G. Let

$$W = \sum_{k=1}^{p} \sqrt{\lambda_k} \xi_k \zeta_k^{\top},$$

be a singular value decomposition of W, where

$$\lambda_1 \geqslant \Lambda_2 \geqslant \ldots \geqslant \lambda_p \geqslant 0 = \lambda_{p+1} = \ldots = \lambda_n,$$

 $\{\xi_k\}_{k=1}^p$ are the left-singular vectors, and $\{\zeta_k\}_{k=1}^n$ are the right-singular vectors. Then using (C.13), we can get that for $i,j\in\mathcal{I}_1$ and $\mu,\nu\in\mathcal{I}_2$,

$$G_{ij} = \sum_{k=1}^{p} \frac{\xi_{k}(i)\xi_{k}^{\top}(j)}{\lambda_{k} - z}, \quad G_{\mu\nu} = z \sum_{k=1}^{p} \frac{\zeta_{k}(\mu)\zeta_{k}^{\top}(\nu)}{\lambda_{k} - z} - \sum_{k=p+1}^{n} \zeta_{k}(\mu)\zeta_{k}^{\top}(\nu),$$

$$G_{i\mu} = \sum_{k=1}^{p} \frac{\sqrt{\lambda_{k}}\xi_{k}(i)\zeta_{k}^{\top}(\mu)}{\lambda_{k} - z}, \quad G_{\mu i} = \sum_{k=1}^{p} \frac{\sqrt{\lambda_{k}}\zeta_{k}(\mu)\xi_{k}^{\top}(i)}{\lambda_{k} - z}.$$
(C.15)

We now define the deterministic limit of $\mathcal{G}(z)$. We first define the deterministic limits of $(m_2(z), m_3(z))$, that is $(m_{2c}(z), m_{3c}(z))$, as the (unique) solution to the following system of self-consistent equations

$$\frac{1}{m_{2c}} = -1 + \frac{\gamma_n}{p} \sum_{i=1}^p \frac{\sigma_i^2}{z + \sigma_i^2 r_1 m_{2c} + r_2 m_{3c}}, \quad \frac{1}{m_{3c}} = -1 + \frac{\gamma_n}{p} \sum_{i=1}^p \frac{1}{z + \sigma_i^2 r_1 m_{2c} + r_2 m_{3c}},$$
(C.16)

such that $(m_{2c}(z), m_{3c}(z)) \in \mathbb{C}_+$ for $z \in \mathbb{C}_+$, where, for simplicity, we introduce the parameters

$$\gamma_n := \frac{p}{n}, \quad r_1 \equiv r_1(n) := \frac{n_1}{n}, \quad r_2 \equiv r_2(n) := \frac{n_2}{n}.$$
(C.17)

We then define the matrix limit of G(z) as

$$\Pi(z) := \begin{pmatrix}
-(z + r_1 m_{2c} \Lambda^2 + r_2 m_{3c})^{-1} & 0 & 0 \\
0 & m_{2c}(z) I_{n_1} & 0 \\
0 & 0 & m_{3c}(z) I_{n_2}
\end{pmatrix}.$$
(C.18)

In particular, the matrix limit of $\mathcal{G}(z)$ is given by $-(z + r_1 m_{2c} \Lambda^2 + r_2 m_{3c})^{-1}$.

If z = 0, then the equations (C.16) is reduced to

$$r_1b_2 + r_2b_3 = 1 - \gamma_n, \quad b_2 + \frac{1}{n} \sum_{i=1}^p \frac{\sigma_i^2 b_2}{\sigma_i^2 r_1 b_2 + (1 - \gamma_n - r_1 b_2)} = 1.$$
 (C.19)

where $b_2 := -m_{2c}(0)$ and $b_3 := -m_{3c}(0)$. Note that the function

$$f(b_2) := b_2 + \frac{1}{n} \sum_{i=1}^{p} \frac{\sigma_i^2 b_2}{\sigma_i^2 b_2 + (1 - \gamma_n - r_1 b_2)}$$

is a strictly increasing function on $[0,r_1^{-1}(1-\gamma_n)]$, and f(0)=0<1, $f(r_1^{-1}(1-\gamma_n))=1+\gamma_n>1$. Hence there exists a unique solution (b_2,b_3) to (C.19). Moreover, it is easy to check that $f'(a)=\mathrm{O}(1)$ for $a\in[0,r_1^{-1}(1-\gamma_n)]$, and f(1)>1 if $1\leqslant r_1^{-1}(1-\gamma_n)$. Hence there exists a constant $\tau>0$, such that

$$r_1 \tau \leqslant r_1 b_2 < \min\{(1 - \gamma_n) - r_1 \tau, r_1 (1 - \tau)\}, \quad \tau < r_3 b_3 \leqslant 1 - \gamma_n - r_1 \tau.$$
 (C.20)

For general z around z=0, the existence and uniqueness of the solution $(m_{2c}(z), m_{3c}(z))$ is given by the following lemma. Moreover, we will also include some basic estimates on it. (say something about the previous work)

Lemma C.5. There exist constants c_0 , $C_0 > 0$ depending only on τ in (C.3), (C.5), (C.7) and (C.20) such that the following statements hold. There exists a unique solution to (C.16) under the conditions

$$|z| \le c_0, \quad |m_{2c}(z) - m_{2c}(0)| + |m_{3c}(z) - m_{3c}(0)| \le c_0.$$
 (C.21)

Moreover, the solution satisfies

$$\max_{al=2}^{3} |m_{\alpha c}(z) - m_{\alpha c}(0)| \le C_0 |z|.$$
 (C.22)

The proof is a standard application of the contraction principle. For reader's convenience, we will gives its proof in Appendix D.4. As a byproduct of the contraction mapping argument there, we also obtain the following stability result that will be useful for our proof of Theorem C.7 below.

Lemma C.6. There exist constants c_0 , $C_0 > 0$ depending only on τ in (C.3), (C.5), (C.7) and (C.20) such that the self-consistent equations in (C.16) are stable in the following sense. Suppose $|z| \le c_0$ and $m_\alpha : \mathbb{C}_+ \mapsto \mathbb{C}_+$, $\alpha = 2, 3$, are analytic functions of z such that

$$|m_2(z) - m_{2c}(0)| + |m_3(z) - m_{3c}(0)| \le c_0.$$

Suppose they satisfy the system of equations

$$\frac{1}{m_2} + 1 - \frac{1}{n} \sum_{i=1}^{p} \frac{\sigma_i^2}{z + \sigma_i^2 r_1 m_2 + r_2 m_3} = \mathcal{E}_2, \quad \frac{1}{m_3} + 1 - \frac{1}{n} \sum_{i=1}^{p} \frac{1}{z + \sigma_i^2 r_1 m_2 + r_2 m_3} = \mathcal{E}_3,$$
(C.23)

for some (random) errors satisfying

$$\max_{\alpha=2}^{3} |\mathcal{E}_{\alpha}| \leq \delta(z),$$

where $\delta(z)$ is any deterministic z-dependent function $\delta(z) \leq (\log n)^{-1}$. Then we have

$$\max_{\alpha=2}^{3} |m_{\alpha}(z) - m_{\alpha c}(z)| \leqslant C_0 \delta(z). \tag{C.24}$$

In the following proof, we choose a sufficiently small constants $c_0 > 0$ such that Lemma C.5 and Lemma C.6 hold. Then we define a domain of the spectral parameter z as

$$\mathbf{D} := \left\{ z = E + i\eta \in \mathbb{C}_+ : |z| \le (\log n)^{-1} \right\}. \tag{C.25}$$

The following theorem gives almost optimal estimates on the resolvent G, which are conventionally called local laws.

Theorem C.7. Suppose Assumption C.1 holds, and Z_1, Z_2 satisfy the bounded support condition (C.8) for some deterministic parameter $q \equiv q(n)$ satisfying $n^{-1/2} \leqslant q \leqslant n^{-\phi}$ for some (small) constant $\phi > 0$. Then there exists a sufficiently small constant $c_0 > 0$ such that the following anisotropic local law holds uniformly for all $z \in \mathbf{D}$. For any deterministic unit vectors $\mathbf{u}, \mathbf{v} \in \mathbb{R}^{\mathcal{I}}$, we have

$$\left|\mathbf{u}^{\top}(G(z) - \Pi(z))\mathbf{v}\right| \prec q.$$
 (C.26)

The proof of this theorem will be given in Section D. Using a simple cutoff argument, it is easy to obtain the following corollary under certain moment assumptions.

Corollary C.8. Suppose Assumption C.1 holds. Moreover, assume that the entries of Z_1 and Z_2 are i.i.d. random variables satisfying (C.1) and

$$\max_{i,j} \mathbb{E}|\sqrt{n}z_{ij}^{(\alpha)}|^a = \mathcal{O}(1), \quad \alpha = 1, 2, \tag{C.27}$$

for some fixed a > 4. Then (C.26) holds for $q = n^{2/a - 1/2}$ on an event with probability 1 - o(1).

Proof of Corollary C.8. Fix any sufficiently small constant $\varepsilon > 0$. We then choose $q = n^{-c_a + \varepsilon}$ with $c_a = 1/2 - 2/a$. Then we introduce the truncated matrices \widetilde{Z}_1 and \widetilde{Z}_2 , with entries

$$\widetilde{z}_{ij}^{(\alpha)} := \mathbf{1} \left\{ |\widetilde{z}_{ij}^{(\alpha)}| \leqslant q \right\} \cdot z_{ij}^{(\alpha)}, \quad \alpha = 1, 2.$$

By the moment conditions (C.27) and a simple union bound, we have

$$\mathbb{P}(\widetilde{Z}_1 = Z_1, \widetilde{Z}_2 = Z_2) = 1 - \mathcal{O}(n^{-a\varepsilon}). \tag{C.28}$$

Using (C.27) and integration by parts, it is easy to verify that

$$\mathbb{E}|z_{ij}^{(\alpha)}|1_{|z_{ij}^{(\alpha)}|>q} = \mathcal{O}(n^{-2-\varepsilon}), \quad \mathbb{E}|z_{ij}^{(\alpha)}|^2 1_{|z_{ij}^{(\alpha)}|>q} = \mathcal{O}(n^{-2-\varepsilon}), \quad \alpha = 1, 2,$$

which imply that

$$|\mathbb{E}\tilde{z}_{ij}^{(\alpha)}| = \mathcal{O}(n^{-2-\varepsilon}), \quad \mathbb{E}|\tilde{z}_{ij}^{(\alpha)}|^2 = n^{-1} + \mathcal{O}(n^{-2-\varepsilon}), \quad \alpha = 1, 2, . \tag{C.29}$$

Moreover, we trivially have

$$\mathbb{E}|\tilde{z}_{ij}^{(\alpha)}|^4 \leqslant \mathbb{E}|z_{ij}^{(\alpha)}|^4 = O(n^{-2}), \quad \alpha = 1, 2.$$

Then we centralize and rescale \widetilde{Z}_1 and \widetilde{Z}_2 as

$$\widehat{Z}_{\alpha} := \frac{\widetilde{Z}_{\alpha} - \mathbb{E}\widetilde{Z}_{\alpha}}{(\mathbb{E}|\widetilde{z}_{11}^{(\alpha)}|^2)^{1/2}}, \quad \alpha = 1, 2.$$

Now \widehat{Z}_1 and \widehat{Z}_2 satisfy the assumptions in Theorem C.7 with $q=n^{-c_a+\varepsilon}$, and (C.26) gives that

$$\left|\mathbf{u}^{\top}(G(\widehat{Z}_1,\widehat{Z}_2,z)-\Pi(z))\mathbf{v}\right| \prec q.$$

Then using (C.29) and (D.4) below, we can easily get that

$$\left|\mathbf{u}^{\mathsf{T}}(G(\widehat{Z}_1,\widehat{Z}_2,z)-G(\widetilde{Z}_1,\widetilde{Z}_2,z))\mathbf{v}\right| \prec n^{-1-\varepsilon},$$

where we also used the bound $\|\mathbb{E}\widetilde{Z}_{\alpha}\| = \mathrm{O}(n^{-1-\varepsilon})$. This shows that (C.26) also holds for $G(\widetilde{Z}_1,\widetilde{Z}_2,z)$ with $q=n^{-c_a+\varepsilon}$, and hence concludes the proof by (C.28).

Using Corollary C.8, we can complete the proof of Lemma A.2 and Lemma A.6.

Proof of Lemma A.2. In the setting of Lemma A.2, we can write

$$\mathcal{R} := (w^2 X_1^\top X_1 + X_2^\top X_2)^{-1} = n^{-1} \left(\widetilde{\Sigma}_1^{1/2} Z_1^\top Z_1 \widetilde{\Sigma}_1^{1/2} + \Sigma_2^{1/2} Z_2^\top Z_2 \Sigma_2^{1/2} \right)^{-1},$$

where $\widetilde{\Sigma}_1:=w^2\Sigma_1, \Sigma_2, Z_1$ and Z_2 satisfy Assumption C.1. Here the extra n^{-1} is due to the choice of the variances—in the setting of Lemma A.2 the variances of the entries of $Z_{1,2}$ are equal to 1, while in (C.1) they are taken to be n^{-1} . As in (C.6), we assume that $M:=\widetilde{\Sigma}_1^{1/2}\Sigma_2^{-1/2}$ has singular value decomposition

$$M = U\Lambda V^{\top}, \quad \Lambda = \operatorname{diag}(\sigma, \dots, \sigma_n).$$
 (C.30)

Then as in (C.9), we can write

$$\mathcal{R} = \Sigma_2^{-1/2} V \mathcal{G}(0) V^{\top} \Sigma_2^{-1/2}, \quad \mathcal{G}(0) = \left(\Lambda U^{\top} Z_1^{\top} Z_1 U \Lambda + V^{\top} Z_2^{\top} Z_2 V\right)^{-1}.$$

Now by Corollary C.8, we obtain that for any small constant $\varepsilon > 0$, with probability 1 - o(1),

$$\max_{1 \le i \le p} |(\Sigma_2 \mathcal{R} - \Sigma_2^{1/2} V \Pi(0) V^{\top} \Sigma_2^{-1/2})_{ii}| \le n^{\varepsilon} q, \quad q = n^{2/a - 1/2}, \tag{C.31}$$

where by (C.18),

$$\Pi(0) = -(r_1 m_{2c}(0) \Lambda^2 + r_2 m_{3c}(0))^{-1} = (r_1 b_2 V^{\top} M^{\top} M V + r_2 b_3)^{-1},$$

with (b_2, b_3) satisfying (C.19). Thus from (C.31) we get that

$$n^{-1} \operatorname{Tr}(\Sigma_2 \mathcal{R}) = n^{-1} \operatorname{Tr}(r_1 b_2 M^{\top} M + r_2 b_3)^{-1} + O(n^{\epsilon} q)$$

with probability 1 - o(1). This concludes (??) if we rename $r_1b_2 \to a_1$ and $r_2b_3 \to a_2$. For (??), it is a well-known result for inverse Whishart matrices (add some references). In fact, if we set $n_1 = 0$ and $n_2 = n$, then it is easy to check that $a_1 = 0$ and $a_2 = (n_2 - p)/n_2$ is the solution to (A.7). This gives (??) by (??).

Proof of Lemma A.6. In the setting of Lemma A.6, we can write

$$\Delta := n^2 \left\| \Sigma_2^{1/2} (\hat{w}^2 X_1^\top X_1 + X_2^\top X_2)^{-1} \Sigma_1 (\beta_s - w \beta_t) \right\|^2$$

= $(\beta_s - w \beta_t) \Sigma_1^{1/2} M \left(M^\top Z_1^\top Z_1 M + Z_2^\top Z_2 \right)^{-2} M^\top \Sigma_1^{1/2} (\beta_s - w \beta_t),$

where $\widetilde{\Sigma}_1 := w^2 \Sigma_1$, Σ_2 , Z_1 and Z_2 satisfy Assumption C.1 and $M := \widetilde{\Sigma}_1^{1/2} \Sigma_2^{-1/2}$. Here again the n^2 factor disappears due to the choice of scaling. Again we assume that M has the singular value decomposition (C.30). Then we can write

$$\Delta := \mathbf{v}^{\top}(\mathcal{G}^2)(0)\,\mathbf{v}, \quad \mathbf{v} := V^{\top}M^{\top}\Sigma_1^{1/2}(\beta_s - w\beta_t).$$

Note that $\mathcal{G}^2(0) = \partial_z \mathcal{G}|_{z=0}$. Now using Cauchy's integral formula and Corollary C.8, we get that with probability 1 - o(1),

$$\mathbf{v}^{\top} \mathcal{G}^{2}(0) \mathbf{v} = \frac{1}{2\pi \mathrm{i}} \oint_{\mathcal{C}} \frac{\mathbf{v}^{\top} \mathcal{G}(z) \mathbf{v}}{z^{2}} dz = \frac{1}{2\pi \mathrm{i}} \oint_{\mathcal{C}} \frac{\mathbf{v}^{\top} \Pi(z) \mathbf{v}}{z^{2}} dz + \mathcal{O}_{\prec}(q) = \mathbf{v}^{\top} \Pi'(0) \mathbf{v} + \mathcal{O}_{\prec}(q),$$
(C.32)

where C is the contour $\{z \in \mathbb{C} : |z| \leq (\log n)^{-1}\}$ and we used (C.26) in the second step. Hence it remains to study the derivatives

$$\mathbf{v}^{\top} \Pi'(0) \mathbf{v} = \mathbf{v} \frac{1 + r_1 m'_{2c}(0) \Lambda^2 + r_2 m'_{3c}(0)}{(r_1 m_{2c}(0) \Lambda^2 + r_2 m_{3c}(0))^2} \mathbf{v}.$$
 (C.33)

It remains to calculate the derivatives $m'_{2c}(0)$ and $m'_{3c}(0)$.

By the implicit differentiation of (C.16), we obtain that

$$\frac{1}{m_{2c}^2(0)}m_{2c}'(0) = \frac{1}{n} \sum_{i=1}^p \frac{\sigma_i^2 \left(1 + \sigma_i^2 r_1 m_{2c}'(0) + r_2 m_{3c}'(0)\right)}{(\sigma_i^2 r_1 m_{2c}(0) + r_2 m_{3c}(0))^2},$$

$$\frac{1}{m_{2c}^2(0)} m_{2c}'(0) = \frac{1}{n} \sum_{i=1}^p \frac{\sigma_i^2 \left(1 + \sigma_i^2 r_1 m_{2c}'(0) + r_2 m_{3c}'(0)\right)}{(\sigma_i^2 r_1 m_{2c}'(0) + r_2 m_{3c}'(0))^2},$$

$$\frac{1}{m_{3c}^2(0)}m_{3c}'(0) = \frac{1}{n}\sum_{i=1}^p \frac{1+\sigma_i^2 r_1 m_{2c}'(0) + r_2 m_{3c}'(0)}{(\sigma_i^2 r_1 m_{2c}(0) + r_2 m_{3c}(0))^2}.$$

If we rename $-r_1m_{2c}(0) \to a_1, -r_2m_{3c}(0) \to a_2, r_2m_{3c}'(0) \to a_3$ and $r_1m_{2c}'(0) \to a_4$, then this equation becomes

$$\left(\frac{r_2}{a_2^2} - \frac{1}{n} \sum_{i=1}^p \frac{1}{(\sigma_i^2 a_1 + a_2)^2}\right) a_3 - \left(\frac{1}{n} \sum_{i=1}^p \frac{\sigma_i^2}{(\sigma_i^2 a_1 + a_2)^2}\right) a_4 = \frac{1}{n} \sum_{i=1}^p \frac{1}{(\sigma_i^2 a_1 + a_2)^2},
\left(\frac{r_1}{a_1^2} - \frac{1}{n} \sum_{i=1}^p \frac{\sigma_i^4}{(\sigma_i^2 a_1 + a_2)^2}\right) a_4 - \left(\frac{1}{n} \sum_{i=1}^p \frac{\sigma_i^2}{(\sigma_i^2 a_1 + a_2)^2}\right) a_3 = \frac{1}{n} \sum_{i=1}^p \frac{\sigma_i^2}{(\sigma_i^2 a_1 + a_2)^2}.$$
(C.34)

Then by (C.32) and (C.33), we get

$$\Delta = (\beta_s - w\beta_t)^{\top} \Sigma_1^{1/2} M V \frac{1 + a_3 + a_4 \Lambda^2}{(a_1 \Lambda^2 + a_2)} V^{\top} M^{\top} \Sigma_1^{1/2} (\beta_s - w\beta_t)$$
$$= (\beta_s - w\beta_t)^{\top} \Sigma_1^{1/2} M \frac{1 + a_3 + a_4 M^{\top} M}{(a_1 M^{\top} M + a_2)} M^{\top} \Sigma_1^{1/2} (\beta_s - w\beta_t)$$

where we used $M^{\top}M = V\Lambda^2V^{\top}$ in the second step. This concludes Lemma A.6.

D Proof of Theorem C.7

The main difficulty for the proof of Theorem C.7 is due to the fact that the entries of $Y_1=Z_1U\Lambda$ and $Y_2=Z_2V$ are not independent. However, notice that if the entries of $Z_1\equiv Z_1^{Gauss}$ and $Z_2\equiv Z_2^{Gauss}$ are i.i.d. Gaussian, then by the rotational invariance of the multivariate Gaussian distribution, we have

$$Z_1^{Gauss}U\Lambda \stackrel{d}{=} Z_1^{Gauss}\Lambda, \quad Z_2^{Gauss}V \stackrel{d}{=} Z_2^{Gauss}.$$

In this case, the problem is reduced to proving the anisotropic local law for G with $U = \operatorname{Id}$ and $V = \operatorname{Id}$, such that the entries of Y_1 and Y_2 are independent. This can be handled using the standard resolvent methods as in e.g. [29, 41, 42]. To go from the Gaussian case to the general X case, we will adopt a continuous self-consistent comparison argument developed in [34].

For the case $U = \operatorname{Id}$ and $V = \operatorname{Id}$, we need to deal with the following resolvent:

$$G_0(z) := \begin{pmatrix} -z & \Lambda Z_1^\top & Z_2^\top \\ Z_1 \Lambda & -I_{n_1} & 0 \\ Z_2 & 0 & -I_{n_2} \end{pmatrix}^{-1}, \quad z \in \mathbb{C}_+,$$
 (D.1)

and prove the following result.

Proposition D.1. Suppose Assumption C.1 holds, and Z_1, Z_2 satisfy the bounded support condition (C.8) with $q = n^{-1/2}$. Suppose U and V are identity. Then the estimate (C.26) holds for $G_0(z)$.

In Section D.1, we will collect some a priori estimates and resolvent identities that will be used in the proof of Theorem C.7 and Proposition D.1. Then in Section D.2 we give the proof of Proposition D.1, which, as discussed above, concludes Theorem C.7 for i.i.d. Gaussian Z_1 and Z_2 . Finally, in Section D.3, we will describe how to extend the result in Theorem C.7 from the Gaussian case to the case with generally distributed entries of Z_1 and Z_2 . In the proof, we always denote the spectral parameter by $z = E + \mathrm{i} \eta$.

D.1 Basic estimates

The estimates in this section work for general G, that is, we do not require U and V to be identity.

First, note that $Z_1^{ op}Z_1$ (resp. $Z_2^{ op}Z_2$) is a standard sample covariance matrix, and it is well-known that its nonzero eigenvalues are all inside the support of the Marchenko-Pastur law $[(1-\sqrt{d_1})^2,(1+\sqrt{d_1})^2]$ (resp. $[(1-\sqrt{d_2})^2,(1+\sqrt{d_2})^2]$) with probability $1-\mathrm{o}(1)$ [43]. In our proof, we shall need a slightly stronger probability bound, which is given by the following lemma. Denote the nonzero eigenvalues of $Z_1^{ op}Z_1$ and $Z_2^{ op}Z_2$ by $\lambda_1(Z_1^{ op}Z_1)\geqslant \cdots \geqslant \lambda_{p\wedge n_1}(Z_1^{ op}Z_1)$ and $\lambda_1(Z_2^{ op}Z_2)\geqslant \cdots \geqslant \lambda_p(Z_2^{ op}Z_2)$.

Lemma D.2. Suppose Assumption C.1 holds, and Z_1, Z_2 satisfy the bounded support condition (C.8) for some deterministic parameter $q \equiv q(n)$ satisfying $n^{-1/2} \leqslant q \leqslant n^{-\phi}$ for some (small) constant $\phi > 0$. Then for any constant $\varepsilon > 0$, we have with high probability,

$$\lambda_1(Z_1^\top Z_1) \leqslant (1 + \sqrt{d_1})^2 + \varepsilon,\tag{D.2}$$

and

$$(1 - \sqrt{d_2})^2 - \varepsilon \leqslant \lambda_p(Z_2^\top Z_2) \leqslant \lambda_1(Z_2^\top Z_2) \leqslant (1 + \sqrt{d_2})^2 + \varepsilon. \tag{D.3}$$

Proof. This lemma essentially follows from [29, Theorem 2.10], although the authors considered the case with $q \prec n^{-1/2}$ only. The results for larger q follows from [44, Lemma 3.12], but only the bounds for the largest eigenvalues are given there in order to avoid the issue with the smallest eigenvalue when d_2 is close to 1. However, under the assumption (C.3), the lower bound for the smallest eigenvalue follows from the exactly the same arguments as in [44]. Hence we omit the details.

With this lemma, we can easily obtain the following a priori estimate on the resolvent G(z) for $z \in \mathbf{D}$.

Lemma D.3. Suppose the assumptions of Lemma D.2 holds. Then there exists a constant C > 0 such that the following estimates hold uniformly in $z, z' \in \mathbf{D}$ with high probability:

$$||G(z)|| \leqslant C,\tag{D.4}$$

and for any deterministic unit vectors $\mathbf{u}, \mathbf{v} \in \mathbb{R}^{\mathcal{I}}$,

$$\left|\mathbf{u}^{\top} \left[G(z) - G(z') \right] \mathbf{v} \right| \leqslant C|z - z'|. \tag{D.5}$$

Proof. As in (C.15), we let $\{\lambda_k\}_{1\leqslant k\leqslant p}$ be the eigenvalues of WW^{\top} . By Lemma D.2 and the assumption (C.3), we obtain that

$$\lambda_p \geqslant \lambda_p(Z_2^\top Z_2) \geqslant \varepsilon > 0 \tag{D.6}$$

for some constant $\varepsilon > 0$. In particular, it implies that

$$\inf_{z \in \mathbf{D}} \min_{1 \leqslant k \leqslant p} |\lambda_k - z| \gtrsim 1.$$

Together with (C.15), it implies the estimates (D.4) and (D.5).

Now we introduce the concept of minors, which are defined by removing certain rows and columns of the matrix H.

Definition D.4 (Minors). For any $(p+n) \times (p+n)$ matrix \mathcal{A} and $\mathbb{T} \subseteq \mathcal{I}$, we define the minor $\mathcal{A}^{(\mathbb{T})} := (\mathcal{A}_{\mathfrak{ab}} : \mathfrak{a}, \mathfrak{b} \in \mathcal{I} \setminus \mathbb{T})$ as the $(p+n-|\mathbb{T}|) \times (p+n-|\mathbb{T}|)$ matrix obtained by removing all rows and columns indexed by \mathbb{T} . Note that we keep the names of indices when defining $\mathcal{A}^{(\mathbb{T})}$, i.e. $(\mathcal{A}^{(\mathbb{T})})_{ab} = \mathcal{A}_{ab}$ for $a, b \notin \mathbb{T}$. Correspondingly, we define the resolvent minor as (recall (C.13))

$$G^{(\mathbb{T})} := \left[\begin{pmatrix} H - \begin{pmatrix} zI_p & 0 \\ 0 & I_n \end{pmatrix} \end{pmatrix}^{(\mathbb{T})} \right]^{-1} = \begin{pmatrix} \mathcal{G}^{(\mathbb{T})} & \mathcal{G}^{(\mathbb{T})}W^{(\mathbb{T})} \\ \left(W^{(\mathbb{T})}\right)^{\top}\mathcal{G}^{(\mathbb{T})} & \mathcal{G}_R^{(\mathbb{T})} \end{pmatrix},$$

and the partial traces (recall (C.14))

$$m^{(\mathbb{T})} := \frac{1}{p} \sum_{i \in \mathcal{I}_1} G_{ii}^{(\mathbb{T})}(z), \quad m_1^{(\mathbb{T})} := \frac{1}{p} \sum_{i \in \mathcal{I}_1}^{(\mathbb{T})} \sigma_i^2 G_{ii}^{(\mathbb{T})}(z),$$

$$m_2^{(\mathbb{T})}(z) := \frac{1}{n_1} \sum_{\mu \in \mathcal{I}_2} G_{\mu\mu}^{(\mathbb{T})}(z), \quad m_3^{(\mathbb{T})}(z) := \frac{1}{n_2} \sum_{\mu \in \mathcal{I}_3} G_{\mu\mu}^{(\mathbb{T})}(z),$$
(D.7)

where we abbreviated that $\sum_{a \notin \mathbb{T}}^{(\mathbb{T})} := \sum_{a \notin \mathbb{T}}$. For convenience, we will adopt the convention that for any minor $\mathcal{A}^{(T)}$ defined as above, $\mathcal{A}^{(T)}_{ab} = 0$ if $a \in \mathbb{T}$ or $b \in \mathbb{T}$. Moreover, we will abbreviate $(\{a\}) \equiv (a)$ and $(\{a,b\}) \equiv (ab)$.

Then we record the following resolvent identities.

Lemma D.5. (Resolvent identities).

(i) For $i \in \mathcal{I}_1$ and $\mu \in \mathcal{I}_2 \cup \mathcal{I}_3$, we have

$$\frac{1}{G_{ii}} = -z - \left(WG^{(i)}W^{\top}\right)_{ii}, \quad \frac{1}{G_{\mu\mu}} = -1 - \left(W^{\top}G^{(\mu)}W\right)_{\mu\mu}. \tag{D.8}$$

(ii) For $i \neq j \in \mathcal{I}_1$ and $\mu \neq \nu \in \mathcal{I}_2 \cup \mathcal{I}_3$, we have

$$G_{ij} = -G_{ii} \left(W G^{(i)} \right)_{ij}, \quad G_{\mu\nu} = -G_{\mu\mu} \left(W^{\top} G^{(\mu)} \right)_{\mu\nu}.$$
 (D.9)

For $i \in \mathcal{I}_1$ and $\mu \in \mathcal{I}_2$, we have

$$G_{i\mu} = -G_{ii} \left(W G^{(i)} \right)_{i\mu}, \quad G_{\mu i} = -G_{\mu\mu} \left(W^{\top} G^{(\mu)} \right)_{\mu i}.$$
 (D.10)

(iii) For $a \in \mathcal{I}$ and $b, c \in \mathcal{I} \setminus \{a\}$,

$$G_{bc}^{(a)} = G_{bc} - \frac{G_{ba}G_{ac}}{G_{aa}}, \quad \frac{1}{G_{bb}} = \frac{1}{G_{bb}^{(a)}} - \frac{G_{ba}G_{ab}}{G_{bb}G_{bb}^{(a)}G_{aa}}.$$
 (D.11)

Proof. All these identities can be proved directly using Schur's complement formula. The reader can also refer to, for example, [34, Lemma 4.4].

The following lemma gives large deviation bounds for bounded supported random variables.

Lemma D.6 (Lemma 3.8 of [45]). Let (x_i) , (y_j) be independent families of centered and independent random variables, and (A_i) , (B_{ij}) be families of deterministic complex numbers. Suppose the entries x_i , y_j have variance at most n^{-1} and satisfy the bounded support condition (C.8) with $q \le n^{-\phi}$ for some constant $\phi > 0$. Then we have the following bound:

$$\left| \sum_{i} A_{i} x_{i} \right| \prec q \max_{i} |A_{i}| + \frac{1}{\sqrt{n}} \left(\sum_{i} |A_{i}|^{2} \right)^{1/2}, \quad \left| \sum_{i,j} x_{i} B_{ij} y_{j} \right| \prec q^{2} B_{d} + q B_{o} + \frac{1}{n} \left(\sum_{i \neq j} |B_{ij}|^{2} \right)^{1/2},$$
(D.12)
$$\left| \sum_{i} \bar{x}_{i} B_{ii} x_{i} - \sum_{i} (\mathbb{E}|x_{i}|^{2}) B_{ii} \right| \prec q B_{d}, \quad \left| \sum_{i \neq j} \bar{x}_{i} B_{ij} x_{j} \right| \prec q B_{o} + \frac{1}{n} \left(\sum_{i \neq j} |B_{ij}|^{2} \right)^{1/2},$$
(D.13)

where $B_d := \max_i |B_{ii}|$ and $B_o := \max_{i \neq j} |B_{ij}|$.

D.2 Entrywise local law

The main goal of this subsection is to prove the following entrywise local law. The anisotropic local law (C.26) then follows from the entrywise local law combined with a polynomialization method as we will explain in next subsection. Recall that in the setting of Proposition D.1, we have $q=n^{-1/2}$ and

$$W = (\Lambda Z_1^{\top}, Z_2^{\top}). \tag{D.14}$$

Lemma D.7. Suppose the assumptions in Proposition D.1 hold. Then the following estimate holds uniformly for $z \in \mathbf{D}$:

$$\max_{\mathfrak{a},\mathfrak{b}} |(G_0)_{\mathfrak{a}\mathfrak{b}}(z) - \Pi_{\mathfrak{a}\mathfrak{b}}(z)| \prec n^{-1/2}. \tag{D.15}$$

Proof. The proof of Lemma D.7 is divided into three steps. For simplicity, we will still denote $G \equiv G_0$ in the following proof, while keeping in mind that W takes the form in (D.14).

Step 1: Large deviations estimates. In this step, we prove some (almost) optimal large deviations estimates on the off-diagonal entries of G, and on the following Z variables. In analogy to [45, Section 3] and [34, Section 5], we introduce the Z variables

$$Z_{\mathfrak{a}}^{(\mathbb{T})} := (1 - \mathbb{E}_{\mathfrak{a}}) \left(G_{\mathfrak{a}\mathfrak{a}}^{(\mathbb{T})} \right)^{-1}, \quad \mathfrak{a} \notin \mathbb{T},$$

where $\mathbb{E}_{\mathfrak{a}}[\cdot] := \mathbb{E}[\cdot \mid H^{(\mathfrak{a})}]$, i.e. it is the partial expectation over the randomness of the \mathfrak{a} -th row and column of H. By (D.8), we have

$$Z_{i} = (\mathbb{E}_{i} - 1) \left(W G^{(i)} W^{\top} \right)_{ii} = \sigma_{i}^{2} \sum_{\mu, \nu \in \mathcal{I}_{2}} G_{\mu\nu}^{(i)} \left(\frac{1}{n} \delta_{\mu\nu} - z_{\mu i} z_{\nu i} \right) + \sum_{\mu, \nu \in \mathcal{I}_{3}} G_{\mu\nu}^{(i)} \left(\frac{1}{n} \delta_{\mu\nu} - z_{\mu i} z_{\nu i} \right), \quad i \in \mathcal{I}_{1},$$
(D.16)

and

$$Z_{\mu} = (\mathbb{E}_{\mu} - 1) \left(W^{\top} G^{(\mu)} W \right)_{\mu\mu} = \sum_{i,j \in \mathcal{I}_{1}} \sigma_{i} \sigma_{j} G_{ij}^{(\mu)} \left(\frac{1}{n} \delta_{ij} - z_{\mu i} z_{\mu j} \right), \quad \mu \in \mathcal{I}_{2},$$

$$Z_{\mu} = (\mathbb{E}_{\mu} - 1) \left(W^{\top} G^{(\mu)} W \right)_{\mu\mu} = \sum_{i,j \in \mathcal{I}_{1}} G_{ij}^{(\mu)} \left(\frac{1}{n} \delta_{ij} - z_{\mu i} z_{\mu j} \right), \quad \mu \in \mathcal{I}_{3}.$$
(D.17)

For simplicity, we introduce the following random error

$$\Lambda_o := \max_{\mathfrak{a} \neq \mathfrak{b}} \left| G_{\mathfrak{a}\mathfrak{a}}^{-1} G_{\mathfrak{a}\mathfrak{b}} \right|. \tag{D.18}$$

The following lemma gives the desired large deviations estimates on the Λ_o and the Z variables.

Lemma D.8. Suppose the assumptions in Proposition D.1 hold. Then the following estimates hold uniformly for all $z \in \mathbf{D}$:

$$\Lambda_o + \max_{\mathfrak{a} \in \mathcal{I}} |Z_{\mathfrak{a}}| \prec n^{-1/2}. \tag{D.19}$$

Proof. Note that for any $\mathfrak{a} \in \mathcal{I}$, $H^{(\mathfrak{a})}$ and $G^{(\mathfrak{a})}$ also satisfies the assumptions for Lemma D.3. Hence (D.4) and (D.5) also hold for $G^{(\mathfrak{a})}$. Now applying Lemma D.6 to (D.16) and (D.17), and using the a priori bound (D.4), we get that for any $i \in \mathcal{I}_1$,

$$|Z_i| \lesssim \sum_{\alpha=2}^{3} \left| \sum_{\mu,\nu \in \mathcal{I}_{\alpha}} G_{\mu\nu}^{(i)} \left(\frac{1}{n} \delta_{\mu\nu} - z_{\mu i} z_{\nu i} \right) \right| \prec n^{-1/2} + \frac{1}{n} \left(\sum_{\mu,\nu \in \mathcal{I}_2 \cup \mathcal{I}_3} \left| G_{\mu\nu}^{(i)} \right|^2 \right)^{1/2} \prec n^{-1/2},$$

where in the last step we used that for any μ ,

$$\sum_{\nu \in \mathcal{I}_2 \cup \mathcal{I}_3} \left| G_{\mu\nu}^{(i)} \right|^2 \leqslant \sum_{\mathfrak{a} \in \mathcal{I}} \left| G_{\mu\mathfrak{a}}^{(i)} \right|^2 = \left[G^{(i)} (G^{(i)})^* \right]_{\mu\mu} = \mathcal{O}(1) \tag{D.20}$$

by (D.4). Similarly, applying Lemma D.6 to Z_{μ} in (D.17) and using (D.4), we obtain the same bound. Then we prove the off-diagonal estimates. For $i \neq j \in \mathcal{I}_1$ and $\mu \neq \nu \in \mathcal{I}_2 \cup \mathcal{I}_3$, using (D.9), Lemma D.6 and (D.4), we obtain that

$$\left| G_{ii}^{-1} G_{ij} \right| \prec n^{-1/2} + \frac{1}{\sqrt{n}} \left(\sum_{\mu \in \mathcal{I}_2 \cup \mathcal{I}_3} \left| G_{\mu j}^{(i)} \right|^2 \right)^{1/2} \prec n^{-1/2},$$

and

$$\left| G_{\mu\mu}^{-1} G_{\mu\nu} \right| \prec n^{-1/2} + \frac{1}{\sqrt{n}} \left(\sum_{i \in \mathcal{I}_1} \left| G_{i\nu}^{(\mu)} \right|^2 \right)^{1/2} \prec n^{-1/2}.$$

For $i \in \mathcal{I}_1 \cup \mathcal{I}_2$ and $\mu \in \mathcal{I}_3$, using (D.10), Lemma D.6 and (D.4), we obtain that

$$\left| G_{ii}^{-1} G_{i\mu} \right| + \left| G_{\mu\mu}^{-1} G_{\mu i} \right| \prec n^{-1/2} + \frac{1}{\sqrt{n}} \left(\sum_{\nu \in \mathcal{I}_2 \cup \mathcal{I}_3} \left| G_{\nu\mu}^{(i)} \right|^2 \right)^{1/2} + \frac{1}{\sqrt{n}} \left(\sum_{j \in \mathcal{I}_1} \left| G_{ji}^{(\mu)} \right|^2 \right)^{1/2} \prec n^{-1/2}.$$

Thus we obtain that $\Lambda_o \prec n^{-1/2}$, which concludes (D.19).

Note that comibining (D.4) and (D.19), we immediately conclude (D.15) for $a \neq b$.

Step 2: Self-consistent equations. This is the key step of the proof for Proposition D.7, which derives approximate self-consistent equations safisfised by $m_2(z)$ and $m_3(z)$. More precisely, we will show that $(m_2(z), m_3(z))$ satisfies (C.23) up to some small error $|\mathcal{E}_{2,3}| \prec n^{-1/2}$. Then applying Lemma C.6 shows that $(m_2(z), m_3(z))$ is close to $(m_{2c}(z), m_{3c}(z))$ —this will discussed in Step 3.

We define the following z-dependent event

$$\Xi(z) := \left\{ |m_2(z) - m_{2c}(z)| + |m_3(z) - m_{3c}(z)| \le (\log n)^{-1/2} \right\}.$$
 (D.21)

Note that by (C.22), we have $|m_{2c}+b_2| \lesssim (\log n)^{-1}$ and $|m_{3c}+b_3| \lesssim (\log n)^{-1}$. Together with (C.16), (C.20) and (C.7), we obtain the following basic estimates

$$|m_{2c}| \sim 1$$
, $|m_{3c}| \sim 1$, $|z + \sigma_i^2 r_1 m_{2c} + r_2 m_{3c}| \sim 1$, $|1 + \gamma_n m_c| \sim 1$, $|1 + \gamma_n m_{1c}| \sim 1$, (D.22)

uniformly in $z \in \mathbf{D}$, where we abbreviate

$$m_c(z) := -\frac{1}{n} \sum_{i \in \mathcal{I}_1} \frac{1}{z + \sigma_i^2 r_1 m_{2c} + r_2 m_{3c}}, \quad m_{1c}(z) := -\frac{1}{n} \sum_{i \in \mathcal{I}_1} \frac{\sigma_i^2}{z + \sigma_i^2 r_1 m_{2c} + r_2 m_{3c}}.$$

Plugging (D.22) into (C.18), we get

$$|\Pi_{\mathfrak{a}\mathfrak{a}}(z)| \sim 1$$
 uniformly in $z \in \mathbf{D}, \ \mathfrak{a} \in \mathcal{I}.$ (D.23)

Then we claim the following result.

Lemma D.9. Suppose the assumptions in Proposition D.1 hold. Then the following estimates hold uniformly in $z \in \mathbf{D}$:

$$\mathbf{1}(\Xi) \left| \frac{1}{m_2} + 1 - \frac{1}{n} \sum_{i=1}^{p} \frac{\sigma_i^2}{z + \sigma_i^2 r_1 m_2 + r_2 m_3} \right| < n^{-1/2},$$

$$\mathbf{1}(\Xi) \left| \frac{1}{m_3} + 1 - \frac{1}{n} \sum_{i=1}^{p} \frac{1}{z + \sigma_i^2 r_2 m_2 + r_2 m_3} \right| < n^{-1/2}.$$
(D.24)

Proof. By (D.8), (D.16) and (D.17), we obtain that for $i \in \mathcal{I}_1$, $\mu \in \mathcal{I}_2$ and $\nu \in \mathcal{I}_3$,

$$\frac{1}{G_{ii}} = -z - \frac{\sigma_i^2}{n} \sum_{\mu \in \mathcal{I}_2} G_{\mu\mu}^{(i)} - \frac{1}{n} \sum_{\mu \in \mathcal{I}_3} G_{\mu\mu}^{(i)} + Z_i = -z - \sigma_i^2 r_1 m_2 - r_2 m_3 + \varepsilon_i, \tag{D.25}$$

$$\frac{1}{G_{\mu\mu}} = -1 - \frac{1}{n} \sum_{i \in \mathcal{I}_1} \sigma_i^2 G_{ii}^{(\mu)} + Z_{\mu} = -1 - \gamma_n m_1 + \varepsilon_{\mu}, \tag{D.26}$$

$$\frac{1}{G_{\nu\nu}} = -1 - \frac{1}{n} \sum_{i \in \mathcal{I}_1} G_{ii}^{(\nu)} + Z_{\nu} = -1 - \gamma_n m + \varepsilon_{\nu}, \tag{D.27}$$

where we recall (D.7), and

$$\varepsilon_i := Z_i + \sigma_i r_1 \left(m_2 - m_2^{(i)} \right) + r_2 (m_3 - m_3^{(i)}), \quad \varepsilon_\mu := \begin{cases} Z_\mu + \gamma_n (m_1 - m_1^{(\mu)}), & \text{if } \mu \in \mathcal{I}_2 \\ Z_\mu + \gamma_n (m - m^{(\mu)}), & \text{if } \mu \in \mathcal{I}_3 \end{cases}.$$

By (D.11) we can bound that

$$|m_2 - m_2^{(i)}| \le \frac{1}{n_1} \sum_{\mu \in \mathcal{T}_2} \left| \frac{G_{\mu i} G_{i\mu}}{G_{ii}} \right| \prec n^{-1},$$

where we used (D.19) in the second step. Similarly, we can get that

$$|m - m^{(\mu)}| + |m_1 - m_1^{(\mu)}| + |m_2 - m_2^{(i)}| + |m_3 - m_3^{(i)}| \prec n^{-1}$$
 (D.28)

for any $i \in \mathcal{I}_1$ and $\mu \in \mathcal{I}_2 \cup \mathcal{I}_3$. Together with (D.19), we obtain that for all i and μ ,

$$|\varepsilon_i| + |\varepsilon_\mu| < n^{-1/2}. \tag{D.29}$$

With (D.22) and the definition of Ξ , we get that $\mathbf{1}(\Xi)|z+\sigma_i^2r_1m_2+r_2m_3|\sim 1$. Hence using (D.25), (D.29) and (D.19), we obtain that

$$\mathbf{1}(\Xi)G_{ii} = \mathbf{1}(\Xi) \left[-\frac{1}{z + \sigma_i^2 r_1 m_2 + r_2 m_3} + \mathcal{O}_{\prec} \left(n^{-1/2} \right) \right]. \tag{D.30}$$

Plugging it into the definitions of m and m_1 in (D.7), we get

$$\mathbf{1}(\Xi)m = \mathbf{1}(\Xi) \left[-\frac{1}{p} \sum_{i \in \mathcal{I}_1} \frac{1}{z + \sigma_i^2 r_1 m_2 + r_2 m_3} + \mathcal{O}_{\prec} \left(n^{-1/2} \right) \right], \tag{D.31}$$

$$\mathbf{1}(\Xi)m_1 = \mathbf{1}(\Xi) \left[-\frac{1}{p} \sum_{i \in \mathcal{I}_1} \frac{\sigma_i^2}{z + \sigma_i^2 r_1 m_2 + r_2 m_3} + \mathcal{O}_{\prec} \left(n^{-1/2} \right) \right]. \tag{D.32}$$

As a byproduct, we obtain from the two estimates that

$$1(\Xi) (|m - m_c| + |m_1 - m_{1c}|) \lesssim (\log n)^{-1/2}$$
, with high probability. (D.33)

Together with (D.22), we get that

$$|1 + \gamma_n m_1| \sim 1$$
, $|1 + \gamma_n m| \sim 1$, with high probability on Ξ . (D.34)

Now using (D.26), (D.27), (D.29), (D.19) and (D.34), we can obtain that with high probability,

$$\mathbf{1}(\Xi)G_{\mu\mu} = \mathbf{1}(\Xi) \left[-\frac{1}{1 + \gamma_n m_1} + \mathcal{O}_{\prec} \left(n^{-1/2} \right) \right], \quad \mu \in \mathcal{I}_2,$$
 (D.35)

$$\mathbf{1}(\Xi)G_{\nu\nu} = \mathbf{1}(\Xi) \left[-\frac{1}{1 + \gamma_n m} + \mathcal{O}_{\prec} \left(n^{-1/2} \right) \right], \quad \nu \in \mathcal{I}_3.$$
 (D.36)

Taking average over $\mu \in \mathcal{I}_2$ and $\nu \in \mathcal{I}_3$, we get that with high probability,

$$\mathbf{1}(\Xi)m_{2} = \mathbf{1}(\Xi) \left[-\frac{1}{1 + \gamma_{n}m_{1}} + O_{\prec} \left(n^{-1/2} \right) \right], \quad \mathbf{1}(\Xi)m_{3} = \mathbf{1}(\Xi) \left[-\frac{1}{1 + \gamma_{n}m} + O_{\prec} \left(n^{-1/2} \right) \right], \quad (D.37)$$

which further implies

$$\mathbf{1}(\Xi)\left(\frac{1}{m_2} + 1 + \gamma_n m_1\right) \prec n^{-1/2}, \quad \mathbf{1}(\Xi)\left(\frac{1}{m_3} + 1 + \gamma_n m\right) \prec n^{-1/2}.$$
 (D.38)

Finally, plugging (D.31) and (D.32) into (D.38), we conclude (D.24).

Step 3: Ξ holds with high probability. In this step, we show that the event $\Xi(z)$ in fact holds with high probability for all $z \in \mathbf{D}$. Once we have proved this fact, then applying Lemma C.6 to (D.24) immediately shows that $(m_2(z), m_3(z))$ is equal to $(m_{2c}(z), m_{3c}(z))$ up to an error of order $n^{-1/2}$.

First we claim that it suffices to show that

$$|m_2(0) - m_{2c}(0)| + |m_3(0) - m_{3c}(0)| < n^{-1/2}.$$
 (D.39)

Once we know (D.39), then by (C.22) and (D.5), we know $\max_{\alpha=2}^3 |m_{\alpha c}(z) - m_{\alpha c}(0)| = O((\log n)^{-1})$ and $\max_{\alpha=2}^3 |m_{\alpha}(z) - m_{\alpha}(0)| = O((\log n)^{-1})$ with high probability for $z \in \mathbf{D}$. Together with (D.39), we obtain that

$$|m_2(z) - m_{2c}(z)| + |m_3(z) - m_{3c}(z)| \lesssim (\log n)^{-1}$$
 with high probability. (D.40)

and

$$\sup_{z \in \mathbf{D}} (|m_2(z) - m_{2c}(0)| + |m_3(z) - m_{3c}(0)|) \lesssim (\log n)^{-1} \quad \text{with high probability.} \tag{D.41}$$

The condition (D.40) shows that Ξ holds with high probability, and the condition (D.41) verifies the condition (C.21) of Lemma C.6. Hence applying Lemma C.6 to (D.24), we obtain that

$$|m_2(z) - m_{2c}(z)| + |m_3(z) - m_{3c}(z)| < n^{-1/2}$$
 (D.42)

for all $z \in D$. Plugging (D.42) into (D.25)-(D.27), we get the diagonal estimate

$$\max_{\mathfrak{a}\in\mathcal{I}}|G_{\mathfrak{a}\mathfrak{a}}(z)-\Pi_{\mathfrak{a}\mathfrak{a}}(z)| \prec n^{-1/2}.$$
 (D.43)

Together with the off-diagonal estimate in (D.19), we conclude (D.15).

Lemma D.10. Under the assumptions in Proposition D.1, the estimate (D.39) holds.

Proof. By (C.15), we get

$$m(0) = \frac{1}{p} \sum_{i \in \mathcal{I}_1} G_{ii}(0) = \frac{1}{p} \sum_{k=1}^p \frac{|\xi_k(i)|^2}{\lambda_k} \geqslant \lambda_1^{-1} \gtrsim 1.$$

Similarly, we can also get that $m_1(0)$ is positive and has size $m_1(0) \sim 1$. Hence we have

$$1 + \gamma_n m_1(0) \sim 1$$
, $1 + \gamma_n m_1(0) \sim 1$.

Together with (D.26), (D.27) and (D.29), we obtain that (D.37) and (D.38) hold at z=0 even without the indicator function $\mathbf{1}(\Xi)$. Furthermore, it gives that

$$\left|\sigma_i^2 r_1 m_2(0) + r_2 m_3(0)\right| = \left|\frac{\sigma_i^2 r_1}{1 + \gamma_n m_1(0)} + \frac{r_2}{1 + \gamma_n m(0)} + \mathcal{O}_{\prec}(n^{-1/2})\right| \sim 1$$

with high probability. Then using (D.25) and (D.29), we obtain that (D.31) and (D.32) hold at z=0 even without the indicator function $\mathbf{1}(\Xi)$. Finally, plugging (D.31) and (D.32) into (D.38), we conclude (D.24) holds at z=0, that is,

$$\left| \frac{1}{m_2(0)} + 1 - \frac{1}{n} \sum_{i=1}^p \frac{\sigma_i^2}{\sigma_i^2 r_1 m_2(0) + r_2 m_3(0)} \right| < n^{-1/2},
\left| \frac{1}{m_3(0)} + 1 - \frac{1}{n} \sum_{i=1}^p \frac{1}{\sigma_i^2 r_2 m_2(0) + r_2 m_3(0)} \right| < n^{-1/2}.$$
(D.44)

Denoting $\omega_2 = -m_{2c}(0)$ and $\omega_3 = -m_{2c}(0)$. By (D.38), we have

$$\omega_2 = \frac{1}{1 + \gamma_n m_1(0)} + \mathcal{O}_{\prec}(n^{-1/2}), \quad \omega_3 = \frac{1}{1 + \gamma_n m(0)} + \mathcal{O}_{\prec}(n^{-1/2}).$$

Hence there exists a sufficiently small constant c > 0 such that

$$c \le \omega_2 \le 1$$
, $c \le \omega_3 \le 1$, with high probability. (D.45)

Moreover, one can verify from (D.44) that (ω_2, ω_3) satisfy approximately the same equations as in (C.19):

$$r_1\omega_2 + r_2\omega_3 = 1 - \gamma_n + O_{\prec}(n^{-1/2}), \quad f(\omega_2) = 1 + O_{\prec}(n^{-1/2}).$$
 (D.46)

The first equation and (D.45) together implies that $\omega_2 \in [0, r_1^{-1}(1-\gamma_n)]$ with high probability. Since f is strictly increasing and has bounded derivatives on $[0, r_1^{-1}(1-\gamma_n)]$, by basic calculus the second equation in (D.46) gives that $|\omega_2 - b_2| \prec n^{-1/2}$. Together with the first equation in (D.46), we get $|\omega_3 - b_3| \prec n^{-1/2}$. This concludes (D.39).

This lemma concludes (D.39), and as explained above, concludes the proof of Lemma D.7. \Box

With Lemma D.7, we can conclude the proof of Proposition D.1.

Proof of Proposition D.1. With (D.15), one can repeat the polynomialization method in [29, Section 5] to get the anisotropic local law (C.26) for G_0 . The proof is exactly the same, except for some minor notation difference, so we omit the details.

D.3 Anisotropic local law

In this section, we finish the proof of Theorem C.7 for a general X satisfying the bounded support condition (C.8) with $q \leqslant n^{-\phi}$ for some constant $\phi > 0$. The proposition D.1 implies that (C.26) holds for Gaussian Z_1^{Gauss} and Z_2^{Gauss} . Thus the basic idea is to prove that for Z_1 and Z_2 satisfying the assumptions in Theorem C.7,

$$\mathbf{u}^\top \left(G(Z,z) - G(Z^{Gauss},z) \right) \mathbf{v} \prec q$$

for any deterministic unit vectors $\mathbf{u}, \mathbf{v} \in \mathbb{C}^{\mathcal{I}}$ and $z \in \mathbf{D}$. Here we abbreviated $Z := \begin{pmatrix} Z_1 \\ Z_2 \end{pmatrix}$ and

 $Z^{Gauss}:=\begin{pmatrix} Z_1^{Gauss} \\ Z_2^{Gauss} \end{pmatrix}$. We prove the above statement using a continuous comparison argument introduced in [34]. The proof is similar to the ones in Sections 7-8 of [34], so we only give an outline without writing down all the details.

Definition D.11 (Interpolating matrices). We denote Introduce the notations $Z^0 := Z^{Gauss}$ and $Z^1 := Z$. Let $\rho^0_{\mu i}$ and $\rho^1_{\mu i}$ be the laws of $Z^0_{\mu i}$ and $Z^1_{\mu i}$, respectively. For $\theta \in [0,1]$, we define the interpolated law

$$\rho_{\mu i}^{\theta} := (1 - \theta)\rho_{\mu i}^{0} + \theta \rho_{\mu i}^{1}.$$

We shall work on the probability space consisting of triples (Z^0, Z^θ, Z^1) of independent $n \times p$ random matrices, where the matrix $Z^\theta = (Z^\theta_{ui})$ has law

$$\prod_{i \in \mathcal{I}_1} \prod_{\mu \in \mathcal{I}_2 \cup \mathcal{I}_3} \rho_{\mu i}^{\theta} (\mathrm{d} Z_{\mu i}^{\theta}). \tag{D.47}$$

For $\lambda \in \mathbb{R}$, $i \in \mathcal{I}_1$ and $\mu \in \mathcal{I}_2 \cup \mathcal{I}_3$, we define the matrix $Z_{(\mu i)}^{\theta,\lambda}$ through

$$\left(Z_{(\mu i)}^{\theta,\lambda}\right)_{\nu j} := \begin{cases} Z_{\mu i}^{\theta}, & \text{if } (j,\nu) \neq (i,\mu) \\ \lambda, & \text{if } (j,\nu) = (i,\mu) \end{cases}.$$

We also introduce the matrices

$$G^{\theta}(z) := G\left(Z^{\theta}, z\right), \quad G^{\theta, \lambda}_{(\mu i)}(z) := G\left(Z^{\theta, \lambda}_{(\mu i)}, z\right).$$

We shall prove (C.26) through interpolation matrices Z^{θ} between Z^{0} and Z^{1} . We have see that (C.26) holds for Z^{0} by Proposition D.1. Using (D.47) and fundamental calculus, we get the following basic interpolation formula: for $F: \mathbb{R}^{n \times p} \to \mathbb{C}$,

$$\frac{\mathrm{d}}{\mathrm{d}\theta} \mathbb{E}F(Z^{\theta}) = \sum_{i \in \mathcal{I}_1} \sum_{\mu \in \mathcal{I}_2 \cup \mathcal{I}_3} \left[\mathbb{E}F\left(Z_{(\mu i)}^{\theta, Z_{\mu i}^1}\right) - \mathbb{E}F\left(Z_{(\mu i)}^{\theta, Z_{\mu i}^0}\right) \right] \tag{D.48}$$

provided all the expectations exist.

We shall apply (D.48) to $F(Z) := F_{\mathbf{u}\mathbf{v}}^p(Z,z)$ for (large) $p \in 2\mathbb{N}$ and $F_{\mathbf{v}}(Z,z)$ defined as

$$F_{\mathbf{u}\mathbf{v}}(Z,z) := |G_{\mathbf{u}\mathbf{v}}(Z,z) - \Pi_{\mathbf{u}\mathbf{v}}(z)|.$$
 (D.49)

Here for simplicity of notations, we introduce the following notation of generalized entries: for $\mathbf{u}, \mathbf{v} \in \mathbb{R}^{\mathcal{I}}$, we shall denote $G_{\mathbf{u}\mathbf{v}} := \mathbf{u}^{\top}G\mathbf{v}$. Moreover, we shall abbreviate $G_{\mathbf{u}\mathfrak{a}} := G_{\mathbf{u}\mathfrak{e}_{\mathfrak{a}}}$ for $\mathfrak{a} \in \mathcal{I}$, where $\mathbf{e}_{\mathfrak{a}}$ is the standard unit vector along \mathfrak{a} -th axis. Given any vector $\mathbf{u} \in \mathbb{R}^{\mathcal{I}_{1,2,3}}$, we always identify it with its natural embedding in $\mathbb{R}^{\mathcal{I}}$. The exact meanings will be clear from the context. The main work is to show the following self-consistent estimate for the right-hand side of (D.48) for any fixed $p \in 2\mathbb{N}$ and constant $\varepsilon > 0$:

$$\sum_{i \in \mathcal{I}_1} \sum_{\mu \in \mathcal{I}_2 \cup \mathcal{I}_3} \left[\mathbb{E} F_{\mathbf{u}\mathbf{v}}^p \left(Z_{(\mu i)}^{\theta, Z_{\mu i}^1}, z \right) - \mathbb{E} F_{\mathbf{u}\mathbf{v}}^p \left(Z_{(\mu i)}^{\theta, Z_{\mu i}^0}, z \right) \right] = \mathcal{O} \left((n^{\varepsilon} q)^p + \mathbb{E} F_{\mathbf{u}\mathbf{v}}^p \left(Z^{\theta}, z \right) \right)$$
 (D.50)

for all $\theta \in [0, 1]$. If (D.50) holds, then combining (D.48) with a Grönwall's argument we obtain that for any fixed $p \in 2\mathbb{N}$ and constant $\varepsilon > 0$:

$$\mathbb{E}\left|G_{\mathbf{u}\mathbf{v}}(Z^{1},z) - \Pi_{\mathbf{u}\mathbf{v}}(z)\right|^{p} \leqslant (n^{\varepsilon}q)^{p}$$

Together with Markov's inequality, we conclude (C.26). In order to prove (D.50), we compare $Z_{(\mu i)}^{\theta, Z_{\mu i}^0}$ and $Z_{(\mu i)}^{\theta, Z_{\mu i}^1}$ via a common $Z_{(\mu i)}^{\theta, 0}$, i.e. we will prove that

$$\sum_{i \in \mathcal{I}_1} \sum_{\mu \in \mathcal{I}_2 \cup \mathcal{I}_3} \left[\mathbb{E} F_{\mathbf{u} \mathbf{v}}^p \left(Z_{(\mu i)}^{\theta, Z_{\mu i}^a}, z \right) - \mathbb{E} F_{\mathbf{v}}^p \left(Z_{(\mu i)}^{\theta, 0}, z \right) \right] = \mathcal{O} \left((n^{\varepsilon} q)^p + \mathbb{E} F_{\mathbf{u} \mathbf{v}}^p \left(Z^{\theta}, z \right) \right) \quad \text{(D.51)}$$

for all $a \in \{0, 1\}$ and $\theta \in [0, 1]$. Underlying the proof of (D.51) is an expansion approach. We define the $\mathcal{I} \times \mathcal{I}$ matrix $\Delta_{(ui)}^{\lambda}$ as

$$\Delta_{(\mu i)}^{\lambda} := \lambda \begin{pmatrix} 0 & \mathbf{u}_i^{(\mu)} \mathbf{e}_{\mu}^{\mathsf{T}} \\ \mathbf{e}_{\mu} (\mathbf{u}_i^{(\mu)})^{\mathsf{T}} & 0 \end{pmatrix}, \tag{D.52}$$

where we denote $\mathbf{u}_i^{(\mu)} := \Lambda U \mathbf{e}_i$ if $\mu \in \mathcal{I}_2$ and $\mathbf{u}_i^{(\mu)} := V \mathbf{e}_i$ if $\mu \in \mathcal{I}_3$. Then by the definition of H in (C.11)), we have for any $\lambda, \lambda' \in \mathbb{R}$ and $K \in \mathbb{N}$,

$$G_{(i\mu)}^{\theta,\lambda'} = G_{(\mu i)}^{\theta,\lambda} + \sum_{k=1}^{K} G_{(\mu i)}^{\theta,\lambda} \left(\Delta_{(\mu i)}^{\lambda-\lambda'} G_{(\mu i)}^{\theta,\lambda} \right)^{k} + G_{(\mu i)}^{\theta,\lambda'} \left(\Delta_{(\mu i)}^{\lambda-\lambda'} G_{(\mu i)}^{\theta,\lambda} \right)^{K+1}. \tag{D.53}$$

Using this expansion and the a priori bound (D.4), it is easy to prove the following estimate: if y is a random variable satisfying $|y| \prec q$, then

$$G_{(\mu i)}^{\theta,y} = \mathcal{O}(1), \quad i \in \mathcal{I}_1, \ \mu \in \mathcal{I}_2 \cup \mathcal{I}_3,$$
 (D.54)

with high probability.

In the following proof, for simplicity of notations, we introduce $f_{(\mu i)}(\lambda) := F_{\mathbf{v}}^p(Z_{(\mu i)}^{\theta,\lambda})$. We use $f_{(\mu i)}^{(r)}$ to denote the r-th derivative of $f_{(\mu i)}$. By (D.54), it is easy to see that for any fixed $r \in \mathbb{N}$, $f_{(\mu i)}^{(r)}(y) = \mathrm{O}(1)$ with high probability for any random variable y satisfying $|y| \prec q$. Then the Taylor expansion of $f_{(\mu i)}$ gives

$$f_{(\mu i)}(y) = \sum_{r=0}^{p+4} \frac{y^r}{r!} f_{(\mu i)}^{(r)}(0) + \mathcal{O}_{\prec} (q^{p+4}), \qquad (D.55)$$

Therefore we have for $a \in \{0, 1\}$,

$$\mathbb{E}F_{\mathbf{v}}^{p}\left(Z_{(\mu i)}^{\theta,Z_{\mu i}^{a}}\right) - \mathbb{E}F_{\mathbf{v}}^{p}\left(Z_{(\mu i)}^{\theta,0}\right) = \mathbb{E}\left[f_{(\mu i)}\left(Z_{i\mu}^{a}\right) - f_{(\mu i)}(0)\right]$$

$$= \mathbb{E}f_{(\mu i)}(0) + \frac{1}{2n}\,\mathbb{E}f_{(\mu i)}^{(2)}(0) + \sum_{r=4}^{p+4}\frac{1}{r!}\,\mathbb{E}f_{(\mu i)}^{(r)}(0)\,\mathbb{E}\left(Z_{i\mu}^{a}\right)^{r} + \mathcal{O}_{\prec}(q^{p+4}). \tag{D.56}$$

Here to illustrate the idea in a more concise way, we assume the extra condition

$$\mathbb{E}(Z_{\mu i}^1)^3 = 0, \quad 1 \leqslant \mu \leqslant n, \quad 1 \leqslant i \leqslant p. \tag{D.57}$$

Hence the r=3 term in the Taylor expansion vanishes. However, this is not necessary as we will explain at the end of the proof.

By (C.2) and the bounded support condition, we have

$$\left| \mathbb{E} \left(Z_{i\mu}^a \right)^r \right| \prec n^{-2} q^{r-4}, \quad r \geqslant 4. \tag{D.58}$$

Thus to show (D.51), we only need to prove for r = 4, 5, ..., p + 4,

$$n^{-2}q^{r-4}\sum_{i\in\mathcal{I}_1}\sum_{\mu\in\mathcal{I}_2\cup\mathcal{I}_3}\left|\mathbb{E}\,f_{(\mu i)}^{(r)}(0)\right| = \mathcal{O}\left(\left(n^{\varepsilon}q\right)^p + \mathbb{E}F_{\mathbf{u}\mathbf{v}}^p(Z^{\theta},z)\right). \tag{D.59}$$

In order to get a self-consistent estimate in terms of the matrix Z^{θ} on the right-hand side of (D.59), we want to replace $Z^{\theta,0}_{(\mu i)}$ in $f_{(\mu i)}(0) = F^p_{\mathbf{u}\,\mathbf{v}}(Z^{\theta,0}_{(\mu i)})$ with $Z^{\theta} = Z^{\theta,Z^{\theta}_{\mu i}}_{(\mu i)}$.

Lemma D.12. Suppose that

$$n^{-2}q^{r-4} \sum_{i \in \mathcal{I}_1} \sum_{\mu \in \mathcal{I}_2 \cup \mathcal{I}_3} \left| \mathbb{E} f_{(\mu i)}^{(r)}(Z_{i\mu}^{\theta}) \right| = \mathcal{O}\left((n^{\varepsilon}q)^p + \mathbb{E} F_{\mathbf{v}}^p(X^{\theta}, z) \right)$$
(D.60)

holds for r = 4, ..., 4p + 4. Then (D.59) holds for r = 4, ..., 4p + 4.

Proof. The proof is the same as the one for [34, Lemma 7.16].

What remains now is to prove (D.60). For simplicity of notations, we shall abbreviate $Z^{\theta} \equiv Z$. For any $k \in \mathbb{N}$, we denote

$$A_{\mu i}(k) := \left(\frac{\partial}{\partial Z_{\mu i}}\right)^k \left(G_{\mathbf{u}\mathbf{v}} - \Pi_{\mathbf{u}\mathbf{v}}\right).$$

The derivative on the right-hand side can be calculated using the expansion (D.53). In particular, it is easy to verify that it satisfies the following bound

$$|A_{\mu i}(k)| \prec \begin{cases} (\mathcal{R}_i^{(\mu)})^2 + \mathcal{R}_{\mu}^2, & \text{if } k \geqslant 2\\ \mathcal{R}_i^{(\mu)} \mathcal{R}_{\mu}, & \text{if } k = 1 \end{cases}$$
 (D.61)

where for $i \in \mathcal{I}_1$ and $\mu \in \mathcal{I}_2 \cup \mathcal{I}_3$, we denote

$$\mathcal{R}_i^{(\mu)} := |G_{\mathbf{u}\mathbf{u}_i^{(\mu)}}| + |G_{\mathbf{v}\mathbf{u}_i^{(\mu)}}|, \quad \mathcal{R}_{\mu} := |G_{\mathbf{u}\mu}| + |G_{\mathbf{v}\mu}|. \tag{D.62}$$

Then we can calculate the derivative

$$\left(\frac{\partial}{\partial Z_{\mu i}}\right)^r F_{\mathbf{u}\mathbf{v}}^p(Z) = \sum_{k_1 + \dots + k_p = r} \prod_{t=1}^{p/2} \left(A_{\mu i}(k_t) \overline{A_{\mu i}(k_{t+p/2})} \right).$$

Then to prove (D.60), it suffices to show that

$$n^{-2}q^{r-4} \sum_{i \in \mathcal{I}_1} \sum_{\mu \in \mathcal{I}_2 \cup \mathcal{I}_3} \left| \mathbb{E} \prod_{t=1}^{p/2} A_{\mu i}(k_t) \overline{A_{\mu i}(k_{t+p/2})} \right| = O\left((n^{\varepsilon}q)^p + \mathbb{E} F_{\mathbf{u}\mathbf{v}}^p(Z, z) \right)$$
(D.63)

for $4 \leqslant r \leqslant p+4$ and $(k_1, \cdots, k_p) \in \mathbb{N}^p$ satisfying $k_1 + \cdots + k_p = r$. Treating zero k's separately (note $A_{\mu i}(0) = (G_{\mathbf{uv}} - \Pi_{\mathbf{uv}})$ by definition), we find that it suffices to prove

$$n^{-2}q^{r-4}\sum_{i\in\mathcal{I}_1}\sum_{\mu\in\mathcal{I}_2\cup\mathcal{I}_3}\mathbb{E}|A_{\mu i}(0)|^{p-l}\prod_{t=1}^l|A_{\mu i}(k_t)| = O\left((n^{\varepsilon}q)^p + \mathbb{E}F_{\mathbf{u}\mathbf{v}}^p(Z,z)\right)$$
(D.64)

for $4 \leqslant r \leqslant p+4$ and $1 \leqslant l \leqslant p$. Here without loss of generality, we assume that $k_t=0$ for $l+1 \leqslant t \leqslant p$, and $\sum_{t=1}^{l} k_t = r$ with $k_t \geqslant 1$ for $t \leqslant l$.

Now we first consider the case $r \le 2l - 2$. Then by pigeonhole principle, there exist at least two k_t 's with $k_t = 1$. Therefore by (D.61) we have

$$\prod_{t=1}^{l} |A_{\mu i}(k_t)| < \mathbf{1}(r \geqslant 2l - 1) \left[(\mathcal{R}_i^{(\mu)})^2 + \mathcal{R}_{\mu}^2 \right] + \mathbf{1}(r \leqslant 2l - 2) (\mathcal{R}_i^{(\mu)})^2 \mathcal{R}_{\mu}^2.$$
 (D.65)

Using (D.4) and a similar argument as in (D.20), we get that

$$\sum_{i \in \mathcal{I}_1} (\mathcal{R}_i^{(\mu)})^2 = \mathcal{O}(1), \quad \sum_{\mu \in \mathcal{I}_2 \cup \mathcal{I}_3} \mathcal{R}_{\mu}^2 = \mathcal{O}(1), \quad \text{with high probability.} \tag{D.66}$$

Using (D.66) and $n^{-1/2} \leq q$, we get that

$$n^{-2}q^{r-4} \sum_{i \in \mathcal{I}_1} \sum_{\mu \in \mathcal{I}_2 \cup \mathcal{I}_3} |A_{\mu i}(0)|^{p-l} \prod_{t=1}^l |A_{\mu i}(k_t)| \prec q^{r-4} F_{\mathbf{u} \mathbf{v}}^{p-l}(Z) \left[\mathbf{1}(r \geqslant 2l-1)n^{-1} + \mathbf{1}(r \leqslant 2l-2)n^{-2} \right]$$

$$\leqslant F_{\mathbf{u} \mathbf{v}}^{p-l}(Z) \left[\mathbf{1}(r \geqslant 2l-1)q^{r-2} + \mathbf{1}(r \leqslant 2l-2)q^r \right].$$

If $r\leqslant 2l-2$, then we get $q^r\leqslant q^l$ using the trivial inequality $r\geqslant l$. On the other hand, if $r\geqslant 4$ and $r\geqslant 2l-1$, then $r\geqslant l+2$ and we get $q^r\leqslant q^{l+2}$. Therefore we conclude that

$$n^{-2}q^{r-4}\sum_{i\in\mathcal{I}_1}\sum_{\mu\in\mathcal{I}_2\cup\mathcal{I}_3}|A_{\mu i}(0)|^{p-l}\prod_{t=1}^l|A_{\mu i}(k_t)|\prec F_{\mathbf{u}\mathbf{v}}^{p-l}(Z)q^l.$$

Now (D.64) follows from Hölder's inequality. This concludes the proof of (D.60), and hence of (D.51), and hence of (C.26).

Finally, if the condition (D.57) does not hold, then there is also an r=3 term in the Taylor expansion (D.56):

$$\frac{1}{6} \mathbb{E} f_{(\mu i)}^{(3)}(0) \mathbb{E} \left(Z_{i\mu}^a \right)^3.$$

Note that $\mathbb{E}\left(Z_{i\mu}^a\right)^3$ is of order $n^{-3/2}$, while the sum over i and μ in (D.51) provides a factor n^2 . In fact, $\mathbb{E}\left(f_{(\mu i)}^{(3)}(0)\right)$ will provide an extra $n^{-1/2}$ to compensate the remaining $n^{1/2}$ factor. This follows from an improved self-consistent comparison argument for sample covariance matrices in [34, Section 8]. The argument for our case is almost the same except for some notational differences, so we omit the details.

D.4 Proof of Lemma C.5 and Lemma C.6

Finally, we give the proof of Lemma C.5 and Lemma C.6 using the contraction principle.

Proof of Lemma C.5. One can check that the equations in (C.16) are equivalent to the following ones:

$$r_1 m_{2c} = -(1 - \gamma_n) - r_2 m_{3c} - z \left(\frac{1}{m_{3c}} + 1\right), \quad g_z(m_{3c}(z)) = 1,$$
 (D.67)

where

$$g_z(m_{3c}) := -m_{3c} + \frac{1}{n} \sum_{i=1}^p \frac{m_{3c}}{z - \sigma_i^2 (1 - \gamma_n) + (1 - \sigma_i^2) r_2 m_{3c} - \sigma_i^2 z \left(m_{3c}^{-1} + 1 \right)}.$$

We first show that there exists a unique solution $m_{3c}(z)$ to the equation $g_z(m_{3c}(z)) = r_2$ under the conditions in (C.21), and the solution satisfies (C.22). Now we abbreviate $\varepsilon(z) := m_{3c}(z) - m_{3c}(0)$, and from (D.67) we can obtain that

$$0 = [g_z(m_{3c}(z)) - g_0(m_{3c}(0)) - g'_z(m_{3c}(0))\varepsilon(z)] + g'_z(m_{3c}(0))\varepsilon(z),$$

which implies

$$g_z'(m_{3c}(0))\varepsilon(z) = -\left[g_z(m_{3c}(0)) - g_0(m_{3c}(0))\right] - \left[g_z(m_{3c}(0) + \varepsilon(z)) - g_z(m_{3c}(0)) - g_z'(m_{3c}(0))\varepsilon(z)\right].$$
(D.68)

Inspired by the above equation, we define iteratively a sequence of vectors $\varepsilon^{(k)} \in \mathbb{C}$ such that $\varepsilon^{(0)} = 0$, and

$$\varepsilon^{(k+1)} = -\frac{g_z(m_{3c}(0)) - g_0(m_{3c}(0))}{g_z'(m_{3c}(0))} - \frac{g_z(m_{3c}(0) + \varepsilon^{(k)}) - g_z(m_{3c}(0)) - g_z'(m_{3c}(0))\varepsilon^{(k)}}{g_z'(m_{3c}(0))}.$$

In other words, the above equation defines a mapping $h: \mathbb{C} \to \mathbb{C}$, which maps $\varepsilon^{(k)}$ to $\varepsilon^{(k+1)} = h(\varepsilon^{(k)})$.

With direct calculation, one can get the derivative

$$g_z'(m_{3c}(0)) = -1 - \frac{1}{n} \sum_{i=1}^p \frac{\sigma_i^2(1 - \gamma_n) - z \left[1 - \sigma_i^2 \left(2r_2 m_{3c}^{-1} + 1\right)\right]}{\left[z - \sigma_i^2(1 - \gamma_n) + (1 - \sigma_i^2)m_{3c} - \sigma_i^2 z \left(r_2 m_{3c}^{-1} + 1\right)\right]^2}.$$

Using (C.20), it is easy to check that there exist constants \widetilde{c} , $\widetilde{C} > 0$ depending only on τ in (C.7) and (C.20) such that

$$|[g_z'(m_{3c}(0))]^{-1}| \leqslant \widetilde{C}, \quad \left| \frac{g_z(m_{3c}(0) + \varepsilon_1) - g_z(m_{3c}(0) + \varepsilon_2) - g_z'(m_{3c}(0))(\varepsilon_1 - \varepsilon_2)}{g_z'(m_{3c}(0))} \right| \leqslant \widetilde{C} |\varepsilon_1 - \varepsilon_2|^2,$$

and

$$\left| \frac{g_z(m_{3c}(0)) - g_0(m_{3c}(0))}{g_z'(m_{3c}(0))} \right| \leqslant \widetilde{C}|z|, \tag{D.71}$$

for all $|z| \le \widetilde{c}$ and $|\varepsilon_1| \le \widetilde{c}$, $|\varepsilon_2| \le \widetilde{c}$. Then with (D.70) and (D.71), it is easy to see that there exists a sufficiently small constant $\delta > 0$ depending only on \widetilde{C} , such that h is a self-mapping

$$h: B_r \to B_r, \quad B_r := \{ \varepsilon \in \mathbb{C} : |\varepsilon| \leqslant r \},$$

as long as $r \leqslant \delta$ and $|z| \leqslant c_{\delta}$ for some constant $c_{\delta} > 0$ depending only on \widetilde{C} and δ . Now it suffices to prove that h restricted to B_r is a contraction, which then implies that $\varepsilon := \lim_{k \to \infty} \varepsilon^{(k)}$ exists and $m_{3c}(0) + \varepsilon$ is a unique solution to the second equation of (D.67) subject to the condition $\|\varepsilon\|_{\infty} \leqslant r$.

From the iteration relation (D.69), using (D.70) one can readily check that

$$\varepsilon^{(k+1)} - \varepsilon^{(k)} = h(\varepsilon^{(k)}) - h(\varepsilon^{(k-1)}) \leqslant \widetilde{C} |\varepsilon^{(k)} - \varepsilon^{(k-1)}|^2. \tag{D.72}$$

Hence as long as r is chosen to be sufficiently small such that $2r\widetilde{C} \leqslant 1/2$, then h is indeed a contraction mapping on B_r , which proves both the existence and uniqueness of the solution $m_{3c}(z) = m_{3c}(0) + \varepsilon$, if we choose c_0 in (C.21) as $c_0 = \min\{c_\delta, r\}$. After obtaining $m_{3c}(z)$, we can then find $m_{2c}(z)$ using the first equation in (D.67).

Note that with (D.71) and $\varepsilon^{(0)} = \mathbf{0}$, we get from (D.69) that

$$|\varepsilon^{(1)}| \leqslant \widetilde{C}|z|.$$

With the contraction mapping, we have the bound

$$|\varepsilon| \leqslant \sum_{k=0}^{\infty} \|\varepsilon^{(k+1)} - \varepsilon^{(k)}\|_{\infty} \leqslant 2\widetilde{C}|z|.$$
 (D.73)

This gives the bound (C.22) for $m_{3c}(z)$. Using the first equation in (D.67), we immediately obtain the bound

$$r_1|m_{2c}(z) - m_{2c}(0)| \le C|z|.$$

This gives (C.22) for $m_{2c}(z)$ as long as if $r_1 \gtrsim 1$. To deal with the small r_1 case, we go back to the first equation in (C.16) and treat $m_{2c}(z)$ as the solution to the following equation:

$$\widetilde{g}_z(m_{2c}(z)) = 1, \quad \widetilde{g}_z(x) := -x + \frac{\gamma_n}{p} \sum_{i=1}^p \frac{\sigma_i^2 x}{z + \sigma_i^2 r_1 x + r_2 m_{3c}(z)}.$$

We can calculate that

$$g_z'(m_{2c}(0)) = -1 + \frac{\gamma_n}{p} \sum_{i=1}^p \frac{\sigma_i^2(z + r_2 m_{3c}(z))}{(z + \sigma_i^2 r_1 m_{2c}(0) + r_2 m_{3c}(z))^2}.$$

At z = 0, we have

$$|g_0'(m_{2c}(0))| = \left| 1 + \frac{\gamma_n}{p} \sum_{i=1}^p \frac{\sigma_i^2 r_2 b_3}{(\sigma_i^2 r_1 b_2 + r_2 b_3)^2} \right| \geqslant 1,$$

where b_2 and b_3 satisfy (C.20). Thus under (C.21) we have $|g_z'(m_{2c}(0))| \sim 1$ as long as c_0 is taken sufficiently small. Then with the above arguments for $m_{3c}(z)$ between (D.67) and (D.73), we can conclude (C.22) for $m_{2c}(z)$. This concludes the proof of Lemma C.5.

Proof of Lemma C.6. Under (C.21), we can obtain equation (D.67) approximately up to some small error

$$r_1 m_{2c} = -(1 - \gamma_n) - r_2 m_{3c} - z \left(\frac{1}{m_{3c}} + 1\right) + \mathcal{E}'_2(z), \quad g_z(m_{3c}(z)) = 1 + \mathcal{E}'_3(z), \quad (D.74)$$

with $|\mathcal{E}_2'(z)| + |\mathcal{E}_3'(z)| = O(\delta(z))$. Then we subtract the equations (D.67) from (D.74), and consider the contraction principle for the functions $\varepsilon(z) := m_3(z) - m_{3c}(z)$. The rest of the proof is exactly the one for Lemma C.5, so we omit the details.