

Supplement Not for Online Publication to “Dynamic Local Average Treatment Effects in Time Series”

ALESSANDRO CASINI

University of Rome Tor Vergata

ADAM McCLOSKEY

University of Colorado at Boulder

LUCA ROLLA

University of Rome Tor Vergata

RAIMONDO PALA

University of Rome Tor Vergata

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Abstract

This supplemental material is structured as follows. Section [N.A](#) describes the collection of articles and specifications for the figure reported in the Introduction of the main article. Section [N.B](#) presents the derivations underlying the heteroskedasticity-based identification results in Section 3. Section [N.C](#) presents a GLS criterion for the estimation of the true sub-population. Section [N.D](#) includes the theoretical results about the estimators of the true sub-population proposed in Section 5. Section [N.E](#) presents additional results on identification of compliers under a continuous instrument, including illustrative examples, a test for full-sample identification failure under homoskedasticity, primitive conditions for the assumptions in Section 6, results on consistent covariance matrix estimation for the tests introduced in Section 6 and results about identification-robust inference under strong IV and local or fixed alternatives. Section [N.F](#) presents additional Monte Carlo simulations.

N.A Publication Selection Criterion

We select recent publications from the following five economics journals: American Economic Review, Econometrica, Journal of Political Economy, Quarterly Journal of Economics and Review of Economic Studies. We first identify articles published in these journals between January 2019 and December 2022 that contain the keyword “instrument” in their text. We then exclude articles that do not estimate linear instrumental variables (IV) models or that are based solely on cross-sectional data. This results in 18 articles, listed in Table [1](#). From these 18 articles, we collect all IV specifications reported in their main text. Articles 1–3, 6–7, 9–12, and 14–16 use time series data, while articles 4–5, 8, 13, and 17–18 use panel data. Since the F^* statistic requires time series data, for panel data applications we treat each cross-sectional unit separately. In cases where an application includes thousands of cross-sectional units, we select only a subset to prevent a single

panel application from dominating the distribution of the F statistics. The median cross-sectional size across applications is 34, so for panel specifications with more than 34 units we randomly select 34 units, while for those with fewer than 34 units we include all available units. When a specification involves multiple endogenous regressors, we run separate first-stage regressions for each endogenous regressor. This yields a total of 199 time series specifications and 1,361 panel data specifications.

Table 1: Selected Publications

Article ID	Year, Vol.(Issue)	Title	# of time series specifications	# of cross-sections in panel, # of specifications (# of cross-sections selected)
1	2019, AER	The Social Value of Financial Expertise	1	0
2	2020, AER	Turnover Liquidity and the Transmission of Monetary Policy	2	0
3	2021, AER	The Macroeconomic Effects of Oil Supply News: Evidence from OPEC Announcements	7	0
4	2021, AER	Stock Market Wealth and the Real Economy: A Local Labor Market Approach	0	2092, 2 (68)
5	2022, AER	Convex Supply Curves	0	21, 11 (231)
6	2022, ECMA	Monetary Policy, Redistribution, and Risk Premia	1	0
7	2022, JPE	Instrumental Variable Identification of Dynamic Variance Decompositions	2	0
8	2019, Restud	Innovation and Top Income Inequality	0	50, 8 (272)
9	2019, Restud	The Changing Returns to Crime: Do Criminals Respond to Prices?	8	0
10	2020, Restud	Mergers, Innovation, and Entry-Exit Dynamics: Consolidation of the Hard Disk Drive Industry, 1996–2016	2	0
11	2021, Restud	Bank Capital Redux: Solvency, Liquidity, and Crisis	0	15, 1 (15)
12	2022, Restud	Understanding the Size of the Government Spending Multiplier: It's in the Sign	120	0
13	2021, Restud	Housing Wealth Effects: The Long View	0	4400, 16 (544)
14	2022, Restud	Fiscal Multipliers and Foreign Holdings of Public Debt	48	0
15	2022, Restud	The Real Effects of Monetary Expansions: Evidence from a Large-scale Historical Experiment	5	6, 1 (6)
16	2020, Restud	Identifying Modern Macro Equations with Old Shocks	3	0
17	2022, QJE	Taxation and Innovation in the Twentieth Century	0	52, 1 (34)
18	2022, QJE	The Slope of the Phillips Curve: Evidence from U.S. States	0	34, 4 (191)

N.B Derivation of Heteroskedasticity-Based Identification from Section 3

Given the notation introduced in Section 3, identification can be shown analytically by first solving for the reduced-form of (3.1):

$$\tilde{Y}_t = \frac{1}{1 - a\beta_0} (\eta_t + \beta_0 e_t), \quad \tilde{D}_t = \frac{1}{1 - a\beta_0} (a\eta_t + e_t).$$

Let Σ_i denote the covariance matrix of $[\tilde{Y}_t, \tilde{D}_t]'$ in the subsample $i = P, C$. It follows that

$$\Sigma_i = \frac{1}{(1 - a\beta_0)^2} \begin{bmatrix} \sigma_{\eta,i}^2 + \beta_0^2 \sigma_{e,i}^2 & \beta_0 \sigma_{e,i}^2 + a \sigma_{\eta,i}^2 \\ \beta_0 \sigma_{e,i}^2 + a \sigma_{\eta,i}^2 & \sigma_{e,i}^2 + a^2 \sigma_{\eta,i}^2 \end{bmatrix}, \quad i = P, C.$$

It is typical in the literature to assume within-regime covariance-stationarity, i.e., $\mathbb{E}(e_t^2)$ and $\mathbb{E}(\eta_t^2)$ are constant within each subsample **P** and **C** which is, however, restrictive for economic time series. It turns out that this is not necessary for identification. Volatilities can be time-varying as long as the average volatilities $\sigma_{e,i}$ and $\sigma_{\eta,i}$ ($i = P, C$) satisfy (3.2).

When (3.1) is correctly specified, i.e., the true model is linear with stable parameters, the parameter β_0 can be identified using (3.2) by taking the difference between the covariance matrices in the policy and control samples:

$$\beta_0 = \frac{\Delta \Sigma^{(1,2)}}{\Delta \Sigma^{(2,2)}} = \frac{T_P^{-1} \sum_{t \in \mathbf{P}} \text{Cov}(\tilde{Y}_t, \tilde{D}_t) - T_C^{-1} \sum_{t \in \mathbf{C}} \text{Cov}(\tilde{Y}_t, \tilde{D}_t)}{T_P^{-1} \sum_{t \in \mathbf{P}} \text{Var}(\tilde{D}_t) - T_C^{-1} \sum_{t \in \mathbf{C}} \text{Var}(\tilde{D}_t)}, \quad \text{where} \quad (\text{N.1})$$

$$\Delta \Sigma \triangleq \Sigma_P - \Sigma_C = \frac{\sigma_{e,P}^2 - \sigma_{e,C}^2}{(1 - a\beta_0)^2} \begin{bmatrix} \beta_0^2 & \beta_0 \\ \beta_0 & 1 \end{bmatrix}.$$

To determine which average treatment effect this approach identifies in the general framework, we re-frame this problem in terms of instrumental variables as follows. Let $Z_t = 1$ for $t \in \mathbf{P}$ and $Z_t = 0$ for $t \in \mathbf{C}$. Multiply both sides of (3.1) by $\tilde{D}_t$ to yield $\tilde{D}_t \tilde{Y}_t = \beta_0 \tilde{D}_t^2 + \tilde{D}_t \eta_t$. We can use Z_t as an instrument for $\tilde{D}_t^2$. The first-stage is $\tilde{D}_t^2 = \theta Z_t + \varepsilon_t$, where ε_t is some error term satisfying $\varepsilon_t \geq -\theta Z_t$.

N.C GLS Estimator of $\mathbf{S}_{0,T}$

As a second estimator of the sub-population $\mathbf{S}_{0,T}$, we consider a GLS criterion that minimizes an efficiently weighted combination of the sum of squared residuals of both the structural equation (4.1) and reduced-form equation (4.2). That is, the system of equations (4.1)-(4.2) can be written in reduced-form as

$$\vec{y} = W(\mathbf{S}_{0,T})\xi + \varepsilon, \quad (\text{N.1})$$

where $\vec{y} = (Y', D')'$, $W(\mathbf{S}_{0,T}) = I_2 \otimes [C_{0,T}Z : X]$, $\xi = (\beta\theta', \gamma'_1 + \beta\gamma'_2, \theta', \gamma'_2)'$ and $\varepsilon = (u' + \beta e', e')'$ with $C_{0,T}$ defined as C_T but corresponding to the indices in $\mathbf{S}_{0,T}$. This is a system of two seemingly unrelated regressions. Let

$$\hat{\xi}_{FGLS}(\mathbf{S}_T) = (W(\mathbf{S}_T)' \hat{\Omega}_\varepsilon(\mathbf{S}_T)^{-1} W(\mathbf{S}_T))^{-1} W(\mathbf{S}_T)' \hat{\Omega}_\varepsilon(\mathbf{S}_T)^{-1} \vec{y},$$

denote a feasible GLS estimator of ξ , where $\hat{\Omega}_\varepsilon(\mathbf{S}_T)$ is a consistent estimator of $\mathbb{E}[\varepsilon\varepsilon'|W(\mathbf{S}_T)]$. Our second estimator of $\mathbf{S}_{0,T}$ minimizes the following GLS criterion based upon (N.1):

$$\hat{\mathbf{S}}_{T,FGLS} = \underset{\mathbf{S}_T \in \Xi_{\varepsilon, \pi_0, m_0, T}}{\operatorname{argmin}} \left(\vec{y} - W(\mathbf{S}_T) \hat{\xi}_{FGLS}(\mathbf{S}_T) \right)' \hat{\Omega}_{\varepsilon, \mathbf{S}}^{-1} \left(\vec{y} - W(\mathbf{S}_T) \hat{\xi}_{FGLS}(\mathbf{S}_T) \right).$$

Correspondingly, we estimate β with $\hat{\beta}(\hat{\mathbf{S}}_{T,FGLS})$. In order for $\hat{\beta}(\hat{\mathbf{S}}_{T,FGLS})$ to be provably more efficient than $\hat{\beta}(\hat{\mathbf{S}}_{T,OLS})$, $\hat{\Omega}_{\varepsilon, \mathbf{S}}$ must be a consistent estimator of $\mathbb{E}[\varepsilon\varepsilon'|W(\mathbf{S}_{0,T})]$. When ε_t does not exhibit conditional serial correlation or heteroskedasticity, i.e., $\mathbb{E}[\varepsilon\varepsilon'|W(\mathbf{S}_{0,T})] = \Sigma_\varepsilon \otimes I_T$, this is feasible since one could simply use $\hat{\Omega}_{\varepsilon, \mathbf{S}} = \hat{\Sigma}_\varepsilon \otimes I_T$, where $\hat{\Sigma}_{\varepsilon, i, j} = (T - q - p)^{-1} \hat{\varepsilon}^{i'} \hat{\varepsilon}^j$ for $i, j = 1, 2$ with $\hat{\varepsilon}^1$ ($\hat{\varepsilon}^2$) equal to the first (last) T elements of $\vec{y} - W(\hat{\mathbf{S}}_{T,OLS}) \hat{\xi}_{OLS}(\hat{\mathbf{S}}_{T,OLS})$, as is standard in seemingly unrelated regression. For serially dependent ε_t , consistent estimation of $\mathbb{E}[\varepsilon\varepsilon'|W(\mathbf{S}_{0,T})]$ requires a correctly-specified model for the dependence in ε_t , a strong assumption in some empirical applications. In Section (N.D) we present the consistency results about $\hat{\mathbf{S}}_{T,OLS}$, $\hat{\beta}(\hat{\mathbf{S}}_{T,OLS})$, $\hat{\mathbf{S}}_{T,FGLS}$ and $\hat{\beta}(\hat{\mathbf{S}}_{T,FGLS})$.

N.D Theoretical Results on Estimation of LATE and Sub-populations

In this section we establish theoretical results about the estimator $\hat{\mathbf{S}}_{T,FGLS}$ from which we deduce the same results for $\hat{\mathbf{S}}_{T,OLS}$ as a special case with $\vec{y} = D$ and $\hat{\Omega}_{\varepsilon, \mathbf{S}} = I_T$. We first rewrite the model in matrix format as follows. We have 2 equations and T observations, excluding the initial conditions if lagged dependent variables are included among the regressors. The number of regimes

that define the π sub-population is m , while the total number of regimes in the full sample is $\widetilde{m} \geq m$. For example, if $\pi = 1$ then $\widetilde{m} = m$, else $\widetilde{m} > m$. The break dates are denoted by the $\widetilde{m}$ vector $(T_1, \dots, T_{\widetilde{m}})$ and we use the usual convention that $T_0 = 1$ and $T_{\widetilde{m}+1} = T$. A subscript i indexes a regime ($i = 1, \dots, \widetilde{m} + 1$), a subscript t indexes a temporal observation ($t = 1, \dots, T$) and a subscript j indexes the equation ($j = 1, 2$) to which a scalar dependent variable y_{jt} belongs. According to our model in Section 5 $y_{1t} = Y_t$ and $y_{2t} = D_t$. $q + p$ is the number of regressors and z_t is the set that includes the regressors from all equations $z_t = (z_{1,t}, \dots, z_{q+p,t})' = (Z_t', X_t')'$. The model considered in (5.1) can be written as

$$y_t = (I_2 \otimes z_t') \alpha_i + v_t, \quad (\text{N.1})$$

where v_t has mean zero and covariance matrix Σ . The parameters in regime i are the $p + q$ vector $\alpha_i = (\beta\theta_i', \gamma_1' + \gamma_2'\beta, \theta_i', \gamma_2')'$, where $\theta_i = \theta$ for $T_{i-1} + 1 \leq t \leq T_i$ with $t \in \mathbf{S}_{0,T}$ and $\theta_i = 0$ for $T_{i-1} + 1 \leq t \leq T_i$ with $t \notin \mathbf{S}_{0,T}$. Let $\alpha = (\alpha_1', \dots, \alpha_{\widetilde{m}+1}')'$.

To ease notation, define the $2(q + p) \times 2$ matrix x_t by $x_t' = (I \otimes z_t')$ and rewrite (N.1) as

$$y_t = x_t' \alpha_i + v_t, \quad (\text{N.2})$$

for $T_{i-1} + 1 \leq t \leq T_i$ ($i = 1, \dots, \widetilde{m} + 1$). We now express the model in matrix form. Let $\vec{Y} = (y_1', \dots, y_T')'$ be the $2T$ vector of dependent variables, let $V = (v_1', \dots, v_T')'$ be the error vector, and let the $2T \times 2(q + p)$ matrix of regressors be $\vec{X} = (x_1, \dots, x_T)'$. For a given partition $\mathbf{S}$ with associated breaks $(T_1, \dots, T_{\widetilde{m}})$, we define the block partition of the matrix $\vec{X}$ as the $2T \times 2(q + p)(\widetilde{m} + 1)$ matrix $\overline{X}(\mathbf{S}) = \text{diag}(\overline{X}_1, \dots, \overline{X}_{\widetilde{m}+1})$, where $\overline{X}_i$ ($i = 1, \dots, \widetilde{m} + 1$) is the $2(T_i - T_{i-1}) \times 2(q + p)$ subset of $\vec{X}$ that corresponds to observations in regime i . We also define the subvector V_i of V similarly. Then the regression (N.2) can be written as $\vec{Y} = \overline{X}(\mathbf{S})\alpha + V$. The true values of the parameters are denoted with a 0 superscript so that the data generating process is assumed to be $\vec{Y} = \overline{X}(\mathbf{S}_{0,T})\alpha_0 + V$, where $\overline{X}(\mathbf{S}_{0,T})$ is the diagonal partition of $\vec{X}$ using the partition $\mathbf{S}_{0,T}$, i.e. $(T_1^0, \dots, T_{\widetilde{m}}^0)$. Let $\widehat{\Omega}_{\mathbf{S}}$ be the rearrangement of $\widehat{\Omega}_{\varepsilon, \mathbf{S}}$ in the main text corresponding to the rearrangement $\vec{Y}$ of $\vec{y}$. We make the following assumptions.

Assumption N.D.1. $\sup_{\mathbf{S}} \|\overline{X}(\mathbf{S})' \widehat{\Omega}_{\mathbf{S}}^{-1/2}\| = O_{\mathbb{P}}(T^{1/2})$, $\sup_{\mathbf{S}, \mathbf{S}'} \|\overline{X}(\mathbf{S})' \widehat{\Omega}_{\mathbf{S}}^{-1} \overline{X}(\mathbf{S}')/T\| = O_{\mathbb{P}}(1)$, $\inf_{\mathbf{S}} \lambda_{\min}(\overline{X}(\mathbf{S})' \widehat{\Omega}_{\mathbf{S}}^{-1} \overline{X}(\mathbf{S})/T) > 0$ with probability approaching one, and $\sup_{\mathbf{S}} \|V' \widehat{\Omega}_{\mathbf{S}}^{-1} \overline{X}(\mathbf{S})\| = O_{\mathbb{P}}(T^{1/2})$. For some $C > 0$,

$$\sup_{\mathbf{S}} \max_{1 \leq t \leq T} \sum_{s=1}^T \left\| \left[\widehat{\Omega}_{\mathbf{S}}^{-1} \right]_{(t,s)} \right\| = O_{\mathbb{P}}(1).$$

Assumption N.D.2. *There exists an $l_0 > 0$ such that for all $l > l_0$, the minimum eigenvalue of*

$$M_{\Omega,T} = \begin{bmatrix} M_{11,\Omega,T} & M_{12,\Omega,T} \\ M_{21,\Omega,T} & M_{22,\Omega,T} \end{bmatrix}$$

is bounded away from zero uniformly over $i = 1, \dots, \widetilde{m}$ and $\mathbf{S}$ where

$$\begin{aligned} M_{11,\Omega,T} &= l^{-1} \sum_{t=T_i^0-l+1}^{T_i^0} \sum_{s=T_i^0-l+1}^{T_i^0} x_t [\widehat{\Omega}_{\mathbf{S}}^{-1}]_{(t,s)} x'_s \\ M_{22,\Omega,T} &= l^{-1} \sum_{t=T_i^0+1}^{T_i^0+l} \sum_{s=T_i^0+1}^{T_i^0+l} x_t [\widehat{\Omega}_{\mathbf{S}}^{-1}]_{(t,s)} x'_s \\ M_{12,\Omega,T} &= l^{-1} \sum_{t=T_i^0-l+1}^{T_i^0} \sum_{s=T_i^0+1}^{T_i^0+l} x_t [\widehat{\Omega}_{\mathbf{S}}^{-1}]_{(t,s)} x'_s \end{aligned}$$

$M_{21,\Omega,T} = M'_{12,\Omega,T}$ and $[\widehat{\Omega}_{\mathbf{S}}^{-1}]_{(t,s)}$ denotes the (t, s) -th element of $\widehat{\Omega}_{\mathbf{S}}^{-1}$.

Assumption N.D.3. *The matrix $\sum_{t=k}^l \sum_{s=k}^l x_t [\widehat{\Omega}_{\mathbf{S}}^{-1}]_{(t,s)} x'_s$ is invertible for $l - k \geq k_0$ for some $0 < k_0 < \infty$.*

Assumption N.D.4. *We have $0 < \lambda_1^0 < \dots < \lambda_{\widetilde{m}}^0 < 1$ with $T_i^0 = \lfloor T \lambda_i^0 \rfloor$.*

Assumption N.D.5. *The minimization search is taken over all partitions that satisfy $|\lambda_{i+1} - \lambda_i| \geq \epsilon$, $|\lambda_1| \geq \epsilon$, $|\lambda_{\widetilde{m}}| \leq 1 - \epsilon$.*

Assumption N.D.6. *For $\Omega = \mathbb{E}[VV' | \overline{X}(\mathbf{S}_{0,T})]$, $\sup_{\mathbf{S}, \mathbf{S}'} T^{-1} \overline{X}(\mathbf{S}')' (\widehat{\Omega}_{\mathbf{S}}^{-1} - \Omega^{-1}) \overline{X}(\mathbf{S}) \xrightarrow{\mathbb{P}} 0$, $\sup_{\mathbf{S}} T^{-1} \overline{X}(\mathbf{S})' (\widehat{\Omega}_{\mathbf{S}}^{-1} - \Omega^{-1}) V \xrightarrow{\mathbb{P}} 0$ and $T^{-1} V' (\widehat{\Omega}_{\mathbf{S}}^{-1} - \Omega^{-1}) V \xrightarrow{\mathbb{P}} 0$.*

N.D.1 Consistency Under Fixed Shifts

Let $\widehat{\lambda}$ be the estimate of the break fractions $\lambda_0 = (\lambda_1^0, \lambda_2^0, \dots, \lambda_{\widetilde{m}}^0)$ that corresponds to $\widehat{\mathbf{S}}_{T,FGLS}$. The following proposition states the consistency of $\widehat{\lambda}$ for λ_0 .

Proposition N.D.1. *Let Assumptions **N.D.1-N.D.6** hold. Then, $\widehat{\lambda}_i \xrightarrow{\mathbb{P}} \lambda_i^0$, $i = 1, \dots, \widetilde{m}$.*

We now consider the rate of convergence of $\widehat{\lambda}$.

Proposition N.D.2. *Let Assumptions **N.D.1-N.D.6** hold, for every $\eta > 0$, there exists a $C < \infty$, such that for all large T ,*

$$\mathbb{P} \left(\left| T (\widehat{\lambda}_i - \lambda_i^0) \right| > C \right) < \eta, \quad (i = 1, \dots, \widetilde{m}).$$

Let $\hat{\alpha}(\cdot)$ be defined in analogy with $\hat{\xi}_{FGLS}(\cdot)$ in the main text upon rearrangement of $\vec{y}$ to $\vec{Y}$. The T rate of convergence of $\hat{\lambda}_i$ allows us to obtain the asymptotic equivalence between the estimated slope coefficients with the estimated subpopulation $\hat{\alpha}(\hat{\mathbf{S}}_{T,FGLS})$ and the estimated slope coefficients with known subpopulation $\hat{\alpha}(\mathbf{S}_{0,T})$ so that standard results feasible generalized least squares results implying $\sqrt{T}$ asymptotic normality for the latter also immediately apply to the former.

Proposition N.D.3. *Let Assumptions N.D.1-N.D.6 hold. We have $\sqrt{T}(\hat{\alpha}(\hat{\mathbf{S}}_{T,FGLS}) - \alpha_0) = O_{\mathbb{P}}(1)$.*

N.D.2 Proofs

N.D.2.1 Proof of Proposition N.D.1

We first outline the main steps of the proof using a few lemmas that are proved below. By the definition of $\hat{\mathbf{S}}_{T,FGLS}$ and Assumption N.D.6,

$$\frac{1}{T} \hat{V}(\hat{\mathbf{S}})' \hat{\Omega}_{\hat{\mathbf{S}}}^{-1} \hat{V}(\hat{\mathbf{S}}) \leq \frac{1}{T} V' \Omega^{-1} V, \quad (\text{N.3})$$

with probability approaching one, where $\hat{V}(\hat{\mathbf{S}}) = \vec{Y} - \bar{X}(\hat{\mathbf{S}}) \hat{\alpha}(\hat{\mathbf{S}})$ with $\hat{\mathbf{S}} = \hat{\mathbf{S}}_{T,FGLS}$. Note that

$$\begin{aligned} & \hat{V}(\hat{\mathbf{S}})' \hat{\Omega}_{\hat{\mathbf{S}}}^{-1} \hat{V}(\hat{\mathbf{S}}) \\ &= (V - \bar{X}(\hat{\mathbf{S}})(\hat{\alpha}(\hat{\mathbf{S}}) - \alpha_0) - (\bar{X}(\hat{\mathbf{S}}) - \bar{X}(\mathbf{S}_{0,T}))\alpha_0)' \hat{\Omega}_{\hat{\mathbf{S}}}^{-1} \\ & \quad \times (V - \bar{X}(\hat{\mathbf{S}})(\hat{\alpha}(\hat{\mathbf{S}}) - \alpha_0) - (\bar{X}(\hat{\mathbf{S}}) - \bar{X}(\mathbf{S}_{0,T}))\alpha_0) \\ &= V' \Omega^{-1} V + (V' \hat{\Omega}_{\hat{\mathbf{S}}}^{-1} V - V' \Omega^{-1} V) + (\hat{\alpha}(\hat{\mathbf{S}}) - \alpha_0)' \bar{X}(\hat{\mathbf{S}})' \hat{\Omega}_{\hat{\mathbf{S}}}^{-1} \bar{X}(\hat{\mathbf{S}}) (\hat{\alpha}(\hat{\mathbf{S}}) - \alpha_0) \\ & \quad + \alpha_0' (\bar{X}(\hat{\mathbf{S}}) - \bar{X}(\mathbf{S}_{0,T}))' \hat{\Omega}_{\hat{\mathbf{S}}}^{-1} (\bar{X}(\hat{\mathbf{S}}) - \bar{X}(\mathbf{S}_{0,T})) \alpha_0 \\ & \quad + 2(\hat{\alpha}(\hat{\mathbf{S}}) - \alpha_0)' \bar{X}(\hat{\mathbf{S}})' \hat{\Omega}_{\hat{\mathbf{S}}}^{-1} (\bar{X}(\hat{\mathbf{S}}) - \bar{X}(\mathbf{S}_{0,T})) \alpha_0 \\ & \quad - 2V' \hat{\Omega}_{\hat{\mathbf{S}}}^{-1} \bar{X}(\hat{\mathbf{S}}) (\hat{\alpha}(\hat{\mathbf{S}}) - \alpha_0) - 2V' \hat{\Omega}_{\hat{\mathbf{S}}}^{-1} (\bar{X}(\hat{\mathbf{S}}) - \bar{X}(\mathbf{S}_{0,T})) \alpha_0 \\ &\equiv V' \Omega^{-1} V + \sum_{j=1}^6 E_j. \end{aligned} \quad (\text{N.4})$$

The proof of Proposition N.D.1 uses (N.3)-(N.4) and the limit of $E_1, \dots, E_6$. By Assumption N.D.6, $T^{-1}E_1 \xrightarrow{\mathbb{P}} 0$ and we show that $T^{-1}E_j \xrightarrow{\mathbb{P}} 0$ for $j = 5$ and 6 , in Lemma N.D.1 below. These results combined with (N.3) imply that $T^{-1}(E_2 + E_3 + E_4) \xrightarrow{\mathbb{P}} 0$. The proof follows by showing that the latter imply $\hat{\lambda} \xrightarrow{\mathbb{P}} \lambda_0$ via Lemma N.D.2. We proceed with a couple of lemmas.

Lemma N.D.1. *Let Assumptions N.D.1 and N.D.3 hold. We have $T^{-1}E_j \xrightarrow{\mathbb{P}} 0$ for $j = 5$ and 6.*

Proof of Lemma N.D.1. To prove the lemma, it suffices to show that

$$\sup_{\mathbf{S}} \frac{1}{T} \left| V' \hat{\Omega}_{\mathbf{S}}^{-1} \bar{X}(\mathbf{S}) (\hat{\alpha}(\mathbf{S}) - \alpha_0) \right| = O_{\mathbb{P}}(T^{-1/2}) = o_{\mathbb{P}}(1), \quad (\text{N.5})$$

$$\sup_{\mathbf{S}} \frac{1}{T} \left| V' \hat{\Omega}_{\mathbf{S}}^{-1} (\bar{X}(\mathbf{S}) - \bar{X}(\mathbf{S}_{0,T})) \alpha_0 \right| = O_{\mathbb{P}}(T^{-1/2}) = o_{\mathbb{P}}(1). \quad (\text{N.6})$$

First consider (N.5). We can rewrite

$$\begin{aligned} & V' \hat{\Omega}_{\mathbf{S}}^{-1} \bar{X}(\mathbf{S}) \hat{\alpha}(\mathbf{S}) - V' \hat{\Omega}_{\mathbf{S}}^{-1} \bar{X}(\mathbf{S}) \alpha_0 \\ &= V' \hat{\Omega}_{\mathbf{S}}^{-1} \bar{X}(\mathbf{S}) \left(\bar{X}(\mathbf{S})' \hat{\Omega}_{\mathbf{S}}^{-1} \bar{X}(\mathbf{S}) \right)^{-1} \bar{X}(\mathbf{S})' \hat{\Omega}_{\mathbf{S}}^{-1} \bar{X}(\mathbf{S}_{0,T}) \alpha_0 \\ & \quad + V' \hat{\Omega}_{\mathbf{S}}^{-1} \bar{X}(\mathbf{S}) \left(\bar{X}(\mathbf{S})' \hat{\Omega}_{\mathbf{S}}^{-1} \bar{X}(\mathbf{S}) \right)^{-1} \bar{X}(\mathbf{S})' \hat{\Omega}_{\mathbf{S}}^{-1} V \\ & \quad - V' \hat{\Omega}_{\mathbf{S}}^{-1} \bar{X}(\mathbf{S}) \alpha_0. \end{aligned}$$

Using Assumptions N.D.1 and N.D.3, the first term on the right-hand side is $O_{\mathbb{P}}(T^{1/2}) O_{\mathbb{P}}(T^{-1}) O_{\mathbb{P}}(T) = O_{\mathbb{P}}(T^{1/2})$ uniformly over all partitions. The second term is $O_{\mathbb{P}}(T^{1/2}) O_{\mathbb{P}}(T^{-1}) O_{\mathbb{P}}(T^{1/2}) = O_{\mathbb{P}}(1)$ and the third term is $O_{\mathbb{P}}(T^{1/2})$, both uniformly over all partitions. Then, (N.5) follows. Next, consider (N.6). Using Assumption N.D.1, we have

$$\begin{aligned} V' \hat{\Omega}_{\mathbf{S}}^{-1} (\bar{X}(\mathbf{S}) - \bar{X}(\mathbf{S}_{0,T})) \alpha_0 &= V' \hat{\Omega}_{\mathbf{S}}^{-1} \bar{X}(\mathbf{S}) \alpha_0 - V' \hat{\Omega}_{\mathbf{S}}^{-1} \bar{X}(\mathbf{S}_{0,T}) \alpha_0 \\ &= O_{\mathbb{P}}(T^{1/2}) + O_{\mathbb{P}}(T^{1/2}). \end{aligned}$$

This implies (N.6). $\square$

Lemma N.D.2. *Let Assumptions N.D.2-N.D.5 hold. If $\hat{\lambda}_i \xrightarrow{\mathbb{P}} \lambda_i^0$ for some i , then*

$$\liminf_{T \rightarrow \infty} \mathbb{P} \left(T^{-1} (E_2 + E_3 + E_4) > c \right) > \epsilon_0$$

for some $c > 0$ and $\epsilon_0 > 0$.

Proof of Lemma N.D.2. We have for $T_{i-1} + 1 \leq t \leq T_i$,

$$\begin{aligned} \hat{v}_t(\hat{\mathbf{S}}) &= \begin{bmatrix} Y_t \\ D_t \end{bmatrix} - \begin{bmatrix} Z_t' \hat{\theta}_{\beta, \hat{i}} \\ Z_t' \hat{\theta}_{\hat{i}} \end{bmatrix} - \begin{bmatrix} X_t' \hat{\gamma}_{\beta} \\ X_t' \hat{\gamma}_2 \end{bmatrix} = \begin{bmatrix} Z_t' \theta_{\beta, i} \\ Z_t' \theta_i \end{bmatrix} + \begin{bmatrix} X_t' \gamma_{\beta} \\ X_t' \gamma_2 \end{bmatrix} - \begin{bmatrix} Z_t' \hat{\theta}_{\beta, \hat{i}} \\ Z_t' \hat{\theta}_{\hat{i}} \end{bmatrix} - \begin{bmatrix} X_t' \hat{\gamma}_{\beta} \\ X_t' \hat{\gamma}_2 \end{bmatrix} + v_t \\ &= \begin{bmatrix} Z_t' (\theta_{\beta, i} - \hat{\theta}_{\beta, \hat{i}}) \\ Z_t' (\theta_i - \hat{\theta}_{\hat{i}}) \end{bmatrix} + \begin{bmatrix} X_t' (\gamma_{\beta} - \hat{\gamma}_{\beta}) \\ X_t' (\gamma_2 - \hat{\gamma}_2) \end{bmatrix} + v_t, \end{aligned}$$

where $\alpha_i = (\theta_{\beta,i}, \gamma_\beta, \theta'_i, \gamma'_2)'$, $\hat{\alpha}_i = (\hat{\theta}_{\beta,\hat{i}}, \hat{\gamma}_\beta, \hat{\theta}'_i, \hat{\gamma}'_2)'$ and $\hat{i}$ corresponds to a regime in $\hat{\mathbf{S}}$.

By Assumptions [N.D.4](#) and [N.D.5](#), if there exists a break, say λ_i^0 , which cannot be consistently estimated, then with some probability $\epsilon_0 > 0$ there exists a $\eta > 0$ such that no estimated break falls in the interval $[T(\lambda_j^0 - \eta), T(\lambda_j^0 + \eta)]$ for a subsequence of T . Suppose this interval is classified into the k -th regime, i.e., $\hat{T}_{k-1} \leq T(\lambda_i^0 - \eta)$ and $T(\lambda_i^0 + \eta) \leq \hat{T}_k$. Let d_t denote the difference between the fitted residuals and true errors. Then,

$$d_t = \begin{cases} \begin{bmatrix} Z'_t(\theta_{\beta,i} - \hat{\theta}_{\beta,k}) \\ Z'_t(\theta_i - \hat{\theta}_k) \end{bmatrix} + \begin{bmatrix} X'_t(\gamma_\beta - \hat{\gamma}_\beta) \\ X'_t(\gamma_2 - \hat{\gamma}_2) \end{bmatrix} & \text{for } t \in [T(\lambda_i^0 - \eta), T\lambda_i^0] \\ \begin{bmatrix} Z'_t(\theta_{\beta,i+1} - \hat{\theta}_{\beta,k}) \\ Z'_t(\theta_{i+1} - \hat{\theta}_k) \end{bmatrix} + \begin{bmatrix} X'_t(\gamma_\beta - \hat{\gamma}_\beta) \\ X'_t(\gamma_2 - \hat{\gamma}_2) \end{bmatrix} & \text{for } t \in [T\lambda_i^0, T(\lambda_i^0 + \eta)]. \end{cases}$$

For $t \in [T(\lambda_i^0 - \eta), T\lambda_i^0]$,

$$d_t = x'_t(\alpha_i - \hat{\alpha}_k) = x'_t a_{\beta,i,k}$$

while for $t \in [T\lambda_i^0, T(\lambda_i^0 + \eta)]$,

$$d_t = x'_t(\alpha_{i+1} - \hat{\alpha}_k) = x'_t b_{\beta,i,k}.$$

Under Assumption [N.D.2](#) we have

$$\begin{aligned} E_2 + E_3 + E_4 &= (\hat{V}(\hat{\mathbf{S}}) - V)' \hat{\Omega}_{\mathbf{S}}^{-1} (\hat{V}(\hat{\mathbf{S}}) - V) \\ &= (\bar{X}(\hat{\mathbf{S}}) \hat{\alpha}(\hat{\mathbf{S}}) - \bar{X}(\mathbf{S}_{0,T}) \alpha_0)' \hat{\Omega}_{\mathbf{S}}^{-1} (\bar{X}(\hat{\mathbf{S}}) \hat{\alpha}(\hat{\mathbf{S}}) - \bar{X}(\mathbf{S}_{0,T}) \alpha_0) \\ &= D(\hat{\mathbf{S}})' \hat{\Omega}_{\mathbf{S}}^{-1} D(\hat{\mathbf{S}}) \\ &\geq \begin{bmatrix} a_{\beta,i,k} \\ b_{\beta,i,k} \end{bmatrix}' (T\eta M_{\Omega,T}) \begin{bmatrix} a_{\beta,i,k} \\ b_{\beta,i,k} \end{bmatrix} \\ &\geq \gamma_T [\|a_{\beta,i,k}\|^2 + \|b_{\beta,i,k}\|^2] \\ &\geq 2^{-1} \gamma_T (\|\theta_{\beta,i} - \theta_{\beta,i+1}\|^2 + \|\theta_i - \theta_{i+1}\|^2), \end{aligned}$$

where $D(\hat{\mathbf{S}}) = [d'_1 d'_2 \cdots d'_T]'$, γ_T is the smallest eigenvalue of $T\eta M_{\Omega,T}$, and the last inequality follows

from

$$(x - a)' A (x - a) + (x - b)' A (x - b) \geq \frac{1}{2} (a - b)' A (a - b)$$

for an arbitrary positive definite matrix A and for all x . Since γ_T is of the order $(T\eta)$ we yield

$$\sum_{j=2}^4 E_j > T\eta c_1 \left(\|\theta_{\beta,i} - \theta_{\beta,i+1}\|^2 + \|\theta_i - \theta_{i+1}\|^2 \right) = TC \left(\|\theta_{\beta,i} - \theta_{\beta,i+1}\|^2 + \|\theta_i - \theta_{i+1}\|^2 \right),$$

for some $C = \eta c_1 > 0$ with probability no less than $\epsilon_0 > 0$ as $T \rightarrow \infty$. $\square$

Proof of Proposition N.D.1. Using (N.4), $T^{-1}E_1 \xrightarrow{\mathbb{P}} 0$ and Lemmas N.D.1-N.D.2, and under the supposition that some break date is not consistently estimated, we have the inequality

$$\frac{1}{T} \widehat{V} (\widehat{\mathbf{S}})' \widehat{\Omega}_{\mathbf{S}}^{-1} \widehat{V} (\widehat{\mathbf{S}}) \geq \frac{1}{T} V' \Omega^{-1} V + C + o_{\mathbb{P}}(1)$$

for some $C > 0$ holding with probability no less than some ϵ_0 as $T \rightarrow \infty$. This is in contradiction with (N.3). Hence, all break fractions are consistently estimated. $\square$

N.D.2.2 Proof of Proposition N.D.2

Without loss of generality, we assume there are only three regimes ($\widetilde{m} = 3$) and provide an explicit proof of T -consistency for $\widehat{\lambda}_2$ only. The analysis for $\widehat{\lambda}_1$ and $\widehat{\lambda}_3$ is virtually the same and is omitted.

By Proposition N.D.1, for each $\epsilon > 0$ and T large, we have $|\widehat{T}_i - T_i^0| \leq \epsilon T$ with probability approaching one. For each $\epsilon > 0$, let $\mathbf{T}_\epsilon = \{(T_1, T_2, T_3) : |T_i - T_i^0| \leq \epsilon T \text{ for } i = 1, \dots, 3\}$ so that $\mathbb{P}(\{\widehat{T}_1, \widehat{T}_2, \widehat{T}_3\} \in \mathbf{T}_\epsilon) \rightarrow 1$. Therefore we only need to examine the behavior of the objective function, $Q_T(T_1, T_2, T_3) = \widehat{V}(\mathbf{S})' \widehat{\Omega}_{\mathbf{S}}^{-1} \widehat{V}(\mathbf{S})$, for those T_i corresponding to $\mathbf{S}$ that are close to the true breaks such that $|T_i - T_i^0| < \epsilon T$ for all i . Also using an argument of symmetry, we can without loss of generality, restrict attention to the case $T_2 < T_2^0$. For $C > 0$, define

$$\mathbf{T}_\epsilon(C) = \{(T_1, T_2, T_3) : |T_i - T_i^0| < \epsilon T, 1 \leq i \leq 3, T_2 - T_2^0 < -C\}.$$

Note that $\mathbf{T}_\epsilon(C) \subset \mathbf{T}_\epsilon$. Because $Q_T(\widehat{T}_1, \widehat{T}_2, \widehat{T}_3) \leq Q_T(\widehat{T}_1, T_2^0, \widehat{T}_3)$ with probability 1, to prove the proposition it is enough to show that for each $\eta > 0$, there exist $C > 0$ and $\epsilon > 0$ such that for large T ,

$$\mathbb{P} \left(\min_{\mathbf{T}_\epsilon(C)} \{Q_T(T_1, T_2, T_3) - Q_T(T_1, T_2^0, T_3)\} \leq 0 \right) < \eta, \quad (\text{N.7})$$

or equivalently,

$$\mathbb{P} \left(\min_{\mathbf{T}_\epsilon(C)} \left\{ \left[Q_T(T_1, T_2, T_3) - Q_T(T_1, T_2^0, T_3) \right] / (T_2^0 - T_2) \right\} \leq 0 \right) < \eta. \quad (\text{N.8})$$

That would imply that for a large C , global minimization cannot be achieved on $\mathbf{T}_\epsilon(C)$. Thus with probability approaching one, $|\hat{T}_2 - T_2^0| \leq C$. Now denote

$$\begin{aligned} Q_{1,T} &= Q_T(T_1, T_2, T_3) \\ Q_{2,T} &= Q_T(T_1, T_2^0, T_3) \\ Q_{3,T} &= \hat{V}(\mathbf{S}_{3,T})' \hat{\Omega}_{\mathbf{S}}^{-1} \hat{V}(\mathbf{S}_{3,T}) \end{aligned}$$

where $\mathbf{S}_{3,T}$ is the partition based on (T_1, T_2, T_2^0, T_3) . Subtracting and adding $Q_{3,T}$, we have

$$Q_{1,T} - Q_{2,T} = Q_{1,T} - Q_{3,T} - (Q_{2,T} - Q_{3,T}).$$

This latter relation is useful because it allows us to perform the analysis in terms of two problems involving a single break. Indeed, $Q_{1,T} - Q_{3,T}$ is the difference in the objective function allowing an additional fourth break at time T_2^0 between the breaks T_2 and T_3 . Similarly, $Q_{2,T} - Q_{3,T}$ is the difference in the objective function allowing an additional fourth break at time T_2 between the breaks T_1 and T_2^0 . Consider $Q_{1,T} - Q_{3,T}$ first. Let $(\hat{\alpha}_1^*, \hat{\alpha}_2^*, \hat{\alpha}_\Delta^*, \hat{\alpha}_3^*, \hat{\alpha}_4^*)$ denote the estimator of $(\alpha_1^0, \alpha_2^0, \alpha_3^0, \alpha_4^0)$. In particular, $\hat{\alpha}_2^*$ is an estimate of α_2^0 associated with the regressors $(0, \dots, 0, x_{T_1+1}, \dots, x_{T_2}, 0, \dots, 0)'$, $\hat{\alpha}_\Delta^*$ is the vector of estimated coefficients associated with the regressors $X_\Delta = (0, \dots, 0, x_{T_2+1}, \dots, x_{T_2^0}, 0, \dots, 0)'$.

From the argument on p. 31 in [Amemiya \(1985\)](#),

$$\begin{aligned} Q_{1,T} - Q_{3,T} &= \hat{V}(\mathbf{S}_{1,T})' \hat{\Omega}_{\mathbf{S}}^{-1} \hat{V}(\mathbf{S}_{1,T}) - \hat{V}(\mathbf{S}_{3,T})' \hat{\Omega}_{\mathbf{S}}^{-1} \hat{V}(\mathbf{S}_{3,T}) \\ &= (\hat{\alpha}_3^* - \hat{\alpha}_\Delta^*)' X_\Delta' \hat{\Omega}_{\mathbf{S}}^{-1/2} M_{\hat{\Omega}_{\mathbf{S}}^{-1/2} \bar{X}(\mathbf{S}_{1,T})} \hat{\Omega}_{\mathbf{S}}^{-1/2} X_\Delta (\hat{\alpha}_3^* - \hat{\alpha}_\Delta^*), \end{aligned}$$

where $M_X = I - X(X'X)^{-1}X'$ for a matrix X and $\hat{\alpha}_3^*$ is the vector of estimated coefficients associated with the regressors $(0, \dots, 0, x_{T_2^0+1}, \dots, x_{T_3}, 0, \dots, 0)'$. Similarly, we have for $Q_{2,T} - Q_{3,T}$,

$$Q_{2,T} - Q_{3,T} = (\hat{\alpha}_2^* - \hat{\alpha}_\Delta^*)' X_\Delta' \hat{\Omega}_{\mathbf{S}}^{-1/2} M_{\hat{\Omega}_{\mathbf{S}}^{-1/2} \bar{X}(\mathbf{S}_{2,T})} \hat{\Omega}_{\mathbf{S}}^{-1/2} X_\Delta (\hat{\alpha}_2^* - \hat{\alpha}_\Delta^*).$$

Thus,

$$\begin{aligned} Q_{1,T} - Q_{2,T} &\geq (\hat{\alpha}_3^* - \hat{\alpha}_\Delta^*)' X'_\Delta \hat{\Omega}_S^{-1/2} M_{\hat{\Omega}_S^{-1/2} \bar{X}(\mathbf{S}_{1,T})} \hat{\Omega}_S^{-1/2} X_\Delta (\hat{\alpha}_3^* - \hat{\alpha}_\Delta^*) \\ &\quad - (\hat{\alpha}_2^* - \hat{\alpha}_\Delta^*)' X'_\Delta \hat{\Omega}_S^{-1} X_\Delta (\hat{\alpha}_2^* - \hat{\alpha}_\Delta^*), \end{aligned} \quad (\text{N.9})$$

where we used

$$X'_\Delta \hat{\Omega}_S^{-1/2} M_{\hat{\Omega}_S^{-1/2} \bar{X}(\mathbf{S}_{2,T})} \hat{\Omega}_S^{-1/2} X_\Delta \leq X'_\Delta \hat{\Omega}_S^{-1} X_\Delta.$$

From the definition of $M_{\hat{\Omega}_S^{-1/2} \bar{X}(\mathbf{S}_{1,T})}$, we have

$$\begin{aligned} \frac{Q_{1,T} - Q_{2,T}}{T_2^0 - T_2} &\geq (\hat{\alpha}_3^* - \hat{\alpha}_\Delta^*)' \frac{X'_\Delta \hat{\Omega}_S^{-1} X_\Delta}{T_2^0 - T_2} (\hat{\alpha}_3^* - \hat{\alpha}_\Delta^*) \\ &\quad - (\hat{\alpha}_3^* - \hat{\alpha}_\Delta^*)' \frac{X'_\Delta \hat{\Omega}_S^{-1} \bar{X}(\mathbf{S}_{1,T})}{T_2^0 - T_2} \left[\frac{\bar{X}(\mathbf{S}_{1,T})' \hat{\Omega}_S^{-1} \bar{X}(\mathbf{S}_{1,T})}{T} \right]^{-1} \\ &\quad \times \frac{\bar{X}(\mathbf{S}_{1,T})' \hat{\Omega}_S^{-1} X_\Delta}{T} (\hat{\alpha}_3^* - \hat{\alpha}_\Delta^*) \\ &\quad - (\hat{\alpha}_2^* - \hat{\alpha}_\Delta^*)' \frac{X'_\Delta \hat{\Omega}_S^{-1} X_\Delta}{T_2^0 - T_2} (\hat{\alpha}_2^* - \hat{\alpha}_\Delta^*) + Q_{2,T} - Q_{2,T} \\ &\equiv L_1 - L_2 - L_3. \end{aligned} \quad (\text{N.10})$$

Consider the limiting behavior of L_1 . Note first that for small ϵ , the estimates $\hat{\alpha}_i^*$ will be close to α_i^0 with high probability in the sequential limit as $T \rightarrow \infty$ followed by $C \rightarrow \infty$, since on the set $\mathbf{T}_\epsilon(C)$, the distance between T_i and T_i^0 can be made small by choosing a small ϵ . Further, $\hat{\alpha}_\Delta^*$ is estimated using observations from the second true regime only and it is close to α_2^0 in probability on $\mathbf{T}_\epsilon(C)$ for a large enough C . Hence, for large C , large T and small ϵ , L_1 is larger than

$$(\hat{\alpha}_3^* - \hat{\alpha}_\Delta^*)' \frac{X'_\Delta \hat{\Omega}_S^{-1} X_\Delta}{T_2^0 - T_2} (\hat{\alpha}_3^* - \hat{\alpha}_\Delta^*) \geq \frac{1}{2} (\alpha_3^0 - \alpha_2^0)' \frac{X'_\Delta \hat{\Omega}_S^{-1} X_\Delta}{T_2^0 - T_2} (\alpha_3^0 - \alpha_2^0)$$

with high probability.

Next consider L_2 . It is easy to see that on $\mathbf{T}_\epsilon(C)$, $\hat{\alpha}_3^*$ and $\hat{\alpha}_\Delta^*$ are $O_{\mathbb{P}}(1)$ uniformly. Also on $\mathbf{T}_\epsilon(C)$,

$$\frac{\bar{X}(\mathbf{S}_{1,T})' \hat{\Omega}_S^{-1} \bar{X}(\mathbf{S}_{1,T})}{T} = O_{\mathbb{P}}(1),$$

and

$$\frac{X'_\Delta \hat{\Omega}_{\mathbf{S}}^{-1} \bar{X}(\mathbf{S}_{1,T})}{T_2^0 - T_2} = O_{\mathbb{P}}(1)$$

by Assumption [N.D.1](#) since $X'_\Delta \hat{\Omega}_{\mathbf{S}}^{-1} \bar{X}(\mathbf{S}_{1,T})$ involves less than $T_2^0 - T_2$ observations. Furthermore,

$$\left\| \frac{\bar{X}(\mathbf{S}_{1,T})' \hat{\Omega}_{\mathbf{S}}^{-1} X_\Delta}{T} \right\| = \left\| \frac{\bar{X}(\mathbf{S}_{1,T})' \hat{\Omega}_{\mathbf{S}}^{-1} X_\Delta}{T_2^0 - T_2} \frac{T_2^0 - T_2}{T} \right\| \leq \epsilon O_{\mathbb{P}}(1).$$

Thus L_2 is no larger than $\epsilon O_{\mathbb{P}}(1)$. Consider finally L_3 . Because both $\hat{\alpha}_2^*$ and $\hat{\alpha}_\Delta^*$ are close to α_2^0 , $\|\hat{\alpha}_2^* - \hat{\alpha}_\Delta^*\| < \rho$ with probability approaching one for any given small number $\rho > 0$. We also have

$$\left\| \frac{X'_\Delta \hat{\Omega}_{\mathbf{S}}^{-1} X_\Delta}{T_2^0 - T_2} \right\| = O_{\mathbb{P}}(1),$$

uniformly on $\mathbf{T}_\epsilon(C)$. Thus L_3 is no larger than $\rho O_{\mathbb{P}}(1)$. In summary, the following inequality holds with probability approaching one on $\mathbf{T}_\epsilon(C)$:

$$\frac{Q_{1,T} - Q_{2,T}}{T_2^0 - T_2} \geq \frac{1}{2} (\alpha_3^0 - \alpha_2^0)' \frac{X'_\Delta \hat{\Omega}_{\mathbf{S}}^{-1} X_\Delta}{T_2^0 - T_2} (\alpha_3^0 - \alpha_2^0) - \epsilon O_{\mathbb{P}}(1) - \rho O_{\mathbb{P}}(1). \quad (\text{N.11})$$

By Assumption [N.D.2](#),

$$\frac{X'_\Delta \hat{\Omega}_{\mathbf{S}}^{-1} X_\Delta}{T_2^0 - T_2}$$

has its minimum eigenvalue bounded away from zero on $\mathbf{T}_\epsilon(C)$. Thus, the first term on the right-hand side of [\(N.11\)](#) is positive and larger in absolute value than the other two terms. Thus, we have

$$\frac{Q_{1,T} - Q_{2,T}}{T_2^0 - T_2} > 0$$

with probability approaching one. This proves [\(N.8\)](#) and the proposition. $\square$

N.D.3 Proof of Proposition [N.D.3](#)

Begin by noting that

$$\hat{\Omega}_{\mathbf{S}}^{-1/2} \left(\bar{X}(\hat{\mathbf{S}}) - \bar{X}(\mathbf{S}_{0,T}) \right)$$

involves $\sum_{i=1}^{\tilde{m}} |\hat{T}_i - T_i^0| = O_{\mathbb{P}}(\tilde{m})$ nonzero observations by Proposition [N.D.2](#) so that, after applying Assumption [N.D.1](#),

$$T^{-1} \bar{X}(\hat{\mathbf{S}})' \hat{\Omega}_{\mathbf{S}}^{-1} \bar{X}(\hat{\mathbf{S}}) = T^{-1} \bar{X}(\mathbf{S}_{0,T})' \hat{\Omega}_{\mathbf{S}}^{-1} \bar{X}(\mathbf{S}_{0,T}) + O_{\mathbb{P}}(T^{-1/2}).$$

Similarly,

$$T^{-1/2} \bar{X}(\hat{\mathbf{S}})' \hat{\Omega}_{\mathbf{S}}^{-1} (\bar{X}(\mathbf{S}_{0,T}) - \bar{X}(\hat{\mathbf{S}})) \alpha_0 = O_{\mathbb{P}}(T^{-1/2})$$

and

$$T^{-1/2} \bar{X}(\hat{\mathbf{S}})' \hat{\Omega}_{\mathbf{S}}^{-1} V - T^{-1/2} \bar{X}(\mathbf{S}_{0,T})' \hat{\Omega}_{\mathbf{S}}^{-1} V = T^{-1/2} (\bar{X}(\hat{\mathbf{S}}) - \bar{X}(\mathbf{S}_{0,T}))' \hat{\Omega}_{\mathbf{S}}^{-1} V = O_{\mathbb{P}}(T^{-1/2})$$

so that another application of Assumption [N.D.1](#) yields

$$\begin{aligned} & \sqrt{T} (\hat{\alpha}(\hat{\mathbf{S}}_{T,FGLS}) - \hat{\alpha}(\mathbf{S}_{0,T})) = \left((T^{-1} \bar{X}(\mathbf{S}_{0,T})' \hat{\Omega}_{\mathbf{S}}^{-1} \bar{X}(\mathbf{S}_{0,T}))^{-1} + o_{\mathbb{P}}(1) \right) \\ & \times T^{-1/2} \left(\bar{X}(\hat{\mathbf{S}})' \hat{\Omega}_{\mathbf{S}}^{-1} (\bar{X}(\mathbf{S}_{0,T}) - \bar{X}(\hat{\mathbf{S}})) \alpha_0 + \bar{X}(\hat{\mathbf{S}})' \hat{\Omega}_{\mathbf{S}}^{-1} V - \bar{X}(\mathbf{S}_{0,T})' \hat{\Omega}_{\mathbf{S}}^{-1} V \right) = O_{\mathbb{P}}(T^{-1/2}). \quad \square \end{aligned}$$

N.E Additional Results

N.E.1 Identification of Compliers for Continuous Instrument and Examples

We extend our results to the case of a continuous instrument and illustrate their applicability in two settings: the identification of fiscal multipliers using [Ramey's \(2011\)](#) defense news shocks, and the identification of monetary policy shocks using [Romer and Romer's \(2004\)](#) narrative measure.

N.E.1.1 Identification of Compliers for Continuous Instrument

We say that observation $t \in \mathbf{S}_{0,T}$ is a complier if and only if $\mathbb{E}(D_t(z))$ is strictly increasing in z . In the present case of a continuous instrument, the policy and control samples $\mathbf{P}$ and $\mathbf{C}$ need to be redefined relative to the simpler case of a binary instrument in the main text. Let $\mathbf{P} \subset \{1, \dots, T\}$ and $\mathbf{C} \subset \{1, \dots, T\}$ satisfy $\mathbf{P} \cap \mathbf{C} = \emptyset$ and $\min_{t \in \mathbf{P}} Z_t > \max_{t \in \mathbf{C}} Z_t$. Let $\mathbf{Z}_{\mathbf{P}}$ denote the values in $\mathbf{Z}$ such that $Z_t \in \mathbf{Z}_{\mathbf{P}}$ is equivalent to $t \in \mathbf{P}$ and similarly for $\mathbf{Z}_{\mathbf{C}}$. Construct $\mathbf{P}$ and $\mathbf{C}$ such that $\underline{z}_{\mathbf{P}} \equiv \inf(\mathbf{Z}_{\mathbf{P}}) > \sup(\mathbf{Z}_{\mathbf{C}}) \equiv \bar{z}_{\mathbf{C}}$. For example, a simple way to define the policy and control samples is $\mathbf{P} = \{t \in \{1, \dots, T\} : Z_t \geq \tilde{z} + \epsilon\}$ and $\mathbf{C} = \{t \in \{1, \dots, T\} : Z_t \leq \tilde{z} - \epsilon\}$ for some $\tilde{z}$ and small $\epsilon > 0$. With these definitions, we impose the continuous instrument-analogs of Assumptions [2.7-2.8](#) in the main text.

Assumption N.E.1. (i) For any $t \in \mathbf{C}$, $\overline{D}_{C,t,n} \xrightarrow{\mathbb{P}} \mathbb{E}(D_t(Z_t)|Z_t \in \mathbf{Z}_\mathbf{C})$ as $n \rightarrow \infty$ with $n/|\mathbf{C}| \rightarrow 0$.
(ii) For $t \in \mathbf{P}$, $\mathbb{E}[D_{t-1}(Z_{t-1})|Z_{t-1} \in \mathbf{Z}_\mathbf{C}] = \mathbb{E}[D_t(Z_t)|Z_t \in \mathbf{Z}_\mathbf{C}]$.

Assumption N.E.2. (i) For any $t \in \mathbf{P}$, $\overline{D}_{P,t,n} \xrightarrow{\mathbb{P}} \mathbb{E}(D_t(Z_t)|Z_t \in \mathbf{Z}_\mathbf{P})$ as $n \rightarrow \infty$ with $n/|\mathbf{P}| \rightarrow 0$.
(ii) For $t \in \mathbf{C}$, $\mathbb{E}[D_t(Z_t)|Z_t \in \mathbf{Z}_\mathbf{P}] = \mathbb{E}[D_{s^*(t)}(Z_{s^*(t)})|Z_{s^*(t)} \in \mathbf{Z}_\mathbf{P}]$, where $s^*(t) = \arg \min_{s \in \mathbf{P}} |t - s|$.

Let $F_{\tilde{V}_t}$ denote the distribution function of $\tilde{V}_t$.

Proposition N.E.1. Let Assumptions 2.1, N.E.1 and N.E.2 hold, $n_0, n_1 \rightarrow \infty$ with $n_0/|\mathbf{C}|, n_1/|\mathbf{P}| \rightarrow 0$, and for $z' > z$, $\mathbb{E}(D_t(z')|\tilde{V}_t = \tilde{v}) \geq \mathbb{E}(D_t(z)|\tilde{V}_t = \tilde{v})$ for $F_{\tilde{V}_t}$ -almost every $\tilde{v}$ with strict inequality on a set of positive $F_{\tilde{V}_t}$ -measure. We have:

- (i) if $t \in \mathbf{P}$ is a complier, then $\overline{D}_{P,t,n_1} - \overline{D}_{C,t-1,n_0} \xrightarrow{\mathbb{P}} c$ where $c > 0$.
- (ii) if $t \in \mathbf{C}$ is a complier, then $\overline{D}_{P,s^*(t),n_1} - \overline{D}_{C,t,n_0} \xrightarrow{\mathbb{P}} \tilde{c}$ where $\tilde{c} > 0$.

Proof of Proposition N.E.1. Consider first the policy sample. Suppose $t \in \mathbf{P}$ is a complier. Then, by Assumptions N.E.1(i) and N.E.2(i), as $n_0, n_1 \rightarrow \infty$,

$$\begin{aligned}
\overline{D}_{P,t,n_1} - \overline{D}_{C,t-1,n_0} &\xrightarrow{\mathbb{P}} \mathbb{E}(D_t(Z_t)|Z_t \in \mathbf{Z}_\mathbf{P}) - \mathbb{E}(D_{t-1}(Z_{t-1})|Z_{t-1} \in \mathbf{Z}_\mathbf{C}) \\
&= \mathbb{E}(D_t(Z_t)|Z_t \in \mathbf{Z}_\mathbf{P}) - \mathbb{E}(D_t(Z_t)|Z_t \in \mathbf{Z}_\mathbf{C}) \\
&= \int_{\tilde{v}} \mathbb{E}(D_t(Z_t)|Z_t \in \mathbf{Z}_\mathbf{P}, \tilde{V}_t = \tilde{v}) dF_{\tilde{V}_t}(\tilde{v}) - \int_{\tilde{v}} \mathbb{E}(D_t(Z_t)|Z_t \in \mathbf{Z}_\mathbf{C}, \tilde{V}_t = \tilde{v}) dF_{\tilde{V}_t}(\tilde{v}) \\
&= \int_{\tilde{v}} \int_{z \in \mathbf{Z}_\mathbf{P}} \mathbb{E}(D_t(z)|\tilde{V}_t = \tilde{v}) dF_{Z_t|\tilde{V}_t=\tilde{v}, Z_t \in \mathbf{Z}_\mathbf{P}}(z) dF_{\tilde{V}_t}(\tilde{v}) \\
&\quad - \int_{\tilde{v}} \int_{z \in \mathbf{Z}_\mathbf{C}} \mathbb{E}(D_t(z)|\tilde{V}_t = \tilde{v}) dF_{Z_t|\tilde{V}_t=\tilde{v}, Z_t \in \mathbf{Z}_\mathbf{C}}(z) dF_{\tilde{V}_t}(\tilde{v}) \\
&\geq \int_{\tilde{v}} \int_{z \in \mathbf{Z}_\mathbf{P}} \mathbb{E}(D_t(\underline{z}_\mathbf{P})|\tilde{V}_t = \tilde{v}) dF_{Z_t|\tilde{V}_t=\tilde{v}, Z_t \in \mathbf{Z}_\mathbf{P}}(z) dF_{\tilde{V}_t}(\tilde{v}) \\
&\quad - \int_{\tilde{v}} \int_{z \in \mathbf{Z}_\mathbf{C}} \mathbb{E}(D_t(\bar{z}_\mathbf{C})|\tilde{V}_t = \tilde{v}) dF_{Z_t|\tilde{V}_t=\tilde{v}, Z_t \in \mathbf{Z}_\mathbf{C}}(z) dF_{\tilde{V}_t}(\tilde{v}) \\
&= \mathbb{E}(D_t(\underline{z}_\mathbf{P})) - \mathbb{E}(D_t(\bar{z}_\mathbf{C})) > 0,
\end{aligned}$$

where $F_{Z_t|\tilde{V}_t=\tilde{v}, Z_t \in \mathbf{Z}_\mathbf{P}}(\cdot)$ is the conditional distribution function of Z_t given $\tilde{V}_t = \tilde{v}$ and $Z_t \in \mathbf{Z}_\mathbf{P}$, $F_{Z_t|\tilde{V}_t=\tilde{v}, Z_t \in \mathbf{Z}_\mathbf{C}}(\cdot)$ is the conditional distribution function of Z_t given $\tilde{V}_t = \tilde{v}$ and $Z_t \in \mathbf{Z}_\mathbf{C}$, the first equality follows from Assumption N.E.1(ii), the third equality follows from Assumption 2.1 and the inequalities follow from conditional monotonicity and the separation of $\mathbf{Z}_\mathbf{P}$ and $\mathbf{Z}_\mathbf{C}$. The proof for the control sample is entirely analogous and therefore omitted. $\square$

Assumption N.E.3. (Continuous Case Monotonicity) $\mathbb{E}(D_t(z))$ is monotonic in z .

Under Assumption [N.E.3](#), assume without loss of generality that $\mathbb{E}(D_t(z))$ is increasing in z . We obtain the following characterization of compliers under monotonicity for the continuous instrument case.

Corollary N.E.1. *Let Assumptions [2.3](#) without conditioning on $\tilde{V}_t$, and [N.E.3](#) hold. If, in addition, the first-stage function is strictly increasing in z for each $t \in \mathbf{S}_{0,T}$, the set of compliers coincides with $\mathbf{S}_{0,T}$.*

Proof of Corollary [N.E.1](#). Assumption [N.E.3](#) rules out defiers, i.e., $\mathbb{E}(D_t(z))$ being strictly decreasing in z , so that non-compliers are characterized by $\mathbb{E}(D_t(z))$ being constant in z . Therefore, a non-complier cannot belong to $\mathbf{S}_{0,T}$ by definition. And any complier belongs to $\mathbf{S}_{0,T}$ by definition. $\square$

N.E.1.2 Identification of fiscal Multipliers [cf. [Ramey \(2011\)](#)]

[Ramey \(2011\)](#) studies how to correctly identify exogenous government spending shocks in order to estimate the fiscal multiplier. The author constructs a new external instrument: the defense news shock, based on narrative information extracted from Business Week and other sources. For each quarter, she reads all articles discussing military contracts, wars, geopolitical events, or defense policies that are likely to affect future U.S. defense spending. She then constructs a measure of the expected present value of changes in future defense spending induced by those events.

Let Z_t denote the defense news shock at time t . Specifically, Z_t is an observable proxy for exogenous revisions in expectations about future fiscal policy, assumed to be independent of contemporaneous shocks to output and other endogenous variables. Values $Z_t > 0$ correspond to anticipated military buildups (e.g., the U.S. entry into WWII), $Z_t < 0$ correspond to anticipated military drawdowns (e.g., the Vietnam drawdown in the 1970s), and $Z_t = 0$ indicates the absence of defense-related news.

Let D_t denote the average of current and future government spending from time t to $t+m$, for some integer $m > 0$. Ramey (2011) documents that defense spending begins to rise approximately 4–6 quarters after a news shock, with the peak response occurring roughly 6–8 quarters thereafter. Accordingly, we define $D_t = \frac{1}{m+1} \sum_{h=0}^m G_{t+h}$ where G_{t+h} denotes government spending at time $t+h$. This delayed response reflects well-known institutional features of military spending—such as procurement, planning, and production lags—which generate a long implementation horizon.

A complier is defined as an observation t where $\mathbb{E}(D_t(z))$ is strictly increasing in z . Formally, date t is a complier if the expected average of government spending over the next m periods increases (respectively, decreases) if and only if Z_t is positive (respectively, negative). In

this continuous-instrument case, we define t to be a complier whenever the corresponding signed conditional-mean contrast is strictly positive.

A noncomplier is a time period in which an expansionary (respectively, contractionary) defense news shock does not translate into higher (respectively, lower) government spending in the subsequent m periods. These episodes reflect situations where the defense news shock is dominated by other general equilibrium forces or non-defense news in expectation.

Let $\mathbf{P}_+ \subset \{1, \dots, T\}$ denote the set of dates for which $Z_t > 0$, and let $\mathbf{P}_- \subset \{1, \dots, T\}$ denote the set of dates for which $Z_t < 0$. Let $\mathbf{C} \subset \{1, \dots, T\}$ denote the set of dates for which $Z_t = 0$. Thus, $\mathbf{P}_+$ corresponds to the subsample of expansionary defense news shocks, $\mathbf{P}_-$ corresponds to the subsample of contractionary defense news shocks, and $\mathbf{C}$ corresponds to the subsample with no defense news shocks. Note that $\mathbf{P} \cap \mathbf{C} = \emptyset$ where $\mathbf{P} = \mathbf{P}_+ \cup \mathbf{P}_-$ and the ordering of shock realizations implies $\min_{t \in \mathbf{P}_+} Z_t > \max_{t \in \mathbf{C}} Z_t$ and $\min_{t \in \mathbf{C}} Z_t > \max_{t \in \mathbf{P}_-} Z_t$. Define $\mathbf{Z}_{\mathbf{P}_+}$ and $\mathbf{Z}_{\mathbf{P}_-}$ as the set of values in $\mathbf{Z}$ such that $Z_t \in \mathbf{Z}_{\mathbf{P}_+}$ for all $t \in \mathbf{P}_+$ and $Z_t \in \mathbf{Z}_{\mathbf{P}_-}$ for all $t \in \mathbf{P}_-$, respectively.

Let $\bar{D}_{P_+, t_0, n_1} = n_1^{-1} \sum_{s \in N_+(t_0)} D_s$ and $N_+(t_0)$ denotes the n_1 largest indices $s \in \mathbf{P}_+$ such that $s \leq t_0$ and $\bar{D}_{P_-, t_0, n_1} = n_1^{-1} \sum_{s \in N_-(t_0)} D_s$ and $N_-(t_0)$ denotes the n_1 largest indices $s \in \mathbf{P}_-$ such that $s \leq t_0$.

Replace Assumptions [N.E.1-N.E.2](#) by the following two assumptions.

Assumption N.E.4. (i) For any $t \in \mathbf{C}$, $\bar{D}_{C, t, n} \xrightarrow{\mathbb{P}} \mathbb{E}(D_t(Z_t) | Z_t \in \mathbf{Z}_{\mathbf{C}})$ as $n \rightarrow \infty$ with $n/|\mathbf{C}| \rightarrow 0$. (ii) For any $t \in \mathbf{P}_+ \cup \mathbf{P}_-$, $\mathbb{E}[D_t(Z_t) | Z_t \in \mathbf{Z}_{\mathbf{C}}] = \mathbb{E}[D_{s_C(t)}(Z_{s_C(t)}) | Z_{s_C(t)} \in \mathbf{Z}_{\mathbf{C}}]$, where $s_C(t) = \arg \min_{s \in \mathbf{C}} |t - s|$ denotes the closest index in $\mathbf{C}$ to t .

Assumption N.E.5. (i) For any $t \in \mathbf{P}_+$, $\bar{D}_{P_+, t, n} \xrightarrow{\mathbb{P}} \mathbb{E}(D_t(Z_t) | Z_t \in \mathbf{Z}_{\mathbf{P}_+})$ as $n \rightarrow \infty$ with $n/|\mathbf{P}_+| \rightarrow 0$. (ii) For any $t \in \mathbf{P}_-$, $\bar{D}_{P_-, t, n} \xrightarrow{\mathbb{P}} \mathbb{E}(D_t(Z_t) | Z_t \in \mathbf{Z}_{\mathbf{P}_-})$ as $n \rightarrow \infty$ with $n/|\mathbf{P}_-| \rightarrow 0$. (iii) For $t \in \mathbf{C}$, $\mathbb{E}[D_t(Z_t) | Z_t \in \mathbf{Z}_{\mathbf{P}_+}] = \mathbb{E}[D_{s_+(t)}(Z_{s_+(t)}) | Z_{s_+(t)} \in \mathbf{Z}_{\mathbf{P}_+}]$ where $s_+(t) = \arg \min_{s \in \mathbf{P}_+} |t - s|$, and $\mathbb{E}[D_t(Z_t) | Z_t \in \mathbf{Z}_{\mathbf{P}_-}] = \mathbb{E}[D_{s_-(t)}(Z_{s_-(t)}) | Z_{s_-(t)} \in \mathbf{Z}_{\mathbf{P}_-}]$, where $s_-(t) = \arg \min_{s \in \mathbf{P}_-} |t - s|$.

Assumption [N.E.4](#)(i) requires that the local average $\bar{D}_{C, t, n}$ consistently estimates the conditional mean of the potential treatment path on the support $\mathbf{Z}_{\mathbf{C}}$. This is a standard smoothness requirement: as the averaging window n grows—but remains small relative to the size of $\mathbf{C}$ —the local neighborhood around each $t \in \mathbf{C}$ becomes dense enough to recover the underlying conditional expectation.

Assumption [N.E.4](#)(ii) ensures that the counterfactual mean $\mathbb{E}[D_t(Z_t) | Z_t \in \mathbf{Z}_{\mathbf{C}}]$ for any expansionary-shock date $t \in \mathbf{P}_+$ can be reliably approximated using the nearest no-shock date $s_C(t)$. This is a smoothness or continuity restriction across adjoining regions of the support: the conditional

mean of the potential treatment path at a shock date should not jump discontinuously relative to its value at nearby no-shock dates. Economically, it imposes that the latent drivers of $D_t(z)$ evolve smoothly over time, so that adjacent no-shock observations provide a credible proxy for the counterfactual treatment path at a shock date.

Together, these conditions allow the no-shock subsample $\mathbf{C}$ to function as a consistent local benchmark for imputing the counterfactual behavior of D_t at dates affected by defense news shocks.

Assumption N.E.5 imposes mild smoothness and local-averaging conditions that ensure the conditional expectations of the potential treatment path $D_t(Z_t)$ are well-behaved across the supports corresponding to expansionary shocks ($\mathbf{P}_+$), contractionary shocks ($\mathbf{P}_-$), and no-shock periods ($\mathbf{C}$).

Parts (i) and (ii) require that the local averages $\bar{D}_{P_+,t,n}$ and $\bar{D}_{P_-,t,n}$ consistently estimate the corresponding conditional mean functions on the supports $\mathbf{Z}_{P_+}$ and $\mathbf{Z}_{P_-}$. These conditions mirror standard requirements in nonparametric and local-smoothing arguments: the window size n must grow, but slowly relative to the size of the relevant subsample, so that the averaging neighborhood becomes dense on the support while remaining asymptotically local.

Part (iii) ensures comparability between no-shock observations and nearby shock observations. For any $t \in \mathbf{C}$, the conditional expectation must coincide with the corresponding conditional expectation at the closest expansionary (or contractionary) shock date. This requirement is a smoothness or stability restriction: moving from a no-shock date to a nearby shock date should not induce a discontinuous jump in the conditional mean of the potential treatment path. It formalizes the idea that counterfactual treatment levels at no-shock dates can be imputed reliably using shock dates in their immediate neighborhood.

Corollary N.E.2. *Let Assumptions 2.1, N.E.4-N.E.5 hold and $n_0, n_1 \rightarrow \infty$ with $n_0/|\mathbf{C}|, n_1/|\mathbf{P}| \rightarrow 0$. Then:*

- (i) $t \in \mathbf{P}_+$ is a complier if and only if $\bar{D}_{P_+,t,n_1} - \bar{D}_{C,s_C(t),n_0} \xrightarrow{\mathbb{P}} c$ where $c > 0$.
- (ii) $t \in \mathbf{P}_-$ is a complier if and only if $\bar{D}_{P_-,t,n_1} - \bar{D}_{C,s_C(t),n_0} \xrightarrow{\mathbb{P}} c_-$ where $c_- < 0$.
- (iii) $t \in \mathbf{C}$ is a complier if and only if $\bar{D}_{P_+,s_+(t),n_1} - \bar{D}_{C,t,n_0} \xrightarrow{\mathbb{P}} \tilde{c}_+$ where $\tilde{c}_+ > 0$ and $\bar{D}_{P_-,s_-(t),n_1} - \bar{D}_{C,t,n_0} \xrightarrow{\mathbb{P}} \tilde{c}_-$ where $\tilde{c}_- < 0$.

The proof of the corollary is analogous to that of Proposition N.C.1. This result shows that compliers can be characterized entirely in terms of local deviations of the realized treatment path D_t , from its counterfactual evolution around no-shock dates. For expansionary shock dates $t \in \mathbf{P}_+$, compliers correspond to observing a strictly positive asymptotic gap between the local average of D_t within the expansionary-shock support and the corresponding local counterfactual average drawn

from the closest no-shock date. Analogously, being compliers at contractionary shock dates $t \in \mathbf{P}_-$ requires that this asymptotic gap be strictly negative.

For dates with no defense news $t \in \mathbf{C}$, compliers amount to the existence of a nearby expansionary (respectively, contractionary) shock date whose local average treatment level differs from the local no-shock average by a strictly positive (respectively, strictly negative) amount. Intuitively, the behavior of D_t at $t \in \mathbf{C}$ must be consistent with what would be observed if a small expansionary or contractionary news shock were applied in the neighborhood of that date.

Taken together, the three cases establish that compliance is equivalent to observing a persistent, sign-consistent divergence between the realized treatment path and its smoothly evolving counterfactual counterpart. This links the structural definition of compliers to a set of estimable conditions based on local averages of the treatment variable across the supports $\mathbf{P}_+$, $\mathbf{P}_-$ and $\mathbf{C}$.

N.E.1.3 Identification of Monetary Policy based on [Romer and Romer \(2004\)](#)

Romer and Romer (2004) construct a narrative monetary policy shock series using a two-step procedure that isolates the component of the federal funds rate target change that is orthogonal to the Federal Reserve's internal forecasts of macroeconomic conditions. The key goal is to obtain a measure of unexpected changes in monetary policy that is not contaminated by systematic responses of the Fed to expected economic developments.

Let Z_t denote the Romer-Romer shock at time t . The Romer-Romer shock is never exactly zero. Thus, values of the Romer-Romer shock that are close to zero are considered as no monetary surprise. Let $\delta > 0$ denote some threshold. Values $Z_t > \delta$ correspond to unexpected contractionary monetary policy news, $Z_t < -\delta$ correspond to unexpected expansionary monetary policy news, and $Z_t \in [-\delta, \delta]$ indicates the absence of monetary policy news. The threshold δ could be set equal to the 25th quantile of the distribution of $|Z_t|$, for example. For the original Romer-Romer shocks series, the 25th quantile is 0.063. Recall that the Romer-Romer shock is constructed at FOMC meeting-frequency and then transformed into monthly frequency.

Z_t is believed to be a proxy for the true monetary policy shock. Hence, D_t could be the change in a short-term Treasury yield or the VAR residual of the Fed fund rate equation. Note that for what concerns the identification of the compliers, Z_t need not satisfy the exclusion restriction. Thus, D_t could be the change in the Fed fund rate. Indeed, the identification of the compliers is then use to test the exclusion restriction.

A complier is defined as an observation t where $\mathbb{E}(D_t(z))$ is strictly increasing in z . Formally, date t is a complier if the short-term Treasury yield increases (respectively, decreases) in expectation if and only if $Z_t > \delta$ (respectively, $Z_t < -\delta$).

A noncomplier is a time period in which an unexpected expansionary (respectively, contractionary) monetary policy shock does not translate into higher (respectively, lower) short-term Treasury yield. These episodes reflect situations where the monetary policy news is dominated by other general equilibrium forces or non-monetary news on average.

Let $\mathbf{P}_+ \subset \{1, \dots, T\}$ denote the set of dates for which $Z_t > \delta$, and let $\mathbf{P}_- \subset \{1, \dots, T\}$ denote the set of dates for which $Z_t < -\delta$. Let $\mathbf{C} \subset \{1, \dots, T\}$ denote the set of dates for which $Z_t \in [-\delta, \delta]$. Thus, $\mathbf{P}_+$ corresponds to the subsample of expansionary monetary news, $\mathbf{P}_-$ corresponds to the subsample of contractionary monetary news, and $\mathbf{C}$ corresponds to the subsample with no monetary news. Note that $\mathbf{P} \cap \mathbf{C} = \emptyset$ where $\mathbf{P} = \mathbf{P}_+ \cup \mathbf{P}_-$ and the ordering of shock realizations implies $\min_{t \in \mathbf{P}_+} Z_t > \max_{t \in \mathbf{C}} Z_t$ and $\min_{t \in \mathbf{C}} Z_t > \max_{t \in \mathbf{P}_-} Z_t$. Define $\mathbf{Z}_{\mathbf{P}_+}$ and $\mathbf{Z}_{\mathbf{P}_-}$ as the set of values in $\mathbf{Z}$ such that $Z_t \in \mathbf{Z}_{\mathbf{P}_+}$ for all $t \in \mathbf{P}_+$ and $Z_t \in \mathbf{Z}_{\mathbf{P}_-}$ for all $t \in \mathbf{P}_-$, respectively.

Let $\overline{D}_{C,t,n}$, $\overline{D}_{P_+,t_0,n_1}$ and $\overline{D}_{P_-,t_0,n_1}$ be defined as in Section N.E.1.2. Under the same assumptions used to identify compliers in the fiscal policy setting from Section N.E.1.2, we present the complier identification result based on Romer-Romer shocks.

Corollary N.E.3. *Let Assumptions 2.1, N.E.4-N.E.5 hold and $n_0, n_1 \rightarrow \infty$ with $n_0/|\mathbf{C}|, n_1/|\mathbf{P}| \rightarrow 0$. Then:*

- (i) $t \in \mathbf{P}_+$ is a complier if and only if $\overline{D}_{P_+,t,n_1} - \overline{D}_{C,s_C(t),n_0} \xrightarrow{\mathbb{P}} c$ where $c > 0$.
- (ii) $t \in \mathbf{P}_-$ is a complier if and only if $\overline{D}_{P_-,t,n_1} - \overline{D}_{C,s_C(t),n_0} \xrightarrow{\mathbb{P}} c_-$ where $c_- < 0$.
- (iii) $t \in \mathbf{C}$ is a complier if and only if $\overline{D}_{P_+,s_+(t),n_1} - \overline{D}_{C,t,n_0} \xrightarrow{\mathbb{P}} \tilde{c}$ where $\tilde{c}_+ > 0$ and $\overline{D}_{P_-,s_-(t),n_1} - \overline{D}_{C,t,n_0} \xrightarrow{\mathbb{P}} \tilde{c}_-$ where $\tilde{c}_- < 0$.

N.E.2 Primitive Conditions on IVs, Exogenous Regressors and Errors for the Assumptions of Section 6.2

Assumptions 6.1-6.3 of Section 6.2 are implied by any one of the following assumptions:

Assumption N.E.6. $\{(v_t, w_t) : t \geq 1\}$ are i.i.d., $\mathbb{E}(v_t \otimes w_t) = 0$, $\mathbb{E}\|v_t\|^2 + \mathbb{E}\|w_t\|^2 + \mathbb{E}\|v_t \otimes w_t\|^2 < \infty$, $\Sigma_v = \mathbb{E}(v_t v_t')$ is positive definite, and uniformly in $\mathbf{S}_T, \mathbf{S}_T' \in \mathcal{S}$, for $\mathbf{S} = \lim_{T \rightarrow \infty} T^{-1} \mathbf{S}_T$ and $\mathbf{S}' = \lim_{T \rightarrow \infty} T^{-1} \mathbf{S}_T'$, $T^{-1} \sum_{t=1}^T \mathbb{E}(w_t(\mathbf{S}_T) w_t(\mathbf{S}_T')') \rightarrow Q(\mathbf{S}, \mathbf{S}')$ for some $(q+p) \times (q+p)$ matrix $Q(\mathbf{S}, \mathbf{S}')$ for which $Q(\mathbf{S}, \mathbf{S})$ is positive definite and $T^{-1} \sum_{t=1}^T \mathbb{E}((v_t \otimes w_t(\mathbf{S}_T)) (v_t \otimes w_t(\mathbf{S}_T'))') \rightarrow \Psi(\mathbf{S}, \mathbf{S}')$ for some $2(q+p) \times 2(q+p)$ matrix $\Psi(\mathbf{S}, \mathbf{S}')$.

Assumption N.E.7. $\{(v_t, w_t) : t \geq 1\}$ are independent, $\mathbb{E}(v_t \otimes w_t) = 0$ for all $t \geq 1$, $\sup_{t \geq 1} (\mathbb{E}\|v_t\|^{2+\varsigma} + \mathbb{E}\|w_t\|^{2+\varsigma} + \mathbb{E}\|v_t \otimes w_t\|^{2+\varsigma}) < \infty$ for some $\varsigma > 0$, $T^{-1} \sum_{t=1}^T \mathbb{E}(v_t v_t') \rightarrow \Sigma_v$ for some positive definite

2×2 matrix Σ_v , and uniformly in $\mathbf{S}_T, \mathbf{S}'_T \in \mathcal{S}$, for $\mathbf{S} = \lim_{T \rightarrow \infty} T^{-1} \mathbf{S}_T$ and $\mathbf{S}' = \lim_{T \rightarrow \infty} T^{-1} \mathbf{S}'_T$, $T^{-1} \sum_{t=1}^T \mathbb{E}(w_t(\mathbf{S}_T) w_t(\mathbf{S}'_T)') \rightarrow Q(\mathbf{S}, \mathbf{S}')$ for some $(q+p) \times (q+p)$ matrix $Q(\mathbf{S}, \mathbf{S}')$ for which $Q(\mathbf{S}, \mathbf{S})$ is positive definite and $T^{-1} \sum_{t=1}^T \mathbb{E}((v_t \otimes w_t(\mathbf{S}_T)) (v_t \otimes w_t(\mathbf{S}'_T))') \rightarrow \Psi(\mathbf{S}, \mathbf{S}')$ for some $2(q+p) \times 2(q+p)$ matrix $\Psi(\mathbf{S}, \mathbf{S}')$.

Assumption N.E.8. $\{(v_t \otimes w_t, \mathcal{F}_t) : t \geq 1\}$ is a martingale difference sequence, where $\mathcal{F}_t = \sigma(v_t, w_t, v_{t-1}, w_{t-1}, \dots)$, $\{(v_t \otimes w_t) : t \geq 1\}$ is an ergodic sequence, $\sup_{t \geq 1} (\mathbb{E}||v_t||^2 + \mathbb{E}||w_t||^2 + \mathbb{E}||v_t \otimes w_t||^2) < \infty$, $\Sigma_v = \mathbb{E}(v_t v_t')$ is positive definite, and uniformly in $\mathbf{S}_T, \mathbf{S}'_T \in \mathcal{S}$, for $\mathbf{S} = \lim_{T \rightarrow \infty} T^{-1} \mathbf{S}_T$ and $\mathbf{S}' = \lim_{T \rightarrow \infty} T^{-1} \mathbf{S}'_T$, $T^{-1} \sum_{t=1}^T \mathbb{E}(w_t(\mathbf{S}_T) w_t(\mathbf{S}'_T)') \rightarrow Q(\mathbf{S}, \mathbf{S}')$ for some $(q+p) \times (q+p)$ matrix $Q(\mathbf{S}, \mathbf{S}')$ for which $Q(\mathbf{S}, \mathbf{S})$ is positive definite and $T^{-1} \sum_{t=1}^T \mathbb{E}((v_t \otimes w_t(\mathbf{S}_T)) (v_t \otimes w_t(\mathbf{S}'_T))') \rightarrow \Psi(\mathbf{S}, \mathbf{S}')$ for some $2(q+p) \times 2(q+p)$ matrix $\Psi(\mathbf{S}, \mathbf{S}')$.

Assumption N.E.9. $\{(v_t, w_t) : t = \dots, 0, 1, \dots\}$ is a doubly infinite ergodic sequence with $\mathbb{E}(v_t \otimes w_t) = 0$, $\sup_{t \geq 1} (\mathbb{E}||v_t||^2 + \mathbb{E}||w_t||^2 + \mathbb{E}||v_t \otimes w_t||^2) < \infty$, $\sup_{t \geq 1} \sum_{j=1}^{\infty} (\mathbb{E}||\mathbb{E}(v_t \otimes w_t | \mathcal{F}_{t-j})||^2)^{1/2} < \infty$ where $\mathcal{F}_t = \sigma(v_t, w_t, v_{t-1}, w_{t-1}, \dots)$, $T^{-1} \sum_{t=1}^T \mathbb{E}(v_t v_t') \rightarrow \Sigma_v$ for some positive definite 2×2 matrix Σ_v , and uniformly in $\mathbf{S}_T, \mathbf{S}'_T \in \mathcal{S}$, for $\mathbf{S} = \lim_{T \rightarrow \infty} T^{-1} \mathbf{S}_T$ and $\mathbf{S}' = \lim_{T \rightarrow \infty} T^{-1} \mathbf{S}'_T$, $T^{-1} \sum_{t=1}^T \mathbb{E}(w_t(\mathbf{S}_T) w_t(\mathbf{S}'_T)') \rightarrow Q(\mathbf{S}, \mathbf{S}')$ for some $(q+p) \times (q+p)$ matrix $Q(\mathbf{S}, \mathbf{S}')$ for which $Q(\mathbf{S}, \mathbf{S})$ is positive definite and $T^{-1} \sum_{t=1}^T \sum_{j=-\infty}^{\infty} \mathbb{E}((v_t \otimes w_t(\mathbf{S}_T)) (v_{t-j} \otimes w_{t-j}(\mathbf{S}'_T))') \rightarrow \Psi(\mathbf{S}, \mathbf{S}') = \int_0^1 \Psi_u(\mathbf{S}, \mathbf{S}') du$, where $\Psi_u(\mathbf{S}, \mathbf{S}')$ is the local long-run covariance matrix of $v_t \otimes w_t(\mathbf{S}_T)$ and $v_t \otimes w_t(\mathbf{S}'_T)$.

The random vectors $\{(v_t, w_t) : t = \dots, 0, 1, \dots\}$ are uncorrelated under Assumptions N.E.6-N.E.8, but are (possibly) correlated under Assumption N.E.9. Assumptions N.E.7-N.E.9 allow for nonstationarity (i.e., time-varying moments). In particular, they are satisfied by segmented local stationarity [see Casini (2024, 2023)].

If the errors are conditionally homoskedastic and $\{(v_t, w_t) : t \geq 1\}$ are uncorrelated, the following assumption holds.

Assumption N.E.10. $\Psi(\mathbf{S}, \mathbf{S}') = \Sigma_v \otimes Q(\mathbf{S}, \mathbf{S}')$, where $Q(\cdot)$ is defined in Assumption 6.1 and $\Psi(\cdot)$ is defined in Assumption 6.3.

This assumption is implied by any one of Assumptions N.E.6, N.E.7, and N.E.8 plus the following.

Assumption N.E.11. $\mathbb{E}((v_t v_t') \otimes (w_t(\mathbf{S}_T) w_t(\mathbf{S}'_T)')) = \Sigma_v \otimes \tilde{Q}(\mathbf{S}_T, \mathbf{S}'_T)$ for all $t \geq 1$ and $\mathbf{S}_T, \mathbf{S}'_T \in \mathcal{S}$.

By iterated expectations, a sufficient condition for Assumption N.E.11 is $\mathbb{E}(v_t v_t' | w_t(\mathbf{S}_T), w_t(\mathbf{S}'_T)) = \mathbb{E}(v_t v_t') = \Sigma_v$ a.s. for all $\mathbf{S}_T, \mathbf{S}'_T \in \mathcal{S}$ and all $t \geq 1$. Note that Assumptions N.E.8 and N.E.9 allow for intertemporal conditional heteroskedasticity even when Assumption N.E.11 holds. The following lemma summarizes the relations between the assumptions.

Lemma N.E.1. (i) Any one of Assumptions [N.E.6](#), [N.E.7](#), [N.E.8](#) and [N.E.9](#) implies Assumptions [6.1-6.3](#);
 (ii) Any one of Assumptions [N.E.6](#), [N.E.7](#), [N.E.8](#) plus Assumption [N.E.11](#) imply Assumption [N.E.10](#).

Proof of Lemma N.E.1. Although $w(\mathbf{S}_T)$ is a function of the partition $\mathbf{S}_T$, Assumptions [6.1-6.3](#) are stated uniformly over the admissible partition class $\mathcal{S}$. Therefore, the required law of large numbers and weak convergence hold uniformly in $\mathbf{S}_T$, so the partition dependence does not affect the argument. The result then follows immediately from Lemma 4 in [Andrews, Moreira, and Stock \(2004\)](#), which establishes the corresponding conclusion for full-sample averages. Their proof relies only on the convergence statements in Assumptions [6.1-6.3](#) and not on stationarity per se. $\square$

N.E.3 Consistent Covariance Matrix Estimation

Let $V_{b,t}(\mathbf{S}_T) = v_t' b_0 \bar{Z}_t(C_T)$ and $V_{a,t}(\mathbf{S}_T) = v_t' \Sigma_v^{-1} a_{0,\beta} \bar{Z}_t(C_T)$. We have

$$\begin{aligned}\Sigma_{N_1}(\mathbf{S}) &= \lim_{T \rightarrow \infty} T^{-1} \sum_{t=1}^T \sum_{r=1}^T \mathbb{E} \left(V_{b,t}(\mathbf{S}_T) V_{b,r}(\mathbf{S}_T)' \right), \\ \Sigma_{N_1 N_2}(\mathbf{S}) &= \lim_{T \rightarrow \infty} T^{-1} \sum_{t=1}^T \sum_{r=1}^T \mathbb{E} \left(V_{b,t}(\mathbf{S}_T) V_{a,r}(\mathbf{S}_T)' \right), \\ \Sigma_{N_2}^*(\mathbf{S}) &= \lim_{T \rightarrow \infty} T^{-1} \sum_{t=1}^T \sum_{r=1}^T \mathbb{E} \left(V_{a,t}(\mathbf{S}_T) V_{a,r}(\mathbf{S}_T)' \right).\end{aligned}$$

Let $\hat{V}_{b,t}(\mathbf{S}_T) = \hat{v}_t(\mathbf{S}_T)' b_0 \bar{Z}_t(C_T)$ and $\hat{V}_{a,t}(s) = \hat{v}_t(\mathbf{S}_T)' \hat{\Sigma}_v^{-1} a_{0,\beta} \bar{Z}_t(C_T)$. We consider both HAC and double-kernel HAC (DK-HAC) estimators of $\Sigma_{v\bar{Z}}(\mathbf{S})$. Here we discuss the HAC estimators of [Newey and West \(1987\)](#) and [Andrews \(1991\)](#). The DK-HAC estimator was recently proposed by [Casini \(2023\)](#). It is consistent under both the null and the alternative so that tests based on it do not suffer from power losses induced by nonstationarity [cf. [Casini, Deng, and Perron \(2025\)](#)].

The HAC estimators are defined as

$$\begin{aligned}\hat{\Sigma}_{N_1}(\mathbf{S}_T) &= \frac{T}{T - q - p} \sum_{k=-T+1}^{T-1} K_1(b_{1,T}k) \hat{\Gamma}_{bb}(k, \mathbf{S}_T), \\ \text{with} \quad \hat{\Gamma}_{bb}(k, \mathbf{S}_T) &= \begin{cases} T^{-1} \sum_{t=k+1}^T \hat{V}_{b,t}(\mathbf{S}_T) \hat{V}_{b,t-k}'(\mathbf{S}_T), & k \geq 0 \\ T^{-1} \sum_{t=-k+1}^T \hat{V}_{b,t+k}(\mathbf{S}_T) \hat{V}_{b,t}'(\mathbf{S}_T), & k < 0 \end{cases},\end{aligned}$$

and $\hat{\Sigma}_{N_2}^*(\mathbf{S}_T)$ and $\hat{\Sigma}_{N_1 N_2}(\mathbf{S}_T)$ are defined analogously to $\hat{\Sigma}_{N_1}(\mathbf{S}_T)$ after replacing $\hat{\Gamma}_{bb}(k, \mathbf{S}_T)$ with

analogous quantities $\hat{\Gamma}_{aa}(k, \mathbf{S}_T)$ and $\hat{\Gamma}_{ab}(k, \mathbf{S}_T)$, respectively. We consider the following class of kernels

$$\mathbf{K}_1 = \left\{ K_1(\cdot) : \mathbb{R} \rightarrow [-1, 1] : K_1(0) = 1, K_1(x) = K_1(-x), \forall x \in \mathbb{R}, \int_{-\infty}^{\infty} |K_1(x)| dx < \infty \right. \quad (\text{N.1})$$

$$\left. \int_{-\infty}^{\infty} K_1^2(x) dx < \infty, K_1(\cdot) \text{ is continuous at 0 and at all but a finite number of points} \right\}.$$

The class $\mathbf{K}_1$ was considered by [Andrews \(1991\)](#) and [Casini \(2023\)](#). Examples of kernels in $\mathbf{K}_1$ include the Truncated, Bartlett, Parzen, Quadratic Spectral (QS) and Tukey-Hanning kernels.

The DK-HAC estimators are defined as

$$\begin{aligned} \hat{\Sigma}_{N_1}(\mathbf{S}_T) &= \frac{T}{T - q - p} \sum_{k=-T+1}^{T-1} K_1(b_{1,T}k) \hat{\Gamma}_{\text{DK}}(k, \mathbf{S}_T), \quad \text{with} \\ \hat{\Gamma}_{\text{DK}}(k, \mathbf{S}_T) &= \frac{n_T}{T - n_T} \sum_{r=0}^{\lfloor (T-n_T)/n_T \rfloor} \hat{c}_{bb}(rn_T/T, k, \mathbf{S}_T), \end{aligned}$$

where $n_T \rightarrow \infty$ satisfies the conditions given below, and

$$\hat{c}_{bb}(rn_T/T, k, \mathbf{S}_T) = \begin{cases} (Tb_{2,T})^{-1} \sum_{t=k+1}^T K_2\left(\frac{((r+1)n_T - (t-k/2))/T}{b_{2,T}}\right) \hat{V}_{b,t}(\mathbf{S}_T) \hat{V}'_{b,t-k}(\mathbf{S}_T), & k \geq 0 \\ (Tb_{2,T})^{-1} \sum_{t=-k+1}^T K_2\left(\frac{((r+1)n_T - (t+k/2))/T}{b_{2,T}}\right) \hat{V}_{b,t+k}(\mathbf{S}_T) \hat{V}'_{b,t}(\mathbf{S}_T), & k < 0 \end{cases}, \quad (\text{N.2})$$

with K_2 being a kernel and $b_{2,T}$ is a bandwidth sequence. The DK-HAC estimators $\hat{\Sigma}_{N_2}^*(\mathbf{S}_T)$ and $\hat{\Sigma}_{N_1 N_2}(\mathbf{S}_T)$ are defined analogously to $\hat{\Sigma}_{N_1}(\mathbf{S}_T)$ after replacing $\hat{c}_{bb}(rn_T/T, k, \mathbf{S}_T)$ with analogous quantities $\hat{c}_{aa}(rn_T/T, k, \mathbf{S}_T)$, and $\hat{c}_{ab}(rn_T/T, k, \mathbf{S}_T)$, respectively. [Casini \(2023\)](#) considered the following class of kernels

$$\begin{aligned} \mathbf{K}_2 &= \left\{ K_2(\cdot) : \mathbb{R} \rightarrow [0, \infty] : K_2(x) = K_2(1-x), \int K_2(x) dx = 1, \right. \quad (\text{N.3}) \\ &\quad \left. K_2(x) = 0 \text{ for } x \notin [0, 1], \int_{-\infty}^{\infty} |K_2(x)| dx < \infty, K_2(\cdot) \text{ is continuous} \right\}. \end{aligned}$$

The QS kernel was shown to be optimal in the class $\mathbf{K}_1$ for HAC estimators under the mean-squared error (MSE) criterion by [Andrews \(1991\)](#) and for DK-HAC estimators under a sequential and global MSE criterion by [Casini \(2023\)](#) and [Belotti, Casini, Catania, Grassi, and Perron \(2023\)](#).

The QS kernel is defined as

$$K_1^{\text{QS}}(x) = \frac{25}{12\pi^2 x^2} \left(\frac{\sin(6\pi x/5)}{6\pi x/5} - \cos(6\pi x/5) \right).$$

Casini (2023) showed that the optimal kernel in the class $\mathbf{K}_2$ is a quadratic-type kernel, $K_2^{\text{opt}}(x) = 6x(1-x)$, $0 \leq x \leq 1$.

For both HAC and DK-HAC estimators, define

$$\begin{aligned} \hat{\Sigma}_{v\bar{Z}}(\mathbf{S}_T) &= \begin{bmatrix} \hat{\Sigma}_{N_1}(\mathbf{S}_T) & \hat{\Sigma}_{N_1 N_2}(\mathbf{S}_T)' \\ \hat{\Sigma}_{N_1 N_2}(\mathbf{S}_T) & \hat{\Sigma}_{N_2}^*(\mathbf{S}_T) \end{bmatrix} \\ \hat{\Sigma}_{N_2}(\mathbf{S}_T) &= \hat{\Sigma}_{N_2}^*(\mathbf{S}_T) - \hat{\Sigma}_{N_1 N_2}(\mathbf{S}_T) \hat{\Sigma}_{N_1}^{-1}(\mathbf{S}_T) \hat{\Sigma}_{N_1 N_2}'(\mathbf{S}_T), \end{aligned}$$

where throughout this section, we adopt the convention $\hat{\Sigma}_{N_1 N_2}(\mathbf{S}_T) = \hat{\Sigma}_{N_2 N_1}'(\mathbf{S}_T)$. We now provide sufficient conditions under which the HAC and DK-HAC estimators are uniformly consistent for $\Sigma_{v\bar{Z}}(\mathbf{S})$, and therefore $\Sigma_{N_2}(\mathbf{S})$, over $\mathbf{S} \in \mathcal{S}$. Let $V_t(\mathbf{S}_T) = v_t \otimes w_t(\mathbf{S}_T)$ where $w_t(\mathbf{S}_T)$ is the t th row of $w = [C_T Z : X]$ written as a column vector.

Assumption N.E.12. *((i) $\{V_t(\mathbf{S}_T)\}$ satisfies*

$$\sum_{j=-\infty}^{\infty} \sup_{t \geq 1} \sup_{\mathbf{S}_T \in \mathcal{S}} \|\mathbb{E}(V_t(\mathbf{S}_T) V_{t-j}'(\mathbf{S}_T))\| < \infty,$$

and for all conformable $a_1, a_2, a_3, a_4 \in \mathbb{Z}_+$, $\sum_{n=1}^{\infty} \sum_{j=1}^{\infty} \sum_{m=1}^{\infty} \sup_{t \geq 1} |\kappa_{V,t}^{(a_1, a_2, a_3, a_4)}(n, j, m, \mathbf{S}_T)| < \infty$ where $\kappa_{V,t}^{(a_1, a_2, a_3, a_4)}(n, j, m, \mathbf{S}_T)$ is the time- t fourth-order cumulant of

$$(V_t^{(a_1)}(\mathbf{S}_T), V_{t+n}^{(a_2)}(\mathbf{S}_T), V_{t+j}^{(a_3)}(\mathbf{S}_T), V_{t+m}^{(a_4)}(\mathbf{S}_T)).$$

$$(ii) \sup_{t \geq 1} \sup_{\mathbf{S}_T \in \mathcal{S}} \mathbb{E} \|V_t(\mathbf{S}_T)\|^2 < \infty.$$

Assumption N.E.13. $b_{1,T} \rightarrow 0$ with $Tb_{1,T}^2 \rightarrow \infty$ and $K_1(\cdot) \in \mathbf{K}_1$.

Assumption N.E.12 imposes conditions on the temporal dependence of the instruments and errors. It is a standard assumption in the literature, see Andrews (1991) and Casini (2023). Note that Assumption N.E.12 allows for nonstationary random variables (i.e., time-varying moments). The condition on the bandwidth in Assumption N.E.13 is from Andrews (1991). Assumptions N.E.12-N.E.13 are sufficient for the consistency of HAC estimators.

Assumption N.E.14. $b_{1,T}, b_{2,T} \rightarrow 0$, $n_T \rightarrow \infty$, $n_T/T \rightarrow 0$, $(Tb_{1,T}b_{2,T})^{-1} \rightarrow 0$, $\sqrt{T}b_{1,T} \rightarrow \infty$, $K_1(\cdot) \in \mathbf{K}_1$ and $K_2(\cdot) \in \mathbf{K}_2$.

The conditions on the bandwidths $b_{1,T}, b_{2,T}$ and on n_T are from Casini (2023). Assumptions N.E.12 and N.E.14 are sufficient for the consistency of the DK-HAC estimators.

Lemma N.E.2. *Let Assumptions 6.1-6.3 hold. We have: $\hat{\Sigma}_{v\bar{Z}}(\mathbf{S}_T) \xrightarrow{\mathbb{P}} \Sigma_{v\bar{Z}}(\mathbf{S})$ for $\mathbf{S} = T^{-1}\mathbf{S}_T$ uniformly in $\mathbf{S}_T \in \mathcal{S}$ under Assumptions N.E.12-N.E.13 for the HAC estimator and under Assumptions N.E.12 and N.E.14 for the DK-HAC estimator.*

The proof of Lemma N.E.2 is omitted. For the HAC estimator the proof follows from the discussion in Section 8 in Andrews (1991) who extended the consistency result in Theorem 1 in Andrews (1991) to nonstationary random variables. See also Casini (2022) who provided a solution to some issues in Section 8 in Andrews (1991). The proof for the DK-HAC estimator follows from Theorem 4.2 in Casini and Perron (2024). $\square$

N.E.4 Strong IV and Local Alternative (SIV-LA) Asymptotics for Identification-Robust Tests

We analyze the strong IV asymptotic properties of the tests considered above for local alternatives. Under strong IV asymptotics, $\theta \neq 0$ is fixed. For local alternatives, β is local to β_0 .

Assumption N.E.15. (SIV-LA) (i) $\beta = \beta_0 + r/T^{1/2}$ for some constant $r \in \mathbb{R}$; (ii) θ is a fixed non-zero q -vector; (iii) There exists an estimator $\hat{\mathbf{S}}_T$ such that $T^{-1}\hat{\mathbf{S}}_T \xrightarrow{\mathbb{P}} \mathbf{S}_0$.

Under strong IVs, part (iii) is satisfied by, for example, $\hat{\mathbf{S}}_T$ in (6.5), $\hat{\mathbf{S}}_{T,OLS}$ and $\hat{\mathbf{S}}_{T,FGLS}$ where the optimization is over $\Xi_{\epsilon, \pi_0, m_0, T}$. Under SIV-LA asymptotics, $N_{1,T}(\mathbf{S}_T)$ and $N_{2,T}(\mathbf{S}_T)$ depend asymptotically on $\zeta_{N_1}(\mathbf{S}) \sim \mathcal{N}(\alpha_{N_1}(\mathbf{S}), I_q)$, $\alpha_{N_1}(\mathbf{S}) = \Sigma_{N_1}^{-1/2}(\mathbf{S})\Sigma_{\bar{Z}}(\mathbf{S}, \mathbf{S}_0)\theta r$, and $\alpha_{N_2}(\mathbf{S}) = \Sigma_{N_2}^{-1/2}(\mathbf{S})\Sigma_{\bar{Z}}(\mathbf{S}, \mathbf{S}_0)\theta(a'_{0,\beta}\Sigma_v^{-1}a_{0,\beta})$.

We now determine the asymptotic distributions of the LR, LM and AR test statistics.

Theorem N.E.1. *Let Assumptions 6.1-6.4 and N.E.15 hold. We have: (i) $AR_T(\hat{\mathbf{S}}_T) \xrightarrow{d} \zeta_{N_1}(\mathbf{S}_0)' \zeta_{N_1}(\mathbf{S}_0) \sim \chi_q^2(\alpha_{N_1}(\mathbf{S}_0)' \alpha_{N_1}(\mathbf{S}_0))$; (ii) $LM_T(\hat{\mathbf{S}}_T) \xrightarrow{d} (\alpha_{N_2}(\mathbf{S}_0)' \zeta_{N_1}(\mathbf{S}_0))^2 / \|\alpha_{N_2}(\mathbf{S}_0)\|^2 \sim \chi_1^2((\alpha_{N_2}(\mathbf{S}_0)' \alpha_{N_1}(\mathbf{S}_0))^2 / \|\alpha_{N_2}(\mathbf{S}_0)\|^2)$; (iii) $LR_T(\hat{\mathbf{S}}_T) = LM_T(\mathbf{S}_0) + o_{\mathbb{P}}(1) \xrightarrow{d} (\alpha_{N_2}(\mathbf{S}_0)' \zeta_{N_1}(\mathbf{S}_0))^2 / \|\alpha_{N_2}(\mathbf{S}_0)\|^2 \sim \chi_1^2((\alpha_{N_2}(\mathbf{S}_0)' \alpha_{N_1}(\mathbf{S}_0))^2 / \|\alpha_{N_2}(\mathbf{S}_0)\|^2)$.*

Since $T^{-1}\hat{\mathbf{S}}_T \xrightarrow{\mathbb{P}} \mathbf{S}_0$ under strong IVs, the test statistics above are evaluated at $\mathbf{S}_0$ asymptotically. Akin to the case of known partition, the LM and LR test statistics are asymptotically

equivalent under SIV-LA asymptotics for any value of q . When $q = 1$, $AR_T(\mathbf{S}_0)$, $LM_T(\mathbf{S}_0)$ and $LR_T(\mathbf{S}_0)$ are the same and so the three tests are asymptotically equivalent.

Under SIV-LA asymptotics and i.i.d. normal errors with unknown covariance matrix Σ_v and known $\mathbf{S}_0$, the model for y is a regular parametric model in the sense of standard likelihood theory. Hence, $LR_T(\mathbf{S}_0)$ and $LM_T(\mathbf{S}_0)$ are asymptotically efficient. This means they have standard large-sample optimality properties such as uniformly maximizing asymptotic power among asymptotically unbiased tests. Adapting the proof of Theorem 7 in [Andrews, Moreira, and Stock \(2006\)](#) while using $T^{-1}\hat{\mathbf{S}}_T \xrightarrow{\mathbb{P}} \mathbf{S}_0$ it follows that $LM_T(\hat{\mathbf{S}}_T)$ and $LR_T(\hat{\mathbf{S}}_T)$ are asymptotically efficient under SIV-LA asymptotics and i.i.d. normal errors.

N.E.5 Strong IV and Fixed Alternative (SIV-FA) Asymptotics for Identification-Robust Tests

We now consider strong IV-fixed alternative (SIV-FA) asymptotics to determine the consistency, or lack thereof, of the tests.

Assumption N.E.16. *SIV-FA. (i) $\beta \neq \beta_0$ is fixed; (ii) θ is a fixed non-zero q -vector; (iii) There exists an estimator $\hat{\mathbf{S}}_T$ such that $T^{-1}\hat{\mathbf{S}}_T \xrightarrow{\mathbb{P}} \mathbf{S}_0$.*

Let $\Sigma_{\bar{Z}}(\mathbf{S}_0) = \Sigma_{\bar{Z}}(\mathbf{S}_0, \mathbf{S}_0)$. Define $\varphi_{N_1}(\mathbf{S}_0) = \Sigma_{N_1}^{-1/2}(\mathbf{S}_0) \Sigma_{\bar{Z}}(\mathbf{S}_0) \theta (\beta - \beta_0)$,

$$\varphi_{N_2}(\mathbf{S}_0) = \Sigma_{N_2}^{-1/2}(\mathbf{S}_0) \left(\Sigma_{\bar{Z}}(\mathbf{S}_0) \theta a'_{\beta} \Sigma_v^{-1} a_{0,\beta} - \Sigma_{N_1 N_2}(\mathbf{S}_0, \mathbf{S}_0) \Sigma_{N_1}^{-1/2}(\mathbf{S}_0) \varphi_{N_1}(\mathbf{S}_0) \right), \quad \varsigma_q \sim \mathcal{N}(0, I_q). \quad (\text{N.4})$$

We now determine the asymptotic behavior of the test statistics under SIV-FA asymptotics.

Theorem N.E.2. *Let Assumptions 6.1-6.4 and N.E.16. We have: (i) $AR_T(\hat{\mathbf{S}}_T)/T \xrightarrow{\mathbb{P}} \varphi_{N_1}(\mathbf{S}_0)' \varphi_{N_1}(\mathbf{S}_0) > 0$, (ii) $LM_T(\hat{\mathbf{S}}_T)/T \xrightarrow{\mathbb{P}} (\varphi_{N_1}(\mathbf{S}_0)' \varphi_{N_2}(\mathbf{S}_0))^2 / \varphi_{N_2}(\mathbf{S}_0)' \varphi_{N_2}(\mathbf{S}_0) > 0$ provided $\varphi_{N_1}(\mathbf{S}_0)' \varphi_{N_2}(\mathbf{S}_0) \neq 0$; (iii)*

$$2LR_T(\hat{\mathbf{S}}_T)/T \xrightarrow{\mathbb{P}} \varphi_{N_1}(\mathbf{S}_0)' \varphi_{N_1}(\mathbf{S}_0) - \varphi_{N_2}(\mathbf{S}_0)' \varphi_{N_2}(\mathbf{S}_0) + \sqrt{(\varphi_{N_1}(\mathbf{S}_0)' \varphi_{N_1}(\mathbf{S}_0) - \varphi_{N_2}(\mathbf{S}_0)' \varphi_{N_2}(\mathbf{S}_0))^2 + 4(\varphi_{N_1}(\mathbf{S}_0)' \varphi_{N_2}(\mathbf{S}_0))^2}.$$

The theorem shows that the test $AR_T(\hat{\mathbf{S}}_T)$ is consistent against any alternative $\beta \neq \beta_0$, $LM_T(\hat{\mathbf{S}}_T)$ is consistent against any alternative $\beta \neq \beta_0$ such that $\varphi_{N_2}(\mathbf{S}_0) \neq 0$, and $LR_T(\hat{\mathbf{S}}_T)$ is consistent against any alternative for which the limit value given in the theorem is non-zero.

N.E.6 Proofs of Section N.E.4 and N.E.5

N.E.6.1 Proof of Theorem N.E.1

We begin with the following lemma.

Lemma N.E.3. *Let Assumptions 6.1-6.4 and N.E.15 hold. We have: (i) $(N_{1,T}(\hat{\mathbf{S}}_T), T^{-1/2}N_{2,T}(\hat{\mathbf{S}}_T)) \xrightarrow{d} (\zeta_{N_1}(\mathbf{S}_0), \alpha_{N_2}(\mathbf{S}_0))$ and (ii) $(M_{1,T}(\hat{\mathbf{S}}_T), T^{-1/2}M_{1,2,T}(\hat{\mathbf{S}}_T), T^{-1}M_{2,T}(\hat{\mathbf{S}}_T)) \xrightarrow{d} (\zeta_{N_1}(\mathbf{S}_0)' \zeta_{N_1}(\mathbf{S}_0), \alpha_{N_2}(\mathbf{S}_0)' \zeta_{N_1}(\mathbf{S}_0), \alpha_{N_2}(\mathbf{S}_0)' \alpha_{N_2}(\mathbf{S}_0))$.*

Proof of Lemma N.E.3. Part (i) for $N_{1,T}(\hat{\mathbf{S}}_T)$ follows from

$$\begin{aligned}
N_{1,T}(\hat{\mathbf{S}}_T) &= \hat{\Sigma}_{N_1}^{-1/2}(\hat{\mathbf{S}}_T) T^{-1/2} \bar{Z}(\hat{C}_T)' y b_0 \\
&= \hat{\Sigma}_{N_1}^{-1/2}(\hat{\mathbf{S}}_T) T^{-1/2} \bar{Z}(\hat{C}_T)' (\bar{Z}(C_{0,T}) \theta a'_{\beta} + v) b_0 \\
&= \hat{\Sigma}_{N_1}^{-1/2}(\hat{\mathbf{S}}_T) T^{-1} \bar{Z}(\hat{C}_T)' \bar{Z}(C_{0,T}) \theta r + \hat{\Sigma}_{N_1}^{-1/2}(\hat{\mathbf{S}}_T) T^{-1/2} \bar{Z}(\hat{C}_T)' v b_0 \\
&= \hat{\Sigma}_{N_1}^{-1/2}(\mathbf{S}_0) T^{-1} \bar{Z}(C_{0,T})' \bar{Z}(C_{0,T}) \theta r + \hat{\Sigma}_{N_1}^{-1/2}(\mathbf{S}_0) T^{-1/2} \bar{Z}(C_{0,T})' v b_0 + o_{\mathbb{P}}(1) \\
&= \Sigma_{N_1}^{-1/2}(\mathbf{S}_0) \Sigma_{\bar{Z}}(\mathbf{S}_0) \theta r + \Sigma_{N_1}^{-1/2}(\mathbf{S}_0) T^{-1/2} \bar{Z}(C_{0,T})' v b_0 + o_{\mathbb{P}}(1) \\
&\Rightarrow \Sigma_{N_1}^{-1/2}(\mathbf{S}_0) \Sigma_{\bar{Z}}(\mathbf{S}_0) \theta r + \Sigma_{N_1}^{-1/2}(\mathbf{S}_0) [I_q : -Q_{12}(\mathbf{S}_0) Q_{22}^{-1}] (b'_0 \otimes I_{q+p}) \mathcal{G}(\mathbf{S}_0) \\
&\sim \zeta_{N_1}(\mathbf{S}_0),
\end{aligned}$$

where the third and fourth equalities hold by Assumption N.E.15, the final equality holds by Assumptions 6.1 and 6.4 and the convergence holds by Assumption 6.3. Under Assumptions 6.2 and N.E.15(iii), $\hat{\Sigma}_v(\hat{\mathbf{S}}_T) \xrightarrow{\mathbb{P}} \Sigma_v$ by the same arguments as when $\mathbf{S}_0$ is known. Then, part (i) for $N_{2,T}(\hat{\mathbf{S}}_T)$ holds using

$$\begin{aligned}
T^{-1/2}N_{2,T}(\hat{\mathbf{S}}_T) &= \hat{\Sigma}_{N_2}^{-1/2}(\hat{\mathbf{S}}_T) \left(T^{-1} \bar{Z}(\hat{C}_T)' y \hat{\Sigma}_v^{-1}(\hat{\mathbf{S}}_T) a_{0,\beta} - T^{-1/2} \hat{\Sigma}_{N_1 N_2}(\hat{\mathbf{S}}_T) \hat{\Sigma}_{N_1}^{-1/2}(\hat{\mathbf{S}}_T) N_{1,T}(\hat{\mathbf{S}}_T) \right) \\
&\quad \quad \quad (N.5) \\
&= \hat{\Sigma}_{N_2}^{-1/2}(\hat{\mathbf{S}}_T) \left(T^{-1} \bar{Z}(\hat{C}_T)' y \Sigma_v^{-1} a_{0,\beta} - T^{-1/2} \hat{\Sigma}_{N_1 N_2}(\hat{\mathbf{S}}_T) \hat{\Sigma}_{N_1}^{-1/2}(\hat{\mathbf{S}}_T) N_{1,T}(\hat{\mathbf{S}}_T) \right) + o_{\mathbb{P}}(1) \\
&= \hat{\Sigma}_{N_2}^{-1/2}(\mathbf{S}_0) \left(T^{-1} \bar{Z}(C_{0,T})' y \Sigma_v^{-1} a_{0,\beta} - T^{-1/2} \hat{\Sigma}_{N_1 N_2}(\mathbf{S}_0) \hat{\Sigma}_{N_1}^{-1/2}(\mathbf{S}_0) N_{1,T}(\mathbf{S}_0) \right) + o_{\mathbb{P}}(1) \\
&= \Sigma_{N_2}^{-1/2}(\mathbf{S}_0) \left(T^{-1} \bar{Z}(C_{0,T})' y \Sigma_v^{-1} a_{0,\beta} - T^{-1/2} \Sigma_{N_1 N_2}(\mathbf{S}_0) \Sigma_{N_1}^{-1/2}(\mathbf{S}_0) N_{1,T}(\mathbf{S}_0) \right) + o_{\mathbb{P}}(1) \\
&= \Sigma_{N_2}^{-1/2}(\mathbf{S}_0) \left(T^{-1} \bar{Z}(C_{0,T})' (\bar{Z}(C_{0,T}) \theta a'_{\beta} + v) \Sigma_v^{-1} a_{0,\beta} \right) + o_{\mathbb{P}}(1) \\
&= \Sigma_{N_2}^{-1/2}(\mathbf{S}_0) \Sigma_{\bar{Z}}(\mathbf{S}_0) \theta a'_{\beta} \Sigma_v^{-1} a_{0,\beta} + o_{\mathbb{P}}(1),
\end{aligned}$$

where the third equality follows from Assumption N.E.15(iii), the fourth by Assumption 6.4, the

fifth by part (i) for $N_{1,T}(\hat{\mathbf{S}}_T)$ and the final equality holds by Assumption 6.1. Part (ii) holds by part (i) and the continuous mapping theorem. $\square$

Proof of Theorem N.E.1. Parts (i) and (ii) of the theorem follow immediately from Lemma N.E.3(ii). Part (iii) of the theorem is established as follows. Following the argument based on a mean-value expansion of the LR_T statistic in the proof of Theorem 9 in Andrews, Moreira, and Stock (2004) (see eq. (14.50)-(14.53)) with references to Lemma 9-(b) there replaced by references to Lemma N.E.3, we have

$$\begin{aligned} LR_T(\hat{\mathbf{S}}_T) &= \frac{1}{2} \left(2M_{1,T}(\hat{\mathbf{S}}_T) - 2 \left(M_{1,T}(\hat{\mathbf{S}}_T) - M_{1,2,T}(\hat{\mathbf{S}}_T)^2 / M_{2,T}(\hat{\mathbf{S}}_T) \right) \right) + o_{\mathbb{P}}(1) \\ &= M_{1,2,T}^2(\mathbf{S}_0) / M_{2,T}(\mathbf{S}_0) + o_{\mathbb{P}}(1) \\ &= LM_T(\mathbf{S}_0) + o_{\mathbb{P}}(1), \end{aligned} \quad (\text{N.6})$$

where we used Assumption N.E.15(iii). $\square$

N.E.6.2 Proof of Theorem N.E.2

We begin with the following lemma.

Lemma N.E.4. (i) Under Assumptions 6.1-6.4 and N.E.16, (i) $(N_{1,T}(\hat{\mathbf{S}}_T)/T^{1/2}, N_{2,T}(\hat{\mathbf{S}}_T)/T^{1/2}) \xrightarrow{\mathbb{P}} (\varphi_{N_1}(\mathbf{S}_0), \varphi_{N_2}(\mathbf{S}_0))$ and (ii) $(M_{1,T}(\hat{\mathbf{S}}_T)/T, M_{1,2,T}(\hat{\mathbf{S}}_T)/T, M_{2,T}(\hat{\mathbf{S}}_T)/T) \xrightarrow{\mathbb{P}} (\varphi_{N_1}(\mathbf{S}_0)' \varphi_{N_1}(\mathbf{S}_0), \varphi_{N_1}(\mathbf{S}_0)' \varphi_{N_2}(\mathbf{S}_0), \varphi_{N_2}(\mathbf{S}_0)' \varphi_{N_2}(\mathbf{S}_0))$.

Proof of Lemma N.E.4. Part (i) of the lemma is established as follows:

$$\begin{aligned} T^{-1} \bar{Z}(\hat{C}_T)' y b_0 &= T^{-1} \bar{Z}(\hat{C}_T)' \left(\bar{Z}(C_{0,T}) \theta a'_{\beta} + X \eta + v \right) b_0 \\ &= T^{-1} \bar{Z}(\hat{C}_T)' \bar{Z}(C_{0,T}) \theta a'_{\beta} b_0 + T^{-1} \bar{Z}(\hat{C}_T)' v b_0 \xrightarrow{\mathbb{P}} \Sigma_{\bar{Z}}(\mathbf{S}_0) \theta a'_{\beta} b_0, \end{aligned}$$

using Assumptions 6.1, 6.3, N.E.16(iii) and $\bar{Z}(\hat{C}_T)' X = 0$. Hence, by Assumptions 6.1, 6.4 and

N.E.16(iii), we have

$$\begin{aligned}
N_{1,T}(\hat{\mathbf{S}}_T)/T^{1/2} &= \hat{\Sigma}_{N_1}^{-1/2}(\hat{\mathbf{S}}_T) T^{-1} \bar{Z}(\hat{C}_T)' y b_0 \\
&= \hat{\Sigma}_{N_1}^{-1/2}(\hat{\mathbf{S}}_T) T^{-1} \bar{Z}(\hat{C}_T)' (\bar{Z}(C_{0,T}) \theta a'_\beta + v) b_0 \\
&= \hat{\Sigma}_{N_1}^{-1/2}(\hat{\mathbf{S}}_T) T^{-1} \bar{Z}(\hat{C}_T)' \bar{Z}(C_{0,T}) \theta a'_\beta b_0 + \hat{\Sigma}_{N_1}^{-1/2}(\hat{\mathbf{S}}_T) T^{-1} \bar{Z}(\hat{C}_T)' v b_0 \\
&= \hat{\Sigma}_{N_1}^{-1/2}(\mathbf{S}_0) T^{-1} \bar{Z}(C_{0,T})' \bar{Z}(C_{0,T}) \theta a'_\beta b_0 + \hat{\Sigma}_{N_1}^{-1/2}(\mathbf{S}_0) T^{-1} \bar{Z}(C_{0,T})' v b_0 + o_{\mathbb{P}}(1) \\
&= \Sigma_{N_1}^{-1/2}(\mathbf{S}_0) \Sigma_{\bar{Z}}(\mathbf{S}_0) \theta a'_\beta b_0 + \Sigma_{N_1}^{-1/2}(\mathbf{S}_0) T^{-1} \bar{Z}(C_{0,T})' v b_0 + o_{\mathbb{P}}(1) \\
&= \Sigma_{N_1}^{-1/2}(\mathbf{S}_0) \Sigma_{\bar{Z}}(\mathbf{S}_0) \theta (\beta - \beta_0) + o_{\mathbb{P}}(1).
\end{aligned} \tag{N.7}$$

Similarly,

$$\begin{aligned}
T^{-1/2} N_{2,T}(\hat{\mathbf{S}}_T) &= \hat{\Sigma}_{N_2}^{-1/2}(\hat{\mathbf{S}}_T) \left(T^{-1} \bar{Z}(\hat{C}_T)' y \hat{\Sigma}_v^{-1}(\hat{\mathbf{S}}_T) a_{0,\beta} - T^{-1/2} \hat{\Sigma}_{N_1 N_2}(\hat{\mathbf{S}}_T) \hat{\Sigma}_{N_1}^{-1/2}(\hat{\mathbf{S}}_T) N_{1,T}(\hat{\mathbf{S}}_T) \right) \\
&= \hat{\Sigma}_{N_2}^{-1/2}(\hat{\mathbf{S}}_T) \left(T^{-1} \bar{Z}(\hat{C}_T)' y \Sigma_v^{-1} a_{0,\beta} - T^{-1/2} \hat{\Sigma}_{N_1 N_2}(\hat{\mathbf{S}}_T) \hat{\Sigma}_{N_1}^{-1/2}(\hat{\mathbf{S}}_T) N_{1,T}(\hat{\mathbf{S}}_T) \right) + o_{\mathbb{P}}(1) \\
&= \hat{\Sigma}_{N_2}^{-1/2}(\mathbf{S}_0) \left(T^{-1} \bar{Z}(C_{0,T})' y \Sigma_v^{-1} a_{0,\beta} - T^{-1/2} \hat{\Sigma}_{N_1 N_2}(\mathbf{S}_0) \hat{\Sigma}_{N_1}^{-1/2}(\mathbf{S}_0) N_{1,T}(\mathbf{S}_0) \right) + o_{\mathbb{P}}(1) \\
&= \Sigma_{N_2}^{-1/2}(\mathbf{S}_0) \left(T^{-1} \bar{Z}(C_{0,T})' y \Sigma_v^{-1} a_{0,\beta} - T^{-1/2} \Sigma_{N_1 N_2}(\mathbf{S}_0) \Sigma_{N_1}^{-1/2}(\mathbf{S}_0) N_{1,T}(\mathbf{S}_0) \right) + o_{\mathbb{P}}(1) \\
&= \Sigma_{N_2}^{-1/2}(\mathbf{S}_0) \left(T^{-1} \bar{Z}(C_{0,T})' (\bar{Z}(C_{0,T}) \theta a'_\beta + v) \Sigma_v^{-1} a_{0,\beta} \right) \\
&\quad - \Sigma_{N_2}^{-1/2}(\mathbf{S}_0) \Sigma_{N_1 N_2}(\mathbf{S}_0) \Sigma_{N_1}^{-1/2}(\mathbf{S}_0) \varphi_{N_1}(\mathbf{S}_0) + o_{\mathbb{P}}(1) \\
&= \Sigma_{N_2}^{-1/2}(\mathbf{S}_0) \Sigma_{\bar{Z}}(\mathbf{S}_0) \theta a'_\beta \Sigma_v^{-1} a_{0,\beta} - \Sigma_{N_2}^{-1/2}(\mathbf{S}_0) \Sigma_{N_1 N_2}(\mathbf{S}_0) \Sigma_{N_1}^{-1/2}(\mathbf{S}_0) \varphi_{N_1}(\mathbf{S}_0) + o_{\mathbb{P}}(1) \\
&= \varphi_{N_2}(\mathbf{S}_0) + o_{\mathbb{P}}(1).
\end{aligned} \tag{N.8}$$

Part (ii) of the lemma follows from part (i) and Slutsky's Theorem. $\square$

Proof of Theorem N.E.2. Parts (i)-(iii) of the theorem hold by Lemma N.E.4 and simple calculations. In the case of $LM_T(\hat{\mathbf{S}}_T)$, the convergence only holds if β is such that $\varphi_{N_2}(\mathbf{S}_0) \neq 0$ because $\varphi_{N_2}(\mathbf{S}_0)$ appears in the denominator. $\square$

N.F Additional Monte Carlo Simulations

N.F.1 Finite-Sample Properties of $\hat{S}_{T,OLS}$, $\hat{S}_{T,FGLS}$, $\hat{\beta}(\hat{S}_{T,OLS})$ and $\hat{\beta}(\hat{S}_{T,FGLS})$

We study the finite-sample bias and mean-squared error (MSE) of the proposed estimators of $S_{0,T}$ and β . For the latter, we compare them with $\hat{\beta}_{FS}$, the full sample IV estimator of β . We consider

the same DGP as in (N.1)-(N.2) where $X_t \sim \text{i.i.d. } \mathcal{N}(1, 1)$. We set $\beta = 1$.¹ The number of simulations is 5,000.

Table 2 reports the MSE of $\hat{\mathbf{S}}_{T,OLS}$ and $\hat{\mathbf{S}}_{T,FGLS}$ for the case $\theta_1 = \theta_3 = dT^{-1/2}$ and $\theta_2 = 0$, where $\hat{\mathbf{S}}_{T,FGLS}$ is constructed using $\hat{\Omega}_\varepsilon(\mathbf{S}_T) = \hat{\Sigma}_\varepsilon \otimes I_T$, which is misspecified in the presence of serial correlation. $\hat{\mathbf{S}}_{T,OLS}$ performs better when d is small, whereas $\hat{\mathbf{S}}_{T,FGLS}$ performs better when d is large. These results continue to hold when the errors are serially correlated.

Table 2: MSE of $\hat{S}_{T,OLS}$ and $\hat{S}_{T,FGLS}$

$\rho = 0.25, \pi_0 = 0.6$	$\rho_e = \rho_u = 0$			$\rho_e = \rho_u = 0.50$		
	$d = 4$	$d = 16$	$d = 32$	$d = 4$	$d = 16$	$d = 32$
MSE($\hat{S}_{T,OLS}$)	31.59	2.31	0.78	37.30	4.03	1.11
MSE($\hat{S}_{T,FGLS}$)	38.11	2.38	0.72	40.69	4.69	1.07
$\rho = 0.75, \pi_0 = 0.6$	$\rho_e = \rho_u = 0$			$\rho_e = \rho_u = 0.50$		
	$d = 4$	$d = 16$	$d = 32$	$d = 4$	$d = 16$	$d = 32$
MSE($\hat{S}_{T,OLS}$)	31.59	2.31	0.78	37.30	4.03	1.11
MSE($\hat{S}_{T,FGLS}$)	32.45	1.16	0.41	38.94	2.15	0.55
$\rho = 0.25, \pi_0 = 0.8$	$\rho_e = \rho_u = 0$			$\rho_e = \rho_u = 0.50$		
	$d = 4$	$d = 16$	$d = 32$	$d = 4$	$d = 16$	$d = 32$
MSE($\hat{S}_{T,OLS}$)	22.58	2.25	1.00	26.61	3.93	1.26
MSE($\hat{S}_{T,FGLS}$)	27.59	2.45	0.97	29.40	6.16	1.31
$\rho = 0.75, \pi_0 = 0.8$	$\rho_e = \rho_u = 0$			$\rho_e = \rho_u = 0.50$		
	$d = 4$	$d = 16$	$d = 32$	$d = 4$	$d = 16$	$d = 32$
MSE($\hat{S}_{T,OLS}$)	22.58	2.25	1.00	26.61	3.93	1.26
MSE($\hat{S}_{T,FGLS}$)	24.18	1.39	0.75	28.29	2.42	0.87

To facilitate readability each value is multiplied by 10^2 .

Next, we examine the bias and MSE of $\hat{\beta}(\hat{\mathbf{S}}_{T,OLS})$, $\hat{\beta}(\hat{\mathbf{S}}_{T,FGLS})$ and $\hat{\beta}_{FS}$ for the case $\theta_1 = \theta_3 = dT^{-1/2}$ and $\theta_2 = 0$, where $\hat{\mathbf{S}}_{T,FGLS}$ is again constructed using $\hat{\Omega}_\varepsilon(\mathbf{S}_T) = \hat{\Sigma}_\varepsilon \otimes I_T$. Table 3 reports the results. When the instrument is weak throughout the sample ($d = 4$) no estimator uniformly dominates in terms of bias or MSE, and the results vary considerably across settings. In contrast, when the instrument is strong in parts of the sample ($d = 16, 32$), $\hat{\beta}(\hat{\mathbf{S}}_{T,OLS})$ and $\hat{\beta}(\hat{\mathbf{S}}_{T,FGLS})$ exhibit lower bias and MSE than $\hat{\beta}_{FS}$, with $\hat{\beta}(\hat{\mathbf{S}}_{T,OLS})$ performing slightly better than $\hat{\beta}(\hat{\mathbf{S}}_{T,FGLS})$. The reduction in bias and MSE can be substantial—often at least 50% in many configurations. The results for the cases $\theta_2 = -0.5dT^{-1/2}$ and $\theta_2 = -0.5$ are similar and omitted.

We now evaluate the performance of the estimators of β in a setting where the instrument is strong in the subsample $S_{0,T}$ but weak in the remainder of the sample. This allows us to assess whether the full sample estimator—which combines the strongly identified subsample with

¹The results do not change with other values of β .

the weakly identified subsample—outperforms estimators based solely on the strongly identified subsample. We consider the DGP specified in (N.1)-(N.2), setting $\theta_2 = d_2/\sqrt{T}$ with $d_2 \in \{4, 8\}$ and $\theta_1 = \theta_3 = d/\sqrt{T}$ with $d \in \{16, 24, 32\}$.

The results are reported in Table 4. When the instrument is weak in the remaining part of the sample ($d_2 = 4$), $\hat{\beta}(\hat{\mathbf{S}}_{T,OLS})$ and $\hat{\beta}(\hat{\mathbf{S}}_{T,FGLS})$ consistently yield lower MSE for both $\rho = 0.25$ and 0.75 . The bias of the full sample estimator is smaller than that of $\hat{\beta}(\hat{\mathbf{S}}_{T,OLS})$ and $\hat{\beta}(\hat{\mathbf{S}}_{T,FGLS})$ when $d = 16$. For larger values of d , $\hat{\beta}(\hat{\mathbf{S}}_{T,OLS})$ and $\hat{\beta}(\hat{\mathbf{S}}_{T,FGLS})$ generally exhibit lower bias.

When the instrument has intermediate identification strength in the remaining sample ($d_2 = 8$), $\hat{\beta}_{FS}$ delivers lower bias and MSE than $\hat{\beta}(\hat{\mathbf{S}}_{T,OLS})$ and $\hat{\beta}(\hat{\mathbf{S}}_{T,FGLS})$ when $d = 16$. However, for $d = 24$ and 32 , $\hat{\beta}(\hat{\mathbf{S}}_{T,OLS})$ exhibits lower bias and MSE than both $\hat{\beta}_{FS}$ and $\hat{\beta}(\hat{\mathbf{S}}_{T,FGLS})$.

Overall, the results indicate that when the instrument is weak in the remaining part of the sample, $\hat{\beta}(\hat{\mathbf{S}}_{T,OLS})$ and $\hat{\beta}(\hat{\mathbf{S}}_{T,FGLS})$ outperform the full sample estimator that also incorporates the weakly identified subsample, provided that the instrument is sufficiently strong in the strongly identified subsample.

Table 3: Bias and MSE of $\hat{\beta}_{FS}$, $\hat{\beta}(\hat{S}_{T,OLS})$ and $\hat{\beta}(\hat{S}_{T,FGLS})$

$\theta_2 = 0$	$\rho_e = \rho_u = 0$			$\rho_e = \rho_u = 0.50$		
$\rho = 0.25, \pi_0 = 0.6$	$d = 4$	$d = 16$	$d = 32$	$d = 4$	$d = 16$	$d = 32$
Bias($\hat{\beta}_{FS}$)	-12.00	-0.32	-0.16	-46.77	-0.60	-0.02
Bias($\hat{\beta}(\hat{S}_{T,OLS})$)	24.89	0.06	-0.05	19.34	0.10	-0.02
Bias($\hat{\beta}(\hat{S}_{T,FGLS})$)	-24.86	-0.05	-0.05	-40.44	-0.26	-0.04
MSE($\hat{\beta}_{FS}$)	2802.00	1.18	0.28	114150.00	1.66	0.39
MSE($\hat{\beta}(\hat{S}_{T,OLS})$)	44649.00	0.70	0.17	30190.00	0.96	0.23
MSE($\hat{\beta}(\hat{S}_{T,FGLS})$)	45405.00	0.71	0.17	52584.00	0.96	0.23
$\rho = 0.75, \pi_0 = 0.6$	$\rho_e = \rho_u = 0$			$\rho_e = \rho_u = 0.50$		
	$d = 4$	$d = 16$	$d = 32$	$d = 4$	$d = 16$	$d = 32$
Bias($\hat{\beta}_{FS}$)	-18.00	-0.82	-0.22	29.23	-1.13	-0.29
Bias($\hat{\beta}(\hat{S}_{T,OLS})$)	30.88	0.19	-0.07	54.22	0.22	-0.11
Bias($\hat{\beta}(\hat{S}_{T,FGLS})$)	126.40	-0.64	-0.18	9.22	-1.04	-0.29
MSE($\hat{\beta}_{FS}$)	6879.05	1.27	0.29	114152.48	1.68	0.40
MSE($\hat{\beta}(\hat{S}_{T,OLS})$)	12408.54	0.71	0.17	98469.13	0.98	0.23
MSE($\hat{\beta}(\hat{S}_{T,FGLS})$)	397451.09	0.79	0.19	38287.46	1.18	0.25
$\rho = 0.25, \pi_0 = 0.8$	$\rho_e = \rho_u = 0$			$\rho_e = \rho_u = 0.50$		
	$d = 4$	$d = 16$	$d = 32$	$d = 4$	$d = 16$	$d = 32$
Bias($\hat{\beta}_{FS}$)	-16.50	-0.18	-0.11	0.58	-0.40	0.01
Bias($\hat{\beta}(\hat{S}_{T,OLS})$)	0.44	0.07	-0.04	-9.14	-0.06	0.04
Bias($\hat{\beta}(\hat{S}_{T,FGLS})$)	-11.13	-0.30	-0.04	-7.35	-0.23	0.03
MSE($\hat{\beta}_{FS}$)	7589.60	0.62	0.15	1969.60	0.85	0.20
MSE($\hat{\beta}(\hat{S}_{T,OLS})$)	96.21	0.50	0.12	2185.10	0.70	0.16
MSE($\hat{\beta}(\hat{S}_{T,FGLS})$)	519.38	0.51	0.12	7360.50	0.73	0.16
$\rho = 0.75, \pi_0 = 0.8$	$\rho_e = \rho_u = 0$			$\rho_e = \rho_u = 0.50$		
	$d = 4$	$d = 16$	$d = 32$	$d = 4$	$d = 16$	$d = 32$
Bias($\hat{\beta}_{FS}$)	-37.82	-0.41	-0.18	-6.43	-0.58	-0.16
Bias($\hat{\beta}(\hat{S}_{T,OLS})$)	-1.07	0.24	-0.06	81.81	0.16	-0.04
Bias($\hat{\beta}(\hat{S}_{T,FGLS})$)	-3.94	-0.14	-0.09	11.11	-0.48	-0.07
MSE($\hat{\beta}_{FS}$)	37454.00	0.63	0.15	3808.60	0.87	0.21
MSE($\hat{\beta}(\hat{S}_{T,OLS})$)	264.07	0.50	0.12	272130.00	0.69	0.17
MSE($\hat{\beta}(\hat{S}_{T,FGLS})$)	653.99	0.51	0.12	21200.00	0.71	0.17

To facilitate readability each value is multiplied by 10^2 .

Table 4: Bias and MSE of $\hat{\beta}_{FS}$, $\hat{\beta}(\hat{S}_{T,OLS})$ and $\hat{\beta}(\hat{S}_{T,FGLS})$

$d_2 = 4, \pi_0 = 0.6$	$\rho = 0.25$			$\rho = 0.75$		
	$d = 16$	$d = 24$	$d = 32$	$d = 16$	$d = 24$	$d = 32$
Bias($\hat{\beta}_{FS}$)	-0.13	-0.12	-0.14	-0.55	-0.26	-0.25
Bias($\hat{\beta}(\hat{S}_{T,OLS})$)	0.47	0.05	-0.04	1.11	0.15	-0.04
Bias($\hat{\beta}(\hat{S}_{T,FGLS})$)	-0.26	-0.02	-0.04	-0.02	-0.12	-0.10
MSE($\hat{\beta}_{FS}$)	0.87	0.40	0.24	0.89	0.40	0.24
MSE($\hat{\beta}(\hat{S}_{T,OLS})$)	0.68	0.30	0.17	0.70	0.30	0.17
MSE($\hat{\beta}(\hat{S}_{T,FGLS})$)	0.72	0.31	0.17	0.74	0.31	0.17
$d_2 = 8, \pi_0 = 0.6$	$\rho = 0.25$			$\rho = 0.75$		
	$d = 16$	$d = 24$	$d = 32$	$d = 16$	$d = 24$	$d = 32$
Bias($\hat{\beta}_{FS}$)	-0.09	-0.10	-0.13	-0.40	-0.21	-0.22
Bias($\hat{\beta}(\hat{S}_{T,OLS})$)	1.05	0.12	-0.01	2.88	0.39	0.03
Bias($\hat{\beta}(\hat{S}_{T,FGLS})$)	0.78	-0.02	-0.04	1.33	-0.04	-0.20
MSE($\hat{\beta}_{FS}$)	0.65	0.33	0.20	0.67	0.33	0.20
MSE($\hat{\beta}(\hat{S}_{T,OLS})$)	0.71	0.30	0.17	0.76	0.30	0.17
MSE($\hat{\beta}(\hat{S}_{T,FGLS})$)	0.81	0.31	0.17	0.74	0.31	0.17

To facilitate readability each value is multiplied by 10^2 .

N.F.2 Identification-Robust Tests

We consider the performance of the identification-robust tests under serial correlation in the errors. Specifically, we examine DGP (S.C.2)–(S.C.3) and model the error terms as $u_t = \rho_u u_{t-1} + v_{u,t}$ and $e_t = \rho_e e_{t-1} + v_{e,t}$, where $v_{u,t}$ and $v_{e,t}$ are jointly normally distributed with mean zero and covariance matrix Σ_{ue} as in (S.C.4) with $\rho \in \{0, 0.25, 0.5, 0.75\}$ and $\rho_e = \rho_u \in \{0.25, 0.5, 0.75\}$. We set the significance level to 5% and number of Monte Carlo replications to 10,000. Table 5 and Figure 1 report the null rejection frequencies and size-adjusted power of the tests, respectively. Under strong serial dependence ($\rho_e = \rho_u = 0.75$) all tests exhibit rejection rates that exceed the nominal significance level. Specifically, $LM_T(\hat{\mathbf{S}}_T)$ and $CLR_T(\hat{\mathbf{S}}_T)$ are a bit more oversized than LM_T and CLR_T but similar to qLL-S. Under weak serial dependence ($\rho_e = \rho_u = 0.25$), the proposed tests $LM_T(\hat{\mathbf{S}}_T)$ and $CLR_T(\hat{\mathbf{S}}_T)$ are only slightly more oversized than their full sample counterparts, LM_T and CLR_T . Figure 1 shows that the size-adjusted power of the proposed tests is higher than that of the existing tests, similar to the i.i.d. case.

Finally, we consider a model with multiple instruments:

$$Y_t = \beta D_t + u_t, \quad (\text{N.1})$$

where

$$D_t = \begin{cases} \theta_1 (Z_{1,t} + Z_{2,t}) + \tilde{\theta}_1 Z_{3,t} + e_t, & t \leq \lfloor T/4 \rfloor \\ \theta_2 (Z_{1,t} + Z_{2,t} + Z_{3,t}) + e_t, & \lfloor T/4 \rfloor + 1 \leq t \leq \lfloor T/4 \rfloor + \lfloor (1 - \pi_0) T \rfloor \\ \theta_3 (Z_{1,t} + Z_{2,t}) + \tilde{\theta}_3 Z_{3,t} + e_t, & \lfloor T/4 \rfloor + \lfloor (1 - \pi_0) T \rfloor + 1 \leq t \leq T, \end{cases} \quad (\text{N.2})$$

$Z_{i,t} \sim \text{i.i.d. } \mathcal{N}(1, 1)$ for $i = 1, 2, 3$, and u_t and e_t are i.i.d. jointly normal with mean zero and covariance

$$\Sigma_{ue} = \begin{bmatrix} 1 & \rho \\ \rho & 1 \end{bmatrix}, \quad (\text{N.3})$$

with $\rho \in \{0.25, 0.75\}$. We set $\theta_1 = \theta_3 = dT^{-1/2}$ with $d \in \{2, 4, 8\}$, $\theta_2 = 0$ and $\tilde{\theta}_1 = \tilde{\theta}_3 = 16/\sqrt{T}$. We set $\pi_0 \in \{0.6, 0.8\}$ and $T = 200$.

Table 6 reports the null rejection frequencies. The AR_T , $AR_T(\hat{\mathbf{S}}_T)$, Split-S, qLL-S, ave-S and exp-S are severely undersized across all values of d and ρ . LM_T , $LM_T(\hat{\mathbf{S}}_T)$, CLR_T , $CLR_T(\hat{\mathbf{S}}_T)$ lead to quite accurate null rejection rates whereas Split-CLR displays null rejection rates substantially beyond the nominal level. Figure 2 plots the power functions. $LM_T(\hat{\mathbf{S}}_T)$ and $CLR_T(\hat{\mathbf{S}}_T)$ are the most powerful, followed by $AR_T(\hat{\mathbf{S}}_T)$ and then by the full sample counterparts of these tests. qLL-S displays the lowest power across all configurations. The power gains of $LM_T(\hat{\mathbf{S}}_T)$ and $CLR_T(\hat{\mathbf{S}}_T)$ are substantial across all configurations.

Table 5: Finite-Sample Null Rejection Frequencies of Tests

$\rho = 0.50$	$\rho_e = \rho_u = 0.25$			$\rho_e = \rho_u = 0.50$			$\rho_e = \rho_u = 0.75$		
$T = 200, \pi_0 = 0.8$	$d = 10$	$d = 16$	$d = 24$	$d = 10$	$d = 16$	$d = 24$	$d = 10$	$d = 16$	$d = 24$
LM_T	0.073	0.073	0.073	0.089	0.089	0.089	0.158	0.158	0.158
CLR_T	0.079	0.077	0.076	0.096	0.094	0.093	0.179	0.169	0.166
$LM_T(\hat{\mathbf{S}}_T)$	0.086	0.081	0.078	0.118	0.107	0.101	0.204	0.185	0.176
$CLR_T(\hat{\mathbf{S}}_T)$	0.088	0.082	0.078	0.126	0.108	0.101	0.214	0.191	0.179
split – S	0.054	0.055	0.055	0.071	0.074	0.076	0.139	0.148	0.155
split – CLR	0.146	0.149	0.150	0.174	0.186	0.187	0.278	0.291	0.300
qqL – S	0.028	0.028	0.028	0.047	0.052	0.052	0.191	0.198	0.199
ave – S	0.047	0.046	0.047	0.068	0.071	0.070	0.153	0.161	0.171
exp – S	0.019	0.020	0.020	0.034	0.034	0.034	0.109	0.109	0.107
$\rho = 0.50$	$\rho_e = \rho_u = 0.25$			$\rho_e = \rho_u = 0.50$			$\rho_e = \rho_u = 0.75$		
$T = 400, \pi_0 = 0.6$	$d = 10$	$d = 16$	$d = 24$	$d = 10$	$d = 16$	$d = 24$	$d = 10$	$d = 16$	$d = 24$
LM_T	0.061	0.061	0.061	0.074	0.074	0.074	0.124	0.124	0.124
CLR_T	0.070	0.067	0.067	0.088	0.083	0.083	0.153	0.142	0.136
$LM_T(\hat{\mathbf{S}}_T)$	0.076	0.068	0.068	0.109	0.091	0.091	0.174	0.158	0.142
$CLR_T(\hat{\mathbf{S}}_T)$	0.081	0.070	0.069	0.106	0.093	0.093	0.190	0.166	0.145
split – S	0.046	0.044	0.043	0.060	0.059	0.059	0.113	0.113	0.113
split – CLR	0.133	0.134	0.134	0.150	0.154	0.154	0.234	0.240	0.243
qqL – S	0.042	0.043	0.039	0.042	0.063	0.063	0.169	0.177	0.180
ave – S	0.048	0.050	0.048	0.048	0.073	0.068	0.130	0.133	0.137
exp – S	0.0243	0.023	0.024	0.023	0.034	0.037	0.098	0.098	0.095

Table 6: Finite-Sample Null Rejection Frequencies of Tests for the model (N.1)-(N.2)

$T = 200, \pi_0 = 0.6$	$\rho = 0.25$			$\rho = 0.50$			$\rho = 0.75$		
	$d = 2$	$d = 4$	$d = 8$	$d = 2$	$d = 4$	$d = 8$	$d = 2$	$d = 4$	$d = 8$
AR_T	0.000	0.000	0.001	0.008	0.000	0.001	0.001	0.000	0.001
LM_T	0.063	0.062	0.061	0.059	0.064	0.062	0.060	0.069	0.060
CLR_T	0.070	0.070	0.056	0.070	0.073	0.068	0.066	0.076	0.067
$AR_T(\hat{\mathbf{S}}_T)$	0.002	0.001	0.001	0.001	0.000	0.001	0.001	0.001	0.001
$LM_T(\hat{\mathbf{S}}_T)$	0.074	0.067	0.072	0.067	0.064	0.073	0.064	0.072	0.070
$CLR_T(\hat{\mathbf{S}}_T)$	0.075	0.068	0.071	0.067	0.073	0.073	0.065	0.072	0.070
split – S	0.015	0.015	0.015	0.015	0.015	0.016	0.016	0.015	0.016
split – CLR	0.096	0.090	0.086	0.094	0.092	0.082	0.091	0.094	0.087
qqL – S	0.008	0.009	0.008	0.008	0.009	0.010	0.010	0.009	0.009
ave – S	0.020	0.025	0.023	0.020	0.024	0.023	0.021	0.027	0.023
exp – S	0.004	0.034	0.005	0.005	0.004	0.004	0.005	0.004	0.004
$T = 200, \pi_0 = 0.8$	$\rho = 0.25$			$\rho = 0.50$			$\rho = 0.75$		
	$d = 2$	$d = 4$	$d = 8$	$d = 2$	$d = 4$	$d = 8$	$d = 2$	$d = 4$	$d = 8$
AR_T	0.000	0.000	0.000	0.001	0.004	0.001	0.001	0.000	0.001
LM_T	0.061	0.063	0.058	0.056	0.066	0.059	0.060	0.066	0.059
CLR_T	0.048	0.050	0.062	0.059	0.070	0.061	0.063	0.068	0.062
$AR_T(\hat{\mathbf{S}}_T)$	0.002	0.002	0.002	0.002	0.003	0.002	0.002	0.004	0.002
$LM_T(\hat{\mathbf{S}}_T)$	0.064	0.061	0.065	0.067	0.072	0.067	0.070	0.079	0.055
$CLR_T(\hat{\mathbf{S}}_T)$	0.064	0.062	0.065	0.067	0.073	0.068	0.069	0.079	0.054
split – S	0.020	0.017	0.021	0.022	0.019	0.021	0.018	0.018	0.020
split – CLR	0.102	0.098	0.086	0.103	0.100	0.098	0.103	0.108	0.098
qqL – S	0.009	0.074	0.009	0.010	0.008	0.011	0.010	0.008	0.009
ave – S	0.020	0.022	0.018	0.023	0.022	0.019	0.019	0.025	0.018
exp – S	0.004	0.026	0.004	0.005	0.003	0.005	0.005	0.005	0.004

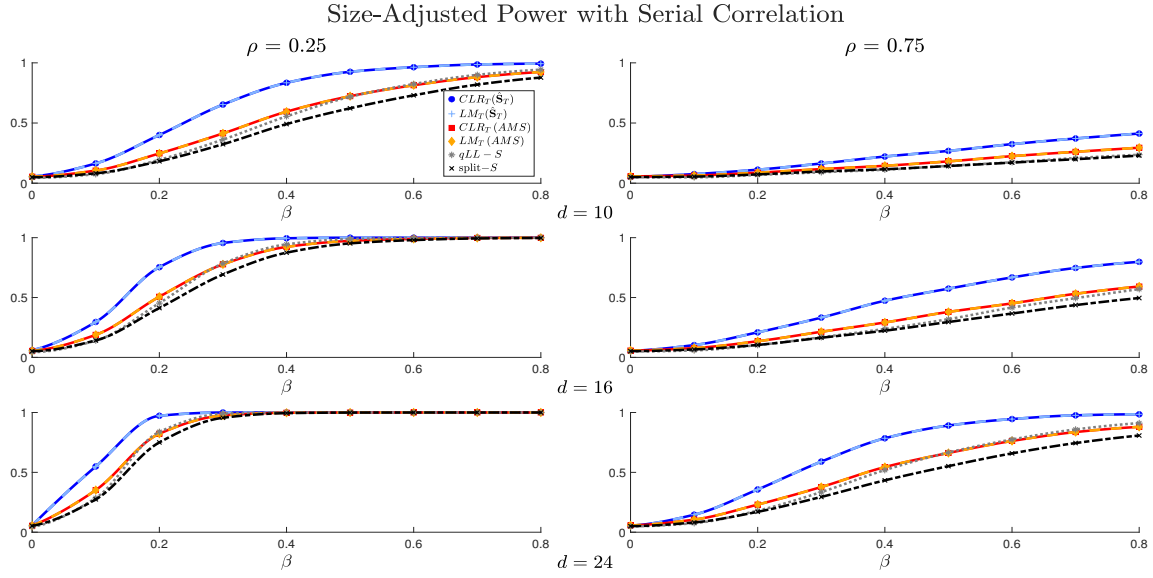


Figure 1: Size-adjusted power of identification robust tests for $T = 400$ and $\pi_0 = 0.6$ in (N.1)-(N.2).

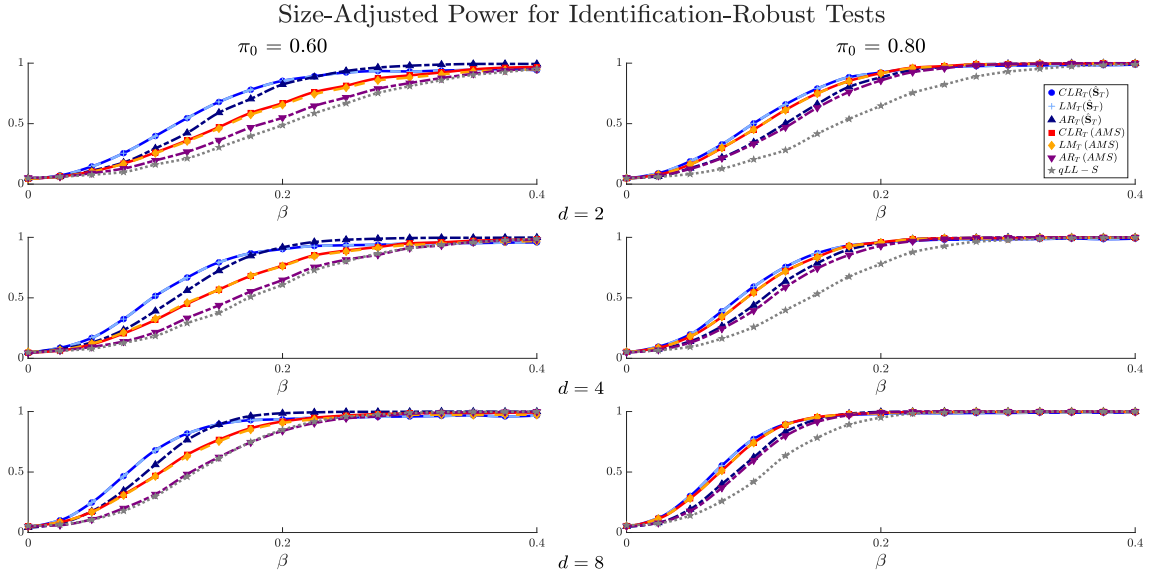


Figure 2: Size-adjusted power of identification robust tests for $T = 200$ in (N.1)-(N.2).

References

- AMEMIYA, T. (1985): *Advanced Econometrics*. Harvard University Press.
- ANDREWS, D. W. K. (1991): “Heteroskedasticity and Autocorrelation Consistent Covariance Matrix Estimation,” *Econometrica*, 59(3), 817–858.
- ANDREWS, D. W. K., M. MOREIRA, AND J. H. STOCK (2004): “Optimal Two-Sided Invariant Similar Tests for Instrumental Variables Regression,” Discussion Paper 1476, Cowles Foundation, Yale University.
- (2006): “Optimal Two-Sided Invariant Similar Tests for Instrumental Variables Regression,” *Econometrica*, 74(3), 715–752.
- BELOTTI, F., A. CASINI, L. CATANIA, S. GRASSI, AND P. PERRON (2023): “Simultaneous Bandwidths Determination for Double-Kernel HAC Estimators and Long-Run Variance Estimation in Nonparametric Settings,” *Econometric Reviews*, 42(3), 281–306.
- CASINI, A. (2022): “Comment on Andrews (1991) ”Heteroskedasticity and Autocorrelation Consistent Covariance Matrix Estimation”,” *Econometrica*, 90(4), 1–2.
- (2023): “Theory of Evolutionary Spectra for Heteroskedasticity and Autocorrelation Robust Inference in Possibly Misspecified and Nonstationary Models,” *Journal of Econometrics*, 235(2), 372–392.
- (2024): “The Fixed-b Limiting Distribution and the ERP of HAR Tests Under Nonstationarity,” *Journal of Econometrics*, 238(2), 105625.
- CASINI, A., T. DENG, AND P. PERRON (2025): “Theory of Low Frequency Contamination from Nonstationarity and Misspecification: Consequences for HAR Inference,” *Econometric Theory*, forthcoming.
- CASINI, A., AND P. PERRON (2024): “Prewhitened Long-Run Variance Estimation Robust to Nonstationarity,” *Journal of Econometrics*, 242(1), 105794.
- NEWKEY, W. K., AND K. D. WEST (1987): “A Simple Positive Semidefinite, Heteroskedastic and Autocorrelation Consistent Covariance Matrix,” *Econometrica*, 55(3), 703–708.
- RAMEY, V. A. (2011): “Identifying Government Spending Shocks: It’s All in the Timing,” *Quarterly Journal of Economics*, 126(1), 1–50.
- ROMER, C. D., AND D. H. ROMER (2004): “A New Measure of Monetary Shocks: Derivation and Implications,” *American Economic Review*, 94(4), 1055–1084.