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# Estimation and Inference for Varying-coefficient Models with Nonstationary Regressors using Penalized Splines

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Estimation and Inference for Varying-coefficient Models with

Nonstationary Regressors using Penalized Splines \*

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Abstract

This paper considers estimation and inference for varying-coefficient models with nonstation-

ary regressors. We propose a nonparametric estimation method using penalized splines, which

achieves the same optimal convergence rate as kernel-based methods, but enjoys computation

advantages. Utilizing the mixed model representation of penalized splines, we develop a likeli-

hood ratio test statistic for checking the stability of the regression coefficients. We derive both

the exact and the asymptotic null distributions of this test statistic. We also demonstrate its op-

timality by examining its local power performance. These theoretical findings are well supported

by simulation studies.

JEL Classification: C12, C14, C22

Keywords: Nonstationary Time Series; Varying-coefficient Model; Likelihood Ratio Test;

Penalized Splines

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## 1 Introduction

Regression models with nonstationary regressors have received much attention in the literature of theoretical and applied econometrics since the seminal work by Nelson and Plosser (1982). Amongst the popular research in this field has been the study of cointegration. The traditional framework of Engle and Granger (1987), assuming constant cointegrating coefficients, provides an appealing analytical framework to characterize the long-run equilibrium relationship. However, very few empirical studies support the presence of cointegration with constant coefficients. Such an empirical frustration is due to the lack of flexibility of traditional cointegration models in accommodating the time-varying long-run equilibrium relationship.

There are many empirical examples in economics and finance demonstrating time-varying features in cointegrating relationships. For example, Goldfajn and Baig (1998) argue that, during the Asian currency crisis, the cointegrating relation between the spot exchange rate and the interest rate differential is not fixed but depends on the level of the interest rate. Another example is in the literature of stock return predictability, where one of the theoretical and practical issues is to answer whether we could predict the asset return from fundamental variables such as the dividend-price ratio and the earning-price ratio, which are well known nonstationary time series variables (Campbell and Yogo 2006). Although linear prediction models have been extensively used, Lettau and Ludvigsson (2001), Goyal and Welch (2003) and Paye and Timmermann (2006) find empirical evidence that the cointegrating stock return forecasting models might be instable.

Many studies adopt nonparametric methods to capture the time-varying relation with nonstationary variables. The latest research include Wang and Phillips (2009a, 2009b) and Cai, Li and Park (2009), among others. Wang and Phillips (2009b) construct asymptotic theories for the local time density estimation in nonparametric cointegrating regression. They find that the so called ill-posed inverse problem does not exist in nonparametric regression

with nonstationary endogenous regressors. Cai et al. (2009) investigate the asymptotic property of local linear estimators in a varying-coefficient model when the smoothing variable is nonstationary but the covariates are either stationary.

Testing the stability of varying coefficients becomes another important issue in this literature. For example, Park and Hahn (1999) construct two residual-based statistics to test the constancy of the cointegrating coefficients based on the series estimation. Kasparis (2008) develops two residual-based statistics for testing the functional form misspecification in cointegrating relations. Bierens and Martins (2010) propose a vector error correction model with cointegration coefficients estimated by Chebyshev polynomials, and conduct a likelihood ratio test on the stability by examining whether all Chebyshev polynomial coefficients are jointly zero. Juhl (2005) and Xiao (2009) further extend the studies to the case where cointegration coefficients are general smooth stochastic functions depending on some stationary covariates. Xiao (2009) considers both kernel and local polynomial estimators and establish their asymptotics. Moreover, he proposes a test statistic by comparing the functional-coefficient estimates to a fixed value estimated under the null. <sup>1</sup>

In this paper, we consider varying-coefficient regression models with nonstationary regressors. Our model setting is similar to Cai et al. (2009) and Xiao (2009). However, we propose to estimate the varying coefficients by penalized splines and construct a likelihood ratio test (LRT) for the stability of the varying coefficients. The basic idea of spline methods is to approximate the unknown regression function by splines, which are piecewise polynomials, and then estimate the spline coefficients by the least squares method. In order to maintain a good balance between the goodness of fit and the model variability, a large number of basis func-

<sup>&</sup>lt;sup>1</sup>Other contributions to the literature on the stability tests in cointegrating regression includes Hansen (1992), Hao (1996), Quintos (1997), Kuo (1998), Hansen and Johansen (1999), Johansen, Mosconi and Nielsen (2000), Harris, McCabe and Leybourne (2002), and among others. Furthermore, Hong and Phillips (2010) propose a modified RESET test for testing linearity in the cointegration model. Gao, King, Lu and Tjøstheim (2009) consider a nonparametric specification test for a nonparametric time series model with nonstationary variables.

tions are employed and a penalty term is imposed to avoid overfitting (Eilers and Marx, 1996). There are several prominent features of such a penalized splines approach. First, this method is simple and easy to implement. Its computation is usually less time-consuming compared to other nonparametric methods such as kernel or local polynomials. Second, it could easily incorporate correlation structure to improve the efficiency of estimator or to account for longitudinal and spatial effects (Ruppert, Wand and Carroll 2003). Third, it has close connections with Ridge regression, Bayesian methods and mixed model representation, thus allowing fitting and testing to be conducted through the paradigm of likelihood (Crainiceanu, Ruppert, Claeskens and Wand 2005). However, theoretical explorations of penalized splines were less well developed until recently. Li and Ruppert (2008) establish the asymptotic normality of the penalized splines estimation. Claeskens, Krivobokova and Opsomer (2009) systematically compare the asymptotics of penalized splines, regression splines and smoothing splines. All these studies are under the univariate nonparametric model assumption  $y_t = \theta(z_t) + v_t$  for stationary covariate  $z_t$ 's.

Our studies contribute to the literature through the following aspects. First, we propose the penalized spline estimation method for varying-coefficient models with nonstationary regressors. We establish the consistency as well as the optimal convergence rate of the penalized splines estimators. In our study, the choice of the spline basis is not crucial, but the penalty parameter plays the key role of smoothing. To our best knowledge, this is the first work in establishing the asymptotics of penalized splines estimators for varying-coefficient models with nonstationary regressors. Second, we consider testing the stability of the regression coefficients. By utilizing the mixed model representation of penalized splines, we relate this problem to testing zero variance component of the spline coefficients. We then adopt the likelihood ratio test (LRT) procedure and derive the exact and the asymptotic null distribution. Since the exact null distribution is non-standard, we provide a fast algorithm to simulate its critical values when the sample size is small. By assuming both the sample size and the number of spline functions grow to infinity, we, for the first time, show that

the limiting null distribution of the proposed LRT statistic follows a simple  $\chi^2$  distribution rather than a mixture of  $\chi^2$  distributions. We also study the local power of the proposed LRT by deriving the asymptotic distribution under the local alternative. Simulations show that our method works very well.

The rest of the paper is organized as follows. In Section 2, we introduce the varying-coefficient regression model with nonstationary regressors and discuss some regularity assumptions. The penalized splines estimation of the varying coefficients and its asymptotics are presented in Section 3. In Section 4, we construct the LRT for the stability of the varying coefficients. Both the exact distribution and the asymptotic null distribution are derived. The local power property is examined as well. Simulations are conducted in Section 5 to demonstrate the finite sample performance of our procedure, while Section 6 concludes.

In matters of notations,  $\stackrel{D}{=}$  denotes equality in distribution, =: denotes definition,  $\Rightarrow$  denotes convergence in distribution,  $\stackrel{a.s.}{\to}$  denotes almost sure convergence,  $a \sim b$  denotes that a and b has the same order, [Ts] denotes the integer part of Ts and  $\int$  denotes the integration from 0 to 1.

# 2 The Model and Assumptions

Consider the following varying-coefficient regression model without intercept and time trends

$$y_t = x_t \theta(z_t) + u_t, \tag{1}$$

where  $\theta(\cdot)$  is a smooth function of unknown form,  $y_t$  can be either stationary or nonstationary,  $z_t$  and  $u_t$  are stationary, and  $x_t$  is an integrated process of order one, whose generating mechanism is given by

$$x_t = x_{t-1} + v_t,$$
 for  $t = 1, 2, ...T,$ 

with  $v_t$  being stationary. We set  $x_0 = 0$  to avoid some unnecessary complications in exposition, although  $x_0 = o_{a.s.}(\sqrt{T})$  is sufficient for the asymptotic results.<sup>2</sup>

<sup>&</sup>lt;sup>2</sup>To save notations, we only consider the simple case when  $z_t$  is univariate.

Compared to traditional varying-coefficient models, which usually deals with independent and identically distributed (iid) or stationary time series, Model (1) allows the regressors to be highly persistent variables, such as interest rate, GDP growth rate and unemployment rates. On the other hand, compared to traditional linear cointegration models, which are widely used in the literature to capture the long term equilibrium among highly persistent economic variables, Model (1) affords more flexibility as it allows the relationship to be varying according to some state variable  $z_t$ .

Before describing our estimating and testing procedures, we first discuss some regularity assumptions for our Model (1).

**Assumption 1:** The sequence  $\{v_t\}$  is stationary  $\alpha - mixing$  with  $E(|v_t|^{\rho}) < \infty$  for some  $\rho > 2 + r_1$  with  $0 < r_1 \le 2$  and the mixing coefficients  $\alpha_m$  satisfying  $\sum_{m=1}^{\infty} \alpha_m^{1/(2+r_1)-1/\rho} < \infty$ .

### **Assumption 2:**

i) The error term  $u_t$  is a general linear process satisfying

$$u_t = \sum_{i=0}^{\infty} c_i e_{t-i} = \mathbb{C}_u(L)e_t,$$

where  $\{e_t\}_{t=-\infty}^{\infty}$  are i.i.d  $N(0, \sigma_e^2)$  with  $\sigma_e^2 > 0$  and  $\{c_i\}_{i=0}^{\infty}$  satisfies the summability conditions  $\sum_{i=0}^{\infty} i|c_i| < \infty$  with  $\mathbb{C}_u(1) \neq 0$ .

ii)  $u_t$  is independent of  $z_t$  and  $v_t$ .

### **Assumption 3:**

- i) The sequence  $\{z_t\}$  is strictly stationary, ergodic and  $\alpha mixing$  with mixing coefficients  $\alpha_m$  satisfying  $\sum_{m=1}^{\infty} \alpha_m^{1/2-1/r_2} < \infty$  for some  $r_2 > 2$ .
- ii)  $z_t$  has a marginal density  $f_z(z)$  on a finite support [0,1].  $f_z(z)$  is continuously differentiable and bounded away from 0.
- iii)  $\theta(z)$  belongs to the Sobolev space of the m-th order  $W^m[0,1]$ , i.e.  $\theta(z)$  is (m-1)-th continuously differentiable and  $\int_0^1 \{\theta^{(m)}(z)\}^2 dz < \infty$ .

Assumption 1 and Assumption 2 i) provide sufficient conditions of strong approximations for the partial sum  $(\frac{1}{\sqrt{T}}\sum_{t=1}^{[Ts]}u_t, \frac{1}{\sqrt{T}}\sum_{t=1}^{[Ts]}v_t),$ 

$$\sup_{s \in [0,1]} ||(\frac{1}{\sqrt{T}} \sum_{t=1}^{[Ts]} u_t, \frac{1}{\sqrt{T}} \sum_{t=1}^{[Ts]} v_t) - \{B_u(s), B_v(s)\}|| \stackrel{a.s.}{\to} 0, \text{ as } T \to \infty,$$

where  $\{B_u(s), B_v(s)\}$  are two Brownian motions defined on D[0, 1], the space of cadlag functions defined in the unit interval [0, 1]. Note that the strong approximation is stronger than the multivariate invariance principle, but it is commonly used in the literature of nonlinear regression model with nonstationary regressors, including Park and Hahn (1999), Park and Phillips (2001), Kasparis (2008), Wang and Phillips (2009a, 2009b), Cai et al. (2009), Shi and Phillips (2012) among others. Sufficient conditions to derive strong approximations for dependent random variables are also well established in the literature. For example, Lemma 1 in Park and Hahn (1999) establish conditions of strong approximations for a general linear process  $u_t$  and Theorem 4.1 in Shao and Lu (1987) give conditions of strong approximations for an  $\alpha - mixing$  process  $v_t$ .

Assumption 2 i) also defines  $u_t$  as an invertible Gaussian moving average process. The normality assumption is somewhat restrictive but it is for the purpose of employing likelihood principles. Define the vector  $u = (u_1, ..., u_T)^T$  and denote  $var(u) = \sigma^2 \sum_u$ . At the current point, we assume that  $\sum_u$  is known so that the full likelihood function could be constructed. In practice,  $\sum_u$  can be posited to be of a particular form  $\sum_u(\phi)$ , where  $\phi$  is a vector of parameters that could be estimated from the data. A simple example is to treat  $u_t$  as an AR(1) process. Then  $\sum_u$  is a function of the first order coefficient  $\rho$ . One could apply a two-step procedure to obtain the estimate  $\hat{\rho}$  and replace  $\sum_u$  by its estimate  $\sum_u(\hat{\rho})$ .

To simplify the derivation of the LRT, Assumption 2 ii) assumes the independence condition between the error term  $u_t$  and  $x_t$ , though this might be further relaxed. Following Saikkonen (1991) and Saikkonen and Choi (2004), we might remove the endogeneity between  $u_t$  and  $x_t$  by adding leads and lags of the term  $\{v_t\}_{t=1}^{\infty}$  in the regression. On the other hand, we assume the independence between  $u_t$  and  $z_t$ , which rules out the ill-posed inverse problems

in the nonparametric estimation with stationary smoothing variables.

Assumption 3 i) guarantees that  $z_t$  is strictly stationary and imposes some conditions on its dependency and moments. In this paper, we do not consider the case when  $z_t$  is nonstationary. Assumption 3 ii) requires that  $z_t$  has a bounded support. In practice, one can always conduct some transformations, such as the probability integral transform, to satisfy this assumption. We also assume that the marginal density of z is continuously differentiable and bounded away from 0, thus ensuring there are enough observations for estimation. Finally, Assumption 3 iii) imposes some smoothness conditions on the unknown function  $\theta(z)$ , which is a standard assumption in nonparametric regression analysis.

# 3 Penalized Splines Estimation

We employ the penalized splines approach to estimate the varying-coefficient regression model with nonstationary regressors. First, we approximate the unknown varying coefficient by splines basis. A popular choice is the uniform B-splines family defined by a set of equally spaced knots  $\kappa_k = k/K$ , for  $k = 0, \dots, K$ . The simplest case is the zero degree B-splines, which are indicator functions between  $\kappa_{k-1}$  and  $\kappa_k$ . In general, we could use the iterative algorithm proposed by de Boor (1978) to calculate the p-th degree B-splines and express  $\theta(z)$  as

$$\theta(z) = \sum_{k=1}^{K+p} \Psi_k^{[p]}(z) \beta_k + O(K^{-1}).$$

Following the idea of penalized least squares, we could estimate the spline coefficients  $\beta$  by minimizing the following criterion:

$$\sum_{t=1}^{T} \left( y_t - x_t \sum_{k=1}^{K+p} \beta_k \Psi_k^{[p]}(z) \right)^2 + \tilde{\lambda}^{-1} \int \{ \theta^{(m)}(z) \}^2 dz.$$
 (2)

In a discrete version, this could be written as

$$\sum_{t=1}^{T} \left( y_t - x_t \sum_{k=1}^{K+p} \beta_k \Psi_k^{[p]}(z) \right)^2 + \tilde{\lambda}^{-1} K^{-1} \sum_{k=m+1}^{K} \left( \frac{\Delta^m \beta_k}{K^{-m}} \right)^2, \tag{3}$$

where  $\Delta$  is the differencing operator such that  $\Delta \beta_k = \beta_k - \beta_{k-1}$ , m is a positive integer indicating the order of differencing with  $\Delta^m = \Delta(\Delta^{m-1})$ . Let  $Y = (y_1, \dots, y_T)^T$  and X be the diagonal matrix whose (i, i)th element is  $x_i$ , and  $\Psi$  be the matrix<sup>3</sup> whose (i, j)-th element is  $\Psi_i^{[p]}(z_j)$ . Define  $D_m$  as the differencing matrix such that the j-th element of  $D_m\beta$  is  $\Delta^m\beta_j$ . Then the above minimization criterion could be written in a matrix form as

$$(Y - X\Psi\beta)^T (Y - X\Psi\beta) + \tilde{\lambda}^{-1} K^{2m-1} \beta^T D_m^T D_m \beta.$$

In general, if we take into account the correlation among  $u_t$  and the fact that  $var(u) = \sigma^2 \Sigma_u$ , we could incorporate the weighted penalized splines approach and estimate  $\hat{\beta}$  by minimizing

$$(Y - X\Psi\beta)^T \Sigma_u^{-1} (Y - X\Psi\beta) + \tilde{\lambda}^{-1} K^{2m-1} \beta^T D_m^T D_m \beta. \tag{4}$$

A direct calculation shows that the solution to (4) is

$$\hat{\beta} = (\Psi^T X \Sigma_u^{-1} X \Psi + \tilde{\lambda}^{-1} K^{2m-1} D_m^T D_m)^{-1} \Psi^T X \Sigma_u^{-1} Y.$$
 (5)

Then the penalized spline estimator of  $\theta(z)$  for model (1) is defined as

$$\hat{\theta}(z) = \sum_{k=1}^{K+p} \Psi_k^{[p]}(x) \hat{\beta}_k.$$

The methodology and applications of penalized splines are discussed extensively in Ruppert et al. (2003), but its theoretical studies had been largely absent until recently. For the univariate nonparametric model, Hall and Opsomer (2005) establish the consistency of the penalized splines estimators. Li and Ruppert (2008) derive the asymptotic normality and they were the first to obtain explicit formula for the asymptotic bias and variance. Claeskens et al. (2009) study the convergence rate of the penalized spline estimation and discussed the impact of the number of knots. However, all of these results are not directly applicable for varying-coefficient models with nonstationary regressors.

The following theorem establishes the consistency of the penalized spline estimator. Please note that all proofs of the theorems are relegated to the appendix.

<sup>&</sup>lt;sup>3</sup>The dimensions of  $\Psi$  and  $\beta$  both depend on the degree of splines p and the number of knots K. For notation simplicity, we suppress the subscripts p and K.

### Theorem 3.1. Suppose that Assumptions 1-3 hold. In addition, assume that

- i) The pth degree uniform B-splines are used to model  $\theta(z)$ . The number of knots satisfies  $K \sim T^{r_1}$  with  $\frac{2m}{2m+1} < r_1 < 1$ .
- ii) The mth order penalty is used and the penalty parameter  $\tilde{\lambda}$  satisfies that  $\tilde{\lambda} \to 0$  and  $T^2\tilde{\lambda} \to \infty$ .

Then for  $z \in (0,1)$ , the penalized spline estimator  $\hat{\theta}(z)$  satisfies

$$\hat{\theta}(z) - \theta(z) = O_p(T^{-1}\tilde{\lambda}^{-1/2}) + O_p(T^{-1+1/(2m)}\tilde{\lambda}^{1/(4m)}) + O_p(K^{-1}).$$
(6)

REMARK 3.1. Theorem 3.1 establishes the consistency of the penalized splines estimator for varying-coefficient models with nonstationary regressors. The term  $O_p(T^{-1}\tilde{\lambda}^{-1/2})$  reflects the order of asymptotic bias due to smoothing. The term  $O_p(K^{-1})$  is the design bias due to the use of splines in approximating the smooth functions  $\theta(z)$ . These results are standard. However, the order of asymptotic variance becomes  $O_p(T^{-2+1/m}\tilde{\lambda}^{1/(2m)})$  when  $x_t$  is an integrated process of order one, compared to a slower rate  $O_p(T^{-1+1/2m}\tilde{\lambda}^{1/(2m)})$  for stationary  $x_t$ . Correspondingly, when  $\tilde{\lambda}$  satisfies  $\tilde{\lambda} \sim T^{-2/(2m+1)}$ , our estimator achieves the optimal convergence rate  $T^{\frac{2m}{2m+1}}$ , faster than the  $T^{m/(2m+1)}$  convergence rate with  $\tilde{\lambda} \sim T^{-1/(2m+1)}$  for stationary  $x_t$ . In particular, when  $\theta(z) \in W^2[0,1]$ , i.e. m=2, the optimal convergence rate of  $\hat{\theta}(z)$  is  $T^{4/5}$  when  $x_t$  is integrated with order one. Such a result is consistent with Cai et al. (2009) and Xiao (2009).

REMARK 3.2. Penalized splines allow one to flexibly select the degree of splines p, the number of knots K, and the amount of penalty  $\tilde{\lambda}^{-1}$ . Our results have the following implications. First, the degree of splines p has no impact in the convergence rate of the estimator. Second, the number of knots K is not crucial as long as it exceeds a certain minimum bound. Third, the penalty parameter  $\tilde{\lambda}$  could serve as the key smoother and it determines the convergence rate of the estimator. These three conclusions are consistent with the results obtained

in Li and Ruppert (2008). However, when  $x_t$  is stationary, the term  $(\tilde{\lambda}^{-1}/T)^{1/(2m)}$  serves as the equivalent bandwidth used in a Nadaraya-Watson kernel estimator. In contrast, when  $x_t$  is integrated with order 1, the term  $(\tilde{\lambda}^{-1}/T^2)^{1/(2m)}$  serves as the equivalent bandwidth used in a Nadaraya-Watson kernel estimator.

Besides the B-spline family, another popular choice of the basis is the p-th degree truncated power polynomial basis (TPS) defined as

$$\{1, z, \cdots, z^p, (z - \kappa_1)_+^p, \cdots, (z - \kappa_K)_+^p\},\$$

where  $(z-a)_+^p = {\max(0, z-a)}^p$ . Note that the *p*-th degree TPS and the *p*-th degree B-splines span the same linear space. For any given  $\beta_k$ 's, there exists  $\gamma_j$ 's and  $b_k$ 's such that

$$\sum_{k=1}^{K+p} \Psi_k(z) \beta_k = \sum_{j=0}^p \gamma_j z^j + \sum_{k=1}^K (z - \kappa_k)_+^p b_k.$$

Moreover, a direct calculation shows that the coefficients of TPS and B-splines satisfy  $b_{k+1} = (-K)^p \Delta^{p+1} \beta_k$  (de Boor, 1978). Hence imposing penalty on  $\sum_k b_k^2$  is equivalent to imposing the (p+1)-th order differencing penalty on the B-splines coefficients  $\beta_k$ 's. In general, we could define the penalty matrix<sup>4</sup>  $\Lambda$  such that  $\tilde{\lambda}^{-1} K^{2m-1} \beta^T D_m^T D_m \beta = \tilde{\lambda}^{-1} K^{2m-2p-1} b^T \Lambda b$ . Equivalently, we could rewrite the minimization criterion (4) for TPS as

$$(Y - XZ_1\gamma - XZ_2b)^T \Sigma_u^{-1} (Y - XZ_1\gamma - XZ_2b) + \tilde{\lambda}^{-1} K^{2m-2p-1} b^T \Lambda b, \tag{7}$$

where  $Z_1$  and  $Z_2$  are matrices whose *i*-th row are  $(1, z, \dots, z^p)$  and  $\{(z-\kappa_1)_+^p, \dots, (z-\kappa_K)_+^p\}$  respectively, and  $\gamma = (\gamma_0, \dots, \gamma_p)^T$ . Because of the equivalence between TPS and B-splines of the same degree, the penalized spline estimator based on TPS could achieve the same optimal convergence rate  $T^{2m/(2m+1)}$  when the penalty parameter satisfies  $\tilde{\lambda} \sim T^{-2/(2m+1)}$ . Since the choice of p will not affect the convergence rate of the spline estimator, a conventional choice is to let p = 1 for  $\theta(z) \in W^2[0,1]$ , i.e. m = 2. In this case,  $\Lambda$  becomes the identity matrix

<sup>&</sup>lt;sup>4</sup>The choice of  $\Lambda$  depends on both the degree of splines p and the order of penalty m. For notation simplicity, we suppress the subscripts p and m.

 $I_K$ . If we denote  $\lambda^{-1} = \tilde{\lambda}^{-1}K$ , then the minimization criterion (7) could be written as

$$(Y - XZ_1\gamma - XZ_2b)^T \Sigma_u^{-1} (Y - XZ_1\gamma - XZ_2b) + \lambda^{-1}b^Tb,$$
(8)

and the optimal rate of  $\lambda$  is of the order  $T^{-2/5}K^{-1}$ .

# 4 Inference using Likelihood Ratio Tests

In this section, we consider testing whether the functional coefficients  $\theta(z)$  is time-invariant. The null hypothesis is  $H_0$ :  $\theta(z) \equiv \theta_0$ . Under the alternative,  $\theta(z)$  is a smooth function of unknown form. Such a stability test is of both theoretical and empirical importance. For example, when a linear cointegration model is misspecified, the resulting estimation of  $\hat{\theta}_0$  would not be consistent and neither of the equilibrium residuals. As a result, the traditional cointegration tests might fail to detect the cointegrating relationship.

In the literature of nonparametric regression, there are also lots of discussions on checking whether there is enough evidence to support the use of the general nonparametric method rather than a simple linear cointegration model. In general, traditional approaches often rely on i) comparing the discrepancy measures between the estimates obtained under the null and the alternative, see Härdle and Mammen (1993); or ii) constructing the F-test statistic based on the sum of residuals, see Hong and White (1995); or iii) conducting the generalized likelihood ratio test using a reasonable smooth estimate under the alternative, see Fan, Zhang and Zhang (2001). In any of these methods, it is crucial to select the smoothing parameter under the nonparametric alternative. In practice, the power of the test is likely to be affected by the smoothing parameter, especially when it is chosen by some ad hoc methods.

In contrast, we are going to propose a likelihood ratio test procedure that could circumvent this difficulty as we use maximum likelihood principles for both estimation and inference. First, we model  $\theta(z)$  by the p-th degree splines in order to define a general nonparametric alternative. As we show in the section above, there is not much difference to estimate  $\theta(z)$ by using either the B-splines family or the TPS family. Moreover, the choice of p is not important. Therefore, we mainly focus on using the linear TPS family in this section. Since we could view the spline coefficients  $b_k$ 's, associated with  $(z - \kappa_k)_+$ 's, as the deviations from the linear function. Hence testing the stability of  $\theta(z)$  is equivalent to testing both the linear coefficient and the spline coefficients being 0, i.e.

$$H_0: \gamma_1 = 0$$
 and  $b_1 = \dots = b_K = 0$ ,

against

$$H_A: \gamma_1 \neq 0$$
 or  $\exists k, s.t. b_k \neq 0$ .

Note that this is a multiple testing problem and the number of restrictions under  $H_0$  grows as the sample size does. To circumvent this difficulty, a new idea is to utilize the mixed model representation for spline estimates based on TPS by treating  $b_k$ 's as random coefficients with a common variance component, and then relate the null hypothesis above to the significance test of zero variance. More details are given below.

Note that minimizing (7) is equivalent to solving a system of equations

$$\begin{pmatrix} A_1^T \Sigma_u^{-1} A_1 & A_1^T \Sigma_u^{-1} A_2 \\ A_2^T \Sigma_u^{-1} A_1 & A_2^T \Sigma_u^{-1} A_2 + \lambda^{-1} \Lambda \end{pmatrix} \begin{pmatrix} \hat{\gamma} \\ \hat{b} \end{pmatrix} = \begin{pmatrix} A_1^T \Sigma_u^{-1} Y \\ A_2^T \Sigma_u^{-1} Y \end{pmatrix},$$

where  $A_1 = XZ_1$  and  $A_2 = XZ_2$ . The above equation is essentially Henderson's mixed model equations, which motivates us to utilize the mixed model representation to obtain  $\hat{\gamma}$  and  $\hat{b}$  as the best linear unbiased predictors (BLUP) in the following model. To be specific, let  $Y = A_1\gamma + A_2b + u$ , where  $\gamma$  is the  $2 \times 1$  vector of fixed effect coefficients and b is the  $K \times 1$  vector of random effect coefficients with mean 0 and variance  $\lambda \sigma^2 \Lambda^{-1}$ , with  $\Lambda = I_K$  when m = 2 and p = 1. The parameter  $\lambda$  controls the amount of smoothing and it could be viewed as the signal to noise ratio. Following Crainiceanu and Ruppert (2004), we could treat Y as

$$Y|b, A_1, A_2 \stackrel{D}{=} N(A_1\gamma + A_2b, \sigma^2\Sigma_u), \quad b \stackrel{D}{=} N(0_{K\times 1}, \lambda\sigma^2\Lambda^{-1}).$$

Note that  $E(Y) = A_1 \gamma$  and  $var(Y) = \sigma^2(\Sigma_u + \lambda A_2 \Lambda^{-1} A_2^T) =: \sigma^2 \Omega_{\lambda}$ . Hence we could define

a twice of the log-likelihood of Y as

$$2l(\gamma, \lambda, \sigma^2) = -(Y - A_1 \gamma)(\sigma^2 \Omega_\lambda)^{-1} (Y - A_1 \gamma) - \log|\sigma^2 \Omega_\lambda| - T \log(2\pi). \tag{9}$$

By maximizing (9), we could estimate the variance components by  $\hat{\sigma}^2$  and  $\hat{\lambda}$ . Define  $\hat{\Omega}_{\lambda} = \Sigma_u + \hat{\lambda} A_2 \Lambda^{-1} A_2^T$ . The BLUP of  $\gamma$  and b are then obtained as

$$\hat{\gamma} = (A_1^T \hat{\Omega}_{\lambda}^{-1} A_1)^{-1} A_1^T \hat{\Omega}_{\lambda}^{-1} Y, \quad \text{and} \quad \hat{b} = \hat{\lambda} \Lambda^{-1} A_2^T \hat{\Omega}_{\lambda}^{-1} (Y - A_1 \hat{\gamma}), \quad (10)$$

and we could estimate  $\theta(z)$  by

$$\hat{\theta}(z) = \sum_{k=0}^{p} \hat{\gamma}_k z^k + \sum_{k=1}^{K} \hat{b}_k (z - \kappa_k)_+^p.$$

For the same  $\lambda$ , minimizing equation (7) yields the same solution as (10). However, the use of the mixed model representation allows us to adopt the maximum likelihood principle to make estimation as well as inference on  $\lambda$ . In particular,  $\lambda = 0$  implies  $b_k = 0$  for all k. Hence the hypothesis test of  $\theta(z)$  being constant is equivalent to testing

$$H_0: \gamma_1 = 0$$
 and  $\lambda = 0$ .

against

$$H_A: \gamma_1 \neq 0$$
 or  $\lambda \neq 0$ .

Then it is straightforward to rely on the LRT statistic for inference, where

$$LRT_T = \sup_{H_A} 2 \log l(\gamma, \lambda, \sigma^2) - \sup_{H_0} 2 \log l(\gamma, \lambda, \sigma^2).$$

Note that the null distribution of the LRT statistic is not standard as the parameter  $\lambda$  is always non-negative and it lies on the boundary of the parameter space under  $H_0$ . Therefore, we derive the exact and the limiting null distributions of our test statistic below.

First we consider the exact case, where both T and K are relatively small and could be treated as fixed. Let P be the projection matrix  $P = I_T - \Sigma_u^{-1/2} A_1 (A_1^T \Sigma_u^{-1} A_1)^{-1} A_1^T \Sigma_u^{-1/2}$ . Define  $\xi_{s,T}$  and  $\eta_{s,T}$  as the s-th eigenvalues of the  $K \times K$  matrices  $\Lambda^{-1/2} A_2^T \Sigma_u^{-1} A_2 \Lambda^{-1/2}$  and  $\Lambda^{-1/2} A_2^T \Sigma_u^{-1/2} P \Sigma_u^{-1/2} A_2 \Lambda^{-1/2}$  respectively. We have the following results.

THEOREM 4.1. Suppose that Assumptions 1-3 hold and the linear TPS with equi-spaced knots are used. Then under  $H_0: \theta(z) \equiv \theta_0$ ,

$$LRT_{T} \stackrel{D}{=} \sup_{\lambda \ge 0} \left[ T \log \left\{ 1 + \frac{N_{T}(\lambda)}{D_{T}(\lambda)} \right\} - \sum_{s=1}^{K} \log \left( 1 + \lambda \xi_{s,T} \right) \right] + T \log \left\{ 1 + \frac{w_{T-1}^{2}}{\sum_{s=1}^{T-2} w_{s}^{2}} \right\}, \tag{11}$$

where 
$$N_T(\lambda) = \sum_{s=1}^K \frac{\lambda \eta_{s,T}}{1 + \lambda \eta_{s,T}} w_s^2$$
,  $D_T = \sum_{s=1}^K \frac{w_s^2}{1 + \lambda \eta_{s,T}} + \sum_{s=K+1}^{T-2} w_s^2$  and  $w_s \stackrel{D}{=} iidN(0,1)$ .

Theorem 4.1 derives the exact null distribution of the LRT statistic when the sample size T is finite. Although equation (11) does not have a close form, we could efficiently simulate this distribution using the following Algorithm A.

Step 1. define a grid  $0 = \lambda_1 < \lambda_2 < \cdots < \lambda_L$  of possible values for  $\lambda$ .

Step 2. simulate K independent  $\chi^2_1$  random variables  $w^2_1, \cdots, w^2_K$ .

Step 3. simulate a random variable  $\nu_1$  that follows  $\chi^2_{T-2-K}$ .

Step 4. simulate a random variable  $\nu_0$  that follows  $\chi_1^2$ .

Step 5. for every 
$$\lambda_i$$
, compute  $f_T(\lambda_i) = T \log\{1 + \frac{N_T(\lambda_i)}{D_T(\lambda_i)}\} - \sum_{s=1}^K (1 + \lambda \xi_{s,T})$ .

Step 6. determine  $\lambda_{\text{max}}$  which maximizes  $f_T(\lambda_i)$  over  $\lambda_i$ 's.

Step 7. compute 
$$f_T(\lambda_{\max}) + T \log\{1 + \frac{\nu_0}{\sum_{i=1}^K w_s^2 + \nu_1}\}$$
.

Step 8. repeat steps 2–7.

If we treat K as fixed and let T grow to infinity, we have the following results.

THEOREM 4.2. Suppose that Assumptions 1-3 hold and the linear TPS with equi-spaced knots are used. Then there exist  $\eta_s$  and  $\xi_s$ , for  $s = 1, \dots, K$ , such that

$$T^{-2}\eta_{s,T} \Rightarrow \eta_s, \qquad T^{-2}\xi_{s,T} \Rightarrow \xi_s, \qquad \text{as} \quad T \to \infty.$$
 (12)

Moreover, under  $H_0: \theta(z) \equiv \theta_0$ ,

$$LRT_T \Rightarrow \sup_{d \ge 0} \{ \sum_{s=1}^K \frac{d\eta_s}{1 + d\eta_s} w_s^2 - \sum_{s=1}^K \log(1 + d\xi_s) \} + \chi_1^2.$$
 (13)

REMARK 4.1. In Theorem 4.2, we have explicitly derived that the convergence rate of the eigenvalues  $\eta_{s,T}$  and  $\xi_{s,T}$  is  $T^2$ , which is faster than the T convergence rate due to nonstationarity. When K is fixed, the part corresponding to testing  $\lambda = 0$  converges to the term  $\sup_{d\geq 0} \{\sum_{s=1}^K \frac{d\eta_s}{1+d\eta_s} w_s^2 - \sum_{s=1}^K \log(1+d\xi_s)\}$ . Since this limiting distribution is nonstandard, one could simulate it by modifying Algorithm A described above.

Furthermore, if we assume K and T both grow to infinity, the null distribution of  $LRT_T$  approaches to a simple  $\chi^2$  distribution.

THEOREM 4.3. Suppose that Assumptions 1-3 hold and the linear TPS with equi-spaced knots are used. Let the number of knots K satisfying  $K \sim T^r$  with 4/5 < r < 1. Then there exist  $\bar{\eta}$  and  $\bar{\xi}$  such that

$$K^{-1}T^{-2}\sum_{k=1}^{K}\eta_{k,T} \Rightarrow \bar{\eta}, \qquad K^{-1}T^{-2}\sum_{k=1}^{K}\xi_{k,T} \Rightarrow \bar{\xi}.$$
 (14)

Under  $H_0: \theta(z) \equiv \theta_0$ , we have,

$$LRT_T \Rightarrow \chi_1^2$$
. (15)

REMARK 4.2. Theorem 4.3 assumes that K grows as T does. Compared to the fixed K case, the amount of penalty  $\lambda^{-1}$  is expected to be larger, and the probability of obtaining the maximum likelihood estimate (MLE) of  $\lambda$  at its actual value 0 approaches to 1 provided that  $H_0$  is true. Therefore, the part corresponding to testing  $\lambda = 0$  degenerates and we have a simple  $\chi^2$  distribution.

For the local alternatives, we assume that  $\theta(z) = \theta_0 + T^{-\alpha}\theta_1(z)$ , where  $\theta_1(z)$  is a nonzero smooth function that belongs to  $W^2[0,1]$ . Suppose we span  $\theta_1(z)$  with the first degree TPS as  $\bar{\gamma}_0 + \bar{\gamma}_1 z + \sum_{k=1}^K \bar{b}_k (z - \kappa_k)_+$ . Utilizing the mixed model representation, we treat the spline coefficients  $\bar{b} = (\bar{b}_1, \dots, \bar{b}_K)^T$  as random with mean 0 and variance  $\bar{\lambda}_0 \sigma^2 \Lambda^{-1}$ . It has been shown in Section 3 that  $\bar{\lambda}_0$  converges to 0 at the rate of  $T^{2/5}K$ . Therefore, we denote  $\bar{\lambda}_0 = \bar{d}_0 T^{-2/5}K^{-1}$  for some constant  $\bar{d}_0 \geq 0$ . Recall that our LRT test will examine both the fixed part  $\bar{\gamma}_1$  and the variance part  $\bar{d}_0$ . Therefore, we will consider two different cases in the

local alternatives. In Case 1,  $\theta_1(z)$  is a linear function with nonzero slope. i.e.  $\bar{\gamma}_1 \neq 0$  but  $\bar{d}_0 = 0$ . The local alternative is then set as  $H_{01}: \theta(z) = \theta_0 + T^{-1}\theta_1(z)$ . In Case 2,  $\theta_1(z)$  has  $\bar{\gamma}_1 = 0$  but  $\bar{d}_0 \neq 0$ . The local alternative is set as  $H_{02}: \theta(z) = \theta_0 + T^{-4/5}\theta_1(z)$ .

THEOREM 4.4. Suppose that Assumptions 1-3 hold and the linear TPS with equi-spaced knots are used. Let the number of knots K satisfying  $K \sim T^r$  with 4/5 < r < 1.

Under the local alternative  $H_{01}$ , the LRT statistic converges to a noncentral  $\chi_1^2$ , i.e.

$$LRT_T \Rightarrow (w_1 + \bar{\gamma}_1 \bar{\pi}_2)^2$$
,

where  $w_1 \stackrel{D}{=} N(0,1)$  and  $\bar{\pi}_2$  is defined right before equation (52).

Under the local alternative  $H_{02}$ ,

$$LRT_T \Rightarrow \max_{d \in [0,\bar{d}_0]} \{\bar{d}_0\bar{\eta} + \bar{g}_3(d) - \bar{g}_2(d)\} + (1 + \bar{d}_0\varrho)w_1^2,$$

where  $w_1 \stackrel{D}{=} N(0,1)$ ,  $\bar{\eta}$  and  $\bar{g}_2(d)$  are defined as in Theorem 4.3,  $\bar{g}_3(d)$  is defined right after equation (53) and  $\varrho$  is defined right before equation (55).

REMARK 4.3. Strictly speaking, we should also consider Case 3 where neither  $\bar{\gamma}_1$  or  $\bar{d}_0$  equals 0. For this case, the local alternative could be set as  $H_{03}:\theta(z)=\theta_0+T^{-1}\theta_1(z)$ . Notice that such a local alternative converges with a rate faster than  $T^{4/5}$ . The nonzero variance component will not affect the asymptotic distribution. Hence in Case 3, the LRT statistic still converges to a noncentral  $\chi^2$  distribution as in Case 1. To save the length of this paper, we omit the detailed discussions of Case 3.

REMARK 4.4. Under  $H_{02}$ , the asymptotic distribution of  $LRT_T$  has two components, where the first part is nonnegative and the second part is a scaled  $\chi_1^2$ . In summary, our penalized spline estimator of  $\theta(z)$  has the  $T^{4/5}$  convergence rate, while our test statistics could detect an alternative whose convergence rate is not faster than  $T^{4/5}$ . On the other hand, for any sequence such that  $H_{0A}: \theta(z) \equiv \theta_0 + T^{-\alpha}\theta_1(z)$  and  $\alpha < 4/5$ ,  $LRT_T$  diverges and the power function satisfies

$$P(LRT_T > \chi^2_{1,\alpha/2}) \to 1,$$

where  $\chi^2_{1,\alpha/2}$  is the upper  $\alpha/2$  quantile of  $\chi^2_1$  distribution. Hence the proposed likelihood ratio test could achieve the optimality.

# 5 Finite Sample Performance

Monte Carlo simulations are conducted in this section to examine the finite sample performance of the proposed LRT test. The data generating process is

$$y_t = \theta(z_t)x_t + u_t,$$

where  $x_t = x_{t-1} + v_t$ ,  $u_t = \rho u_{t-1} + \varepsilon_t$ ,  $v_t$ 's and  $\varepsilon_t$ 's are iid N(0,1), and they are independent of each other. The initial values are set to be zero. In particular, four cases for the parameter values are considered: i)  $\theta(z_t) = 0.25$ ,  $\rho = 0$ ; ii)  $\theta(z_t) = 0.25$ ,  $\rho = 0.5$ ; iii)  $\theta(z_t) = (z_t - 0.5)^2$ ,  $\rho = 0$ ; iv)  $\theta(z_t) = (z_t - 0.5)^2$ ,  $\rho = 0.5$ , where the first two cases are related to calculating the size of the test and the last two are related to calculating the power of the test. The simulation designs above are similar to those in Xiao (2009). The sample sizes we consider are T = 100 and 300. In particular, we would like to examine the impact of the number of knots K. Hence when T = 100, we consider three situations, K = 10, K = 20 and K = 40; when T = 300, we consider K = 20, K = 40, and K = 80. All reported results are based on 2000 replications.

Table 1 report the size of the proposed likelihood ratio test when  $\rho$  is given, i.e. the true covariance matrix  $\Sigma_u$  is known. The five columns on the left use the critical values based on the finite distribution derived in Theorem 4.1, while the five columns on the right use the critical values based on the asymptotic distribution  $\chi_1^2$  as indicated in Theorem 4.3. From Table 1, we find that both the finite distribution and the  $\chi^2$  limiting distribution work very well. For example, consider Panel A1 with K chosen as 10. Even though the sample size is just 100, the actual rejection rates based on the asymptotic  $\chi^2$  distribution are 0.1995, 0.1490, 0.0995, 0.0465 and 0.0100, very close to the nominal sizes 0.2, 0.15, 0.10, 0.05 and 0.01 respectively. Moreover, we find that the number of knots does not have much impact on

the size performance. For a given nominal level in any reported panel, the absolute differences in the rejection rates associated with different K are not greater than 0.005. This is consistent with the empirical conclusions that the number of knots is not important, provided that it is above some minimum threshold (Ruppert, 2002).

Table 2 repeats all designs in Table 1, except that the covariance matrix  $\Sigma_u$  is treated as unknown and replaced by an estimate. We find that the our LRT procedure still performs well and is less likely to be affected by the fact that the covariance is unknown.

Table 3 reports the power of our test statistic. Once again, we find that the choice of the number of knots K is not important and the procedure is robust against the use of an estimated covariance. When the sample size increase from 100 to 300, the rejection rates are all greater than 0.98, implying very good power performance of our testing procedure.

# 6 Conclusions

Varying-coefficient regression models with nonstationary regressors have received heated interests in recent years. This paper proposes a penalized splines approach to estimate the varying coefficients. Compared to kernel-based methods, penalized splines estimation not only achieves the same optimal convergence rate, but also enjoys the advantage of fast computation. Utilizing the mixed model representation of penalized splines, we construct a likelihood ratio test for the stability of the varying coefficient. We derive the exact and limiting distributions of the proposed test statistic. When the number of knots is treated as fixed, the null distribution is non-standard, but could be simulated via a proposed algorithm using spectral decomposition. When the number of knots grows as the sample size does, the limiting null distribution converges to a simple  $\chi^2$  distribution. Our test is less likely to be suffered from the mis-selection of the smoothing parameters. Simulations show that the asymptotic distribution works very well even for the finite sample case.

There are some issues worth of future studies. One potential analysis is to extend the current setting to the case allowing for dependence between  $u_t$  and  $v_t$ . Another natural extension is to consider a more general varying-coefficient cointegrating regression model which includes both the stochastic and the deterministic functional coefficients in the cointegrating relationship.

# Appendix A: Proofs

Proof of Theorem 3.1: Note that our model could be written as  $Y = X\Psi\beta + u$ . If  $\Sigma_u \neq I_T$ , we could always multiply  $\Sigma_u^{-1/2}$  and consider instead  $\tilde{Y} = \tilde{X}\Psi\beta + \tilde{u}$ , where  $\tilde{Y} = \Sigma_u^{-1/2}Y$ ,  $\tilde{X} = \Sigma_u^{-1/2}X$  is an integrated process and the elements of  $\tilde{u} = \Sigma_u^{-1/2}u$  are uncorrelated. Hence without loss of generality, we only need to show that equation (6) holds when  $u_t$ 's are uncorrelated, i.e.  $\Sigma_u$  is the identity matrix.

Recall that the penalized spline estimator  $\hat{\theta}(z)$  could be written as in equation (5), i.e.

$$\hat{\theta}(z) = \Psi_z (\Psi^T X^2 \Psi + \tilde{\lambda}^{-1} K^{2m-1} D_m^T D_m)^{-1} \Psi^T X Y, \tag{16}$$

where  $\Psi_z = \{\Psi_1^{[p]}(z), \dots, \Psi_{K+p}^{[p]}(z)\}$ . First consider the (i, j)th element of the term  $\Psi^T X^2 \Psi$ . Define  $R_0 =: T^{-2} \sum_{t=1}^T x_t^2 [\Psi_i^{[p]}(z_t) \Psi_j^{[p]}(z_t) - E\{\Psi_i^{[p]}(z_t) \Psi_j^{[p]}(z_t)\}]$ . By subtracting and adding the mean, we have,

$$T^{-2}(\Psi^T X^2 \Psi)_{i,j} = T^{-2} \sum_{t=1}^T x_t^2 \Psi_i^{[p]}(z_t) \Psi_j^{[p]}(z_t) = R_0 + T^{-2} \sum_{t=1}^T E\{\Psi_i^{[p]}(z_t) \Psi_j^{[p]}(z_t)\} x_t^2.$$

Recall that  $\Psi_i^{[p]}(z)$  is nonzero only in a small interval of length (p+1)/K. For example, when zero degree splines are used,  $\Psi_i^{[0]}(z)$  is the indicator function  $I_{(i-1)/K < z \le i/K}$ . Hence  $E\{\Psi_i^{[p]}(z_t)\Psi_j^{[p]}(z_t)\} = O(K^{-1})$  and  $var\{\Psi_i^{[p]}(z_t)\Psi_j^{[p]}(z_t)\} = O(K^{-1})$ . Moreover,

$$T^{-2} \sum_{t=1}^{T} E\{\Psi_i^{[p]}(z_t)\Psi_j^{[p]}(z_t)\}x_t^2 = E\{\Psi_i^{[p]}(z_t)\Psi_j^{[p]}(z_t)\}T^{-1} \sum_{t=1}^{T} (\frac{x_t}{\sqrt{T}})^2 = O_p(K^{-1}),$$

and hence

$$R_0 = T^{-1} \sum_{t=1}^{T} \left( \frac{x_t}{\sqrt{T}} \right)^2 \left[ \Psi_i^{[p]}(z_t) \Psi_j^{[p]}(z_t) - E\{\Psi_i^{[p]}(z_t) \Psi_j^{[p]}(z_t) \} \right] = O_p\{(TK)^{-1/2}\} = o_p(K^{-1}),$$

where the last equality holds as K = o(T). Denote  $q_{ij}$  as the limit of  $KE\{\Psi_i^{[p]}(z_t)\Psi_j^{[p]}(z_t)\}$ . Then

$$KT^{-2}(\Psi^T X X \Psi)_{ij} = K[T^{-2} \sum_{t=1}^{T} x_t^2 E\{\Psi_i^{[p]}(z_t) \Psi_j^{[p]}(z_t)\} + R_0] \Rightarrow q_{ij} \int B_v^2(s) ds. \quad (17)$$

Therefore,

$$\hat{\theta}(z) = \Psi_z (\Psi^T X^2 \Psi + \tilde{\lambda}^{-1} K^{2m-1} D_m^T D_m)^{-1} \Psi^T X Y, 
= \Psi_z (\Psi^T X^2 \Psi + \tilde{\lambda}^{-1} K^{2m-1} D_m^T D_m)^{-1} \Psi^T X^2 \Psi \beta 
+ \Psi_z (\Psi^T X^2 \Psi + \tilde{\lambda}^{-1} K^{2m-1} D_m^T D_m)^{-1} \Psi^T X u + O_p (K^{-1}) 
=: R_1 + R_2 + O_p (K^{-1}),$$
(18)

where the term  $O_p(K^{-1})$  comes from the bias due to splines approximation,

$$R_{1} =: \Psi_{z} \left\{ \frac{K}{T^{2}} (\Psi^{T} X^{2} \Psi + \tilde{\lambda}^{-1} K^{2m-1} D_{m}^{T} D_{m}) \right\}^{-1} \frac{K}{T^{2}} \Psi^{T} X^{2} \Psi \beta$$
$$= \Psi_{z} \left\{ \left( \frac{K}{T^{2}} \Psi^{T} X^{2} \Psi \right) + \bar{\lambda} K^{2m} D_{m}^{T} D_{m} \right\}^{-1} \left( \frac{K}{T^{2}} \Psi^{T} X^{2} \Psi \right) \beta,$$

$$R_{2} =: \Psi_{z} \left\{ \frac{K}{T^{2}} (\Psi^{T} X^{2} \Psi + \tilde{\lambda}^{-1} K^{2m-1} D_{m}^{T} D_{m}) \right\}^{-1} \frac{K}{T^{2}} \Psi^{T} X u$$

$$= \frac{1}{\sqrt{T}} \Psi_{z} \left\{ (\frac{K}{T^{2}} \Psi^{T} X^{2} \Psi) + \bar{\lambda} K^{2m} D_{m}^{T} D_{m} \right\}^{-1} \left\{ \frac{\Psi^{T} (T^{-1/2} X u)}{T/K} \right\},$$

and  $\bar{\lambda} = \tilde{\lambda}^{-1}/T^2$ . By equation (17), the term  $\frac{K}{T^2}\Psi^TX^2\Psi$  converges. Using a similar technique as in Li and Ruppert (2008), we can show that the term  $\bar{\lambda}^{1/(2m)}$  serves equivalently as the bandwidth h used in a Nadaraya-Watson kernel estimator. Therefore,

$$R_1 - \Psi_z \beta = O_p(h^m) = O_p\{\bar{\lambda}^{m/(2m)}\} = O_p(T^{-1}\tilde{\lambda}^{-1/2}).$$
(19)

Now consider the second term  $R_2$ . Note that  $ER_2 = 0$  and the *i*-th element of  $\Psi^T X u$  satisfies

$$\frac{\Psi^T(T^{-1/2}Xu)_i}{T/K} = \frac{\sum_{t=1}^T \{\Psi_i^{[p]}(z_t) \frac{x_t}{\sqrt{T}} u_t\}}{T/K} = O_p(\frac{1}{\sqrt{T/K}}). \tag{20}$$

By the fact that  $\bar{\lambda} = \tilde{\lambda}^{-1}/T^2$  and  $\bar{\lambda}^{1/(2m)}$  serves as the bandwidth used in Nadaraya-Watson kernel estimate, we have

$$var(\sqrt{T}R_2) = O(\frac{1}{Kh}\frac{1}{T/K}) = O(\frac{1}{Th}) = O(\frac{1}{T\bar{\lambda}^{1/(2m)}}) = O(T^{-1+1/m}\tilde{\lambda}^{1/(2m)}),$$

and hence

$$R_2 = \frac{1}{\sqrt{T}} O_p(T^{-1/2 + 1/(2m)} \tilde{\lambda}^{1/(4m)}) = O_p(T^{-1 + 1/(2m)} \tilde{\lambda}^{1/(4m)}). \tag{21}$$

Together with equation (18), (19) and (21),

$$\hat{\theta}(z) - \theta(z) = O_p(T^{-1}\tilde{\lambda}^{-1/2}) + O_p(T^{-1+1/(2m)}\tilde{\lambda}^{1/(4m)}) + O_p(K^{-1}),$$

and Theorem 3.1 holds.  $\Box$ 

Note that the proofs of Theorem 4.1–4.4 share lots of similarity. Thus we first provide a general description and four useful propositions which could be applied to all these theorems.

Recall that twice of the log-likelihood of Y is defined as in equation (9), i.e.

$$2l(\gamma, \lambda, \sigma^2) = -(Y - A_1 \gamma)(\sigma^2 \Omega_{\lambda})^{-1}(Y - A_1 \gamma) - \log|\sigma^2 \Omega_{\lambda}| - T\log(2\pi).$$

where  $\Omega_{\lambda} = \Sigma_{u} + \lambda A_{2} \Lambda^{-1} A_{2}^{T}$  and  $\sigma^{2} \Omega_{\lambda} = var(Y)$ . Instead of maximizing  $2l(\gamma, \lambda, \sigma^{2})$  over the parameter space  $(\gamma, \lambda, \sigma^{2})$ , we consider maximizing the profile log-likelihood function  $2 \log L(\lambda) = 2l(\hat{\gamma}_{\lambda}, \lambda, \hat{\sigma}_{\lambda}^{2})$  over the parameter space  $\lambda \geq 0$ , where  $\hat{\gamma}_{\lambda}$  and  $\hat{\sigma}_{\lambda}^{2}$  are the profile maximum likelihood estimates (MLE) for  $\gamma$  and  $\sigma^{2}$  when  $\lambda$  is given. Specifically, they satisfy

$$\hat{\gamma}_{\lambda} = (A_1^T \Omega_{\lambda}^{-1} A_1)^{-1} A_1^T \Omega_{\lambda}^{-1} Y \quad \text{and} \quad \hat{\sigma}_{\lambda}^2 = T^{-1} (Y - A_1 \hat{\gamma}_{\lambda})^T \Omega_{\lambda}^{-1} (Y - A_1 \hat{\gamma}_{\lambda}).$$

By plugging  $\hat{\gamma}_{\lambda}$  and  $\hat{\sigma}_{\lambda}^2$  into equation (9), we could simplify the profile log-likelihood as

$$2\log L(\lambda) = -T - \log|\hat{\sigma}_{\lambda}^{2}\Omega_{\lambda}| - T\log(2\pi). \tag{22}$$

Denote  $\log L^0$  as the maximum log-likelihood under the null hypothesis. Then we can decompose  $2LRT_T$  into two parts by adding and subtracting  $2\log L(0)$ , i.e.

$$2LRT_T = \sup_{\lambda > 0} \{2\log L(\lambda) - 2\log L(0)\} + \{2\log L(0) - 2\log L^0\},$$

where the first part corresponds to testing  $\lambda = 0$  and the second part corresponds to testing the linear coefficient  $\gamma_1 = 0$  given that  $\lambda = 0$ . The following propositions establish some useful preliminary results. In particular, Proposition 6.1–Proposition 6.3 studies the property related to the first part, while Proposition 6.4 studies the property of the second part.

Proposition 6.1. Denote  $\tilde{Y} = \Sigma_u^{-1/2} Y$ ,  $\tilde{A}_1 = \Sigma_u^{-1/2} A_1$ ,  $\tilde{A}_2 = \Sigma_u^{-1/2} A_2$ . We have

$$2\log L(\lambda) - 2\log L(0) = T\log(1 + \frac{\lambda R_3/\sigma^2}{T\hat{\sigma}_{\lambda}^2/\sigma^2}) - \sum_{k=1}^{K} \log(1 + \lambda \xi_{k,T}).$$
 (23)

where  $R_3$  is given right below equation (29), and  $\xi_{k,T}$ 's are defined right before Theorem 4.1.

Proof of Proposition 6.1: Define  $V_{\lambda} =: I_T + \lambda \tilde{A}_2 \Lambda^{-1} \tilde{A}_2^T$ . With the notations of  $\tilde{Y}$ ,  $\tilde{A}_1$  and  $\tilde{A}_2$ , we have  $\Omega_{\lambda} = \Sigma_u^{1/2} V_{\lambda} \Sigma_u^{1/2}$ . Correspondingly, the profile MLE can be rewritten as

$$\hat{\gamma}_{\lambda} = (\tilde{A}_1^T V_{\lambda}^{-1} \tilde{A}_1)^{-1} \tilde{A}_1 V_{\lambda}^{-1} \tilde{Y} \text{ and } \hat{\sigma}_{\lambda}^2 = T^{-1} (\tilde{Y} - \tilde{A}_1 \hat{\gamma}_{\lambda})^T V_{\lambda}^{-1} (\tilde{Y} - \tilde{A}_1 \hat{\gamma}_{\lambda}).$$

By Patterson and Thompson (1971), it is well-known that there exists a matrix W satisfying

$$WW^{T} = P = I_{T} - \tilde{A}_{1}(\tilde{A}_{1}^{T}\tilde{A}_{1})^{-1}\tilde{A}_{1}^{T}, \qquad W^{T}W = I_{T-2}, \tag{24}$$

and that

$$T\hat{\sigma}_{\lambda}^{2} = (\tilde{Y} - \tilde{A}_{1}\hat{\gamma}_{\lambda})^{T}V_{\lambda}^{-1}(\tilde{Y} - \tilde{A}_{1}\hat{\gamma}_{\lambda}) = \tilde{Y}^{T}W(W^{T}V_{\lambda}W)^{-1}W^{T}\tilde{Y}.$$

By equation (22) and that  $|\Omega_{\lambda}| = |\Sigma_u^{1/2}||V_{\lambda}||\Sigma_u^{1/2}| = |V_{\lambda}||\Sigma_u|$ , we have

$$-2\log L(\lambda) = \log|\hat{\sigma}_{\lambda}^{2}\Omega_{\lambda}| + T\{1 + \log(2\pi)\}$$

$$= T\log(\hat{\sigma}_{\lambda}^{2}) + \log|V_{\lambda}| + \log|\Sigma_{u}| + T\{1 + \log(2\pi)\}$$

$$= T\log(\frac{T\hat{\sigma}_{\lambda}^{2}}{\sigma^{2}}) + \log|V_{\lambda}| + C_{0}.$$
(25)

where  $C_0 = \log |\Sigma_u| + T\{1 + \log(2\pi)\} - T\log(T/\sigma^2)$ .

Note that for any  $p_1 \times p_2$  matrices A and B, it holds that  $|I_{p_1} + A^T B| = |I_{p_2} + B^T A|$ . Moreover, |A| equals the product of its eigenvalues. Hence

$$\log |V_{\lambda}| = \log |I_T + \lambda \tilde{A}_2 \Lambda^{-1} \tilde{A}_2^T| = \log |I_K + \lambda \Lambda^{-1/2} \tilde{A}_2^T \tilde{A}_2 \Lambda^{-1/2}| = \sum_{k=1}^K \log(1 + \lambda \xi_{k,T}).$$

Together with equation (25), we have

$$-2\log L(\lambda) = T\log(\frac{T\hat{\sigma}_{\lambda}^2}{\sigma^2}) + \sum_{k=1}^K \log(1 + \lambda \xi_{k,T}) + C_0.$$
 (26)

When  $\lambda = 0$ , we have  $V_0 = I_T$ ,  $W^T V_0 W = W^T W = I_{T-2}$  and  $T \hat{\sigma}_0^2 = \tilde{Y}^T W W^T \tilde{Y}$ . Thus,

$$-2\log L(0) = T\log(\frac{\tilde{Y}^T W W^T \tilde{Y}}{\sigma^2}) + C_0. \tag{27}$$

It follows that

$$2\log L(\lambda) - 2\log L(0) = T\log\{1 + \frac{(\tilde{Y}^T W W^T \tilde{Y} - T\hat{\sigma}_{\lambda}^2)/\sigma^2}{T\hat{\sigma}_{\lambda}^2/\sigma^2}\} - \sum_{k=1}^{K} (1 + \lambda \xi_{k,T}). \quad (28)$$

Note that  $(W^T V_{\lambda} W)^{-1} = (I_{T-2} + \lambda W^T \tilde{A}_2 \tilde{\Lambda}^{-1} A_2^T W)^{-1} = I_{T-2} - \lambda \Delta_1$ , where, by Woodbury Matrix Identity,

$$\Delta_1 =: W^T \tilde{A}_2 \tilde{\Lambda}^{-1/2} (I_K + \lambda \Lambda^{-1/2} \tilde{A}_2^T W W^T \tilde{A}_2 \Lambda^{-1/2})^{-1} \tilde{\Lambda}^{-1/2} \tilde{A}_2^T W. \tag{29}$$

Hence  $T\hat{\sigma}_{\lambda}^2 = \tilde{Y}^T W (W^T V_{\lambda} W)^{-1} W^T \tilde{Y} = \tilde{Y}^T W W^T \tilde{Y} - \lambda R_3$ , where  $R_3 = \tilde{Y}^T W \Delta_1 W^T \tilde{Y}$ . Together with equation (28), we conclude that equation (23) holds. 

Proposition 6.2. For any number of knots K, the following results hold.

Under 
$$H_0$$
 or  $H_{01}$ ,  $\lambda R_3/\sigma^2 = \sum_{k=1}^K \frac{\lambda \eta_{k,T}}{1+\lambda \eta_{k,T}} w_{k1}^2 \stackrel{D}{=} N_T(\lambda);$   
Under  $H_{02}$ ,  $\lambda R_3/\sigma^2 = \lambda \sum_{k=1}^K \frac{\eta_{k,T} + \bar{d}_0(KT^2)^{-1} \eta_{k,T}^2}{1+\lambda \eta_{k,T}} w_{k2}^2 \stackrel{D}{=} N_T'(\lambda);$ 

Under 
$$H_{02}$$
,  $\lambda R_3/\sigma^2 = \lambda \sum_{k=1}^K \frac{\eta_{k,T} + d_0(KT^2)^{-1} \eta_{k,T}^2}{1 + \lambda \eta_{k,T}} w_{k2}^2 \stackrel{D}{=} N_T'(\lambda);$ 

where  $N_T(\lambda)$  is defined in Theorem 4.1,  $N_T'(\lambda) =: \lambda \sum_{k=1}^K \frac{\eta_{k,T} + \bar{d}_0(KT^2)^{-1} \eta_{k,T}^2}{1 + \lambda \eta_{k,T}} w_k^2$  with  $w_k \stackrel{D}{=}$ iidN(0,1), and  $\eta_{k,T}$ 's are defined right before Theorem 4.1.

Proof of Proposition 6.2: Under  $H_0$  or  $H_{01}$ , the spline coefficients are all 0. Therefore,  $\tilde{Y} \stackrel{D}{=} N(\tilde{A}_1 \gamma, \sigma^2 I_T)$ . Recall that the matrix W satisfies equation (24). Hence

$$W^T \tilde{A}_1 = (W^T W) W^T \tilde{A}_1 = W^T \{ I_T - \tilde{A}_1 (\tilde{A}_1^T \tilde{A}_1)^{-1} \tilde{A}_1^T \} \tilde{A}_1 = \mathbf{0}_{(T-2) \times 2},$$

and  $W^T \tilde{Y} / \sigma \stackrel{D}{=} N(0_{T-2}, I_{T-2})$ , where  $0_{T-2}$  is the  $(T-2) \times 1$  vector whose components are 0. Recall that the matrix  $\Delta_1$  is defined as in equation (29). Suppose its eigen decomposition is  $U_1S_1U_1$ , where  $S_1$  is the diagonal matrix whose (i,i)th element  $\phi_{i1}$  is also the ith eigenvalue of  $\Delta_1$ . Then

$$\frac{R_3}{\sigma^2} = (\frac{W^T \tilde{Y}}{\sigma})^T \Delta_1(\frac{W^T \tilde{Y}}{\sigma}) = (\frac{U_1^T W^T \tilde{Y}}{\sigma})^T S_1(\frac{U_1^T W^T \tilde{Y}}{\sigma}) = \sum_{i=1}^{T-2} \phi_{i1} w_{i1}^2, \tag{30}$$

where  $w_{i1}$ 's are the elements of  $U_1^T W^T \tilde{Y} / \sigma$ .

Now we show that the eigenvalues  $\phi_{i1}$ 's satisfy  $\phi_{i1} = \eta_{i,T}/(1 + \lambda \eta_{i,T})$  for  $= 1, \dots, K$  and  $\phi_{i1} = 0$  for i > K. Note that for any  $p_1 \times p_2$  matrices A and B, AB and BA shares the same nonzero eigenvalues. Hence the nonzero eigenvalues of  $\Delta_1$  are the same as those of the matrix  $(I_K + \lambda \Lambda^{-1/2} \tilde{A}_2^T W W^T \tilde{A}_2 \Lambda^{-1/2})^{-1} \Lambda^{-1/2} \tilde{A}_2^T W W^T \tilde{A}_2 \Lambda^{-1/2}$ . Suppose the eigen decomposition of the matrix  $\Lambda^{-1/2} \tilde{A}_2^T W W^T \tilde{A}_2 \Lambda^{-1/2}$  is  $U_2 S_2 U_2^T$ , where  $S_2$  is the  $K \times K$  diagonal matrix whose (k, k)th element is the kth eigenvalue  $\eta_{k,T}$ . Then

$$(I_K + \lambda \Lambda^{-1/2} \tilde{A}_2^T W W^T \tilde{A}_2 \Lambda^{-1/2})^{-1} \Lambda^{-1/2} \tilde{A}_2^T W W^T \tilde{A}_2 \Lambda^{-1/2}$$

$$= (U_2 U_2^T + \lambda U_2 S_2 U_2^T)^{-1} U_2 S_2 U_2^T$$

$$= U_2 \{ (I_K + \lambda S_2)^{-1} S_2 \} U_2^T.$$

Note that the (k, k)th element in the diagonal matrix  $(I_K + \lambda S_2)^{-1})S_2$  is  $\eta_{k,T}/(1 + \lambda \eta_{k,T})$ , which equals the first K eigenvalues  $\phi_{k1}$ 's. Since the rest T-2-K eigenvalues are 0, we have  $R_3/\sigma^2 = \sum_{k=1}^K \frac{\eta_{k,T}}{1+\lambda \eta_{k,T}} w_{k1}^2$ . Note that  $W^T \tilde{Y}/\sigma \stackrel{D}{=} N(0_{T-2}, I_{T-2})$ . Moreover,  $U_1$  is an orthogonal matrix. Hence we conclude that  $U_1^T W^T \tilde{Y}/\sigma \stackrel{D}{=} N(0_{T-2}, I_{T-2})$ , i.e.  $w_{k1} \stackrel{D}{=} iidN(0,1)$ . Therefore,  $\lambda R_3/\sigma^2 \stackrel{D}{=} N_T(\lambda)$ .

Under  $H_{02}$ ,  $\theta(z) = \theta_0 + T^{-4/5}\theta_1(z)$ , where  $\theta_1(z)$  has an associated  $\bar{\lambda}_0$  satisfying  $\bar{\lambda}_0 = \bar{d}_0 T^{-2/5} K^{-1}$  for some positive constant  $\bar{d}_0$ . Let  $\bar{b}$  be the spline coefficients of  $\theta_1(z)$ . Then  $\text{var}(T^{-4/5}\bar{b}) = T^{-8/5}\bar{d}_0 T^{-2/5} K^{-1} \sigma^2 \Lambda^{-1} = \bar{d}_0 (KT^2)^{-1} \sigma^2 \Lambda^{-1}$ . Since  $W^T \tilde{A}_1 = \mathbf{0}_{(T-2)\times 2}$ , we have  $W^T \tilde{Y} = T^{-4/5} W^T \tilde{A}_2 \bar{b} + W^T \tilde{u}$ . Equivalently, we could write that  $W^T \tilde{Y}/\sigma \stackrel{D}{=} N \left(0_{T-2}, I_{T-2} + \bar{d}_0 (KT^2)^{-1} W^T \tilde{A}_2 \Lambda^{-1} \tilde{A}_2^T W\right)$ . Denote

$$\Delta_2 = \{I_{T-2} + \bar{d}_0(KT^2)^{-1}W^T\tilde{A}_2\Lambda^{-1}\tilde{A}_2^TW\}^{1/2}\Delta_1\{I_{T-2} + \bar{d}_0(KT^2)^{-1}W^T\tilde{A}_2\Lambda^{-1}\tilde{A}_2^TW\}^{1/2}.$$

Similar as equation (30), we conclude that  $\frac{R_3}{\sigma^2} = \sum_{i=1}^{T-2} \phi_{i2} w_{i2}^2$ , where  $\phi_{i2}$  is the *i*th eigenvalue of  $\Delta_2$  and  $w_{i2}$  is the *i*th element of  $U_1^T \{I_{T-2} + \bar{d}_0 (KT^2)^{-1} W^T \tilde{A}_2 \Lambda^{-1} \tilde{A}_2^T W\}^{-1/2} W^T \tilde{Y} / \sigma$ . Note the nonzero eigenvalues of  $\Delta_2$  are the same as those of

$$(I_K + \lambda \Lambda^{-1} \tilde{A}_2^T W W^T \tilde{A}_2 \Lambda^{-1/2})^{-1} \Lambda^{-1/2} \tilde{A}_2^T W \{I_{T-2} + \bar{d}_0 (KT^2)^{-1} W^T \tilde{A}_2 \Lambda^{-1} \tilde{A}_2^T W\} W^T \tilde{A}_2 \Lambda^{-1/2}.$$

With the eigen decomposition  $\Lambda^{-1/2}\tilde{A}_2^TWW^T\tilde{A}_2\Lambda^{-1/2}=U_2S_2U_2^T,$  the above matrix equals

$$(I_K + \lambda U_2 S_2 U_2^T)^{-1} \{ U_2 S_2 U_2^T + \bar{d}_0 (KT^2)^{-1} U_2 S_2 U_2^T U_2 S_2 U_2^T \}$$

$$= U_2 (I_K + \lambda S_2)^{-1} \{ S_2 + \bar{d}_0 (KT^2)^{-1} S_2^2 \} U_2^T.$$

Since the (k,k)th element in the diagonal matrix  $(I_K + \lambda S_2)^{-1} \{ S_2 + \bar{d}_0 (KT^2)^{-1} S_2^2 \}$  is  $\frac{\eta_{k,T} + \bar{d}_0 (KT^2)^{-1} \eta_{k,T}^2}{1 + \eta_{k,T}}$ , we have  $\frac{R_3}{\sigma^2} = \sum_{k=1}^K \frac{\eta_{k,T} + \bar{d}_0 (KT^2)^{-1} \eta_{k,T}^2}{1 + \lambda \eta_{k,T}} w_{k2}^2$ . Note that  $w_{k2} \stackrel{D}{=} iidN(0,1)$  because  $var[U_1^T \{ I_{T-2} + \bar{d}_0 (KT^2)^{-1} W^T \tilde{A}_2 \Lambda^{-1} \tilde{A}_2^T W \}^{-1/2} W^T \tilde{Y}] = \sigma^2 I_{T-2}$  under  $H_{02}$ . Therefore, Proposition 6.2 is proved.

Proposition 6.3. Assume K = o(T). Then under  $H_0$ ,  $H_{01}$  or  $H_{02}$ ,  $\tilde{Y}^T W W^T \tilde{Y} / T \Rightarrow \sigma$ .

Proof of Proposition 6.3: Under  $H_0$  or  $H_{01}$ ,  $W^T \tilde{Y} / \sigma \stackrel{D}{=} N(0_{T-2}, I_{T-2})$ . Therefore,

$$\frac{\tilde{Y}^T W W^T \tilde{Y}}{T \sigma^2} = \frac{T - 2}{T} + o_p(1) = 1 + o_p(1).$$

Under  $H_{02} W^T \tilde{Y} / \sigma \stackrel{D}{=} N \left( 0_{T-2}, I_{T-2} + \bar{d}_0 (KT^2)^{-1} W^T \tilde{A}_2 \Lambda^{-1} \tilde{A}_2^T W \right)$ . Therefore,

$$\tilde{Y}^T W W^T \tilde{Y} / (T\sigma^2) = T^{-1} \sum_{i=1}^{T-2} \phi'_{i2} w_{i2}^2 = T^{-1} \sum_{i=1}^{T-2} \phi'_{i2} + o_p(1), \tag{31}$$

where  $\phi'_{i2}$  is the *i*th eigenvalue of the matrix  $I_{T-2} + \bar{d}_0(KT^2)^{-1}W^T\tilde{A}_2\Lambda^{-1}\tilde{A}_2^TW$  and  $w_{i2}$  is defined as in the proof of Proposition 6.2 and satisfying  $w_{i2} \stackrel{D}{=} iidN(0,1)$ . Since  $W^T\tilde{A}_2\Lambda^{-1}\tilde{A}_2^TW$  and  $\Lambda^{-1/2}\tilde{A}_2^TWW^T\tilde{A}_2\Lambda^{-1/2}$  share the same nonzero eigenvalues,

$$T^{-1}\sum_{i=1}^{T-2}\phi'_{i2} = \frac{T-2}{T} + T^{-1}\sum_{i=1}^{K}\bar{d}_0(KT^2)^{-1}\eta_{i,T} = 1 + o_p(1).$$

With equation (31),  $\tilde{Y}^T W W^T \tilde{Y} / T \Rightarrow \sigma^2$ . Hence, Proposition 6.3 is proved.

Proposition 6.4. Let  $\tilde{A}^c_{01}$  be defined right before equation (33). It holds that

$$2\log L(0) - 2\log L^{0} = T\log\{1 + \frac{\tilde{Y}^{T}\tilde{A}_{01}^{c}(\tilde{A}_{01}^{c})^{T}\tilde{Y}/\sigma^{2}}{\tilde{Y}^{T}WW^{T}\tilde{Y}/\sigma^{2}}\}.$$
 (32)

Proof of Proposition 6.4: By equation (27),  $2 \log L(0) = T \log(\tilde{Y}^T W W^T \tilde{Y}/\sigma^2) + C_0$ , where W is the  $T \times (T-2)$  matrix satisfying equation (24). Now we show that there exists a  $T \times (T-1)$  matrix  $W_0$  such that  $2 \log L^0 = T \log(\tilde{Y}^T W_0 W_0^T \tilde{Y}/\sigma^2) + C_0$ .

Partition  $A_1$  as  $A_1 = (A_{01}, A_{r1})$ , where  $A_{01}$  is the  $T \times 1$  vector whose elements are  $x_t$ 's and  $A_{r1}$  is the  $T \times 1$  vector whose elements are  $x_t z_t$ 's. Define  $\tilde{A}_{01} = \Sigma_u^{-1/2} A_{01}$  and  $\tilde{A}_{r1} = \Sigma_u^{-1/2} A_{r1}$ . Note that the maximum likelihood estimates associated with  $\log L^0$  satisfies that

$$\hat{\gamma}_0 = (\tilde{A}_{01}^T \tilde{A}_{01})^{-1} \tilde{A}_{01}^T \tilde{Y}, \qquad \hat{\sigma}_0^2 = T^{-1} (\tilde{Y} - \tilde{A}_{01} \hat{\gamma}_0)^T (\tilde{Y} - \tilde{A}_{01} \hat{\gamma}_0).$$

Similar as equation (24), there exists  $W_0$  such that  $(\tilde{Y} - \tilde{A}_{01}\hat{\gamma}_0)^T(\tilde{Y} - \tilde{A}_{01}\hat{\gamma}_0) = \tilde{Y}^T W_0 W_0^T \tilde{Y}$ ,  $W_0^T W_0 = I_{T-1}$  and  $W_0 W_0^T = P_0 =: I_T - \tilde{A}_{01} (\tilde{A}_{01}^T \tilde{A}_{01})^{-1} \tilde{A}_{01}^T$ . Thus

$$-2\log L^{0} = T\log(\hat{\sigma}_{0}^{2}) + (\tilde{Y} - \tilde{A}_{01}\hat{\gamma}_{0})^{T}(\tilde{Y} - \tilde{A}_{01}\hat{\gamma}_{0}) + T\log(2\pi) = T\log(\frac{\tilde{Y}^{T}W_{0}W_{0}^{T}\tilde{Y}}{\sigma^{2}}) + C_{0},$$

and  $2 \log L(0) - 2 \log L^0 = T \log \{1 + \frac{(\tilde{Y}^T W_0 W_0^T \tilde{Y} - \tilde{Y}^T W W^T \tilde{Y})/\sigma^2}{\tilde{Y}^T W W^T \tilde{Y}/\sigma^2} \}$ . We could project  $\tilde{A}_{r1}$  onto the unit direction  $\frac{\tilde{A}_{01}}{||\tilde{A}_{01}||}$  and the unit direction orthogonal to  $\tilde{A}_{01}$ , i.e.

$$\tilde{A}_{r1} = \pi_1 \frac{\tilde{A}_{01}}{||\tilde{A}_{01}||} + \pi_2 \tilde{A}_{01}^c. \tag{33}$$

By standard linear algebra,  $W_0W_0^T - WW^T = \tilde{A}_{01}^c(\tilde{A}_{01}^c)^T$ . Hence equation (32) holds.  $\square$ 

Now we study Theorem 4.1.

Proof of Theorem 4.1: By Proposition 6.2,  $\lambda R_3/\sigma^2 = \sum_{k=1}^K \frac{\lambda \eta_{k,T}}{1+\lambda \eta_{k,T}} w_{k1}^2 \stackrel{D}{=} N_T(\lambda)$ . Note that  $\tilde{Y}^T W W^T \tilde{Y}/\sigma^2 = \sum_{k=1}^{T-2} w_{k1}^2$ . Hence

$$\frac{T\hat{\sigma}_{\lambda}^{2}}{\sigma^{2}} = \frac{\tilde{Y}^{T}WW^{T}\tilde{Y} - \lambda R_{3}}{\sigma^{2}} = \sum_{k=1}^{K} \frac{w_{k1}^{2}}{1 + \lambda \eta_{k,T}} + \sum_{k=K+1}^{T-2} w_{k1}^{2} \stackrel{D}{=} D_{T}(\lambda).$$

Together with Proposition 6.1, we have

$$\sup_{\lambda \ge 0} \{ 2 \log L(\lambda) - 2 \log L(0) \} \stackrel{D}{=} \sup_{\lambda \ge 0} [T \log \{ 1 + N_T(\lambda) / D_T(\lambda) \} - \sum_{k=1}^K \log (1 + \lambda \xi_{k,T})]. \tag{34}$$

Denote  $(\tilde{A}_{01}^c)^T \tilde{Y} / \sigma = w_{T-1}$ . Proposition 6.4 yields

$$2\log L(0) - 2\log L^0 = T\log\{1 + \frac{\tilde{Y}^T\tilde{A}_{01}^c(\tilde{A}_{01}^c)^T\tilde{Y}/\sigma^2}{\tilde{Y}^TWW^T\tilde{Y}/\sigma^2}\} = T\log\{1 + \frac{w_{T-1}^2}{\sum_{k=1}^{T-2}w_{k1}^2}\}.$$

It remains to show that under  $H_0$ ,  $w_{T-1} \stackrel{D}{=} N(0,1)$  and it is independent of  $w_{k1}$ 's, or equivalently, the vector  $U_1^T W^T \tilde{Y}$ . Under  $H_0$ ,  $\tilde{Y} = \tilde{A}_{1\gamma} = \tilde{A}_{01\gamma_0} + \tilde{u}$ . Recall that  $\tilde{A}_{01}^c$  is the unit direction that is orthogonal to both  $\tilde{A}_{01}$  and W. Since  $(\tilde{A}_{01}^c)^T \tilde{Y} = (\tilde{A}_{01}^c)^T \tilde{u} \stackrel{D}{=} N(0, \sigma^2)$  and

$$cov\{(\tilde{A}_{01}^c)^T\tilde{Y}, U_1^TW^T\tilde{Y}\} = \sigma^2(\tilde{A}_{01}^c)^TWU_1 = \sigma^20_{T-2}U_1 = 0_{T-2},$$

we conclude Theorem 4.1 holds.

Now we study Theorem 4.2.

Proof of Theorem 4.2: First we show equation (12) holds. Recall that  $\xi_{s,T}$  and  $\eta_{s,T}$  are the sth eigenvalues of  $\Lambda^{-1/2}A_2^T\Sigma_u^{-1}A_2\Lambda^{-1/2}$  and  $\Lambda^{-1/2}A_2^T\Sigma_u^{-1/2}P\Sigma_u^{-1/2}A_2\Lambda^{-1/2}$  respectively. Note that  $\Lambda^{-1/2}$  and  $\Sigma_u^{-1/2}$  are bounded deterministic matrices. By continuous mapping theorem, it suffices to show that  $T^{-2}A_2^TA_2$  and  $T^{-2}A_2^TP'A_2$  converge, where  $P' = I_T - A_2(A_2^TA_2)^{-1}A_2^T$ . Define that  $\zeta_t = E(z_t^t)$ ,  $\zeta_t(\kappa_i) = E\{z_t^t(z_t - \kappa_i)_+\}$ , and  $\zeta(\kappa_i, \kappa_j) = E\{(z_t - \kappa_i)_+(z_t - \kappa_j)_+\}$ , for l = 0 or 1,  $i = 1, \dots, K$  and  $j = 1, \dots, K$ . Let  $\Pi_1$  be the  $2 \times 2$  matrices with (i, j)-th element  $\zeta_{i+j-2}$ ,  $\Pi_2$  be the  $K \times K$  matrix with (i, j)-th element  $\zeta(\kappa_i, \kappa_j)$  and  $\Pi_3$  be the  $2 \times K$  matrix with (i, j)-th element  $\zeta_{i-1}(\kappa_j)$ . We first show that

$$T^{-2}A_1^T A_1 \Rightarrow \Pi_1 \int B_v^2(s) ds, \tag{35}$$

$$T^{-2}A_2^T A_2 \Rightarrow \Pi_2 \int B_v^2(s) ds, \tag{36}$$

$$T^{-2}A_1^T A_2 \Rightarrow \Pi_3 \int B_v^2(s) ds. \tag{37}$$

Take the proof of equation (35) as an example. Note that the (i, j)th element of  $A_1^T A_1$  satisfying  $(A_1^T A_1)_{i,j} = \sum x_t^2 z_t^{i+j-2}$ . By subtracting and adding the mean, we have

$$\begin{split} T^{-2} \sum x_t^2 z_t^{i+j-2} &= T^{-2} \sum x_t^2 \left( z_t^{i+j-2} - E(z_t^{i+j-2}) \right) + T^{-2} E(z_t^{i+j-2}) \sum x_t^2 \\ &= o_p(1) + T^{-2} E(z_t^{i+j-2}) \sum x_t^2. \end{split}$$

Hence  $T^{-2}(A_1^T A_1)_{i,j} = T^{-2} \sum x_t^2 z_t^{i+j-2} \Rightarrow \zeta_{i+j-2} \int B_v^2(s) ds = (\Pi_1)_{i,j} \int B_v^2(s) ds$ . Similarly,

we could show that equation (36) and (37) are true. Moreover,

$$T^{-2}A_2^T P' A_2 = T^{-2}A_2^T A_2 - T^{-2}A_2^T A_1 (T^{-2}A_1^T A_1)^{-1}T^{-2}A_1^T A_2$$
  
$$\Rightarrow (\Pi_2 - \Pi_3^T \Pi_1^{-1}\Pi_3) \int B_v^2(s) ds.$$

Therefore, we conclude that equation (12) is true.

Next we prove that equation (13) is valid. Recall that equation (34) holds for any T and K. Let  $d = \lambda T^2$ . Then we have

$$\sup_{\lambda \ge 0} \{ 2 \log L(\lambda) - 2 \log L(0) \} \stackrel{D}{=} \sup_{d \ge 0} [T \log \{ 1 + \frac{N_T(dT^{-2})}{D_T(dT^{-2})} \} - \sum_{k=1}^K \log (1 + dT^{-2}\xi_{k,T}) ].$$
 (38)

Define the right hand side in equation (38) as  $\sup_{d>0} f_T(d)$ . We want to show that

$$\sup_{d \ge 0} f_T(d) \Rightarrow \sup_{d \ge 0} f(d) =: \sup_{d \ge 0} \{ \sum_{s=1}^K \frac{d\eta_s}{1 + d\eta_s} w_s^2 - \sum_{s=1}^K \log(1 + d\xi_s) \}.$$
 (39)

We first establish the finite dimensional convergence of  $f_T(d)$ . Since  $T^{-2}\eta_{s,t} \Rightarrow \eta_s$ , we have  $N_T(dT^{-2})$  converges to  $\sum_{s=1}^K \frac{d\eta_s}{1+d\eta_s} w_s^2$  for every fixed d. By Proposition 6.3,

$$T^{-1}D_T(dT^{-2}) = T^{-1}\{\tilde{Y}^TWW^T\tilde{Y}/\sigma^2 - N_T(dT^{-2})\} = 1 + o_p(1).$$

Therefore,  $N_T(dT^{-2})/D_T(dT^{-2}) = O_p(T^{-1})$  and we have

$$T\log\{1 + \frac{N_T(dT^{-2})}{D_T(dT^{-2})}\} = T\left[\frac{N_T(dT^{-2})}{T\{1 + o_p(1)\}} + O_p(T^{-2})\right] = \sum_{s=1}^K \frac{d\eta_s}{1 + d\eta_s} w_s^2 + o_p(1).$$

Together with the fact  $T^{-2}\xi_{s,T} \Rightarrow \xi_s$ , we conclude  $f_T(d)$  converges to f(d) for every fixed d. Lemma 6.1 below shows that  $f_T(d)$  form a tight sequence and hence  $f_T(d)$  converges to f(d) weakly. Lemma 6.2 further shows that  $\sup_{d\geq 0} f_T(d) \Rightarrow \sup_{d\geq 0} f(d)$  by proving a continuous mapping theorem type results holds. Therefore, equation (38) holds.

By Proposition 6.4,  $2 \log L(0) - 2 \log L^0 = T \log(1 + \frac{w_{T-1}^2}{\tilde{Y}^T W W^T \tilde{Y} / \sigma^2})$ , where  $w_{T-1} = (\tilde{A}_{01}^c)^T \tilde{Y} / \sigma$ . Under  $H_0$ ,  $w_{T-1} \stackrel{D}{=} N(0,1)$ . Recall that  $W^T \tilde{Y} / \sigma \stackrel{D}{=} N(0_{T-2}, I_{T-2})$ . Hence

$$\log(1 + \frac{w_{T-1}^2}{\tilde{Y}^T W W^T \tilde{Y} / \sigma^2}) = \frac{w_{T-1}^2}{T \{1 + o_p(1)\}} + O_p(T^{-2}) = \frac{w_{T-1}^2}{T} + o_p(T^{-1}).$$

Therefore,  $2 \log L(0) - 2 \log L^0 = w_{T-1}^2 + o_p(1)$  and Theorem 4.2 is proved.

LEMMA 6.1. Under the assumptions of Theorem 4.2,  $f_T(d) \Rightarrow f(d)$ , where  $f_T(d)$  and f(d) are defined between equation (38) and equation (39).

Proof of Lemma 6.1: We have already established the finite dimensional convergence of  $f_T(d)$  to f(d). It suffices to show that  $f_T(d)$  form a tight sequence, i.e. for every  $\epsilon$  and  $\eta$  strictly positive, there exist  $\delta$  and  $T_0$  such that for  $T \geq T_0$ ,

$$\delta^{-1} \mathbb{P}\{ \sup_{t \le u \le t+\delta} |f_T(u) - f_T(t)| \ge \epsilon \} \le \eta.$$

By the definition of  $f_T(\cdot)$ , we have  $|f_T(u) - f_T(t)| = T \log\{\frac{D_T(T^{-2}t)}{D_T(T^{-2}u)}\} + \sum_{s=1}^K \log\frac{1+tT^{-2}\xi_{s,T}}{1+uT^{-\alpha}\xi_{s,T}}$ . Since  $\log(1+x) < x$  for every x > 0, it holds that

$$\log\left\{\frac{D_T(T^{-2}t)}{D_T(T^{-2}u)}\right\} \le \frac{D_T(T^{-2}t) - D_T(T^{-2}u)}{D_T(T^{-2}u)} \le (u-t) \frac{\sum_{s=1}^K T^{-2} \eta_{s,T} w_s^2}{\sum_{s=K+1}^{T-2} w_s^2}.$$

Since  $T^{-2}\eta_{s,T} \Rightarrow \eta_s$ , there exists a constant  $C_1$  such that  $T^{-2}\eta_{s,T}w_s^2 \leq w_s^2$  for all s and T. Denote  $R_{K,T} = \frac{\sum_{s=1}^K w_s^2/K}{\sum_{s=K+1}^{T-2} w_s^2/(T-K-2)}$ . Then  $T\log\{\frac{D_T(T^{-2}t)}{D_T(T^{-2}u)}\} \leq (u-t)C_1KR_{K,T}$ . Since  $T^{-2}\xi_{s,T} \Rightarrow \xi_s$ , there exists a constant  $C_2$  such that

$$\sum_{s=1}^{K} \log(\frac{1 + uT^{-2}\xi_{s,T}}{1 + tT^{-2}\xi_{s,T}}) \le (u - t) \sum_{s=1}^{K} \xi_{s,T} T^{-2} \le (u - t) C_2 K.$$

Let  $C_3 = \max(C_1, C_2)$ . Then  $\mathbb{P}\{\sup_{t \leq u \leq t+\delta} |f_T(u) - f_T(t)| \geq \epsilon\} \leq \mathbb{P}(R_{K,T} \geq \frac{\epsilon}{C_3K\delta} - 1)$ , and it reduces to show the cumulative distribution function (c.d.f.)  $H_{K,T}$  of  $R_{K,T}$  satisfies

$$1 - H_{K,T}(\frac{\epsilon}{C_3 K \delta} - 1) \le \eta \delta. \tag{40}$$

Note that  $R_{K,T}$  follows the F-distribution with degrees of freedom K and T-2-K. For every x,  $\lim_{T\to\infty} H_{K,T}(x) = H_K(Kx)$ , where  $H_K$  is the c.d.f of  $\chi_K^2$  random variables. Using L'Hospital rules, we have  $\lim_{\delta\to 0+} \{1-H_K(\frac{\epsilon}{C_3\delta}-K)\}/\{\frac{\eta\delta}{2}\} = 0$ . Therefore, we could find  $\delta = \delta(\epsilon,\eta)$ , with  $\delta < 1$  and  $\frac{\epsilon}{C_\delta} - K > 0$ , such that  $1-H_K(\frac{\epsilon}{C_3\delta}-K) \leq \frac{\eta\delta}{2}$ . Because of the convergence of  $H_{K,T}(x)$  to  $H_K(Kx)$ , we could find  $T_0$  such that for  $T \geq T_0$ , it holds that

 $|H_{K,T}(\frac{\epsilon}{C_3K\delta}-1)-H_K(\frac{\epsilon}{C_3\delta}-K)| \leq \frac{\eta\delta}{2}$ . Thus equation (40) holds and we conclude that it converges to f(d) weakly.

LEMMA 6.2. Under the assumptions of Theorem 4.2,  $\sup_{d>0} f_T(d) \Rightarrow \sup_{d>0} f(d)$ .

Proof of Lemma 6.2: Lemma 6.1 shows that  $f_T(d)$  weakly converges to f(d). Similar as Crainiceanu and Ruppert (2004), we first find a random variable  $F_{K,T}$  such that

$$\sup_{d>0} f_T(d) = \max_{d \in [0, F_{K,T}]} f_T(d). \tag{41}$$

Note that  $f_T(0) = 0$  for all T. It suffices to find  $F_{K,T}$  such that  $f_T(d) < 0$  when  $d > F_{K,T}$ . Recall that  $\log(1+x) \le x$  when  $x \ge 0$ . By definition,  $D_T(dT^{-2}) \ge \sum_{s=K+1}^{T-2} w_s^2$ , and  $N_T(dT^{-2}) \le \sum_{s=1}^K w_s^2$  for all d. Hence

$$T\log\{1 + \frac{N_T(dT^{-2})}{D_T(dT^{-2})}\} \le T\frac{N_T(dT^{-2})}{D_T(dT^{-2})} \le T\frac{\sum_{s=1}^K w_s^2}{\sum_{s=K+1}^{T-2} w_s^2}.$$
(42)

Let  $m_0$  be the positive constant such that all nonzero  $\xi_{s,T}$ 's satisfy  $T^{-2}\xi_{s,T} \geq m_0$ . Then

$$-\sum_{s=1}^{K} \log(1 + T^{-2}\xi_{s,T}) \le -K \log(1 + dm_0). \tag{43}$$

With equations (42) and (43), we establish that

$$f_T(d) \le T \frac{\sum_{s=1}^K w_s^2}{\sum_{s=K+1}^{T-2} w_s^2} - K \log(1 + dm_0).$$

$$\tag{44}$$

Let  $F_{K,T} =: m_0^{-1} \{ \exp(\frac{T}{T-2-K}KR_{K,T}) - 1 \}$ . Mind that  $f_T(F_{K,T}) = 0$ . Since the right hand side of equation (44) is monotonic decreasing in d,  $F_{K,T}$  has the desired property (41).

For any fixed M > 0 and  $t \ge 0$ , we have  $\max_{d \in [0,M]} f_T(d) \le \sup_{d \ge 0} f_T(d)$ . Hence

$$\mathbb{P}\{\sup_{d>0} f_T(d) \le t\} \le \mathbb{P}\{\max_{d\in[0,M]} f_T(d) \le t\}.$$

Taking  $\limsup for T \to \infty$  and applying the Continuous Mapping Theorem,

$$\limsup_{T \to \infty} \mathbb{P}(\sup_{d > 0} f_T(d) \le t) \le \limsup_{T \to \infty} \mathbb{P}(\max_{d \in [0, M]} f_T(d) \le t) = \mathbb{P}(\max_{d \in [0, M]} f(d) \le t).$$

Using the fact that

$$\lim_{M \to \infty} \mathbb{P}\{\max_{d \in [0,M]} f(d) \le t\} = \mathbb{P}\{\sup_{d > 0} f(d) \le t\},\tag{45}$$

we have

$$\lim_{T \to \infty} \sup_{d \ge 0} \mathbb{P}(\sup_{d \ge 0} f_T(d) \le t) \le \mathbb{P}\{\sup_{d \ge 0} f(d) \le t\}. \tag{46}$$

Since  $\mathbb{P}(AB) \geq \mathbb{P}(A) - \mathbb{P}(B^c)$ ,

$$\begin{split} \mathbb{P}(\sup_{d \geq 0} f_T(d) \leq t) & \geq & \mathbb{P}(\sup_{d \geq 0} f_T(d) \leq t, F_{K,T} < M) \\ & = & \mathbb{P}(\max_{d \in [0,M]} f_T(d) \leq t, F_{K,T} < M) \\ & \geq & \mathbb{P}(\max_{d \in [0,M]} f_T(d) \leq t) - \mathbb{P}(F_{K,T} > M). \end{split}$$

Note that  $\mathbb{P}(F_{K,T} > M) \to \mathbb{P}(F_K > KM)$ , where  $F_K$  is a  $\chi_K^2$  distributed random variable. Taking  $\liminf_{T \to \infty} \mathbb{P}(\sup_{d \geq 0} f_T(d) \leq t) \geq \mathbb{P}(\max_{d \in [0,M]} f(d) \leq t) - \mathbb{P}(F_K \geq KM)$ . Using equation (45) and that  $\lim_{M \to \infty} \mathbb{P}(F_K \geq KM) = 0$ , we conclude

$$\liminf_{T \to \infty} \mathbb{P}(\sup_{d \ge 0} f_T(d) \le t) \ge \mathbb{P}\{\sup_{d \ge 0} f(d) \le t\}.$$

Together with equation (46), the limit of  $\mathbb{P}\{\sup_{d>0} f_T(d)\}$  exists and satisfying

$$\lim_{T \to \infty} \mathbb{P} \{ \sup_{d > 0} f_T(d) \le t \} = \mathbb{P} \{ \sup_{d > 0} f(d) \le t \}.$$

Therefore,  $\sup_{d\geq 0} f_T(d) \Rightarrow \sup_{d\geq 0} f(d)$ .

Now we study Theorem 4.3.

Proof of Theorem 4.3: First we show that equation (14) holds. Note that for any matrix A, its trace equals the sum of its eigenvalues. Since  $\Lambda^{-1/2}$  and  $\Sigma_u^{-1/2}$  are bounded deterministic matrices, it suffices to show both  $(KT^2)^{-1}tr(A_2^TA_2)$  and  $(KT^2)^{-1}tr(A_2^TP'A_2)$  converge, where P' is defined the same as in the proof of Theorem 4.2.

Take the term  $(KT^2)^{-1}tr(A_2^TA_2)$  as an example. Let  $\prod_2^* = \prod_2 \int B_v^2(s)ds$ . From Equation (36), the (i,j)-th element of  $T^{-2}A_2^TA_2$  satisfies  $T^{-2}(A_2^TA_2)_{ij} \Rightarrow (\prod_2^*)_{ij}$ . Note that

$$K^{-1} \sum_{i=1}^{K} (\Pi_2^*)_{ii} = K^{-1} \sum_{i=1}^{K} [E\{(z_t - \kappa_i)_+^2\} \int B_v^2(s) ds] \le C \int B_v^2(s) ds,$$

where the last inequality holds as  $E\{(z_t - \kappa_i)_+^2\}$  is bounded. Therefore,  $(KT^2)^{-1}tr(A_2^TA_2)$  converges. Similarly,  $(KT^2)^{-1}tr(A_2^TP'A_2)$  converges and thus equation (14) is true.

Next we prove equation (15). Let  $d = \lambda KT^2$ . By equation (34), we could conclude that  $\sup_{\lambda \geq 0} \{2 \log L(\lambda) - 2 \log L(0)\} \stackrel{D}{=} \sup_{d \geq 0} g_T(d)$ , where

$$\sup_{d\geq 0} g_T(d) =: \sup_{d\geq 0} \left[ T \log\left\{1 + \frac{N_T \left(d(KT^2)^{-1}\right)}{D_T \left(d(KT^2)^{-1}\right)} \right\} - \sum_{k=1}^K \log\left\{1 + d(KT^2)^{-1}\xi_{k,T}\right\} \right]. \tag{47}$$

Let  $g_1(d)$  and  $g_2(d)$  be continuous functions defined respectively as the following limits:

$$\bar{g}_1(d) = \lim_{T \to \infty} \sum_{k=1}^K \frac{dK^{-1}T^{-2}\eta_{k,T}}{1 + dK^{-1}T^{-2}\eta_{k,T}}, \qquad \bar{g}_2(d) = \lim_{T \to \infty} \sum_{k=1}^K \log(1 + dK^{-1}T^{-2}\xi_{k,T}). \tag{48}$$

Lemma 6.3 shows that  $\bar{g}_1(d)$  and  $\bar{g}_2(d)$  exist for every fixed d and established the finite dimensional convergence of  $g_T(d)$  to  $\bar{g}_1(d) - \bar{g}_2(d)$ . Similar as Lemma 6.1, we could show that  $g_T(d)$  form a tight sequence and hence  $g_T(d)$  converges to  $\bar{g}_1(d) - \bar{g}_2(d)$  weakly. Similar as Lemma 6.2, we could establish a continuous mapping theorem type result and conclude that

$$\sup_{d>0} g_T(d) \Rightarrow \sup_{d>0} \{ \bar{g}_1(d) - \bar{g}_2(d) \}, \tag{49}$$

Next we want to prove that

$$\sup_{d \ge 0} \{ \bar{g}_1(d) - \bar{g}_2(d) \} = 0, \tag{50}$$

To prove equation (50), we will show that  $\bar{g}_1(d) - \bar{g}_2(d) \leq 0$  for all d. Note that the first derivative of the partial sum induced by  $\bar{g}_1(d) - \bar{g}_2(d)$  satisfies

$$\left\{ \sum_{k=1}^{K} \frac{dK^{-1}T^{-2}\eta_{k,T}}{1 + dK^{-1}T^{-2}\eta_{k,T}} - \sum_{k=1}^{K} \log(1 + dK^{-1}T^{-2}\xi_{k,T}) \right\}^{(1)}$$

$$= \sum_{k=1}^{K} \frac{K^{-1}T^{-2}\eta_{k,T}}{(1 + dK^{-1}T^{-2}\eta_{k,T})^{2}} - \sum_{k=1}^{K} \frac{K^{-1}T^{-2}\xi_{k,T}}{1 + dK^{-1}T^{-2}\xi_{k,T}}$$

$$=: Q_{1,T} + Q_{2,T},$$

where

$$Q_{1,T} = \sum_{k=1}^{K} \frac{K^{-1}T^{-2}\eta_{k,T}}{(1 + dK^{-1}T^{-2}\eta_{k,T})^2} - \sum_{k=1}^{K} \frac{K^{-1}T^{-2}\eta_{k,T}}{1 + dK^{-1}T^{-2}\eta_{k,T}},$$

and

$$Q_{2,T} = \sum_{k=1}^{K} \frac{K^{-1}T^{-2}\eta_{k,T}}{1 + dK^{-1}T^{-2}\eta_{k,T}} - \sum_{k=1}^{K} \frac{K^{-1}T^{-2}\xi_{k,T}}{1 + dK^{-1}T^{-2}\xi_{k,T}}.$$

Since  $(1 + dK^{-1}T^{-2}\eta_{k,T}) \ge 1$ , we have  $Q_{1,T} \le 0$  for all T. Moreover,  $Q_{2,T} \le 0$  as we explained below. Recall that  $\xi_{k,T}$  and  $\eta_{k,T}$  are the k-th eigenvalues of  $\Lambda^{-1/2}\tilde{A}_2^T\tilde{A}_2\Lambda^{-1/2}$  and  $\Lambda^{-1/2}\tilde{A}_2^TP\tilde{A}_2\Lambda^{-1/2}$  respectively. Moreover,

$$\Lambda^{-1/2} \tilde{A}_2^T \tilde{A}_2 \Lambda^{-1/2} - \Lambda^{-1/2} \tilde{A}_2^T P \tilde{A}_2 \Lambda^{-1/2} = \Lambda^{-1/2} \tilde{A}_2^T \tilde{A}_1 (\tilde{A}_1^T \tilde{A}_1)^{-1} \tilde{A}_1^T \tilde{A}_2 \Lambda^{-1/2}$$

which is a semi-positive definite matrix. Hence  $\eta_{k,T} \leq \xi_{k,T}$  for all k and T. Since  $\frac{x}{1+dx}$  is an increasing function of x, we have

$$Q_{2,T} =: \sum_{k=1}^{K} \frac{K^{-1}T^{-2}\eta_{k,T}}{1 + dK^{-1}T^{-2}\eta_{k,T}} - \sum_{k=1}^{K} \frac{K^{-1}T^{-2}\xi_{k,T}}{1 + dK^{-1}T^{-2}\xi_{k,T}} \le 0.$$
 (51)

Because  $\bar{g}_1(d)$  and  $\bar{g}_2(d)$  are both absolutely summable, we could change the order between summation and derivative. Since  $Q_{1,T} \leq 0$  and  $Q_{2,T} \leq 0$ , we conclude that the first derivative of  $\bar{g}_1(d) - \bar{g}_2(d)$  satisfies

$$\bar{g}_1'(d) - \bar{g}_2'(d) = \lim_{T \to \infty} (Q_{1,T} + Q_{2,T}) \le 0.$$

Recall that  $\bar{g}_1(0) - \bar{g}_2(0) = 0$ . For  $d \geq 0$ ,  $\bar{g}_1(d) - \bar{g}_2(d) = 0 + \int_0^d \{\bar{g}_1'(x) - \bar{g}_2'(x)\} dx \leq 0$ . Therefore, equation (50) holds and thus  $\sup_{\lambda \geq 0} \{2 \log L(\lambda) - 2 \log L(0)\} = \sup_{d \geq 0} g_T(d) \Rightarrow 0$ . Similarly as in Theorem 4.2, we conclude that

$$2\log L(0) - 2\log L^0 = w_{T-1}^2 + o_p(1),$$

where  $w_{T-1} \stackrel{D}{=} N(0,1)$  under  $H_0$ . Therefore,  $LRT_T \Rightarrow \chi_1^2$  under  $H_0$ .

LEMMA 6.3. Assume the conditions in Theorem 4.3 and define  $g_T(d)$  as in equation (47). Then  $g_T(d)$  converges to  $\bar{g}_1(d) - \bar{g}_2(d)$  for every fixed d, where  $\bar{g}_1(d)$  and  $\bar{g}_2(d)$  are defined in equation (48). Proof of Lemma 6.3: First, we prove that  $\bar{g}_1(d)$  and  $\bar{g}_2(d)$  exist for any fixed  $d \geq 0$ . Since  $1 + x \leq \exp(x)$  for  $x \geq 0$ , we have

$$1 + d\sum_{k=1}^{K} K^{-1}T^{-2}\xi_{k,T} \le \prod_{k=1}^{K} (1 + dK^{-1}T^{-2}\xi_{k,T}) \le \exp\{d\sum_{k=1}^{K} K^{-1}T^{-2}\xi_{k,T}\}.$$

If  $d\sum_{k=1}^K K^{-1}T^{-2}\xi_{k,T}$  converges, so does  $\log\{\prod_{k=1}^K (1+dK^{-1}T^{-2}\xi_{k,T})\}$ . We have already proved equation (14), i.e.  $\sum_{k=1}^K K^{-1}T^{-2}\xi_{k,T} \Rightarrow \bar{\xi}$ . Hence  $\sum_{k=1}^K \log(1+dK^{-1}T^{-2}\xi_{k,T})$  converges and its limit  $\bar{g}_2(d)$  exists. Note that

$$0 \le \frac{dK^{-1}T^{-2}\eta_{k,T}}{1 + dK^{-1}T^{-2}\eta_{k,T}} \le dK^{-1}T^{-2}\eta_{k,T} \le dK^{-1}T^{-2}\xi_{k,T}.$$

Since  $\sum_{k=1}^K dK^{-1}T^{-2}\xi_{k,T}$  converges, the limit of  $\sum_{k=1}^K \frac{dK^{-1}T^{-2}\eta_{k,T}}{1+dK^{-1}T^{-2}\eta_{k,T}}$  exists and we could denote it as  $\bar{g}_1(d)$ .

Next we establish the finite dimensional convergence of  $g_T(d)$  to  $\bar{g}_1(d) - \bar{g}_2(d)$ . Lemma 6.4 below shows  $N_T(dK^{-1}T^{-2})$  converges to  $\bar{g}_1(d)$  for every fixed d. By Proposition 6.3,

$$T^{-1}D_T\left(d(KT^2)^{-1}\right) = T^{-1}\{\tilde{Y}^TWW^T\tilde{Y}/\sigma^2 - N_T\left(d(KT^2)^{-1}\right)\} = 1 + o_p(1).$$

Correspondingly  $N_T\left(d(KT^2)^{-1}\right)/D_T\left(d(KT^2)^{-1}\right)=O_p(T^{-1})$  and

$$T\log\{1 + \frac{N_T\left(d(KT^2)^{-1}\right)}{D_T\left(d(KT^2)^{-1}\right)}\} = T\left[\frac{N_T\left(d(KT^2)^{-1}\right)}{T\{1 + o_p(1)\}} + O_p(T^{-2})\right] = \bar{g}_1(d) + o_p(1).$$

With the fact that  $\sum_{k=1}^{K} \log\{1 + d(KT^2)^{-1}\xi_{k,T}\}$  converges to  $\bar{g}_2(d)$  for every fixed d, Lemma 6.3 holds.

Lemma 6.4. Assume the conditions in Theorem 4.3.

Then  $N_T(dK^{-1}T^{-2})$  converges to  $\bar{q}_1(d)$  for every fixed d.

Proof of Lemma 6.4: Notice that  $w_k \stackrel{D}{=} iidN(0,1)$ . It suffices to consider show that  $\sum_{k=1}^K \frac{dK^{-1}T^{-2}\eta_{k,T}}{1+dK^{-1}T^{-2}\eta_{k,T}} w_k^2 - \bar{g}_1(d) = \sum_{k=1}^K \frac{dK^{-1}T^{-2}\eta_{k,T}}{1+dK^{-1}T^{-2}\eta_{k,T}} (w_k^2 - 1) \text{ converges to 0 in finite dimension.}$  For every fixed d,

$$E\sum_{k=1}^{K} \frac{dK^{-1}T^{-2}\eta_{k,T}}{1 + dK^{-1}T^{-2}\eta_{k,T}} (w_k^2 - 1) = 0.$$

Moreover,

$$var\{\sum_{k=1}^{K}\frac{dK^{-1}T^{-2}\eta_{k,T}}{1+dK^{-1}T^{-2}\eta_{k,T}}(w_k^2-1)\}=2\sum_{k=1}^{K}(\frac{dK^{-1}T^{-2}\eta_{k,T}}{1+dK^{-1}T^{-2}\eta_{k,T}})^2=O(K^{-1})=o(1).$$

Hence Lemma 6.4 is valid.

Now we consider Theorem 4.4.

Proof of Theorem 4.4: First we consider the local alternative  $H_{01}$ . Note that all spline coefficients are 0 under  $H_{01}$ . Therefore, it still holds  $\sup_{\lambda \geq 0} \{2 \log L(\lambda) - 2 \log L(0)\} \Rightarrow 0$ . It suffices to show  $2 \log L(0) - 2 \log L^0$  converges to a noncentral  $\chi_1^2$  with parameter  $\bar{\gamma}_1 \bar{\pi}_2$ . As  $W^T \tilde{A}_1 = \mathbf{0}_{(T-2)\times 2}$ , we have  $W^T \tilde{Y} = W^T \tilde{u}$ . By the fact that  $W_0^T \tilde{A}_{r1} = W_0^T (\pi_1 \frac{\tilde{A}_{01}}{||\tilde{A}_{01}||} + \pi_2 \tilde{A}_{01}^c) = \pi_2 W_0^T \tilde{A}_{01}^c$ , it holds  $W_0^T \tilde{Y} = (T^{-1} \bar{\gamma}_1) \pi_2 W_0^T \tilde{A}_{01}^c + W_0^T \tilde{u}$ . Hence

$$\begin{split} &\frac{\tilde{Y}^T W_0 W_0^T \tilde{Y} - \tilde{Y}^T W W^T \tilde{Y}}{\sigma^2} \\ &= \frac{T^{-2} \bar{\gamma}_1^2 \pi_2^2 (\tilde{A}_{01}^c)^T W_0 W_0^T \tilde{A}_{01}^c}{\sigma^2} + \frac{2T^{-1} \bar{\gamma}_1 \pi_2 (\tilde{A}_{01}^c)^T W_0 W_0^T \tilde{u}}{\sigma^2} + \frac{\tilde{u}^T (W_0 W_0^T - W W^T) \tilde{u}}{\sigma^2} \\ &= \frac{T^{-2} \bar{\gamma}_1^2 \pi_2^2}{\sigma^2} + \frac{2T^{-1} \bar{\gamma}_1 \pi_2 (\tilde{A}_{01}^c)^T \tilde{u}}{\sigma^2} + \frac{\tilde{u}^T \tilde{A}_{01}^c (\tilde{A}_{01}^c)^T \tilde{u}}{\sigma^2} \\ &\stackrel{D}{=} (w_1 + T^{-1} \pi_2 \bar{\gamma}_1)^2, \end{split}$$

where  $w_1 = (\tilde{A}_{01}^c)^T \tilde{u}/\sigma \stackrel{D}{=} N(0,1)$ , and the second equation holds as any unit direction orthogonal to  $\tilde{A}_{01}$  is the eigenvector of the projection matrix  $W_0 W_0^T$ , i.e.  $(\tilde{A}_{01}^c)^T W_0 W_0^T = (\tilde{A}_{01}^c)^T$ . Now we show the limit of  $T^{-1}\pi_2$  exists. Note that

$$(T^{-1}\pi_2)^2 = (T^{-1}\pi_2)^2 (\tilde{A}_{01}^c)^T W_0 W_0^T \tilde{A}_{01}^c = T^{-2} \tilde{A}_{r1}^T W_0 W_0^T \tilde{A}_{r1} = T^{-2} \tilde{A}_{r1}^T P_0 \tilde{A}_{r1}.$$

Using the same technique as in Theorem 4.2, we conclude all elements in  $T^{-2}A_{r1}^TA_{r1}$  and  $T^{-2}A_{r1}^TP_0'A_{r1}$  converge, where  $P_0' = I_T - A_{01}(A_{01}^TA_{01})^{-1}A_{01}^T$ . Since  $\Sigma_u^{-1/2}$  and  $\Lambda^{-1/2}$  are bounded deterministic matrices,  $T^{-2}\tilde{A}_{r1}^T\tilde{A}_{r1}$  and  $T^{-2}\tilde{A}_{r1}^TP_0\tilde{A}_{r1}$  converge. Equivalently,  $T^{-1}\pi_2$  converges and we could denote its limit as  $\bar{\pi}_2$ . Hence

$$LRT_T = 2\log L(0) - 2\log L^0 + o_p(1) = (w_1 + \bar{\gamma}_1 \bar{\pi}_2)^2 + o_p(1), \tag{52}$$

i.e. the asymptotic distribution is a non-central chi-square distribution with parameter  $\bar{\gamma}_1\bar{\pi}_2$ .

Next consider  $H_{02}$ . Let  $d = \lambda KT^2$ . By Proposition 6.1 and Proposition 6.2, we have  $\sup_{\lambda \geq 0} \{2 \log L(\lambda) - 2 \log L(0)\} \stackrel{D}{=} \sup_{d \geq 0} h_T(d)$ , where  $h_T(d) =: T \log\{1 + \frac{N'_T \left(d(KT^2)^{-1}\right)}{D_T (d(KT^2)^{-1})}\} - \sum_{k=1}^K \log\{1 + d(KT^2)^{-1}\xi_{k,T}\}$ . We want to show

$$\sup_{d \ge 0} h_T(d) \Rightarrow \sup_{d \ge 0} h(d) =: \sup_{d \ge 0} \{ \bar{d}_0 \bar{\eta} + \bar{g}_3(d) - \bar{g}_2(d) \}, \tag{53}$$

where  $\bar{g}_2(d)$  and  $\bar{\eta}$  are defined as in Theorem 4.3 and  $\bar{g}_3(d)$  is the limit of  $\sum_{k=1}^K \frac{(d-\bar{d}_0)(KT^2)^{-1}\eta_{k,T}}{1+d(KT^2)^{-1}\eta_{k,T}}$ . Note that  $\bar{g}_3(d)$  exists because the term  $\sum_{k=1}^K \frac{(d-\bar{d}_0)(KT^2)^{-1}\eta_{k,T}}{1+d(KT^2)^{-1}\eta_{k,T}}$  is bounded by  $-\bar{d}_0\bar{\eta}$  and  $\bar{g}_1(d)$ .

To prove equation (53), we first establish the finite dimensional convergence of  $h_T(d)$  to h(d). For every fixed d, we could simplify  $N_T'\{d(KT^2)^{-1}\}$  as

$$N_T' \left\{ d(KT^2)^{-1} \right\} = \bar{d}_0 \sum_{k=1}^K (KT^2)^{-1} \eta_{k,T} w_k^2 + \sum_{k=1}^K \frac{(d - \bar{d}_0)(KT^2)^{-1} \eta_{k,T}}{1 + d(KT^2)^{-1} \eta_{k,T}} w_k^2.$$

Apply the same technique as Lemma 6.4, we could show that  $\sum_{k=1}^{K} (KT^2)^{-1} \eta_{k,T}(w_k^2 - 1) \Rightarrow 0$  and thus  $\sum_{k=1}^{K} (KT^2)^{-1} \eta_{k,T} w_k^2 \Rightarrow \bar{\eta}$ . Similarly,  $\sum_{k=1}^{K} \frac{(d-\bar{d}_0)(KT^2)^{-1} \eta_{k,T}}{1+d(KT^2)^{-1} \eta_{k,T}} w_k^2$  converges to  $\bar{g}_3(d)$  for every fixed d. Therefore,  $N_T' \left( d(KT^2)^{-1} \right)$  does so to  $\bar{d}_0 \bar{\eta} + \bar{g}_3(d)$ . Using Proposition 6.3,  $T^{-1}D_T \left( d(KT^2)^{-1} \right) = T^{-1} \{ \tilde{Y}^T W W^T \tilde{Y} / \sigma^2 - N_T' \left( d(KT^2)^{-1} \right) \} = 1 + o_p(1)$ . Correspondingly,  $N_T' \left( d(KT^2)^{-1} \right) / D_T \left( d(KT^2)^{-1} \right) = O_p(T^{-1})$ . Hence

$$T\log\{1 + \frac{N_T'\left(d(KT^2)^{-1}\right)}{D_T\left(d(KT^2)^{-1}\right)}\} = T\left[\frac{N_T'\left(d(KT^2)^{-1}\right)}{T\{1 + o_p(1)\}} + O_p(T^{-2})\right] = \bar{d}_0\bar{\eta} + \bar{g}_3(d) + o_p(1).$$

Since  $\sum_{k=1}^{K} \log\{1 + d(KT^2)^{-1}\xi_{k,T}\}$  converges to  $\bar{g}_2(d)$  for every fixed d, so does  $h_T(d)$  to h(d). Similar as Lemma 6.1 and Lemma 6.2, we could further show that  $h_T(d)$  weakly converges to h(d) and a continuous mapping theorem type results holds. Thus equation (53) holds.

Next we want to show

$$\sup_{d>0} h(d) = \max_{d\in[0,\bar{d}_0]} \{\bar{d}_0\bar{\eta} + \bar{g}_3(d) - \bar{g}_2(d)\}. \tag{54}$$

Note that the first derivative of the partial sum induced by  $\bar{g}_3(d) - \bar{g}_2(d)$  satisfies

$$\begin{split} l_T'(d) &= \sum_{k=1}^K \frac{\{1 + d_0(KT^2)^{-1}\eta_{k,T}\}(KT^2)^{-1}\eta_{k,T}}{\{1 + d(KT^2)^{-1}\eta_{k,T}\}^2} - \sum_{k=1}^K \frac{(KT^2)^{-1}\xi_{k,T}}{1 + d(KT^2)^{-1}\xi_{k,T}} \\ &\leq \sum_{k=1}^K \frac{(KT^2)^{-1}\eta_{k,T}}{1 + d(KT^2)^{-1}\eta_{k,T}} - \sum_{k=1}^K \frac{(KT^2)^{-1}\xi_{k,T}}{1 + d(KT^2)^{-1}\xi_{k,T}} \\ &= \sum_{k=1}^K \frac{(KT^2)^{-1}\eta_{k,T}}{1 + d(KT^2)^{-1}\eta_{k,T}} - \sum_{k=1}^K \frac{(KT^2)^{-1}\xi_{k,T}}{1 + d(KT^2)^{-1}\xi_{k,T}}, \end{split}$$

where the inequality holds when  $d \geq \bar{d}_0$ . By equation (51),

$$\sum_{k=1}^{K} \frac{(KT^2)^{-1} \eta_{k,T}}{1 + d(KT^2)^{-1} \eta_{k,T}} \le \sum_{k=1}^{K} \frac{(KT^2)^{-1} \xi_{k,T}}{1 + d(KT^2)^{-1} \xi_{k,T}}.$$

When  $d \geq \bar{d}_0$ ,  $l'_T(d) \leq 0$  and hence  $h(d) \leq h(\bar{d}_0)$  for  $d \geq \bar{d}_0$ . Thus equation (54) holds.

Finally, we will show that  $2 \log L(0) - 2 \log L^0 \Rightarrow (1 + \bar{d}_0 \varrho) w_1^2$ . Under  $H_{02}$ , we have that  $W^T \tilde{Y} = W^T (T^{-4/5} \tilde{A}_2 \bar{b} + \tilde{u})$  and  $W_0^T \tilde{Y} = W_0^T (T^{-4/5} \tilde{A}_2 \bar{b} + \tilde{u})$ . Hence

$$\frac{\tilde{Y}^T W_0 W_0^T \tilde{Y} - \tilde{Y}^T W W^T \tilde{Y}}{\sigma^2} = \frac{(T^{-4/5} \tilde{A}_2 \bar{b} + \tilde{u})^T \tilde{A}_{01}^c (\tilde{A}_{01}^c)^T (T^{-4/5} \tilde{A}_2 \bar{b} + \tilde{u})}{\sigma^2} \stackrel{D}{=} \sum_{i=1}^T \phi_{i4} w_i^2,$$

where  $w_i \stackrel{D}{=} iidN(0,1)$  and  $\phi_{i4}$  is the *i*th eigenvalue of

$$\Delta_4 = \{I_T + \bar{d}_0(KT^2)^{-1}\tilde{A}_2\Lambda^{-1}\tilde{A}_2^T\}^{1/2}\tilde{A}_{01}^c(\tilde{A}_{01}^c)^T\{I_T + \bar{d}_0(KT^2)^{-1}\tilde{A}_2\Lambda^{-1}\tilde{A}_2^T\}^{1/2}.$$

Note that  $\Delta_4$  share the same nonzero eigenvalue as

$$(\tilde{A}_{01}^c)^T \{ I_T + \bar{d}_0 (KT^2)^{-1} \tilde{A}_2 \Lambda^{-1} \tilde{A}_2^T \} \tilde{A}_{01}^c = 1 + \bar{d}_0 \frac{(\tilde{A}_{01}^c)^T \tilde{A}_2 \Lambda^{-1} \tilde{A}_2^T \tilde{A}_{01}^c}{KT^2}.$$

Moreover, we could show that the limit of  $\frac{(\tilde{A}_{01}^c)^T\tilde{A}_2\Lambda^{-1}\tilde{A}_2^T\tilde{A}_{01}^c}{KT^2}$  exists. Using the same technique as in Theorem 4.2, we conclude that each element in  $T^{-2}A_2^TA_{01}(A_{01}^c)^TA_2$  converges. Since  $\Sigma_u^{-1/2}$  and  $\Lambda^{-1/2}$  are bounded deterministic matrices, each element of the matrix  $T^{-2}\Lambda^{-1/2}\tilde{A}_2^T\tilde{A}_{01}^c(\tilde{A}_{01}^c)^T\tilde{A}_2\Lambda^{-1/2}$  converges and  $K^{-1}T^{-2}tr\{\Lambda^{-1/2}\tilde{A}_2^T\tilde{A}_{01}^c(\tilde{A}_{01}^c)^T\tilde{A}_2\Lambda^{-1/2}\}$  converges. By the fact that  $(\tilde{A}_{01}^c)^T\tilde{A}_2\Lambda^{-1}\tilde{A}_2^c\tilde{A}_{01}^c=tr\{\Lambda^{-1/2}\tilde{A}_2^T\tilde{A}_{01}^c(\tilde{A}_{01}^c)^T\tilde{A}_2\Lambda^{-1/2}\}$ , the limit of  $(KT^2)^{-1}(\tilde{A}_{01}^c)^T\tilde{A}_2\Lambda^{-1}\tilde{A}_2^T\tilde{A}_{01}^c$  exists and we could denote it as  $\varrho$ . Therefore,

$$\frac{\tilde{Y}^T W_0 W_0^T \tilde{Y} - \tilde{Y}^T W W^T \tilde{Y}}{\sigma^2} \stackrel{D}{=} (1 + \bar{d}_0 \varrho) w_1^2, \tag{55}$$

and we could further conclude  $2 \log L(0) - 2 \log L^0 \Rightarrow (1 + \bar{d}_0 \varrho) w_1^2$ . Theorem 4.4 is proved.  $\square$ 

## Appendix B: Tables

Table 1: The empirical size of the proposed LRT when  $\Sigma_u$  is known

	Finite dist.						Asymptotic dist.				
	20%	15%	10%	5%	1%	20%	15%	10%	5%	1%	
Panel A1: $T = 100,  \rho = 0$											
K = 10	0.1915	0.1445	0.0930	0.0430	0.0095	0.1995	0.1490	0.0995	0.0465	0.0100	
K = 20	0.1940	0.1425	0.0935	0.0425	0.0100	0.1985	0.1480	0.0990	0.0465	0.0100	
K = 40	0.1920	0.1425	0.0925	0.0415	0.0085	0.1985	0.1480	0.0990	0.0465	0.0100	
Panel A2: $T = 100,  \rho = 0.5$											
K = 10	0.1940	0.1435	0.0975	0.0550	0.0095	0.1975	0.1495	0.1025	0.0585	0.0105	
K = 20	0.1930	0.1440	0.0990	0.0550	0.0095	0.1970	0.1495	0.1035	0.0580	0.0105	
K = 40	0.1945	0.1435	0.0975	0.0555	0.0090	0.1985	0.1500	0.1035	0.0585	0.0105	
Panel B1: $T = 300,  \rho = 0$											
K = 20	0.2035	0.1560	0.0975	0.0540	0.0160	0.2050	0.1565	0.0995	0.0535	0.0165	
K = 40	0.2035	0.1550	0.0990	0.0530	0.0160	0.2050	0.1565	0.0995	0.0535	0.0165	
K = 80	0.2025	0.1560	0.0980	0.0530	0.0160	0.2045	0.1560	0.0995	0.0535	0.0165	
Panel B2: $T = 300,  \rho = 0.5$											
K = 20	0.1985	0.1550	0.0975	0.0485	0.0105	0.1995	0.1555	0.1010	0.0495	0.0110	
K = 40	0.1985	0.1520	0.1000	0.0485	0.0120	0.1980	0.1550	0.1005	0.0490	0.0110	
K = 80	0.1975	0.1540	0.0975	0.0485	0.0105	0.1975	0.1550	0.1005	0.0495	0.0110	

Note: The model is  $y_t = \theta(z_t)x_t + u_t$  with  $x_t = x_{t-1} + v_t$  and  $u_t = \rho u_{t-1} + \varepsilon_t$ , where  $v_t$ 's and  $\varepsilon_t$ 's are  $iid\ N(0,1)$  and are independent with each other. The initial values are set to be zero. In particular,  $\theta(z_t) = 0.25$  and the true covariance matrix is used. The rejection frequencies are calculated based on 2000 replications.

Table 2: The empirical size of the proposed LRT when  $\Sigma_u$  is unknown

	Finite	dist.				Asymptotic dist.					
	20%	15%	10%	5%	1%	20%	15%	10%	5%	1%	
Panel A1	Panel A1: $T = 100$ , $\rho = 0$										
K = 10	0.2170	0.1595	0.1065	0.0500	0.0150	0.2245	0.1635	0.1120	0.0565	0.0155	
K = 20	0.2190	0.1610	0.1075	0.0530	0.0125	0.2245	0.1630	0.1115	0.0565	0.0150	
K = 40	0.2190	0.1585	0.1060	0.0530	0.0135	0.2245	0.1630	0.1115	0.0565	0.0150	
Panel A2	Panel A2: $T = 100,  \rho = 0.5$										
K = 10	0.2065	0.1565	0.1120	0.0635	0.0125	0.2110	0.1610	0.1155	0.0650	0.0145	
K = 20	0.2045	0.1545	0.1080	0.0630	0.0130	0.2105	0.1610	0.1135	0.0650	0.0145	
K = 40	0.2060	0.1555	0.1105	0.0625	0.0125	0.2105	0.1605	0.1135	0.0650	0.0145	
Panel B1: $T = 300,  \rho = 0$											
K = 20	0.2110	0.1640	0.1075	0.0575	0.0165	0.2125	0.1645	0.1075	0.0585	0.0175	
K = 40	0.2110	0.1635	0.1065	0.0570	0.0170	0.2125	0.1645	0.1075	0.0585	0.0175	
K = 80	0.2105	0.1635	0.1060	0.0565	0.0165	0.2120	0.1640	0.1075	0.0585	0.0175	
Panel B2: $T = 300,  \rho = 0.5$											
K = 20	0.2075	0.1550	0.1000	0.0500	0.0115	0.2090	0.1595	0.1025	0.0500	0.0115	
K = 40	0.2065	0.1555	0.0995	0.0480	0.0115	0.2075	0.1590	0.1020	0.0500	0.0110	
K = 80	0.2065	0.1565	0.1000	0.0485	0.0105	0.2070	0.1590	0.1020	0.0500	0.0110	

Note: The model is  $y_t = \theta(z_t)x_t + u_t$  with  $x_t = x_{t-1} + v_t$  and  $u_t = \rho u_{t-1} + \varepsilon_t$ , where  $v_t$ 's and  $\varepsilon_t$ 's are  $iid\ N(0,1)$  and are independent with each other. The initial values are set to be zero. In particular,  $\theta(z_t) = 0.25$ . The true covariance matrix is unknown and is replaced by its estimate. The rejection frequencies are calculated based on 2000 replications.

Table 3 The power of the proposed LRT

	Finite dist.	<u> </u>	Asymptotic dist.						
	$\Sigma_u$ is known	$\Sigma_u$ is unknown	$\Sigma_u$ is known	$\Sigma_u$ is unknown					
Panel A1: $T = 100,  \rho = 0$									
K = 10	0.6220 (0.5255)	0.6330 (0.5405)	0.6255 (0.5300)	$0.6365 \ (0.5455)$					
K = 20	0.6175 (0.5195)	$0.6275 \ (0.5325)$	0.6210 (0.5240)	$0.6330\ (0.5390)$					
K = 40	0.6140 (0.5150)	$0.6265 \ (0.5290)$	0.6195 (0.5205)	$0.6295 \ (0.5350)$					
Panel A2: $T = 100,  \rho = 0.5$									
K = 10	0.7085 (0.6050)	0.6930 (0.5905)	0.7125 (0.6110)	0.6975 (0.5960)					
K = 20	0.7050 (0.6010)	$0.6855 \ (0.5890)$	0.7080 (0.6055)	$0.6890\ (0.5925)$					
K = 40	0.7025 (0.5960)	$0.6855 \ (0.5875)$	0.7055 (0.6025)	$0.6875 \ (0.5900)$					
Panel B1	Panel B1: $T = 300,  \rho = 0$								
K = 20	0.9920 (0.9855)	0.9915 (0.9855)	0.9920 (0.9855)	0.9915 (0.9850)					
K = 40	0.9915 (0.9855))	0.9905 (0.9850)	0.9915 (0.9855)	0.9910 (0.9850)					
K = 80	0.9910 (0.9855)	0.9910 (0.9850)	0.9910 (0.9855)	0.9910 (0.9850)					
Panel B2: $T = 300,  \rho = 0$									
K = 20	0.994 (0.99)	0.9935 (0.9910)	0.994 (0.99)	0.9940 (0.9910)					
K = 40	0.994 (0.99)	0.9940 (0.9905)	0.994 (0.99)	0.9935 (0.9910)					
K = 80	0.994 (0.99)	0.9935 (0.9905)	0.994 (0.99)	0.9935 (0.9910)					

Note: The model is  $y_t = \theta(z_t)x_t + u_t$  with  $x_t = x_{t-1} + v_t$  and  $u_t = \rho u_{t-1} + \varepsilon_t$ , where  $v_t$ 's and  $\varepsilon_t$ 's are  $iid\ N(0,1)$  and are independent with each other. The initial values are set to be zero. In particular,  $\theta(z_t) = (z_t - 0.5)^2$ . The rejection frequencies are calculated using the critical values associated with  $\alpha = 0.05$  or  $\alpha = 0.01$  as indicated inside the parenthesis. All results are based on 2000 replications.

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