# Online Appendix for <br> The impact of uncertainty on investment: Empirical challenges and a new estimator 

The online Appendix contains five sections. Section A provides details on how to construct orthogonal moment conditions to treat measurement error in $q$. Section B gives the assumptions and limiting results for our proposed estimators. Appendix C explains how to construct nonparametric model specification test statistics. Appendix D describes how to construct the 10 instrumental variables for the endogenous individual stock volatility variable. Section E shows more robust results. Section F provides the mathematical proofs of the theorems in Section B.

## Appendix A. Treating measurement error

In this section, we provide details on how to construct orthogonal moment conditions to treat measurement error in $q$. For ease in reading, we introduce our notation here. (i) $\mathbf{i}_{T}$ denotes a $T \times 1$ vector of ones, $0_{n}$ denotes an $n \times 1$ vector of zeros, $0_{m \times n}$ is an $m \times n$ matrix of zeros, $I_{T}$ is a $T \times T$ identity matrix, and $J_{T}=I_{T}-T^{-1} \mathbf{i}_{T} \mathbf{i}_{T}^{\prime}$; (ii) $\tilde{\mathbf{a}}_{i}=J_{T} \mathbf{a}_{i}$ denotes the demeaned data for any $T \times 1$ vector $\mathbf{a}_{i}$; (iii) $M, M_{1}, M_{2}, \ldots$ are constants that can take different values at different locations. Denoting a $T \times\left(3 k_{n T}\right)$ matrix $\mathbf{E}_{i}=\left[\mathbf{E}_{i, 1}, \ldots, \mathbf{E}_{i, T}\right]^{\prime}$, rewrite model (14) in matrix form as follows

$$
\begin{equation*}
\mathbf{y}_{i} \approx \mu_{i} \mathbf{i}_{T}+\boldsymbol{\lambda}_{0}+\alpha_{0} \mathbf{x}_{i,-1}+\mathbf{E}_{i} \boldsymbol{\vartheta}_{0}+\boldsymbol{\varepsilon}_{i} \tag{A.1}
\end{equation*}
$$

and premultiplying $J_{T}$ to both sides of model (A.1) gives

$$
\begin{align*}
\tilde{\mathbf{y}}_{i} & \approx J_{T} \boldsymbol{\lambda}_{0}+\alpha_{0} \tilde{\mathbf{x}}_{i,-1}+\tilde{\mathbf{E}}_{i} \boldsymbol{\vartheta}_{0}+\tilde{\boldsymbol{\varepsilon}}_{i}  \tag{A.2}\\
& =F_{T} \tilde{\boldsymbol{\lambda}}_{0}+\alpha_{0} \tilde{\mathbf{x}}_{i,-1}+\tilde{\mathbf{E}}_{i} \boldsymbol{\vartheta}_{0}+\tilde{\boldsymbol{\varepsilon}}_{i} \tag{A.3}
\end{align*}
$$

where $\tilde{\mathbf{y}}_{i}=J_{T} \mathbf{y}_{i}, \tilde{\mathbf{E}}_{i}=J_{T} \mathbf{E}_{i}$ and $\tilde{\varepsilon}_{i}=J_{T} \boldsymbol{\varepsilon}_{i}$ are all demeaned data. Because $J_{T}$ does not have full rank, we redefine $J_{T} \boldsymbol{\lambda}_{0}=F_{T} \tilde{\boldsymbol{\lambda}}_{0}$, where $F_{T}=\left[I_{T-1},-\mathbf{i}_{T-1}\right]^{\prime}, \tilde{\boldsymbol{\lambda}}_{0}=\left[\tilde{\lambda}_{1}, \ldots, \tilde{\lambda}_{T-1}\right]^{\prime}$, and $\tilde{\lambda}_{t}=\lambda_{t}-T^{-1} \sum_{s=1}^{T} \lambda_{s}$.

In the following, we construct two blocks of orthogonal moment conditions, in the spirit of Meijer, Spierdijk and Wansbeek (2017). The moments in Section A2 are sufficient to identify our model, and we use them to produce the empirical results in this paper. In contrast, the moments in Section A1 are empirically optional, and we include them in the econometric theory for completeness.

## A1. Intertemporal covariance matrix

These moment conditions are based on the intertemporal covariance matrix of $\boldsymbol{\omega}_{i}$ in model (10), $\Sigma_{\omega}=\mathbb{E}\left(\boldsymbol{\omega}_{i} \boldsymbol{\omega}_{i}^{\prime}\right)$, by exploiting the cross-sectional independence across firms:

$$
\begin{align*}
& \mathbb{E}\left(\tilde{\varepsilon}_{i} \otimes \tilde{\mathbf{y}}_{i}\right) \approx \mathbb{E}\left[\left(\tilde{\boldsymbol{\omega}}_{i}-\alpha_{0} \tilde{\mathbf{e}}_{i,-1}\right) \otimes\left(F_{T} \tilde{\boldsymbol{\lambda}}_{0}+\alpha_{0} \tilde{\mathbf{x}}_{i,-1}^{*}+\tilde{\mathbf{E}}_{i} \boldsymbol{\vartheta}_{0}+\tilde{\boldsymbol{\omega}}_{i}\right)\right] \\
& =\mathbb{E}\left(\tilde{\boldsymbol{\omega}}_{i} \otimes \tilde{\boldsymbol{\omega}}_{i}\right)=\mathbb{E}\left[\operatorname{vec}\left(J_{T} \boldsymbol{\omega}_{i} \boldsymbol{\omega}_{i}^{\prime} J_{T}\right)\right]=\left(J_{T} \otimes J_{T}\right) \mathbb{E}\left[\operatorname{vec}\left(\boldsymbol{\omega}_{i} \boldsymbol{\omega}_{i}^{\prime}\right)\right] \\
& =\left(J_{T} \otimes J_{T}\right) \operatorname{vec}\left(\Sigma_{\omega}\right)=\left(J_{T} \otimes J_{T}\right) D_{T} \pi_{\omega} \tag{A.4}
\end{align*}
$$

where $\pi_{\omega}=\operatorname{vech}\left(\Sigma_{\omega}\right)$, and $D_{T}$ is the duplication matrix of dimension $T^{2} \times m_{0}$, with $m_{0}=$ $T(T+1) / 2$, such that $\operatorname{vec}\left(\Sigma_{\omega}\right)=D_{T} \pi_{\omega} .{ }^{1}$ Let $\left[\left(J_{T} \otimes J_{T}\right) D_{T}\right]_{\perp}$ be the orthogonal complement of $\left(J_{T} \otimes J_{T}\right) D_{T}$, i.e., the $T^{2} \times m_{0}$ matrix with full column rank satisfying

$$
\left[\left(J_{T} \otimes J_{T}\right) D_{T}\right]_{\perp}^{\prime}\left[\left(J_{T} \otimes J_{T}\right) D_{T}\right]=0_{m_{0} \times m_{0}} .
$$

Pre-multiplying both sides of (A.4) by $\left[\left(J_{T} \otimes J_{T}\right) D_{T}\right]_{\perp}^{\prime}$ gives

$$
\left[\left(J_{T} \otimes J_{T}\right) D_{T}\right]_{\perp}^{\prime} \mathbb{E}\left(\tilde{\varepsilon}_{i} \otimes \tilde{\mathbf{y}}_{i}\right) \approx 0_{m_{0}}
$$

which is equivalent to ${ }^{2}$

$$
\begin{equation*}
\left[\left(J_{T} \otimes J_{T}\right) D_{T}\right]_{\perp}^{\prime} \mathbb{E}\left[\left(I_{T} \otimes \tilde{\mathbf{y}}_{i}\right) \tilde{\varepsilon}_{i}\right] \approx 0_{m_{0}} \tag{A.5}
\end{equation*}
$$

As a result, the $\left(T \times m_{0}\right)$ matrix $\tilde{\mathbf{d}}_{1, i}=\left(I_{T} \otimes \tilde{\mathbf{y}}_{i}\right)^{\prime}\left[\left(J_{T} \otimes J_{T}\right) D_{T}\right]_{\perp}$ acts as valid instru-

[^0]ments that are orthogonal to $\tilde{\varepsilon}_{i}$. Moreover, $\tilde{\mathbf{d}}_{1, i}$ are relevant instruments if
\[

$$
\begin{equation*}
\mathbb{E}\left(\tilde{\mathbf{d}}_{1, i}^{\prime} \tilde{\mathbf{x}}_{i,-1}\right) \neq 0_{m_{0}} \tag{A.6}
\end{equation*}
$$

\]

which generally holds true as $\left[\left(J_{T} \otimes J_{T}\right) D_{T}\right]_{\perp}^{\prime} \mathbb{E}\left(\tilde{\mathbf{x}}_{i,-1} \otimes\left(\tilde{\mathbf{E}}_{i} \boldsymbol{\vartheta}_{0}\right)\right) \neq 0_{m_{0}}$ if $\boldsymbol{\vartheta}_{0} \neq 0_{3 k}$ and $\tilde{\mathbf{x}}_{i,-1}$ and $\tilde{\mathbf{E}}_{i}$ are correlated.

## A2. Exogenous regressors

These orthogonal moment conditions are based on the strict exogeneity of time-fixed effects and $\tilde{\mathbf{E}}_{i}$ in model (A.3), and include

$$
\begin{align*}
& \mathbb{E}\left[\left(\tilde{\mathbf{y}}_{i}-\tilde{\mathbf{q}}_{i} \boldsymbol{\theta}_{0}\right) \otimes \tilde{\mathbf{E}}_{i}\right] \approx \mathbb{E}\left(\tilde{\varepsilon}_{i} \otimes \tilde{\mathbf{E}}_{i}\right)=0_{T^{2} \times\left(3 k_{n T}\right)}  \tag{A.7}\\
& \mathbb{E}\left[\left(\tilde{\mathbf{y}}_{i}-\tilde{\mathbf{q}}_{i} \boldsymbol{\theta}_{0}\right) \otimes F_{T}\right] \approx \mathbb{E}\left(\tilde{\varepsilon}_{i} \otimes F_{T}\right)=0_{T^{2} \times(T-1)} \tag{A.8}
\end{align*}
$$

where $\tilde{\mathbf{q}}_{i}=\left[F_{T}, \tilde{\mathbf{x}}_{i,-1}, \tilde{\mathbf{E}}_{i}\right]$ contains all the regressors in model (A.3) and $\boldsymbol{\theta}_{0}=\left[\tilde{\boldsymbol{\lambda}}_{0}^{\prime}, \alpha_{0}, \boldsymbol{\vartheta}_{0}^{\prime}\right]^{\prime}$ includes all the parameters to be estimated. These moment conditions are valid because after the control function approach, the only endogenous variable in model (A.3) is $\tilde{\mathbf{x}}_{i,-1}$. Consequently, denoting a $T \times\left[T^{2}\left(T-1+3 k_{n T}\right)\right]$ matrix

$$
\tilde{\mathbf{d}}_{2, i}=\left[I_{T} \otimes \operatorname{vec}\left(\tilde{\mathbf{E}}_{i}\right)^{\prime} \quad I_{T} \otimes \operatorname{vec}\left(F_{T}\right)^{\prime}\right]
$$

we obtain the following orthogonal moment conditions

$$
\begin{equation*}
\mathbb{E}\left(\tilde{\mathbf{d}}_{2, i}^{\prime} \tilde{\varepsilon}_{i}\right) \approx 0_{T^{2}\left(T-1+3 k_{n T}\right)} \tag{A.9}
\end{equation*}
$$

Moreover, $\tilde{\mathbf{d}}_{2, i}$ are relevant instruments if $\mathbb{E}\left(\tilde{\mathbf{d}}_{2, i}^{\prime} \tilde{\mathbf{x}}_{i,-1}\right) \neq 0$, which can be easily tested in the data. We further select a smaller set of instruments based on the method developed in Belloni, Chen, Chernozhukov and Hansen (2012) to mitigate the concern of many-instruments bias in a finite sample.

## Appendix B. Limiting results

For a sufficiently large $n$ and a fixed $T$, we show three theorems in the following (with proofs in the Online Appendix), ensuring that both the penalty estimator, $\hat{\boldsymbol{\theta}}$, and the postpenalty estimator, $\tilde{\boldsymbol{\theta}}$, converge to the true parameter values $\boldsymbol{\theta}_{0}$. The estimators of the Tobin's $q$ coefficient and the three unknown curves are consistent and have asymptotic normal distributions. Theorem B. 1 proves that if we knew which elements in $\boldsymbol{\theta}_{0}$ were zero,
then dropping the corresponding regressors and conducting traditional GMM estimation by minimizing (18) in Li and Sun (2022) would produce a uniformly consistent estimator for the nonzero elements in $\boldsymbol{\theta}_{0}$. At the same time, the estimators of $\alpha_{0}, f_{0}(z), g_{0}(s)$, and $r_{0}(v)$ possess asymptotic normal distributions. Theorem B. 2 demonstrates that for any given tuning parameter $\psi$, there exists a penalty estimator $\hat{\boldsymbol{\theta}}(\psi)$ that minimizes the objective function (19) in Li and $\operatorname{Sun}$ (2022). For sufficiently large $n$ and with a probability approaching 1, this penalty estimator equals the traditional GMM estimator in Theorem B. 1 as if we knew which elements in $\boldsymbol{\theta}_{0}$ were zero. That is, the penalty estimator is oracle efficient. Theorem B. 3 reveals that the post-penalty estimator $\tilde{\boldsymbol{\theta}}$ that solves (20) in Li and Sun (2022) is same as $\hat{\boldsymbol{\theta}}$ asymptotically with a probability approaching 1 .

In Appendix A, we construct the following instruments to deal with the error-ridden $q$ :

$$
\tilde{\mathbf{d}}_{i}=\left[\tilde{\mathbf{d}}_{1, i} \tilde{\mathbf{d}}_{2, i}\right]=\left[\begin{array}{ll}
\left(I_{T} \otimes \tilde{\mathbf{y}}_{i}\right)^{\prime}\left[\left(J_{T} \otimes J_{T}\right) D_{T}\right]_{\perp} & I_{T} \otimes \operatorname{vec}\left(\tilde{\mathbf{E}}_{i}\right)^{\prime} \quad I_{T} \otimes \operatorname{vec}\left(F_{T}\right)^{\prime} \tag{B.1}
\end{array}\right]
$$

which is a $T \times m_{i v}$ matrix, with the number of instruments equal to

$$
m_{i v}=m_{0}+T^{2}\left(T-1+3 k_{n T}\right) .
$$

We then obtain the following orthogonal moment conditions (equivalent to (17) in Li and Sun (2022)):

$$
\begin{equation*}
\mathbb{E}\left[\tilde{\mathbf{d}}_{i}^{\prime}\left(\tilde{\mathbf{y}}_{i}-\tilde{\mathbf{q}}_{i} \boldsymbol{\theta}_{0}\right)\right] \approx 0_{m_{i v}} \tag{B.2}
\end{equation*}
$$

Decompose $\boldsymbol{\theta}_{0}=\left[\boldsymbol{\theta}_{1,0}^{\prime}, 0_{p-J}^{\prime}\right]^{\prime}$, where $\boldsymbol{\theta}_{1,0}^{\prime}$ is the $J \times 1$ vector containing all the non-zero parameters, and $J$ is the number of elements in $\mathcal{J}=\operatorname{supp}\left(\boldsymbol{\theta}_{0}\right)$. Divide the regressors accordingly into two parts, $\tilde{\mathbf{q}}_{i}=\left[\tilde{\mathbf{q}}_{1, i}, \tilde{\mathbf{q}}_{2, i}\right]$, where $\tilde{\mathbf{q}}_{1, i}$ and $\tilde{\mathbf{q}}_{2, i}$ are $T \times J$ and $T \times(p-J)$ matrix, respectively. Our orthogonal moment conditions in (B.2) can be rewritten in this partition by removing all the regressors with zero coefficients, as $\tilde{\mathbf{q}}_{i} \boldsymbol{\theta}_{0}=\tilde{\mathbf{q}}_{1, i} \boldsymbol{\theta}_{1,0}$. Let $\check{\boldsymbol{\theta}}_{1}$ be the estimator of $\boldsymbol{\theta}_{1,0}$, which solves the following optimization:

$$
\begin{equation*}
\min _{\boldsymbol{\theta}_{1} \in \Theta_{1}} \overline{\boldsymbol{\varphi}}_{n}\left(\boldsymbol{\theta}_{1}\right)^{\prime} \boldsymbol{\Omega}_{n} \overline{\boldsymbol{\varphi}}_{n}\left(\boldsymbol{\theta}_{1}\right) \tag{B.3}
\end{equation*}
$$

where $\Theta_{1}$ is a compact subset of $R^{J}$, and

$$
\begin{equation*}
\bar{\varphi}_{n}\left(\boldsymbol{\theta}_{1}\right)=\frac{1}{n} \sum_{i=1}^{n} \widetilde{\hat{\mathbf{d}}}_{i}^{\prime}\left(\tilde{\mathbf{y}}_{i}-\widetilde{\hat{\mathbf{q}}}_{1, i} \boldsymbol{\theta}_{1}\right) \tag{B.4}
\end{equation*}
$$

where $\widetilde{\hat{\mathbf{q}}}_{1, i}$ equals $\tilde{\mathbf{q}}_{1, i}$ with $\mathbf{P}^{k_{n T}}\left(v_{i, t}\right)$ replaced by $\mathbf{P}^{k_{n T}}\left(\hat{v}_{i, t}\right)$ for all $t$. Let $P_{\mathcal{J}}$ be the $J \times p$
selection matrix that satisfies $\boldsymbol{\theta}_{0}=P_{\mathcal{J}}^{\prime} \boldsymbol{\theta}_{1,0} \cdot{ }^{3}$ Consequently, $\check{\boldsymbol{\theta}}=P_{\mathcal{J}}^{\prime} \check{\boldsymbol{\theta}}_{1}$ is an estimator for $\boldsymbol{\theta}_{0}$, when $\mathcal{J}$ is known. We list additional regularity conditions that support the limiting results in the following two assumptions.

Assumption 3. (i) $\mathbb{E}\left\|T^{-1} \mathbf{W}_{i} \mathbf{W}_{i}^{\prime}\right\| \leq M$; also, $T^{-1} \mathbb{E}\left(\left(\mathbf{i}_{T}, \mathbf{W}_{i}\right)^{\prime}\left(\mathbf{i}_{T}, \mathbf{W}_{i}\right)\right), \Sigma_{v}$, and $T^{-1} \mathbb{E}\left(\mathbf{W}_{i} \Sigma_{v} \mathbf{W}_{i}^{\prime}\right)$ are all positive definite matrix; (ii) There exist two constants, $\varsigma_{0}$ and $\varsigma_{1}$, satisfying

$$
0<\varsigma_{0} \leq \lambda_{\min }\left(\mathbb{E}\left(\widetilde{\mathbf{q}}_{i}^{\prime} \tilde{\mathbf{d}}_{i}\right) \Omega_{n} \mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \widetilde{\mathbf{q}}_{i}\right)\right) \leq \lambda_{\max }\left(\mathbb{E}\left(\widetilde{\mathbf{q}}_{i}^{\prime} \tilde{\mathbf{d}}_{i}\right) \Omega_{n} \mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \widetilde{\mathbf{q}}_{i}\right)\right) \leq \varsigma_{1}<\infty
$$

where $\lambda_{\min }(\cdot)$ and $\lambda_{\max }(\cdot)$ denote the smallest and largest eigenvalues of a matrix, respectively; (iii) $\boldsymbol{\Omega}_{n}$ is a non-stochastic positive definite matrix with $\lambda_{\max }\left(\boldsymbol{\Omega}_{n}\right) \leq M<\infty$.

Assumption 4. As $n \rightarrow \infty$, the following conditions hold: (i) $k_{n T} \rightarrow \infty$; (ii) $\sqrt{n} k_{n T}^{-\zeta} \rightarrow 0$; (iii) $k_{n T}^{2} / n \rightarrow 0$; and (iv) $n^{-1} k_{n T}\left\|\mathbf{P}^{k_{n T}}\right\|_{1}^{2} \rightarrow 0$, where $\left\|\mathbf{P}^{k_{n T}}\right\|_{l}=\max _{0 \leq j \leq l} \sup _{x \in R}\left\|\frac{\partial^{j} \mathbf{P}^{k_{n T}(x)}}{\partial x^{j}}\right\|$ for $l \geq 0$, with $\|\cdot\|$ being the Euclidean norm.

Assumption 3(i) regulates the stochastic property of $\left(\mathbf{w}_{i t}, v_{i t}\right)$ in the time dimension to ensure that $\boldsymbol{\pi}_{0}$ in model (6) in Li and Sun (2022) can be estimated at the root- $n$ convergence rate so that the estimated residuals $\hat{v}_{i, t}$ have stochastic properties mimicking those of the true error terms $v_{i, t}$. Assumption 3 (ii) is standard in the literature on series approximation, ensuring the existence of the proposed estimators with nonsingular variance and covariance matrices. In Assumption 4, (i) and (iii) are standard conditions in series approximation, (ii) is a technical condition, and (iv) is used to remove the asymptotic impact of the first-stage estimation of $\boldsymbol{\pi}_{0}$ on the estimation of model (13) in Li and Sun (2022). For the smoothing parameter, the existing literature commonly sets $k_{n T}=c_{k}(n T)^{r}$, where $c_{k}>0$ and $r>0$ are constants. In this setup, Assumption 4 requires that $(2 \zeta)^{-1}<r<0.25$ if the Hermite series are adopted in the series approximation because $\left\|\mathbf{P}^{k_{n T}}\right\|_{1}=O\left(k_{n T}^{3 / 2}\right)$.

Theorem B. 1 Denoting $a_{n}=k_{n T}^{-\zeta}+\sqrt{J / n}$, under Assumptions 1-4, we have
(i) $\left\|\check{\boldsymbol{\theta}}_{1}-\boldsymbol{\theta}_{1,0}\right\|=O_{p}\left(a_{n}\right)$;
(ii)

$$
\begin{aligned}
& \sup _{z \in \mathcal{S}_{z}}\left\|\check{f}(z)-f_{0}(z)\right\|=O_{p}\left(a_{n} k_{n T}^{1 / 2}\right), \\
& \sup _{s \in \mathcal{S}_{s}}\left\|\check{g}(s)-g_{0}(s)\right\|=O_{p}\left(a_{n} k_{n T}^{1 / 2}\right), \\
& \sup _{v \in \mathcal{S}_{v}}\left\|\check{r}(v)-r_{0}(v)\right\|=O_{p}\left(a_{n} k_{n T}^{1 / 2}\right),
\end{aligned}
$$

where $\check{f}(z)=\check{\boldsymbol{\theta}}_{z}^{\prime} \overrightarrow{\mathbf{P}}_{z}^{k_{n T}}, \check{g}(s)=\check{\boldsymbol{\theta}}_{s}^{\prime} \overrightarrow{\mathbf{P}}_{s}^{k_{n T}}$, and $\check{r}(v)=\check{\boldsymbol{\theta}}_{v}^{\prime} \overrightarrow{\mathbf{P}}_{v}^{k_{n T}}$, with $\overrightarrow{\mathbf{P}}^{k_{n T}}(\cdot)$ equal to $\mathbf{P}^{k_{n T}}(\cdot)$

[^1]after removing the terms with zero coefficients, and $\mathcal{S}_{z}, \mathcal{S}_{s}$, and $\mathcal{S}_{v}$ denote the support of $z_{i t}$, $s_{i t}$, and $v_{i t}$, respectively; (iii) if $\sqrt{n / J} k_{n T}^{-\zeta} \rightarrow 0$, we have $\sqrt{n / \hat{\sigma}_{\alpha}^{2}}\left(\check{\alpha}-\alpha_{0}\right) \xrightarrow{d} N(0,1)$ and
\[

\sqrt{n \hat{\Xi}_{n}^{-1}}\left[$$
\begin{array}{l}
\check{f}(z)-f_{0}(z) \\
\check{g}(s)-g_{0}(s) \\
\check{r}(v)-r_{0}(v)
\end{array}
$$\right] \xrightarrow{d} N\left(0_{3}, I_{3}\right),
\]

provided that $\hat{\sigma}_{\alpha}^{2} \xrightarrow{p} \sigma_{\alpha}^{2}>0$ and that $\hat{\Xi}_{n}^{-1}$ is nonsingular, where $\hat{\sigma}_{\alpha}^{2}, \sigma_{\alpha}^{2}$, and $\hat{\Xi}_{n}$ are denoted in (F.13) and (F.14) in the Online Appendix.

When $\mathcal{J}$ is known, Theorem B.1(i) shows the uniform consistency of $\check{\boldsymbol{\theta}}$. Meanwhile, the estimated coefficient $\check{\alpha}$ is a root- $n$ consistent estimator of $\alpha_{0}$. The pointwise convergence rate of $\check{f}(z), \check{g}(s)$, and $\check{r}(v)$ are of order $O_{p}(\sqrt{J / n})$ because $\left\|\hat{\Xi}_{n}\right\|=O_{p}(J)$.

Theorem B. 2 Under Assumptions 1-4, there exists a local minimizer $\hat{\boldsymbol{\theta}}\left(\psi_{n}\right)$ for (19) that satisfies $\operatorname{Pr}\left(\hat{\boldsymbol{\theta}}\left(\psi_{n}\right)=\check{\boldsymbol{\theta}}\right) \rightarrow 1$ asymptotically, provided that the tuning parameter $\psi_{n}$ meets the following conditions: (i) $\psi_{n} \rightarrow 0$; (ii) $\psi_{n} / a_{n} \rightarrow \infty$; and (iii) $\psi_{n} \leq \min _{l \in \mathcal{J}}\left\{\theta_{l, 0}\right\} /\left(q_{0} c\right)$ for some $c>1$ and $q_{0}<1 / 2$.

Theorem B. 2 implies that the penalty estimator $\hat{\boldsymbol{\theta}}\left(\psi_{n}\right)$ performs as if we knew $\mathcal{J}$, i.e., which parameters in $\boldsymbol{\theta}_{0}$ do not equal zero. This property is known as the estimator being oracle efficient. The next theorem shows that the post-MCP-penalty estimator $\tilde{\boldsymbol{\theta}}$ converges to the true parameters, $\boldsymbol{\theta}_{0}$, at the same convergence rate as the penalty estimator. This implies that the post-MCP estimator is as good as the MCP-penalty estimator.

Theorem B. 3 Under the assumptions in Theorems B. 1 and B.2, we obtain

$$
\left\|\tilde{\boldsymbol{\theta}}-\boldsymbol{\theta}_{0}\right\|=O_{p}\left(k_{n T}^{-2 \zeta}+J / n\right),
$$

and that Theorem $B .1$ (ii) and (iii) hold for the post-MCP estimator, $\tilde{f}(z), \tilde{g}(s)$ and $\tilde{r}(v)$.

## Appendix C. Hypothesis tests

Based on the limiting results, we introduce asymptotically valid inference procedures for two types of hypotheses broadly. First, we conduct standard Wald-type tests for the significance of the levels of the estimates. That is, we test for whether each of the time effects, $\left\{\tilde{\lambda}_{t}\right\}_{t=1}^{T-1}$, and the Tobin's $q$ coefficient, $\alpha_{0}$, are zero, and whether the three unknown functions $f_{0}(z)=0, g_{0}(s)=0$, and $r_{0}(v)=0$, respectively. The tests for $f_{0}(z)=0$ and $g_{0}(s)=0$ shed light on whether cash flow and individual stock volatility contribute to
investment in addition to $q$. The test for $r_{0}(v)=0$ instead indicates whether the variable $s_{i, t}$ is correlated with $u_{i, t}$, as $r_{0}\left(v_{i, t}\right)=\mathbb{E}\left(u_{i, t} \mid v_{i, t}\right)$; put differently, it tests for whether the individual stock volatility is a contemporaneous endogenous variable.

Second, we construct the following group of pointwise tests for whether the first-order derivatives of the three unknown curves, $f_{0}^{\prime}(z), g_{0}^{\prime}(s)$, and, $r_{0}^{\prime}(v)$ are constant. We do so to test for the existence of nonlinearities.

$$
\begin{align*}
& H_{0}^{f}: f_{0}(z)=\beta_{f, 0} z, \text { and } H_{1}^{f}: \operatorname{not} H_{0}^{f}  \tag{C.1}\\
& H_{0}^{g}: g_{0}(s)=\beta_{g, 0} s, \text { and } H_{1}^{g}: \text { not } H_{0}^{g}  \tag{C.2}\\
& H_{0}^{r}: r_{0}(v)=\beta_{v, 0} v, \text { and } H_{1}^{r}: \text { not } H_{0}^{r} . \tag{C.3}
\end{align*}
$$

We impose an additional assumption that is stronger than Assumption 2 to asymptotically remove the series approximation bias of the first-order derivatives.

Assumption 5. There exists $\boldsymbol{\vartheta}_{0}=\left[\boldsymbol{\vartheta}_{z, 0}^{\prime}, \boldsymbol{\vartheta}_{s, 0}^{\prime}, \boldsymbol{\vartheta}_{v, 0}^{\prime}\right]^{\prime}$ such that

$$
\begin{aligned}
& \max _{0 \leq l \leq l_{0}} \sup _{z \in \mathcal{S}_{z}}\left|f_{0}^{(l)}(z)-\boldsymbol{\vartheta}_{z, 0}^{\prime} d^{l} \mathbf{P}^{k_{n T}}(z)\right| \leq M_{1} k_{n T}^{-\zeta} \\
& \max _{0 \leq l \leq l_{0}} \sup _{s \in \mathcal{S}_{s}}\left|g_{0}^{(l)}(z)-\boldsymbol{\vartheta}_{s, 0}^{\prime} d^{l} \mathbf{P}^{k_{n T}}(s)\right| \leq M_{2} k_{n T}^{-\zeta} \\
& \max _{0 \leq l \leq l_{0}} \sup _{v \in \mathcal{S}_{v}}\left|r_{0}^{(l)}(v)-\boldsymbol{\vartheta}_{v, 0}^{\prime} d^{l} \mathbf{P}^{k_{n T}}(v)\right| \leq M_{3} k_{n T}^{-\zeta}
\end{aligned}
$$

for some $\zeta>2$ and a nonnegative integer $l_{0} \leq 1$ as $k_{n T} \rightarrow \infty$, where $f_{0}^{(l)}(\cdot)$ and $d^{l} \mathbf{P}^{k_{n T}}(\cdot)$ denote the $l$-th order partial derivatives of $f_{0}(\cdot)$ and $\mathbf{P}^{k_{n T}}(\cdot)$, respectively.

We use the test of (C.1) as an example to illustrate the basic ideas. Under $H_{0}^{f}$, model (14) in the paper becomes

$$
\begin{equation*}
y_{i, t} \approx \mu_{i}+\lambda_{t}+\alpha_{0} x_{i, t-1}+\beta_{f, 0} z_{i, t-1}+\boldsymbol{\vartheta}_{s, 0}^{\prime} \mathbf{P}^{k}\left(s_{i, t-1}\right)+\boldsymbol{\vartheta}_{v, 0}^{\prime} \mathbf{P}^{k}\left(v_{i, t}\right)+\varepsilon_{i, t}, \tag{C.4}
\end{equation*}
$$

in which we replace the nonparametric Hermite expansion $\boldsymbol{\vartheta}_{z, 0}^{\prime} \mathbf{P}^{k}\left(z_{i, t-1}\right)$ with the linear specification $\beta_{f, 0} z_{i, t-1}$. The equivalent hypothesis is

$$
H_{0}: f_{0}^{\prime}(z)=\beta_{f, 0} \text { for all } z, \text { and } H_{1}: \text { not } H_{0}
$$

The post-penalty estimator for $f_{0}^{\prime}(z)$ is defined as

$$
\tilde{f}^{\prime}(z)=d \overrightarrow{\mathbf{P}}^{k_{n T}}(z)^{\prime} \tilde{\boldsymbol{\vartheta}}_{z, 1}
$$

where $d \overrightarrow{\mathbf{P}}^{k_{n T}}(z)=\frac{\partial}{\partial z} \overrightarrow{\mathbf{P}}^{k_{n T}}(z)$. We introduce some other notations: $\widehat{\mathcal{J}}=\operatorname{supp}(\hat{\boldsymbol{\theta}})$, and $\hat{J}$ is the dimension of $\widehat{\mathcal{J}}$; split $\widetilde{\hat{\mathbf{q}}}_{i}=\left[\widetilde{\hat{\mathbf{q}}}_{\widehat{\mathcal{J}}}^{1, i}, \overline{,}, \widetilde{\hat{\mathbf{q}}}_{\widehat{\mathcal{J}}, 2, i}\right]$ and $\tilde{\boldsymbol{\theta}}=\left[\tilde{\boldsymbol{\theta}}_{1}^{\prime}, 0_{p-\hat{J}}^{\prime}\right]^{\prime}$, where the parameters in front of $\widetilde{\hat{\mathbf{q}}}_{\widehat{\mathcal{J}}, 2, i}$ are all zeros and $\widetilde{\hat{\mathbf{q}}}_{i} \tilde{\boldsymbol{\theta}}=\widetilde{\hat{\mathbf{q}}}_{\widehat{\mathcal{J}}, 1, i} \tilde{\boldsymbol{\theta}}_{1} ; \hat{\sigma}_{n}^{2}(z)=\left(d \overrightarrow{\mathbf{P}}^{k_{n T}}(z)\right)^{\prime} \mathbf{S}_{z} \tilde{\boldsymbol{\Sigma}}_{n} \mathbf{S}_{z}^{\prime} d \overrightarrow{\mathbf{P}}^{k_{n T}}(z) ; \tilde{\boldsymbol{\Sigma}}_{n}$ equals $\hat{\boldsymbol{\Sigma}}_{n}$ with $\widetilde{\hat{\mathbf{q}}}_{1, i}$ replaced with $\widetilde{\hat{\mathbf{q}}}_{\widehat{\mathcal{J}}}^{1, i},{ }^{4} \mathbf{S}_{z}$ is part of $I_{\hat{J}}$ such that $\tilde{\boldsymbol{\vartheta}}_{z, 1}=\mathbf{S}_{z} \tilde{\boldsymbol{\theta}}_{1}$; and $\widehat{\tilde{\boldsymbol{\varepsilon}}}_{i}$ are the residuals calculated under the alternative hypothesis $H_{1}$.

Under the null hypothesis $H_{0}$, when Assumptions 1-5 hold, following the proof of Theorem B.1, we obtain

$$
\sqrt{n / \hat{\sigma}_{n}^{2}(z)}\left(\tilde{f}^{\prime}(z)-\beta_{f, 0}\right) \xrightarrow{d} N(0,1) .
$$

Calculating $\hat{\beta}_{f}$ from model (C.4), we have

$$
\sqrt{n / \hat{\sigma}_{n}^{2}(z)}\left(\tilde{f}^{\prime}(z)-\hat{\beta}_{f}\right)=\sqrt{n / \hat{\sigma}_{n}^{2}(z)}\left(\tilde{f}^{\prime}(z)-\beta_{f, 0}\right)+O_{p}\left(\hat{\sigma}_{n}^{-1}(z)\right)
$$

because $\hat{\beta}_{f}-\beta_{f, 0}=O_{p}\left(n^{-1 / 2}\right)$, and $\hat{\sigma}_{n}^{2}(z) \geq \lambda_{\max }\left(\mathbf{S}_{z} \tilde{\boldsymbol{\Sigma}}_{n} \mathbf{S}_{z}^{\prime}\right)\left\|d \overrightarrow{\mathbf{P}}^{k_{n T}}(z)\right\|^{2} \geq M\left\|d \overrightarrow{\mathbf{P}}^{k_{n T}}(z)\right\|^{2}$ under Assumption 3. This is proved by Lemma F. 3 in the Online Appendix. Hence, we have

$$
\sqrt{n / \hat{\sigma}_{n}^{2}(z)}\left(\tilde{f}^{\prime}(z)-\hat{\beta}_{f}\right) \xrightarrow{d} N(0,1) .
$$

On the other hand, under the alternative hypothesis $H_{1}$,

$$
\sqrt{n / \hat{\sigma}_{n}^{2}(z)}\left(\tilde{f}^{\prime}(z)-\hat{\beta}_{f}\right)=\sqrt{n / \hat{\sigma}_{n}^{2}(z)}\left(f_{0}^{\prime}(z)-\beta_{f}\right)+o_{p}(1) \xrightarrow{p} \infty,
$$

provided that there exists a constant $\beta_{f}$ such that $\hat{\beta}_{f}-\beta_{f}=O_{p}\left(n^{-1 / 2}\right)$. Hence, we can construct a $t$ statistic for any given $z$ point to test $H_{0}$ and $H_{1}$.

Our test statistic can also be extended to simultaneously consider a range of $z$ values. Let $z_{i}^{*}$, where $i=1, \ldots, m$, be $m$ distinct points. Then, under $H_{0}$, we have

$$
\sqrt{n \hat{\Upsilon}_{n}^{-1}}\left(\begin{array}{c}
\tilde{f}^{\prime}\left(z_{1}^{*}\right)-\hat{\beta}_{f}  \tag{C.5}\\
\vdots \\
\tilde{f}^{\prime}\left(z_{m}^{*}\right)-\hat{\beta}_{f}
\end{array}\right) \xrightarrow{d} N\left(0, I_{m}\right),
$$

where

$$
\hat{\Upsilon}_{n}=\left(\begin{array}{c}
d \overrightarrow{\mathbf{P}}^{k_{n T}}\left(z_{1}^{*}\right)^{\prime} \\
\vdots \\
d \overrightarrow{\mathbf{P}}^{k_{n T}}\left(z_{m}^{*}\right)^{\prime}
\end{array}\right) \mathbf{S}_{z} \tilde{\boldsymbol{\Sigma}}_{n} \mathbf{S}_{z}^{\prime}\left[d \overrightarrow{\mathbf{P}}^{k_{n T}}\left(z_{1}^{*}\right) \ldots d \overrightarrow{\mathbf{P}}^{k_{n T}}\left(z_{m}^{*}\right)\right]
$$

[^2]is an $m \times m$ non-singular matrix. We can construct a $\chi^{2}$ statistic as follows: under $H_{0}$,
\[

$$
\begin{equation*}
T_{f, n}=n\left\|\hat{\Upsilon}_{n}^{-1 / 2}\left[\tilde{f}^{\prime}\left(z_{1}^{*}\right)-\hat{\beta}_{f}, \ldots, \tilde{f}^{\prime}\left(z_{m}^{*}\right)-\hat{\beta}_{f}\right]^{\prime}\right\| \xrightarrow{d} \chi^{2}(m) \tag{C.6}
\end{equation*}
$$

\]

for a finite $m<\hat{J}$. The latter condition ensures the non-singularity of $\hat{\Upsilon}_{n}$.

## Appendix D. Instruments for volatility

This section describes how to construct the 10 instrumental variables $\mathbf{w}_{i, t}$ to address the endogeneity in individual stock volatility. The work here largely follows Alfaro, Bloom and Lin (2018) and consists of several steps. First, we collect each individual firm's daily stock returns from CRSP. For each firm-day combination in our sample, we regress the firm's stock returns in excess of the risk-free rate on the four asset-pricing factors in Carhart (1997) using a rolling window of the previous 2,520 trading days, if data are available. ${ }^{5}$ The regression residual is referred to as $r_{i, t}^{r i s k-a d j}$, where $i$ denotes the firm and $t$ represents the day.

Second, for each year $t_{o}$ in 2010-2017, we run the following pooled regression for all the firms in a given industry- $j$ using all the days in the recursive rolling window that starts from the year 2010 until the year $t_{o}$.

$$
\begin{equation*}
r_{i, t}^{r i s k-a d j}=\alpha_{j, t_{o}}+\beta_{j, t_{o}}^{(1)} \times r_{t}^{(1)}+\beta_{j, t_{o}}^{(2)} \times r_{t}^{(2)} \ldots+\beta_{j, t_{o}}^{(10)} \times r_{t}^{(10)}+\epsilon_{i, t}, \tag{D.1}
\end{equation*}
$$

where $\alpha_{j, t_{o}}$ and $\beta_{j, t_{o}}^{(c)}$ are the industry-and-year specific intercept and coefficients; $\epsilon_{i, t}$ is the error term; and $r_{t}^{c}$ for $c \in\{1, \ldots, 10\}$ correspond to the 10 different sources of aggregate shocks in Alfaro, Bloom and Lin (2018): (i) When $c$ represents oil prices, $r_{t}^{c}$ is the daily growth rate of the price of the crude oil futures contract CL1 COMB Comdty on Bloomberg; (ii) when $c$ represents U.S. 10-year Treasury, $r_{t}^{c}$ is the daily first difference of the 10-year Treasury yield multiplied by negative $1 ;{ }^{6}$ (iii) when $c$ represents U.S. policy uncertainty, $r_{t}^{c}$ is the daily growth rate of the policy uncertainty index in Baker, Bloom and Davis (2016); and (iv)-(x) when $c$ represents the exchange rates between the U.S. dollars and seven major currencies around the world, $r_{t}^{c}$ is the daily growth rate of the corresponding exchange rate. The seven major currencies are the Australian dollar, British pound, Canadian dollar, Euro, Japanese yen, Swedish krona, and Swiss franc. The daily growth rate of a variable $a$ is calculated as $\left[a_{t}-a_{t-1}\right] /\left[\left(a_{t}+a_{t-1}\right) / 2\right]$, where $t$ is a day. For each estimated coefficient $\beta_{j, t_{o}}^{c}$, we obtain its $t$-statistic $t_{j, t_{o}}^{c}$ and replace the $t$-statistic with 0 if the absolute value of the $t$-statistic is smaller than 1 , indicating that the coefficient is insignificant. ${ }^{7}$ We generate a weighted

[^3]coefficient $\beta_{j, t_{o}}^{c, \text { wighted }}=\left(\left|t_{j, t_{o}}^{c}\right| / \sum_{c}\left|t_{j, t_{o}}^{c}\right|\right) \cdot \beta_{j, t_{o}}^{c}$.
Third, the instruments for $\Delta s_{i, t_{o}}$ are constructed as $w_{i, t_{o}}^{c}=\left|\beta_{j, t_{o}}^{c, w e i g h t e d}\right| \cdot \Delta \sigma_{t_{o}}^{c}$, where firm- $i$ belongs to industry- $j$, and $\sigma_{t_{o}}^{c}$ represents aggregate uncertainty measures calculated as follows. (i) When $c$ represents oil prices, $\sigma_{t_{o}}^{c}$ is the annual average of daily (30-day) volatility of the crude oil futures contract CL1 COMB Comdty (Bloomberg); (ii) when $c$ represents the U.S. 10-year Treasury rate, $\sigma_{t_{o}}^{c}$ is the annual average of daily TYVIX; (iii) when $c$ represents U.S. policy uncertainty, $\sigma_{t_{o}}^{c}$ is the annual average of daily policy uncertainty index in Baker, Bloom and Davis (2016); and (iv)-(x) when $c$ represents the exchange rates between U.S. dollars and seven abovementioned currencies, $\sigma_{t_{o}}^{c}$ is the annual average of daily (three-month) volatility of each exchange rate (Bloomberg CMPN). The first-stage regression results are reported in Table D.1. ${ }^{8}$

Note that the weighting schemes directly follow Alfaro, Bloom and Lin (2018) and aim to reduce noise in the estimation of the industry exposure to different macro uncertainty shocks. Without the weighting schemes, the coefficients in regression (D.1) would represent the industry exposure in a way that a larger coefficient indicates higher exposure. However, a large coefficient can be statistically insignificant if its standard error is proportionally large. The weighting schemes instead weight each coefficient by its $t$-statistics value, where a less significant coefficient (indicated by a smaller $t$ value) receives a smaller weight, and therefore, can result in a better measure of the industry exposure. That said, for our empirical sample, dropping the weighting schemes leads to similar estimation results, shown in Appendix E.

We further follow Alfaro, Bloom and Lin (2018) to construct 10 first-moment control variables to make sure that our results are not driven by movements in oil prices, Treasury yields, U.S. government policies, and exchange rates themselves but rather by the movements in their volatility (see section section IV.C. 3 in the main text). The 10 control variables are constructed by $\beta_{j, t_{o}}^{c, w e i g h t e d} \cdot r_{t_{o}}^{c}$, where $r_{t_{o}}^{c}$ is the first-moment aggregate shock in source $c$ in year $t_{o}$. When $c$ represents oil prices, U.S. 10-year Treasury, or the seven exchange rates, $r_{t_{o}}^{c}$ is the annual average of the corresponding daily $r_{t}^{c}$ in equation (D.1). When $c$ represents U.S. policy uncertainty, $r_{t_{o}}^{c}$ is the annual growth of the U.S. government expenditure as a share of GDP. In addition, in a robustness check (see section IV.C.3), we use out-of-sample data to construct $\mathbf{w}_{i, t}$ to ensure that these instruments are exogenous. That is, we fit regression (D.1) using the individual stock return data between 1998 and 2009 to obtain $\beta_{j, t_{o}}^{c, w e i g h t e d .}$. However, we still use the aggregate data between 2009 and 2017 to calculate $\Delta \sigma_{t_{o}}^{c}$ because
$t$-value larger than 1 (or smaller than -1 ) means that this coefficient is significant in economic terms in the sense that dropping the variable from the regression model can considerably influence the estimation results.
${ }^{8}$ We further add the cash-flow-to-capital ratio and time dummies to the first-stage regression, because they are exogenous variables. The first-stage regression is estimated using pooled OLS; the constant term is included in the estimation but omitted from the table.
aggregate uncertainty shocks are assumed to be exogenous to individual firms' investment plans.

Table D.1: First-stage pooled regression

|  | $(1)$ |
| :--- | :---: |
|  | Volatility changes |
| W oil | $2.644^{* * *}$ |
|  | $(0.206)$ |
| W 10-year Treasury | 28.92 |
|  | $(37.17)$ |
| W policy uncertainty | -26.48 |
|  | $(29.51)$ |
| W AUD | 2.299 |
|  | $(2.906)$ |
| W CAD | 2.000 |
|  | $(2.069)$ |
| W CHF | $6.438^{* * *}$ |
|  | $(1.281)$ |
| W EUR | 3.177 |
|  | $(2.233)$ |
| W GBP | $4.954^{*}$ |
|  | $(2.640)$ |
| W JPY | $4.339^{*}$ |
|  | $(2.553)$ |
| W SEK | $7.791^{* * *}$ |
| CF/K $(z)$ | $(2.155)$ |
|  | 0.00427 |
| Year dummy 2010 | $(0.00496)$ |
| Year dummy 2011 | $-0.205^{* * *}$ |
|  | $(0.00663)$ |
| Year dummy 2012 | $0.0743^{* * *}$ |
| Year dummy 2013 | $(0.00631)$ |
|  | $-0.0484^{* * *}$ |
|  | $(0.00634)$ |
|  | -0.00855 |
|  | $(0.00672)$ |
|  |  |
|  | Continued on the next page |

Table D. 1 - continued from the previous page
(1)

|  | Volatility changes |
| :--- | :---: |
| Year dummy 2014 | $0.0516^{* * *}$ |
|  | $(0.00622)$ |
| Year dummy 2015 | $0.0437^{* * *}$ |
|  | $(0.00742)$ |
| Year dummy 2016 | $0.0629^{* * *}$ |
|  | $(0.00672)$ |
| Number of firms | 1,025 |
| Number of obs | 8,200 |
| R-squared | 0.353 |
| F-test | 248.2 |
| *** p-val $<0.01 ;{ }^{* *}$ p-val $<0.05 ;^{*} \mathrm{p}$-val $<0.1$. |  |

## Appendix E. Robustness results

Figure E. 1 shows the estimates of the first-order derivative, $d \hat{g}(s) / d s$, of the investmentuncertainty relation for the following robustness checks: panel (a) for $k=6$ (footnote 20 in Li and Sun (2022)); (b) for less strict sample selection, and (c) for stricter sample selection (footnote 25 in Li and Sun (2022)); (d)-(l) for robustness outlined in section IV.C. 3 in Li and Sun (2022), (1)-(8). ${ }^{9}$ The results are similar to Figure 1 in Li and Sun (2022).

Figure E. 2 shows results of more robustness checks: panel (a) for using linear specifications for $f_{0}$ and $r_{0} ;{ }^{10}(\mathrm{~b})$ for dropping all the instruments related to time dummies when treating mismeasured $q$; (c) for dropping all the instruments related to cash flow when treating mismeasured $q$; (d) for dropping all the instruments related to volatility when treating mismeasured $q$; (e) for dropping all the instruments related to the control function when treating mismeasured $q ;{ }^{11}$ (f) for using granular instrumental variables (GIV) as in Gabaix

[^4]Figure E.1: Investment-Uncertainty Relation: Robustness

and Koijen (2020); ${ }^{12}(\mathrm{~g})$ for using $W$-variables without the weighting scheme. ${ }^{13}$ The results are similar to Figure 1 in Li and $\operatorname{Sun}$ (2022).

Figure E. 3 shows the comparison of investment-uncertainty relations between low- and high-irreversibility subsamples, where irreversibility is measured by a firm-level scale inflexibility index calculated from quarterly data in our sample period, following Gu, Hackbarth and Johnson (2018). This comparison is similar to Figure 5 in Li and Sun (2022).

Table E. 1 corresponds to the discussion in footnote 29 in Li and Sun (2022) in which we compare three different estimates of the linear model (1). Column (1) includes the results from our baseline estimator but without series approximation. Column (2) shows the results from a modified Erickson, Jiang and Whited (2014) (EJW) estimator, in which we combine their measurement error remedy with the control function approach in our method to simultaneously account for mismeasurement in $q$ and the regressor endogeneity of individual volatility. Column (3) represents the estimates obtained from the original EJW method and their empirical model (note that their model does not include individual volatility as a regressor). The EJW estimates utilize up-to-4th order cumulants following the suggestion in Erickson, Parham and Whited (2017). We observe that the estimates are close to each other in columns (1) and (2), implying that the measurement error remedy in our estimator and that in EJW perform equally well after the regressor endogeneity of uncertainty is treated. Moreover, the coefficient of cash flow is significant in all the three columns. Thus, cash flow is a crucial factor influencing investment. This result is different from those in Erickson and Whited (2000) and Erickson, Jiang and Whited (2014). The comparison here shows that the difference is likely due to different empirical samples rather than estimation methods.

Table E. 2 reports the percentage of firms in an extended sample for which key variables may contain unit roots. The extended sample includes 3,730 individual firms that have at least 3 years of non-missing data between 1986-2017. For each firm, and for each key variable, we conduct the Augmented Dickey-Fuller (ADF) test and the Phillips-Perron (PP) test for unit roots. The null hypothesis of both tests is that the variable contains a unit root, and the alternative is that the variable is stationary. For each variable, we report the percentage of firms for which we cannot reject the null of unit roots under the usual 5 percent statistical level. For example, the number for " y : I/K" under the ADF test is 68.31 , indicating that for 68.31 percent of the 3,730 firms, the variable of investment-to-capital ratio is likely to contain

[^5]Figure E.2: Investment-Uncertainty Relation: More Robustness


Figure E.3: Low Versus High Irreversibility (Measured by Inflexibility)


Table E.1: Linear estimates

|  | $(1)$ <br> Our estimator | $(2)$ <br> Modified EJW | $(3)$ <br> EJW |
| :--- | :---: | :---: | :---: |
| Tobin's $q$ | $0.01^{* * *}$ | $0.01^{*}$ | $0.01^{*}$ |
|  | $(0.002)$ | $(0.004)$ | $(0.003)$ |
| Cash flow | $0.03^{* * *}$ | $0.05^{* * *}$ | $0.05^{* * *}$ |
|  | $(0.008)$ | $(0.014)$ | $(0.013)$ |
| Individual volatility | $-0.08^{* * *}$ | $-0.08^{* * *}$ |  |
|  | $(0.012)$ | $(0.016)$ |  |
| Control function | $-0.03^{* * *}$ | $-0.03^{* * *}$ |  |
|  | $(0.009)$ | $(0.011)$ |  |

Standard errors in parentheses: ${ }^{* * *} p<0.01,{ }^{* *} p<0.05,{ }^{*} p<0.1$

Table E.2: Unit-root tests

|  | $y: \mathrm{I} / \mathrm{K}$ | $x:$ Tobin's $q$ | $z:$ cash flow | $s:$ volatility |
| :--- | :---: | :---: | :---: | :---: |
| ADF | 68.31 | 77.51 | 74.08 | 76.6 |
| PP | 66.06 | 74.64 | 72.04 | 72.49 |

a unit root according to the ADF test. The test results reported here are obtained without a time trend and under the optimal number of lags for each firm, equal to $\operatorname{int}\left\{4(T / 100)^{2 / 9}\right\}$, where int takes the integer value, and $T$ refers to each firm's time-series length (see official documents of Stata 16). The results are also robust to a uniform setting of lags between 1-3 and both with and without a time trend. ${ }^{14}$

## Appendix F: Mathematical Proofs

In this section, $\mathbb{E}_{n}(\cdot)$ denotes the average $n^{-1} \sum_{i=1}^{n}$ over index $i$ and $\mathbb{E}(\cdot)$ denotes the population expectation. Additionally, (i) $\tilde{p}_{l}\left(v_{i, t}\right)=p_{l}\left(v_{i, t}\right)-T^{-1} \sum_{s=1}^{T} p_{l}\left(v_{i, s}\right)$ and $\tilde{p}_{l}^{\prime}\left(v_{i, t}\right)=$ $p_{l}^{\prime}\left(v_{i, t}\right)-T^{-1} \sum_{s=1}^{T} p_{l}^{\prime}\left(v_{i, s}\right)$; (ii) $\left\|P^{k_{n T}}\right\|_{l}=\max _{0 \leq j \leq l} \sup _{x \in R}\left\|d^{j} P^{k_{n T}}(x) / d x^{j}\right\|$ for $l=0,1,2$; (iii) $\|A\|=\left[\operatorname{tr}\left(A A^{\prime}\right)\right]^{1 / 2}$ and $\|A\|_{s p}=\lambda_{\max }^{1 / 2}\left(A A^{\prime}\right)$ denote the Euclidean and spectral norm of matrix $A$, respectively. In addition, to fit our mathematical equations with page margin, we introduce the symbols $\boldsymbol{\uparrow}, \boldsymbol{\leftrightarrow}, \boldsymbol{\star}, \boldsymbol{\nabla}$, and $\nabla$; they may have different meanings in different places.

Lemma F. 1 Under Assumptions $1-4$ and denoting

$$
\delta_{n T}=(n T)^{-1 / 2}\left\|P^{k_{n T}}\right\|_{1}
$$

we have

$$
\begin{align*}
\left\|\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{i}^{\prime} \tilde{\hat{\mathbf{d}}}_{i}-\widetilde{\mathbf{q}}_{i}^{\prime} \tilde{\mathbf{d}}_{i}\right)\right\| & =O_{p}\left(\delta_{n T} k_{n T}^{1 / 2}\right)  \tag{F.1}\\
\left\|\mathbb{E}_{n}\left[\left(\widetilde{\hat{\mathbf{d}}}_{i}-\tilde{\mathbf{d}}_{i}\right)^{\prime} \tilde{\Delta}_{i}\right]\right\| & =O_{p}\left(\delta_{n T} k_{n T}^{-\zeta}\right)  \tag{F.2}\\
\left\|\mathbb{E}_{n}\left[\left(\widetilde{\hat{\mathbf{d}}}_{i}-\tilde{\mathbf{d}}_{i}\right)^{\prime} \tilde{\varepsilon}_{i}\right]\right\| & =O_{p}\left(\delta_{n T} / \sqrt{n}\right) \tag{F.3}
\end{align*}
$$

where $\tilde{\boldsymbol{\Delta}}_{i}=J_{T} \boldsymbol{\Delta}_{i}, \boldsymbol{\Delta}_{i}=\left[\Delta_{i, 1}, \ldots, \boldsymbol{\Delta}_{i, T}\right]^{\prime}$ and $\Delta_{i, t}=f_{0}\left(z_{i, t-1}\right)+g_{0}\left(s_{i, t-1}\right)+r_{0}\left(v_{i, t}\right)-E_{i, t}^{\prime} \boldsymbol{\vartheta}_{0}$.

[^6]Proof of Lemma F.1: The matrix form of model (6) in the main text is given by

$$
\Delta \mathbf{s}=\mathbf{i}_{n T} \eta_{0}+\mathbf{W} \pi_{0}+\mathbf{v}
$$

where $\mathbf{s}=\left[\mathbf{s}_{1}^{\prime}, \ldots, \mathbf{s}_{n}^{\prime}\right]^{\prime}, \mathbf{v}=\left[\mathbf{v}_{1}^{\prime}, \ldots, \mathbf{v}_{n}^{\prime}\right]^{\prime}$, and $\mathbf{W}=\left[\mathbf{W}_{1}, \ldots, \mathbf{W}_{n}\right]^{\prime}$ is an $(n T) \times d_{w}$ matrix. Applying the least squares estimation gives $\left\|\hat{\pi}-\pi_{0}\right\|=O_{p}\left((n T)^{-1 / 2}\right)$ and $\hat{\eta}-\eta_{0}=$ $O_{p}\left((n T)^{-1 / 2}\right)$ under Assumptions 1(i) and 3(i).

First, we verify (F.1). It is readily seen that

$$
\begin{align*}
& \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{i}^{\prime} \widetilde{\hat{\mathbf{d}}}_{i}-\widetilde{\mathbf{q}}_{i}^{\prime} \tilde{\mathbf{d}}_{i}\right)  \tag{F.4}\\
= & \mathbb{E}_{n}\left[\left(\widetilde{\hat{\mathbf{q}}}_{i}-\widetilde{\mathbf{q}}_{i}\right)^{\prime}\left(\widetilde{\hat{\mathbf{d}}}_{i}-\tilde{\mathbf{d}}_{i}\right)\right]+\mathbb{E}_{n}\left[\left(\widetilde{\hat{\mathbf{q}}}_{i}-\widetilde{\mathbf{q}}_{i}\right)^{\prime} \tilde{\mathbf{d}}_{i}\right]+\mathbb{E}_{n}\left[\widetilde{\mathbf{q}}_{i}^{\prime}\left(\widetilde{\hat{\mathbf{d}}}_{i}-\tilde{\mathbf{d}}_{i}\right)\right] \\
= & \mathbb{E}_{n}\left\{\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n} T}\right)^{\prime} J_{T}\left[I_{T} \otimes \operatorname{vec}\left(J_{T}\left(P_{\hat{v}, i}^{k_{n} T}-P_{v, i}^{k_{n, i}}\right)\right)^{\prime}\right]\right\} \\
+ & \mathbb{E}_{n}\left[\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n} T}\right)^{\prime} J_{T} \tilde{\mathbf{d}}_{i}\right]+\mathbb{E}_{n}\left\{\widetilde{\mathbf{q}}_{i}^{\prime}\left[I_{T} \otimes \operatorname{vec}\left(J_{T}\left(P_{\hat{v}, i}^{k_{n} T}-P_{v, i}^{k_{n T}}\right)\right)^{\prime}\right]\right\},
\end{align*}
$$

where for the first term, we have

$$
\begin{align*}
& \left\|\mathbb{E}_{n}\left\{\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)^{\prime} J_{T}\left[I_{T} \otimes \operatorname{vec}\left(J_{T}\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)\right)^{\prime}\right]\right\}\right\| \\
& \leq \mathbb{E}_{n}\left(\left\|\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)^{\prime} J_{T}\left[I_{T} \otimes \operatorname{vec}\left(J_{T}\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)\right)^{\prime}\right]\right\|\right) \\
& =\mathbb{E}_{n}\left\{\operatorname{tr}\left[\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)^{\prime} J_{T}\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)\right]\right\} \\
& =\mathbb{E}_{n}\left(\sum_{l=1}^{k_{n T}} \sum_{t=1}^{T}\left[\tilde{p}_{l}\left(\hat{v}_{i, t}\right)-\tilde{p}_{l}\left(v_{i, t}\right)\right]^{2}\right)=O_{p}\left(\delta_{n T}^{2}\right) \tag{F.5}
\end{align*}
$$

under Assumptions 1(i) and 2(ii) because $\hat{v}_{i, t}-v_{i, t}=\eta_{0}-\hat{\eta}+\mathbf{w}_{i, t}^{\prime}\left(\boldsymbol{\pi}_{0}-\hat{\boldsymbol{\pi}}\right)$, and applying

Taylor expansion yields

$$
\begin{aligned}
& \tilde{p}_{l}\left(\hat{v}_{i, t}\right)-\tilde{p}_{l}\left(v_{i, t}\right)=p_{l}^{\prime}\left(\stackrel{\circ}{i}_{i, t}\right)\left(\hat{v}_{i, t}-v_{i, t}\right)-T^{-1} \sum_{s=1}^{T} p_{l}^{\prime}\left(\stackrel{\circ}{i}_{i, s}\right)\left(\hat{v}_{i, s}-v_{i, s}\right) \\
& =p_{l}^{\prime}(\stackrel{\circ}{i, t})\left[\eta_{0}-\hat{\eta}+\mathbf{w}_{i, t}^{\prime}\left(\boldsymbol{\pi}_{0}-\hat{\boldsymbol{\pi}}\right)\right] \\
& -T^{-1} \sum_{s=1}^{T} p_{l}^{\prime}\left(\stackrel{\circ}{i}_{i, s}\right)\left(\eta_{0}-\hat{\eta}+\mathbf{w}_{i, s}^{\prime}\left(\boldsymbol{\pi}_{0}-\hat{\boldsymbol{\pi}}\right)\right) \\
& =\tilde{p}_{l}^{\prime}\left(\dot{v}_{i, t}\right)\left(\eta_{0}-\hat{\eta}\right)-\left(\boldsymbol{\pi}_{0}-\hat{\boldsymbol{\pi}}\right)^{\prime}\left[\mathbf{w}_{i, t} p_{l}^{\prime}\left(\stackrel{( }{v}_{i, t}\right)-T^{-1} \sum_{s=1}^{T} \mathbf{w}_{i, s} p_{l}^{\prime}\left(\dot{v}_{i, s}\right)\right]
\end{aligned}
$$

where $\stackrel{\circ}{v}_{i, t}$ lies between $\hat{v}_{i, t}$ and $v_{i, t}$. Additionally, we have

$$
\begin{aligned}
& \left\|\mathbb{E}_{n}\left[\left(P_{\hat{v}, i}^{k_{n} T}-P_{v, i}^{k_{n} T}\right)^{\prime} J_{T} \tilde{\mathbf{d}}_{i}\right]\right\| \leq \mathbb{E}_{n}\left\|\left(P_{\hat{v}, i}^{k_{n} T}-P_{v, i}^{k_{n} T}\right)^{\prime} J_{T} \tilde{\mathbf{d}}_{i}\right\| \\
& \leq 2 \mathbb{E}_{n}\left|\operatorname{tr}\left[\left(P_{\hat{v}, i}^{k_{n} T}-P_{v, i}^{k_{n} T}\right)^{\prime} J_{T}\left(P_{\hat{v}, i}^{k_{n} T}-P_{v, i}^{k_{n} T}\right)\right]\left[\operatorname{tr}\left(\tilde{\mathbf{E}}_{i}^{\prime} \tilde{\mathbf{E}}_{i}\right)+\operatorname{tr}\left(F_{T}^{\prime} F_{T}\right)\right]\right|^{1 / 2} \\
& +2 \mathbb{E}_{n}\left|\operatorname{tr}\left(\left(P_{\hat{v}, i}^{k_{n} T}-P_{v, i}^{k_{n} T}\right)^{\prime} J_{T}\left(I_{T} \otimes \tilde{\mathbf{y}}_{i}\right)^{\prime} \boldsymbol{\phi}\left(I_{T} \otimes \tilde{\mathbf{y}}_{i}\right) J_{T}\left(P_{\hat{v}, i}^{k_{n} T}-P_{v, i}^{k_{n} T}\right)\right)\right|^{1 / 2} \\
& =O_{p}\left(\delta_{n T} k_{n T}^{1 / 2}\right)+O_{p}\left(\delta_{n T}\right)
\end{aligned}
$$

under Assumptions 1(i), 2(ii) and 3(i), where we obtain the last line using Hölder's inequality, (F.5), Lemma F.2(i), $\operatorname{tr}(A B) \leq \operatorname{tr}(A) \operatorname{tr}(B)$ for non-negative definite matrices $A$ and $B$ and $\boldsymbol{\phi}=\left[\left(J_{T} \otimes J_{T}\right) D_{T}\right]_{\perp}\left[\left(J_{T} \otimes J_{T}\right) D_{T}\right]_{\perp}^{\prime}, \operatorname{tr}[\boldsymbol{Q}]=m_{0}$. Similarly, we obtain

$$
\begin{aligned}
& \left\|\mathbb{E}_{n}\left\{\widetilde{\mathbf{q}}_{i}^{\prime}\left[I_{T} \otimes \operatorname{vec}\left(J_{T}\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)\right)^{\prime}\right]\right\}\right\| \\
& \leq \mathbb{E}_{n}\left\|\widetilde{\mathbf{q}}_{i}^{\prime}\left[I_{T} \otimes \operatorname{vec}\left(J_{T}\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)\right)^{\prime}\right]\right\| \\
& =\mathbb{E}_{n}\left|\operatorname{tr}\left[\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)^{\prime} J_{T}\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)\right] \operatorname{tr}\left(\widetilde{\mathbf{q}}_{i}^{\prime} \widetilde{\mathbf{q}}_{i}\right)\right|^{1 / 2} \\
& =O_{p}\left(\delta_{n T} k_{n T}^{1 / 2}\right) .
\end{aligned}
$$

Hence, under Assumptions 1-4, taking together the results above gives (i).

Second, applying a similar proof method used above and by (F.5), we have

$$
\begin{aligned}
& \left\|\mathbb{E}_{n}\left[\left(\tilde{\hat{\mathbf{d}}}_{i}-\tilde{\mathbf{d}}_{i}\right)^{\prime} \tilde{\Delta}_{i}\right]\right\|=\left\|\mathbb{E}_{n}\left\{\left[I_{T} \otimes \operatorname{vec}\left(J_{T}\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)\right)\right] \tilde{\Delta}_{i}\right\}\right\| \\
& \leq \sqrt{T \mathbb{E}_{n}}\left(\tilde{\Delta}_{i}^{\prime} \tilde{\Delta}_{i} \operatorname{tr}\left[\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)^{\prime} J_{T}\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)\right]\right)^{1 / 2}=O_{p}\left(\delta_{n T} k_{n T}^{-\zeta}\right)
\end{aligned}
$$

where $\max _{1 \leq i \leq n}\left|\tilde{\Delta}_{i}^{\prime} \tilde{\Delta}_{i}\right| \leq M k_{n T}^{-2 \zeta}$ under Assumption 2(i), and

$$
\begin{aligned}
& \mathbb{E}\left\|\mathbb{E}_{n}\left[\left(\widetilde{\hat{\mathbf{d}}}_{i}-\tilde{\mathbf{d}}_{i}\right)^{\prime} \tilde{\varepsilon}_{i}\right]\right\|^{2} \\
& =\mathbb{E}\left\|\mathbb{E}_{n}\left\{\left[I_{T} \otimes \operatorname{vec}\left(J_{T}\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n}}\right)\right)\right] \tilde{\varepsilon}_{i}\right\}\right\|^{2} \\
& =n^{-1} \mathbb{E}\left\{\left[I_{T} \otimes \operatorname{vec}\left(J_{T}\left(P_{\hat{v}, i}^{k_{n}}-P_{v, i}^{k_{n T}}\right)\right)\right] \tilde{\varepsilon}_{i} \tilde{\varepsilon}_{i}^{\prime}\left[I_{T} \otimes \operatorname{vec}\left(J_{T}\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)\right)\right]\right\} \\
& \leq n^{-1} \mathbb{E}\left\{\tilde{\varepsilon}_{i}^{\prime} \tilde{\varepsilon}_{i} \operatorname{tr}\left[\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n} T}\right)^{\prime} J_{T}\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)\right]\right\}=O\left(\delta_{n T}^{2} / n\right) .
\end{aligned}
$$

This completes the proof of this lemma.

Lemma F. 2 Under Assumptions 1-4, we have

$$
\begin{aligned}
\left\|\mathbb{E}_{n}\left(\tilde{\mathbf{d}}_{i}^{\prime} \widetilde{\mathbf{q}}_{i}\right)-\mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \widetilde{\mathbf{q}}_{i}\right)\right\| & =O_{p}\left(k_{n T} / \sqrt{n}\right) \\
\left\|\mathbb{E}_{n}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\boldsymbol{\Delta}}_{i}\right)-\mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\Delta}_{i}\right)\right\| & =O_{p}\left(k_{n T}^{1 / 2-\zeta} / \sqrt{n}\right) .
\end{aligned}
$$

Proof of Lemma F.2: First, we consider

$$
\begin{aligned}
& \mathbb{E}\left\|\mathbb{E}_{n}\left(\tilde{\mathbf{d}}_{i}^{\prime} \widetilde{\mathbf{q}}_{i}\right)-\mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \widetilde{\mathbf{q}}_{i}\right)\right\|^{2} \\
& =n^{-1} \mathbb{E}\left\{\operatorname{tr}\left[\left(\tilde{\mathbf{d}}_{i}^{\prime} \widetilde{\mathbf{q}}_{i}-\mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \widetilde{\mathbf{q}}_{i}\right)\right)^{\prime}\left(\tilde{\mathbf{d}}_{i}^{\prime} \widetilde{\mathbf{q}}_{i}-\mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \widetilde{\mathbf{q}}_{i}\right)\right)\right]\right\} \\
& \leq n^{-1} \mathbb{E}\left[\operatorname{tr}\left(\tilde{\mathbf{d}}_{i}^{\prime} \widetilde{\mathbf{q}}_{i} \widetilde{\mathbf{q}}_{i}^{\prime} \tilde{\mathbf{d}}_{i}\right)\right] \leq n^{-1} \mathbb{E}\left[\operatorname{tr}\left(\widetilde{\mathbf{q}}_{i} \widetilde{\mathbf{q}}_{i}^{\prime}\right) \operatorname{tr}\left(\tilde{\mathbf{d}}_{i} \tilde{\mathbf{d}}_{i}^{\prime}\right)\right] \\
& =n^{-1} \mathbb{E}\left[\operatorname{tr}\left(F_{T} F_{T}^{\prime}+\tilde{\mathbf{x}}_{i,-1} \tilde{\mathbf{x}}_{i,-1}^{\prime}+\tilde{\mathbf{E}}_{i} \tilde{\mathbf{E}}_{i}^{\prime}\right)\right. \\
& \left.\times \operatorname{tr}\left(\left(I_{T} \otimes \tilde{\mathbf{y}}_{i}\right)^{\prime} \boldsymbol{\mathbf { A }}\left(I_{T} \otimes \tilde{\mathbf{y}}_{i}\right)+\operatorname{tr}\left(\tilde{\mathbf{E}}_{i} \tilde{\mathbf{E}}_{i}^{\prime}+F_{T} F_{T}^{\prime}\right) I_{T}\right)\right] \\
& =n^{-1} \mathbb{E}\left[\left(2(T-1)+\tilde{\mathbf{x}}_{i,-1}^{\prime} \tilde{\mathbf{x}}_{i,-1}+\operatorname{tr}\left(\tilde{\mathbf{E}}_{i} \tilde{\mathbf{E}}_{i}^{\prime}\right)\right)\right. \\
& \left.\times \operatorname{tr}\left(\left(I_{T} \otimes \tilde{\mathbf{y}}_{i}\right)^{\prime} \mathbf{\phi}\left(I_{T} \otimes \tilde{\mathbf{y}}_{i}\right)+\left(\operatorname{tr}\left(\tilde{\mathbf{E}}_{i} \tilde{\mathbf{E}}_{i}^{\prime}\right)+2(T-1)\right) T\right)\right] \\
& \leq n^{-1} T \mathbb{E}\left[\left(2(T-1)+\tilde{\mathbf{x}}_{i,-1}^{\prime} \tilde{\mathbf{x}}_{i,-1}+\operatorname{tr}\left(\tilde{\mathbf{E}}_{i} \tilde{\mathbf{E}}_{i}^{\prime}\right)\right)\left(\tilde{\mathbf{y}}_{i}^{\prime} \tilde{\mathbf{y}}_{i}+\operatorname{tr}\left(\tilde{\mathbf{E}}_{i} \tilde{\mathbf{E}}_{i}^{\prime}\right)+2(T-1)\right)\right] \\
& =O\left(n^{-1} k_{n T}^{2}\right),
\end{aligned}
$$

where $\boldsymbol{\uparrow}=\left[\left(J_{T} \otimes J_{T}\right) D_{T}\right]_{\perp}\left[\left(J_{T} \otimes J_{T}\right) D_{T}\right]_{\perp}^{\prime}$. Second, we obtain

$$
\begin{aligned}
& \mathbb{E}\left\|\mathbb{E}_{n}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\Delta}_{i}\right)-\mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\Delta}_{i}\right)\right\|^{2} \\
& =n^{-1} \mathbb{E}\left\{\operatorname{tr}\left[\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\Delta}_{i}-\mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\Delta}_{i}\right)\right)^{\prime}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\Delta}_{i}-\mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\Delta}_{i}\right)\right)\right]\right\} \\
& \leq n^{-1} \mathbb{E}\left[\operatorname{tr}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\Delta}_{i} \tilde{\Delta}_{i}^{\prime} \tilde{\mathbf{d}}_{i}\right)\right] \leq n^{-1} \mathbb{E}\left[\tilde{\Delta}_{i}^{\prime} \tilde{\Delta}_{i} \operatorname{tr}\left(\tilde{\mathbf{d}}_{i} \tilde{\mathbf{d}}_{i}^{\prime}\right)\right]=O\left(n^{-1} k_{n T}^{1-2 \zeta}\right)
\end{aligned}
$$

under Assumptions 1-2. This completes the proof of this lemma.
Lemma F. 3 Under Assumptions 1-4, we have

$$
\begin{aligned}
\left\|\Lambda_{n}-\Lambda\right\|_{s p} & =O_{p}\left(k_{n T} / \sqrt{n}+\delta_{n T} J^{1 / 2}\right) \\
\left\|A_{n}-A\right\| & =O_{p}\left(k_{n T}^{-\zeta}\left(k_{n T} / \sqrt{n}+\delta_{n T} J^{1 / 2}\right)\right) \\
\left\|\Lambda_{n}^{-1}-\Lambda^{-1}\right\|_{s p} & =O_{p}\left(k_{n T} / \sqrt{n}+\delta_{n T} J^{1 / 2}\right)
\end{aligned}
$$

where we denote

$$
\begin{align*}
& \Lambda_{n}=\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{1, i}^{\prime} \widetilde{\hat{\mathbf{d}}}_{i}\right) \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\hat{\mathbf{q}}}_{1, i}\right), \Lambda=\mathbb{E}\left(\tilde{\mathbf{q}}_{1, i}^{\prime} \tilde{\mathbf{d}}_{i}\right) \Omega_{n} \mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\mathbf{q}}_{1, i}\right)  \tag{F.6}\\
& A_{n}=\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{1, i}^{\prime} \widetilde{\hat{\mathbf{d}}}_{i}\right) \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\Delta}_{i}\right), A=\mathbb{E}\left(\tilde{\mathbf{q}}_{1, i}^{\prime} \tilde{\mathbf{d}}_{i}\right) \Omega_{n} \mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\Delta}_{i}\right) . \tag{F.7}
\end{align*}
$$

Proof of Lemma F.3: For any comparable matrices $A, B, C$, and $D$, we have $A \Omega_{n} C^{\prime}-$ $B \Omega_{n} D^{\prime}=(A-B) \Omega_{n}(C-D)^{\prime}+(A-B) \Omega_{n} D^{\prime}+B \Omega_{n}(C-D)^{\prime}$, so that

$$
\begin{align*}
\left\|A \Omega_{n} C^{\prime}-B \Omega_{n} D^{\prime}\right\|_{s p} & \leq\|A-B\|_{s p}\left\|\Omega_{n}\right\|_{s p}\|C-D\|_{s p}+\|A-B\|_{s p}\left\|\Omega_{n}\right\|_{s p}\|D\|_{s p} \\
& +\|B\|_{s p}\left\|\Omega_{n}\right\|_{s p}\|C-D\|_{s p} \tag{F.8}
\end{align*}
$$

Hence, by Lemmas F. 1 and F. 2 and under Assumption 3, we have

$$
\begin{aligned}
& \left\|\Lambda_{n}-\Lambda\right\|_{s p} \\
& \leq\left\|\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{1, i}^{\prime} \widetilde{\hat{\mathbf{d}}}_{i}\right)-\mathbb{E}\left(\tilde{\mathbf{q}}_{1, i}^{\prime} \tilde{\mathbf{d}}_{i}\right)\right\|_{s p}^{2}\left\|\Omega_{n}\right\|_{s p} \\
& +2\left\|\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{1, i}^{\prime} \widetilde{\hat{\mathbf{d}}}_{i}\right)-\mathbb{E}\left(\tilde{\mathbf{q}}_{1, i}^{\prime} \tilde{\mathbf{d}}_{i}\right)\right\|_{s p}\left\|\Omega_{n}^{1 / 2}\right\|_{s p}\left\|\Omega_{n}^{1 / 2} \mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\mathbf{q}}_{1, i}\right)\right\|_{s p} \\
& =O_{p}\left(k_{n T} / \sqrt{n}+\delta_{n T} J^{1 / 2}\right)
\end{aligned}
$$

and

$$
\begin{aligned}
& \left\|A_{n}-A\right\|_{s p} \\
& \leq\left\|\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{1, i}^{\prime} \widetilde{\hat{\mathbf{d}}}_{i}\right)-\mathbb{E}\left(\tilde{\mathbf{q}}_{1, i}^{\prime} \tilde{\mathbf{d}}_{i}\right)\right\|_{s p}\left\|\Omega_{n}\right\|_{s p}\left\|\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\Delta}_{i}\right)-\mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\Delta}_{i}\right)\right\|_{s p} \\
& +\left\|\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{1, i}^{\prime} \widetilde{\hat{\mathbf{d}}}_{i}\right)-\mathbb{E}\left(\tilde{\mathbf{q}}_{1, i}^{\prime} \tilde{\mathbf{d}}_{i}\right)\right\|_{s p}\left\|\Omega_{n}\right\|_{s p}\left\|\mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\boldsymbol{\Delta}}_{i}\right)\right\|_{s p} \\
& +\left\|\mathbb{E}\left(\tilde{\mathbf{q}}_{1, i}^{\prime} \tilde{\mathbf{d}}_{i}\right) \Omega_{n}^{1 / 2}\right\|_{s p}\left\|\Omega_{n}^{1 / 2}\right\|_{s p}\left\|\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\Delta}_{i}\right)-\mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\Delta}_{i}\right)\right\|_{s p} \\
& =O_{p}\left(k_{n T}^{-\zeta}\left(k_{n T} / \sqrt{n}+\delta_{n T} J^{1 / 2}\right)\right) .
\end{aligned}
$$

Next, we have

$$
\begin{aligned}
\left\|\Lambda_{n}^{-1}-\Lambda^{-1}\right\|_{s p} & =\left\|\Lambda_{n}^{-1}\left(\Lambda_{n}-\Lambda\right) \Lambda^{-1}\right\|_{s p} \\
& \leq\left\|\Lambda_{n}^{-1}\right\|_{s p}\left\|\Lambda_{n}-\Lambda\right\|_{s p}\left\|\Lambda^{-1}\right\|_{s p}=O_{p}\left(k_{n T} / \sqrt{n}+\delta_{n T} J^{1 / 2}\right)
\end{aligned}
$$

under Assumption 3(ii), where $\left\|\Lambda_{n}^{-1}\right\|_{s p}=\left\|\Lambda^{-1}\right\|_{s p}+O\left(\left\|\Lambda_{n}-\Lambda\right\|_{s p}\right)$. This completes the proof of this lemma.

Proof of Theorem B.1: We first verify (i). Applying simple algebra yields

$$
\begin{aligned}
\check{\boldsymbol{\theta}}_{1} & =\Lambda_{n}^{-1} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{1, i}^{\prime} \tilde{\hat{\mathbf{d}}}_{i}\right) \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\mathbf{y}}_{i}\right) \\
& =\boldsymbol{\theta}_{1,0}+\Lambda_{n}^{-1} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{1, i}^{\prime} \widetilde{\hat{\mathbf{d}}}_{i}\right) \Omega_{n}\left[\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\boldsymbol{\Delta}}_{i}\right)+\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\varepsilon}_{i}\right)\right] \\
& =\boldsymbol{\theta}_{1,0}+\Lambda_{n}^{-1} A_{n}+\Lambda_{n}^{-1} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{1, i}^{\prime} \widetilde{\hat{\mathbf{d}}}_{i}\right) \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\varepsilon}_{i}\right) \\
& =\boldsymbol{\theta}_{1,0}+\Lambda^{-1} A+\left(\Lambda_{n}^{-1} A_{n}-\Lambda^{-1} A\right)+\Lambda_{n}^{-1} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{1, i}^{\prime} \widetilde{\hat{\mathbf{d}}}_{i}\right) \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\varepsilon}_{i}\right),
\end{aligned}
$$

where $\|b\|=\|b\|_{s p}$ for any vector $b$, and by Lemmas F.1-F. 3 we obtain

$$
\begin{aligned}
& \left\|\Lambda_{n}^{-1} A_{n}-\Lambda^{-1} A\right\| \\
& \leq\left\|\left(\Lambda_{n}^{-1}-\Lambda^{-1}\right)\left(A_{n}-A\right)\right\|+\left\|\left(\Lambda_{n}^{-1}-\Lambda^{-1}\right) A\right\|+\left\|\Lambda^{-1}\left(A_{n}-A\right)\right\| \\
& =O_{p}\left(\delta_{n T} k_{n T}^{1 / 2-\zeta}\right)
\end{aligned}
$$

and

$$
\begin{align*}
& \left\|\Lambda_{n}^{-1} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{1, i}^{\prime} \widetilde{\hat{\mathbf{d}}}_{i}\right) \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\varepsilon}_{i}\right)\right\| \\
& \leq\left\|\left(\Lambda_{n}^{-1}-\Lambda^{-1}\right) \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{1, i}^{\prime} \widetilde{\hat{\mathbf{d}}}_{i}\right) \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\varepsilon}_{i}\right)\right\|+\left\|\Lambda^{-1} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{1, i}^{\prime} \widetilde{\hat{\mathbf{d}}}_{i}\right) \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\varepsilon}_{i}\right)\right\| \\
& \leq\left(\left\|\Lambda_{n}^{-1}-\Lambda^{-1}\right\|_{s p}+\left\|\Lambda^{-1}\right\|_{s p}\right)\left\|\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{1, i}^{\prime} \widetilde{\hat{\mathbf{d}}}_{i}\right) \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\varepsilon}_{i}\right)\right\|_{s p} \\
& =O_{p}(1)\left(\|\left(\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{1, i}^{\prime} \widetilde{\hat{\mathbf{d}}}_{i}\right)-\mathbb{E}\left(\tilde{\mathbf{q}}_{1, i}^{\prime} \tilde{\mathbf{d}}_{i}\right)\right) \Omega_{n} \mathbb{E}_{n}\left({\left.\left.\widetilde{\hat{\mathbf{d}}_{i}^{\prime}} \tilde{\varepsilon}_{i}\right)\left\|_{s p}+\right\| \boldsymbol{\|} \|_{s p}\right)}_{=O_{p}\left(\left(\delta_{n T}+\sqrt{J}\right) / \sqrt{n}\right)}\right.\right.
\end{align*}
$$

where $=\mathbb{E}\left(\tilde{\mathbf{q}}_{1, i}^{\prime} \tilde{\mathbf{d}}_{i}\right) \Omega_{n} \mathbb{E}_{n}\left(\tilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\varepsilon}_{i}\right)$, because $\mathbb{E}\left(\left\|\mathbb{E}\left(\tilde{\mathbf{q}}_{1, i}^{\prime} \tilde{\mathbf{d}}_{i}\right) \Omega_{n} \mathbb{E}_{n}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\varepsilon}_{i}\right)\right\|^{2}\right)=n^{-1} \mathbb{E}\left[\tilde{\boldsymbol{\varepsilon}}_{i}^{\prime} \tilde{\varepsilon}_{i} \operatorname{tr}\right.$
$=O(J / n) ;{ }^{15} \boldsymbol{\&}=\mathbb{E}\left(\tilde{\mathbf{q}}_{1, i}^{\prime} \tilde{\mathbf{d}}_{i}\right) \Omega_{n} \tilde{\mathbf{d}}_{i}^{\prime} \tilde{\mathbf{d}}_{i} \Omega_{n} \mathbb{E}\left(\tilde{\mathbf{q}}_{1, i}^{\prime} \tilde{\mathbf{d}}_{i}\right)$. Hence, under Assumption 4, we have $\left\|\check{\boldsymbol{\theta}}_{1}-\boldsymbol{\theta}_{1,0}\right\|=O_{p}\left(k_{n T}^{-\zeta}+\sqrt{J / n}\right)$ because the bias term, $\Lambda^{-1} A$, is of order $O\left(k_{n T}^{-\zeta}\right)$ under

[^7]Assumptions 2(i) and 3. Next, since

$$
\begin{aligned}
& {\left[\begin{array}{l}
\check{f}(z)-f_{0}(z) \\
\check{g}(s)-g_{0}(s) \\
\check{r}(v)-r_{0}(v)
\end{array}\right]=\left[\begin{array}{c}
\check{f}(z)-\vartheta_{z, 1,0}^{\prime} \vec{P}^{k_{n T}}(z) \\
\check{g}(s)-\vartheta_{s, 1,0}^{\prime} \vec{P}_{n T}(s) \\
\check{r}(v)-\vartheta_{v, 1,0}^{\prime} \vec{P}^{k_{n T}}(v)
\end{array}\right]+\left[\begin{array}{c}
\vartheta_{z, 1,0}^{\prime} \vec{P}^{k_{n T}}(z)-f_{0}(z) \\
\vartheta_{s, 1,0}^{\prime} \vec{P}^{k_{n T}}(s)-g_{0}(s) \\
\vartheta_{v, 1,0}^{\prime} \vec{P}^{k_{n T}}(v)-r_{0}(v)
\end{array}\right]} \\
& =\left[\begin{array}{ccc}
\vec{P}^{k_{n T}}(z)^{\prime} & 0 & 0 \\
0 & \vec{P}^{k_{n T}}(s)^{\prime} & 0 \\
0 & 0 & \vec{P}^{k_{n T}}(v)^{\prime}
\end{array}\right]\left[\begin{array}{c}
\check{\boldsymbol{\vartheta}}_{z, 1}-\vartheta_{z, 1,0} \\
\check{\boldsymbol{\vartheta}}_{s, 1}-\vartheta_{s, 1,0} \\
\check{\boldsymbol{\vartheta}}_{v, 1}-\vartheta_{v, 1,0}
\end{array}\right]+O_{p}\left(k_{n T}^{-\zeta}\right)
\end{aligned}
$$

under Assumption 2(i); then, under Assumption 2(ii), we obtain result (ii) of this theorem, where the subscript " 1 " has the same meaning as that in $\check{\boldsymbol{\theta}}_{1}$ and $\boldsymbol{\theta}_{1,0}$.

Under Assumption 4 and by the proof above, we have

$$
\left[\begin{array}{c}
\check{\alpha}-\alpha_{0}  \tag{F.10}\\
\check{\boldsymbol{\vartheta}}_{z, 1}-\vartheta_{z, 1,0} \\
\check{\boldsymbol{\vartheta}}_{s, 1}-\vartheta_{s, 1,0} \\
\check{\boldsymbol{\vartheta}}_{v, 1}-\vartheta_{v, 1,0}
\end{array}\right] \approx\left[\mathbb{E}\left(\tilde{\mathbf{q}}_{1, i}^{\prime} \tilde{\mathbf{d}}_{i}\right) \Omega_{n} \mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\mathbf{q}}_{1, i}\right)\right]^{-1} \mathbb{E}\left(\tilde{\mathbf{q}}_{1, i}^{\prime} \tilde{\mathbf{d}}_{i}\right) \Omega_{n} \mathbb{E}_{n}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\varepsilon}_{i}\right),
$$

where for any $p_{1} \times m_{i v}$ selection matrix $S$ with finite $p_{1}$, applying the multivariate central limit theorem yields

$$
\begin{equation*}
n^{-1 / 2} S \sum_{i=1}^{n} \tilde{\mathbf{d}}_{i}^{\prime} \tilde{\varepsilon}_{i} \xrightarrow{d} N\left(0_{p_{1}}, S \mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\varepsilon}_{i} \tilde{\varepsilon}_{i}^{\prime} \tilde{\mathbf{d}}_{i}\right) S\right) \tag{F.11}
\end{equation*}
$$

Denoting a $J \times J$ matrix

$$
\begin{equation*}
\hat{\Sigma}_{n}=\Lambda_{n}^{-1} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{q}}}_{1, i}^{\prime} \tilde{\hat{\mathbf{d}}}_{i}\right) \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widehat{\tilde{\varepsilon}}_{i} \widehat{\tilde{\varepsilon}}_{i}^{\prime} \tilde{\mathbf{d}}_{i}\right) \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widetilde{\hat{\mathbf{q}}}_{1, i}^{\prime}\right) \Lambda_{n}^{-1} \tag{F.12}
\end{equation*}
$$

and applying the delta method yields

$$
\sqrt{n / \hat{\sigma}_{\alpha}^{2}}\left(\check{\alpha}-\alpha_{0}\right) \xrightarrow{d} N(0,1)
$$

where we have

$$
\begin{equation*}
\hat{\sigma}_{\alpha}^{2}=e_{1, J}^{\prime} \hat{\Sigma}_{n} e_{1, J} \xrightarrow{p} \sigma_{\alpha}^{2}, \tag{F.13}
\end{equation*}
$$

and $\sigma_{\alpha}^{2}$ equals $\hat{\sigma}_{\alpha}^{2}$ with $\mathbb{E}_{n}, \tilde{\hat{\mathbf{q}}}_{1, i}, \tilde{\hat{\mathbf{d}}}_{i}$ and $\widehat{\tilde{\varepsilon}}_{i}$ replaced by $\mathbb{E}, \tilde{\mathbf{q}}_{1, i}, \tilde{\mathbf{d}}_{i}$ and $\tilde{\varepsilon}_{i}$, respectively, and $e_{1, J}$ is the first column of $I_{J}$. In addition, we have

$$
\sqrt{n \hat{\Xi}_{n}^{-1}}\left[\begin{array}{c}
\check{f}(z)-f_{0}(z) \\
\check{g}(s)-g_{0}(s) \\
\check{r}(v)-r_{0}(v)
\end{array}\right] \xrightarrow{d} N\left(0_{3}, I_{3}\right)
$$

where, letting $S_{0}$ equal the last $J-1$ rows of the identity matrix $I_{J}$, we define

$$
\hat{\Xi}_{n}=\left[\begin{array}{ccc}
\vec{P}^{k_{n T}}(z)^{\prime} & 0 & 0  \tag{F.14}\\
0 & \vec{P}^{k_{n T}}(s)^{\prime} & 0 \\
0 & 0 & \vec{P}^{k_{n T}}(v)^{\prime}
\end{array}\right] S_{0} \hat{\Sigma}_{n} S_{0}^{\prime}\left[\begin{array}{ccc}
\vec{P}^{k_{n T}}(z) & 0 & 0 \\
0 & \vec{P}^{k_{n T}}(s) & 0 \\
0 & 0 & \vec{P}^{k_{n T}}(v)
\end{array}\right]
$$

This completes the proof of this theorem.
Proof of Theorem B.2: By Theorem B.1, we have $\left\|\check{\boldsymbol{\theta}}_{1}-\boldsymbol{\theta}_{1,0}\right\|=O_{p}\left(a_{n}\right)$, where $a_{n} \rightarrow 0$ as $n \rightarrow \infty$. It follows that $\left\|\check{\boldsymbol{\theta}}-\boldsymbol{\theta}_{0}\right\|=\left\|P_{\mathcal{J}}^{\prime}\left(\check{\boldsymbol{\theta}}_{1}-\boldsymbol{\theta}_{1,0}\right)\right\|=O_{p}\left(a_{n}\right)$, where $\check{\boldsymbol{\theta}}=P_{\mathcal{J}}^{\prime} \check{\boldsymbol{\theta}}_{1}$. Since $Q_{n}(\check{\boldsymbol{\theta}} ; \psi) \leq Q_{n}\left(\boldsymbol{\theta}_{0} ; \psi\right)$, for any $\varpi \in R^{p}$, we have

$$
\begin{align*}
& Q_{n}\left(\boldsymbol{\theta}_{0}+\varpi a_{n} ; \psi\right)-Q_{n}(\check{\boldsymbol{\theta}} ; \psi) \\
& \geq Q_{n}\left(\boldsymbol{\theta}_{0}+\varpi a_{n} ; \psi\right)-Q_{n}\left(\boldsymbol{\theta}_{0} ; \psi\right) \\
& =\bar{G}_{n}\left(\boldsymbol{\theta}_{0}+\varpi a_{n}\right)^{\prime} \Omega_{n} \bar{G}_{n}\left(\boldsymbol{\theta}_{0}+\varpi a_{n}\right)-\bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right)^{\prime} \Omega_{n} \bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right) \\
& +\sum_{l=1}^{p}\left[p_{c}\left(\left|\theta_{l, 0}+\varpi_{l} a_{n}\right|, \psi\right)-p_{c}\left(\left|\theta_{l, 0}\right|, \psi\right)\right] \\
& =-2 a_{n} \bar{G}_{n}(\overline{\boldsymbol{\theta}})^{\prime} \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widetilde{\hat{\mathbf{q}}}_{i}\right) \varpi \\
& +\sum_{l=1}^{p}\left[a_{n} \varpi_{l} p_{c}^{\prime}\left(\left|\theta_{l, 0}\right|, \psi\right)+\frac{a_{n}^{2} \varpi_{l}^{2}}{2} p_{c}^{\prime \prime}\left(\left|\theta_{l, 0}+v_{l} \varpi_{l} a_{n}\right|, \psi\right)\right] \tag{F.15}
\end{align*}
$$

where $\overline{\boldsymbol{\theta}}=a \boldsymbol{\theta}_{0}+(1-a)\left(\boldsymbol{\theta}_{0}+\varpi a_{n}\right)$ for some $a \in(0,1), v_{l} \in(0,1)$ for $l=1, \ldots, p$, and $\partial \bar{G}_{n}(\boldsymbol{\theta}) / \partial \boldsymbol{\theta}^{\prime}=-\mathbb{E}_{n}\left(\tilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\hat{\mathbf{q}}}_{i}\right)$. First, we have

$$
\begin{aligned}
\bar{G}_{n}(\overline{\boldsymbol{\theta}}) & =\mathbb{E}_{n}\left[\widetilde{\hat{\mathbf{d}}}_{i}^{\prime}\left(\tilde{\mathbf{y}}_{i}-\widetilde{\hat{\mathbf{q}}}_{i} \overline{\boldsymbol{\theta}}\right)\right] \\
& =\mathbb{E}_{n}\left[\widetilde{\hat{\mathbf{d}}}_{i}^{\prime}\left(\tilde{\mathbf{q}}_{i} \boldsymbol{\theta}_{0}-\widetilde{\hat{\mathbf{q}}}_{i} \overline{\boldsymbol{\theta}}\right)\right]+\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\boldsymbol{\Delta}}_{i}\right)+\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\boldsymbol{\varepsilon}}_{i}\right) \\
& =\mathbb{E}_{n}\left[\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} J_{T}\left(P_{v}^{k_{n T}}-P_{\hat{v}}^{k_{n} T}\right)\right] \boldsymbol{\vartheta}_{v, 0}+a_{n}(1-a) \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widetilde{\hat{\mathbf{q}}}_{i}\right) \varpi \\
& +\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\boldsymbol{\Delta}}_{i}\right)+\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\varepsilon}_{i}\right) .
\end{aligned}
$$

It follows that $\bar{G}_{n}(\overline{\boldsymbol{\theta}})^{\prime} \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\hat{\mathbf{q}}}_{i}\right) \varpi=A_{1}+a_{n}(1-a) A_{2}+A_{3}+A_{4}$, where the definition of $A_{1}$ to $A_{4}$ will be clarified in the context below.
(i) We have

$$
\begin{aligned}
\left\|A_{1}\right\| & =\left|\boldsymbol{\vartheta}_{v, 0}^{\prime} \mathbb{E}_{n}\left[\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} J_{T}\left(P_{v}^{k_{n T}}-P_{\hat{\boldsymbol{v}}}^{k_{n T}}\right)\right]^{\prime} \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widetilde{\hat{\mathbf{q}}}_{i}\right) \varpi\right| \\
& \leq\left\|\mathbb{E}_{n}\left[\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} J_{T}\left(P_{v}^{k_{n T}}-P_{\hat{v}}^{k_{n T}}\right)\right] \boldsymbol{\vartheta}_{v, 0} \times\right\|\left\|\Omega_{n}\right\|_{s p}\left\|\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widetilde{\hat{\mathbf{q}}}_{i}\right) \varpi\right\| \\
& =\|\varpi\| O_{p}\left(\delta_{n T}^{2}\right)
\end{aligned}
$$

because we have

$$
\begin{aligned}
& \left\|\mathbb{E}_{n}\left[\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} J_{T}\left(P_{v}^{k_{n T}}-P_{\hat{v}}^{k_{n T}}\right)\right] \boldsymbol{\vartheta}_{v, 0}\right\| \\
& =\left\|\mathbb{E}_{n}\left[I_{T} \otimes \operatorname{vec}\left(J_{T}\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)\right)^{\prime} J_{T}\left(P_{v}^{k_{n T}}-P_{\hat{v}}^{k_{n T}}\right)\right] \boldsymbol{\vartheta}_{v, 0}\right\| \\
& \leq \mathbb{E}_{n}\left\|\left[I_{T} \otimes \operatorname{vec}\left(J_{T}\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)\right)^{\prime} J_{T}\left(P_{v}^{k_{n T}}-P_{\hat{v}}^{k_{n T}}\right)\right] \boldsymbol{\vartheta}_{v, 0}\right\| \\
& =\mathbb{E}_{n}\left(\operatorname{tr}[\boldsymbol{\star}] \boldsymbol{\vartheta}_{v, 0}^{\prime}\left(P_{v}^{k_{n T}}-P_{\hat{v}}^{k_{n T}}\right)^{\prime} J_{T}\left(P_{v}^{k_{n T}}-P_{\hat{v}}^{k_{n T}}\right) \boldsymbol{\vartheta}_{v, 0}\right)^{1 / 2} \\
& \leq\left\|\boldsymbol{\vartheta}_{v, 0}\right\| \mathbb{E}_{n}(\operatorname{tr}[\boldsymbol{\star}])=O_{p}\left(\delta_{n T}^{2}\right)
\end{aligned}
$$

by (F.5), where $\star=\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)^{\prime} J_{T}\left(P_{\hat{v}, i}^{k_{n T}}-P_{v, i}^{k_{n T}}\right)$, and

$$
\begin{aligned}
& \left\|\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\hat{\mathbf{q}}}_{i}\right) \varpi\right\| \\
& \leq\left\|\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\hat{\mathbf{q}}}_{i}-\mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\mathbf{q}}_{i}\right)\right) \varpi\right\|+\left\|\mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\mathbf{q}}_{i}\right) \varpi\right\| \\
& \leq\|\varpi\|\left(\left\|\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\hat{\mathbf{q}}}_{i}-\mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\mathbf{q}}_{i}\right)\right)\right\|_{s p}+\left\|\mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\mathbf{q}}_{i}\right)\right\|_{s p}\right)\left(1+O_{p}\left(k_{n T} / \sqrt{n}\right)\right)
\end{aligned}
$$

by Lemma F.2.
(ii) We have $A_{2}=\varpi^{\prime} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widetilde{\hat{\mathbf{q}}}_{i}\right)^{\prime} \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widetilde{\hat{\mathbf{q}}}_{i}\right) \varpi=\varpi^{\prime} \Lambda \varpi+\varpi^{\prime}\left(\Lambda_{n}-\Lambda\right) \varpi$. As $\varpi^{\prime}\left(\Lambda_{n}-\Lambda\right) \varpi$ $\leq\|\varpi\|\left\|\Lambda_{n}-\Lambda\right\|_{s p}=\|\varpi\| O_{p}\left(k_{n T} / \sqrt{n}+\delta_{n T} J^{1 / 2}\right)$ by Lemma F.3, we obtain $A_{2}=\varpi^{\prime} \Lambda \varpi+$ $\|\varpi\| O_{p}\left(k_{n T} / \sqrt{n}+\delta_{n T} J^{1 / 2}\right)$.
(iii) Under Assumption 2(i) and by Lemma F.3, we have

$$
\left|A_{3}\right|=\left|\varpi^{\prime} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\Delta}_{i}\right) \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widetilde{\hat{\mathbf{q}}}_{i}\right) \varpi\right|=\|\varpi\| O_{p}\left(k_{n T}^{-\zeta}\left(k_{n T} / \sqrt{n}+\delta_{n T} J^{1 / 2}+1\right)\right) .
$$

Following the proof of (F.9) we obtain

$$
\left|\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\varepsilon}_{i}\right)^{\prime} \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widetilde{\hat{\mathbf{q}}}_{i}\right) \varpi\right|=\|\varpi\| O_{p}\left(\left(\delta_{n T}+\sqrt{J}\right) / \sqrt{n}\right) .
$$

Taking (i)-(iii) together gives

$$
\begin{align*}
& \bar{G}_{n}(\overline{\boldsymbol{\theta}})^{\prime} \Omega_{n} \mathbb{E}_{n}\left(\tilde{\hat{\mathbf{d}}}_{i}^{\prime} \tilde{\mathbf{q}}_{i}\right) \varpi-\varpi^{\prime} \Lambda \varpi \\
& =O_{p}\left(\delta_{n T}^{2}+k_{n T} / \sqrt{n}+\delta_{n T} J^{1 / 2}+k_{n T}^{-\zeta}+\sqrt{J / n}\right)\|\varpi\| \tag{F.16}
\end{align*}
$$

so that $\bar{G}_{n}(\overline{\boldsymbol{\theta}})^{\prime} \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widetilde{\hat{\mathbf{q}}}_{i}\right) \varpi=\varpi^{\prime} \Lambda \varpi+o_{p}(1)$ is positive in probability approaching one for any finite $\|\varpi\|$.

Second, the MCP penalty and its first-order derivative are defined as

$$
\begin{align*}
& p_{c}(t, \psi)=\left\{\begin{array}{cc}
\psi|t|-t^{2} c \psi / 2 & \text { if }|t| \leq c \psi \\
c \psi^{2} / 2 & \text { otherwise }
\end{array}\right.  \tag{F.17}\\
& p_{c}^{\prime}(t, \psi)=\left\{\begin{array}{cc}
(\psi-|t| / c) \operatorname{sign}(t) & \text { if }|t| \leq c \psi \\
0 & \text { otherwise }
\end{array}\right.
\end{align*}
$$

for $c>1$ and $p_{c}^{\prime \prime}(t, \psi)=-c^{-1} I(|t| \leq c \psi)$, where $\operatorname{sign}(t)=-1,0,1$ for negative, 0 and positive $t$, respectively. Hence, setting $(c, \psi)$ satisfying $c \psi<q_{0} \min _{l \in \mathcal{J}_{n}}\left\{\theta_{l, 0}\right\}$ for some $q_{0} \leq 1$, we have $\sum_{l=1}^{p} a_{n} \varpi_{l} p_{c}^{\prime}\left(\left|\theta_{l, 0}\right|, \psi\right)=0$ and

$$
\begin{aligned}
& \sum_{l=1}^{p} \varpi_{l}^{2} p_{c}^{\prime \prime}\left(\left|\theta_{l, 0}+v_{l} \varpi_{l} a_{n}\right|, \psi\right) \\
& =\sum_{l \in \mathcal{J}_{n}} \varpi_{l}^{2} p_{c}^{\prime \prime}\left(\left|\theta_{l, 0}+v_{l} \varpi_{l} a_{n}\right|, \psi\right)+\sum_{l \notin \mathcal{J}_{n}} \varpi_{l}^{2} p_{c}^{\prime \prime}\left(\left|v_{l} \varpi_{l} a_{n}\right|, \psi\right)=-c^{-1} \sum_{l \notin \mathcal{J}_{n}} \varpi_{l}^{2}
\end{aligned}
$$

where $p_{c}^{\prime \prime}\left(\left|v_{l} \varpi_{l} a_{n}\right|, \psi\right)=-c^{-1}$ for $l \notin \mathcal{J}_{n}$ if $\max _{l \notin \mathcal{J}_{n}}\left|v_{l} \varpi_{l} a_{n}\right|<a_{n}\|\varpi\|<c \psi$, and for $l \in \mathcal{J}_{n}$, $p_{c}^{\prime \prime}\left(\left|\theta_{l, 0}+v_{l} \varpi_{l} a_{n}\right|, \psi\right)=0$, because $\max _{l \in \mathcal{J}_{n}}\left|\theta_{l, 0}+v_{l} \varpi_{l} a_{n}\right|>\left(q_{0}^{-1}-1\right) c \psi>c \psi$ for $q_{0}<1 / 2$.

Combining this result with (F.15) and (F.16) gives

$$
\begin{aligned}
& Q_{n}\left(\boldsymbol{\theta}_{0}+\varpi a_{n} ; \psi\right)-Q_{n}(\check{\boldsymbol{\theta}} ; \psi)=2 a_{n} \varpi^{\prime} \Lambda \varpi-\frac{a_{n}^{2}}{2 c} \sum_{l \notin \mathcal{J}_{n}} \varpi_{l}^{2} \\
& +O_{p}\left(\delta_{n T}^{2}+k_{n T} / \sqrt{n}+\delta_{n T} J^{1 / 2}+k_{n T}^{-\zeta}+\sqrt{J / n}\right) a_{n}\|\varpi\| .
\end{aligned}
$$

Hence, $Q_{n}\left(\boldsymbol{\theta}_{0}+\varpi a_{n} ; \psi\right) \geq Q_{n}(\check{\boldsymbol{\theta}} ; \psi)$ holds in probability approaching one if the tuning parameter, $\psi$, falls into the interval $\left(a_{n}\|\varpi\| / c, \min _{l \in \mathcal{J}_{n}}\left\{\theta_{l, 0}\right\} /\left(c q_{0}\right)\right]$ for some $c>1, q_{0}<1 / 2$ and sufficiently large $\|\varpi\|$. Thus, there exists a local minimizer, $\hat{\boldsymbol{\theta}}(\psi)$, of $Q_{n}(\boldsymbol{\theta} ; \boldsymbol{\psi})$ in the ball $\left\{\boldsymbol{\theta}_{0}+a_{n}\|\varpi\|:\|\varpi\| \leq M\right\}$, and combining this result with Theorem B. 1 yields

$$
\operatorname{Pr}(\check{\boldsymbol{\theta}}=\hat{\boldsymbol{\theta}}(\psi)) \rightarrow 1 \text { as } n \rightarrow \infty .
$$

This completes the proof of this Theorem.
Proof of Theorem B.3: Let $\widehat{\mathcal{J}}=\operatorname{supp}(\hat{\boldsymbol{\theta}})$ and $\hat{J}=\operatorname{dim}(\widehat{\mathcal{J}})$ and split $\widetilde{\hat{\mathbf{q}}}_{i}=$ $\left[\widetilde{\hat{\mathbf{q}}}_{\widehat{\mathcal{J}}, 1, i}, \widetilde{\hat{\mathbf{q}}}_{\widehat{\mathcal{J}}, 2, i}\right]$ and $\tilde{\boldsymbol{\theta}}=\left[\tilde{\boldsymbol{\theta}}_{1}^{\prime}, 0_{p-\hat{J}}^{\prime}\right]^{\prime}$, where the parameters in front of $\widetilde{\hat{\mathbf{q}}}_{\widehat{\mathcal{J}}, 2, i}$ are all zeros and $\widetilde{\hat{\mathbf{q}}}_{i} \tilde{\boldsymbol{\theta}}=\widetilde{\hat{\mathbf{q}}}_{\widehat{\mathcal{J}}, 1, i} \tilde{\boldsymbol{\theta}}_{1}$. Additionally, we denote a parameter vector $\boldsymbol{\theta}_{\widehat{\mathcal{J}}, 0}=\left[\boldsymbol{\theta}_{\widehat{\mathcal{J}}, 1,0}, 0_{p-\hat{J}}\right]$ for a chosen $\hat{J}$ where $\boldsymbol{\theta}_{\widehat{\mathcal{J}}, 0}$ is the parameter vector to which $\hat{\boldsymbol{\theta}}$ converge. By definition of $\tilde{\boldsymbol{\theta}}$, $\bar{G}_{n}(\tilde{\boldsymbol{\theta}})^{\prime} \Omega_{n} \bar{G}_{n}(\tilde{\boldsymbol{\theta}}) \leq \bar{G}_{n}(\hat{\boldsymbol{\theta}})^{\prime} \Omega_{n} \bar{G}_{n}(\hat{\boldsymbol{\theta}})$ and $\bar{G}_{n}(\tilde{\boldsymbol{\theta}})^{\prime} \Omega_{n} \bar{G}_{n}(\tilde{\boldsymbol{\theta}}) \leq \bar{G}_{n}\left(\boldsymbol{\theta}_{\widehat{\mathcal{J}}, 0}\right)^{\prime} \Omega_{n} \bar{G}_{n}\left(\boldsymbol{\theta}_{\widehat{\mathcal{J}}, 0}\right)$, which implies that

$$
\bar{G}_{n}(\tilde{\boldsymbol{\theta}})^{\prime} \Omega_{n} \bar{G}_{n}(\tilde{\boldsymbol{\theta}})-\bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right)^{\prime} \Omega_{n} \bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right) \leq \min \left(P_{n}, S_{n}\right)
$$

where $P_{n}=\bar{G}_{n}(\hat{\boldsymbol{\theta}})^{\prime} \Omega_{n} \bar{G}_{n}(\hat{\boldsymbol{\theta}})-\bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right)^{\prime} \Omega_{n} \bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right)$ and

$$
S_{n}=\bar{G}_{n}\left(\boldsymbol{\theta}_{\widehat{\mathcal{J}}, 0}\right)^{\prime} \Omega_{n} \bar{G}_{n}\left(\boldsymbol{\theta}_{\widehat{\mathcal{J}}, 0}\right)-\bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right)^{\prime} \Omega_{n} \bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right)
$$

For any vector $\omega \in R^{p}$, applying simple algebra gives

$$
\begin{align*}
& \bar{G}_{n}\left(\boldsymbol{\theta}_{0}+\omega\right)^{\prime} \Omega_{n} \bar{G}_{n}\left(\boldsymbol{\theta}_{0}+\omega\right)-\bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right)^{\prime} \Omega_{n} \bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right) \\
& =\omega^{\prime} \Lambda_{n} \omega-2 \omega^{\prime} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \hat{\mathbf{q}}_{i}\right)^{\prime} \Omega_{n} \bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right) \tag{F.18}
\end{align*}
$$

where we denote $\Lambda_{n}=\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widetilde{\hat{\mathbf{q}}}_{i}\right)^{\prime} \Omega_{n} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widetilde{\hat{\mathbf{q}}}_{i}\right)$ since $\bar{G}_{n}\left(\boldsymbol{\theta}_{0}+\omega\right)=\bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right)-\mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widetilde{\hat{\mathbf{q}}}_{i}\right) \omega$.

Additionally, we have

$$
\begin{equation*}
\bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right)^{\prime} \Omega_{n} \bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right)=O_{p}\left(k_{n T}^{-2 \zeta}+\delta_{n T}^{2} k_{n T}+J / n\right) \tag{F.19}
\end{equation*}
$$

by Lemma F.1, which implies that

$$
\begin{aligned}
& \left|\omega^{\prime} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widetilde{\hat{\mathbf{q}}}_{i}\right)^{\prime} \Omega_{n} \bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right)\right| \leq\left\|\omega^{\prime} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widetilde{\hat{\mathbf{q}}}_{i}\right)^{\prime} \Omega_{n}^{1 / 2}\right\|\left\|\Omega_{n}^{1 / 2} \bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right)\right\| \\
& =\left(\omega^{\prime} \Lambda_{n} \omega\right)^{1 / 2} O_{p}\left(k_{n T}^{-\zeta}+\delta_{n T} k_{n T}^{1 / 2}+\sqrt{J / n}\right) .
\end{aligned}
$$

Then, we have

$$
\begin{aligned}
& \left|\bar{G}_{n}(\tilde{\boldsymbol{\theta}})^{\prime} \Omega_{n} \bar{G}_{n}(\tilde{\boldsymbol{\theta}})-\bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right)^{\prime} \Omega_{n} \bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right)-\left(\tilde{\boldsymbol{\theta}}-\boldsymbol{\theta}_{0}\right)^{\prime} \Lambda_{n}\left(\tilde{\boldsymbol{\theta}}-\boldsymbol{\theta}_{0}\right)\right| \\
& \quad \leq\left[\left(\tilde{\boldsymbol{\theta}}-\boldsymbol{\theta}_{0}\right)^{\prime} \Lambda_{n}\left(\tilde{\boldsymbol{\theta}}-\boldsymbol{\theta}_{0}\right)\right]^{1 / 2} O_{p}\left(k_{n T}^{-\zeta}+\delta_{n T} k_{n T}^{1 / 2}+\sqrt{J / n}\right) .
\end{aligned}
$$

Hence, we obtain

$$
\begin{equation*}
\left[\left(\tilde{\boldsymbol{\theta}}-\boldsymbol{\theta}_{0}\right)^{\prime} \Lambda_{n}\left(\tilde{\boldsymbol{\theta}}-\boldsymbol{\theta}_{0}\right)\right]^{1 / 2}=O_{p}\left(k_{n T}^{-\zeta}+\delta_{n T} k_{n T}^{1 / 2}+\sqrt{J / n}\right)+\sqrt{\min \left(P_{n,+}, S_{n,+}\right)} \tag{F.20}
\end{equation*}
$$

where $P_{n,+}=\max \left(P_{n}, 0\right)$ and $S_{n,+}=\max \left(S_{n}, 0\right)$. By (F.18), we have

$$
\begin{aligned}
& P_{n}=\left(\hat{\boldsymbol{\theta}}-\boldsymbol{\theta}_{0}\right)^{\prime} \Lambda_{n}\left(\hat{\boldsymbol{\theta}}-\boldsymbol{\theta}_{0}\right)-2\left(\hat{\boldsymbol{\theta}}-\boldsymbol{\theta}_{0}\right)^{\prime} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widetilde{\hat{\mathbf{q}}}_{i}\right)^{\prime} \Omega_{n} \bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right) \\
& S_{n}=\left(\boldsymbol{\theta}_{\widehat{\mathcal{J}}, 0}-\boldsymbol{\theta}_{0}\right)^{\prime} \Lambda_{n}\left(\boldsymbol{\theta}_{\widehat{\mathcal{J}}, 0}-\boldsymbol{\theta}_{0}\right)-2\left(\boldsymbol{\theta}_{\widehat{\mathcal{J}}, 0}-\boldsymbol{\theta}_{0}\right)^{\prime} \mathbb{E}_{n}\left(\widetilde{\hat{\mathbf{d}}}_{i}^{\prime} \widetilde{\hat{\mathbf{q}}}_{i}\right)^{\prime} \Omega_{n} \bar{G}_{n}\left(\boldsymbol{\theta}_{0}\right) .
\end{aligned}
$$

By Theorem B.1, we have

$$
\begin{aligned}
P_{n,+} & \leq \boldsymbol{\nabla}+\mathbf{\nabla}^{1 / 2} O_{p}\left(k_{n T}^{-\zeta}+\delta_{n T_{0}} k_{n T}^{1 / 2}+\sqrt{J / n}\right) \\
& =O_{p}\left(k_{n T}^{-2 \zeta}+\frac{J}{n}\right)+O_{p}\left[\left(k_{n T}^{-\zeta}+\sqrt{\frac{J}{n}}\right)\left(k_{n T}^{-\zeta}+\delta_{n T} k_{n T}^{1 / 2}+\sqrt{J / n}\right)\right],
\end{aligned}
$$

where $\boldsymbol{\nabla}=\left(\hat{\boldsymbol{\theta}}-\boldsymbol{\theta}_{0}\right)^{\prime} \Lambda_{n}\left(\hat{\boldsymbol{\theta}}-\boldsymbol{\theta}_{0}\right)$. In addition, denoting the complement set of $\mathcal{J}$ and $\widehat{\mathcal{J}}$ by $\mathcal{J}^{c}=\left\{l \in(1, \ldots, p): \theta_{l, 0}=0\right\}$ and $\widehat{\mathcal{J}}^{c}=\left\{l \in(1, \ldots, p): \theta_{\widehat{\mathcal{J}}, l, 0}=0\right\}$, respectively, we have $0 \leq m=\operatorname{dim}\left(\mathcal{J}^{c} \cap \widehat{\mathcal{J}}^{c}\right) \leq p-J$, where $m=0$ if $\mathcal{J} \cap \widehat{\mathcal{J}}=\varnothing$ and $m=p-J$ if $\mathcal{J}=\widehat{\mathcal{J}}$.

It is readily seen that $\boldsymbol{\theta}_{\widehat{\mathcal{J}}, 0}=\boldsymbol{\theta}_{0}$ so that $S_{n}=0$ if $\mathcal{J} \subseteq \widehat{\mathcal{J}}$. If $\mathcal{J} \subseteq \widehat{\mathcal{J}}$ fails to hold, we have

$$
\begin{aligned}
S_{n,+} & \leq \nabla+\nabla^{1 / 2} O_{p}\left(k_{n T}^{-\zeta}+\delta_{n T} k_{n T}^{1 / 2}\right) \\
& =O_{p}(p-m)+O_{p}\left(\sqrt{p-m}\left(k_{n T}^{-\zeta}+\delta_{n T} k_{n T}^{1 / 2}+\sqrt{J / n}\right)\right)
\end{aligned}
$$

where $\nabla=\left(\boldsymbol{\theta}_{\widehat{\mathcal{J}}, 0}-\boldsymbol{\theta}_{0}\right)^{\prime} \Lambda_{n}\left(\boldsymbol{\theta}_{\widehat{\mathcal{J}}, 0}-\boldsymbol{\theta}_{0}\right)$. Therefore, (F.20) is bounded by

$$
\left(\tilde{\boldsymbol{\theta}}-\boldsymbol{\theta}_{0}\right)^{\prime} \Lambda_{n}\left(\tilde{\boldsymbol{\theta}}-\boldsymbol{\theta}_{0}\right)=O_{p}\left(k_{n T}^{-2 \zeta}+\frac{J}{n}\right) .
$$

By Lemma F.3, we have $\left\|\Lambda_{n}-\Lambda\right\|_{s p}=O_{p}\left(k_{n T} / \sqrt{n}+\delta_{n T} J^{1 / 2}\right)$, where $\Lambda=\mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\mathbf{q}}_{i}\right)^{\prime} \Omega_{n} \mathbb{E}\left(\tilde{\mathbf{d}}_{i}^{\prime} \tilde{\mathbf{q}}_{i}\right)$ is a finite non-singular matrix under Assumption 3. Hence, $\left\|\tilde{\boldsymbol{\theta}}-\boldsymbol{\theta}_{0}\right\|=O_{p}\left(k_{n T}^{-\zeta}+\sqrt{J / n}\right)$. This completes the proof of this theorem.

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[^0]:    ${ }^{1}$ If $A=\left(a_{i j}\right)$ is a $3 \times 3$ symmetric matrix, $\operatorname{vec}(A)=\left[a_{11}, a_{12}, a_{13}, a_{21}, a_{22}, a_{23}, a_{31}, a_{32}, a_{33}\right]^{\prime}$ and $\operatorname{vech}(A)=$ $\left[a_{11}, a_{12}, a_{13}, a_{22}, a_{23}, a_{33}\right]^{\prime}$.
    ${ }^{2}$ For any vectors $a(m \times 1)$ and $b(n \times 1), a \otimes b=\left(I_{m} \otimes b\right) a$.

[^1]:    ${ }^{3} P_{\mathcal{J}}$ is the $J \times p$ sub-matrix of the identity matrix $I_{p}$ that satisfies $\boldsymbol{\theta}_{1,0}=P_{\mathcal{J}} \boldsymbol{\theta}_{0}$ and $P_{\mathcal{J}}^{\prime} P_{\mathcal{J}}=I_{J}$.

[^2]:    ${ }^{4}$ See Online Appendix (F.12).

[^3]:    ${ }^{5}$ At least 252 trading days of non-missing data are required in this estimation.
    ${ }^{6}$ This calculation uses the first-order approximation of duration.
    ${ }^{7}$ The usual criterion of statistical significance at the 10 percent level is 1.645 . We instead use 1 because a

[^4]:    ${ }^{9}$ Note that for panel (f) where control variables from Kim and Kung (2017) and Panousi and Papanikolaou (2012) are included, we follow their work to treat those variables as exogenous variables in the estimation.
    ${ }^{10}$ See discussions in section II in Li and Sun (2022).
    ${ }^{11}$ For panels (b)-(e), see discussions in footnote 18 in Li and Sun (2022).

[^5]:    ${ }^{12}$ We follow the idea of equations (5) and (6) in Gabaix and Koijen (2020) to construct a GIV as the difference between size-weighted average stock volatility and equal-weighted average stock volatility for each business sector (classified by two-digit SIC codes). We choose to let our instruments vary with both sectors and time, rather than just time as in Gabaix and Koijen (2020), to increase estimation efficiency. Moreover, given that both our main regression model and the first-stage regression model contain time dummies as regressors, using GIVs which vary only with time as instruments is not feasible due to multicollinearity.
    ${ }^{13}$ See discussions in Appendix D.

[^6]:    ${ }^{14} \mathrm{~A}$ uniform setting of lags means that we set the same number of lags to all firms, regardless of their time-series length. We consider the maximum number of lags to be 3 because the maximum time-series length is 32 years, and $\operatorname{int}\left\{4(T / 100)^{2 / 9}\right\}=3$.

[^7]:    ${ }^{15}$ This is from Assumption 3.

