## Supplement to

## "Locally Robust Inference for

Non-Gaussian Linear Simultaneous
Equations Models"*

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#### Abstract

In this supplementary material we provide the following additional results. S1: A more general model S2: Supporting results for the main Theorems S3: Additional auxillary results S4: A consistent estimator of the Moore - Penrose pseudoinverse S5: Log density score estimation S6: Power optimality under strong identification S7: Additional simulation results S8: Additional empirical results


[^0]Throughout this document, references to lemmas, equations etc. which start with a " S " are references to this document. Those which consist of just a number refer to the main text.

## S1 A more general model

## S1.1 Model setup, ULAN and the effective score

In this section we extend the approach in the main paper to the more general model:

$$
\begin{equation*}
Y_{i}=B\left(b, X_{i}\right)+A\left(\alpha, \sigma, X_{i}\right)^{-1} \epsilon_{i}, \quad i=1, \ldots, n \tag{S1}
\end{equation*}
$$

under Assumptions S1 and S2 below, which are weakened versions of Assumptions 1 and 2 respectively. This version of the model allows (a) (parametric) conditional heteroskedasticity in the reduced form error $A\left(\alpha, \sigma, X_{i}\right)^{-1} \epsilon_{i}$ and (b) the conditional mean $\mathbb{E}\left[Y_{i} \mid X_{i}\right]=B\left(b, X_{i}\right)$ to be a non-linear function of $X_{i}$, known up to a finite dimensional parameter $b$.

Assumption S1. Suppose that for all $(\alpha, \beta) \in \mathcal{A} \times \mathcal{B}$,

1. $A(\alpha, \sigma, X)$ is non-singular for all $X$;
2. $(\alpha, \sigma) \mapsto A(\alpha, \sigma, X)$ and $b \mapsto B(b, X)$ are continuously differentiable for all $X$.

Define the partial derivative matrices $D_{\alpha, l}(\alpha, \sigma, X)=\partial A(\alpha, \sigma, X) / \partial \alpha_{l}$, for $l=1, \ldots, L_{\alpha}$ $D_{\sigma, l}(\alpha, \sigma, X)=\partial A(\alpha, \sigma, X) / \partial \sigma_{l}$, for $l=1, \ldots, L_{\sigma}$ and $D_{b, l}:=\partial B(b, X) / \partial b_{l}$ for $l=1, \ldots, L_{b}$. Further, for each $k, j \in\{1, \ldots, K\}, l \in\left\{1, \ldots, L_{\alpha}\right\}$ and $m \in\left\{1, \ldots, L_{\sigma}\right\}$ define $\zeta_{l, k, j}^{\alpha}(\alpha, \sigma, X):=$ $e_{k}^{\prime} D_{\alpha, l}(\alpha, \sigma, X) A(\alpha, \sigma, X)^{-1} e_{j}$ and $\zeta_{m, k, j}^{\sigma}:=e_{k}^{\prime} D_{\sigma, m}(\alpha, \sigma, X) A(\alpha, \sigma, X)^{-1} e_{j}$. With this notation, for all $(\alpha, \beta) \in \mathcal{A} \times \mathcal{B}$
3. $(\alpha, \sigma) \rightarrow \zeta_{l, k, j}^{\alpha}(\alpha, \sigma, X)$ and $(\alpha, \sigma) \rightarrow \zeta_{m, k, j}^{\sigma}(\alpha, \sigma, X)$ are locally Lipschitz continuous for all $j, k, l, m$ considered and all $X$.
4. $\|A(\alpha, \sigma, X)\|,\left\|A(\alpha, \sigma, X)^{-1}\right\|,\left\|D_{\alpha, l}(\alpha, \sigma, X)\right\|$ and $\left\|D_{\sigma, l}(\alpha, \sigma, X)\right\|$ are locally (in $\left.(\alpha, \sigma)\right)$ bounded.

Assumption S2. For $\epsilon_{i}=\left(\epsilon_{i, 1}, \ldots, \epsilon_{i, K}\right)^{\prime}$ in model (S1), each component $\epsilon_{i, k}$ has a continuously differentiable root density (with respect to Lebesgue measure on $\mathbb{R}$ ). We write the density as $\eta_{k}$ with log density score $\phi_{k}(x)=\partial \log \eta_{k}(x) / \partial x$. We assume that for all $k=1, \ldots, K$ and some $\delta>0$

1. $\mathbb{E} \epsilon_{i, k}=0, \mathbb{E} \epsilon_{i, k}^{2}=1, \mathbb{E} \epsilon_{i, k}^{4+\delta}<\infty, \mathbb{E}\left(\epsilon_{i, k}^{4}\right)-1>\mathbb{E}\left(\epsilon_{i, k}^{3}\right)^{2}$, and $\mathbb{E} \phi_{k}^{4+\delta}\left(\epsilon_{i, k}\right)<\infty$;
2. $\mathbb{E} \phi_{k}\left(\epsilon_{i, k}\right)=0, \mathbb{E} \phi_{k}\left(\epsilon_{i, k}\right) \epsilon_{i, k}=-1, \mathbb{E} \phi_{k}\left(\epsilon_{i, k}\right) \epsilon_{i, k}^{2}=0$ and $\mathbb{E} \phi_{k}\left(\epsilon_{i, k}\right) \epsilon_{i, k}^{3}=-3$;
3. $\epsilon_{i, k}$ is independent of $\epsilon_{i, l}$ for all $k \neq l$;
4. $\eta_{0} \in \mathscr{Z}$ is a density function (with respect to Lebesgue measure on $\mathbb{R}^{d-1}$ ) such that if $\tilde{X}_{i} \sim \eta_{0}, \mathbb{E}\left[\left\|D_{b, l}\left(b+\varrho, X_{i}\right)\right\|^{4+\delta}\right] \leq \bar{D}_{b, l}(b)<\infty$ for all $b \in \mathcal{B}$, all $\varrho$ in a neighbourhood of zero and all $l=1, \ldots, L_{b}$;
5. $\epsilon_{i}$ and $\tilde{X}_{i}$ are independent.

Remark 1. If $A(\alpha, \sigma, X)=A(\alpha, \sigma)$ and $B(b, X)=\operatorname{vec}^{-1}(b) X$ then Assumptions $S 1$ and $S 2$ are implied by Assumptions 1 and 2 respectively.

Formally, the considered model is the collection

$$
\begin{equation*}
\mathcal{P}_{\Theta}=\left\{P_{\theta}: \theta \in \Theta\right\}, \tag{S2}
\end{equation*}
$$

where each $P_{\theta}$ is the law of the data $W_{i}=\left(Y_{i}, \tilde{X}_{i}\right)$ which lies in $\mathcal{W} \subset \mathbb{R}^{K+d-1}$. The parameter space $\Theta$ has the form $\Theta=\mathcal{A} \times \mathcal{B} \times \mathcal{H}$, where $\mathcal{A} \subset \mathbb{R}^{L_{\alpha}}, \mathcal{B} \subset \mathbb{R}^{L_{\beta}}$. $\mathcal{H}$ has the form $\mathscr{Z} \times \prod_{k=1}^{K} \mathscr{H}$, where $\mathscr{Z}$ is the space of density functions $\eta_{0}$ and $\mathscr{H}$ is the space of density functions $\eta_{k}$ such that if $\tilde{X} \sim \eta_{0}$ and $\epsilon_{k} \sim \eta_{k}$ then Assumption S2 parts 1, 3, 4 and 5 hold. ${ }^{\text {S1 }}$.

We write a typical element of $\Theta$ as $\theta=(\alpha, \beta, \eta)$, where $\beta=\left(b^{\prime}, \sigma^{\prime}\right)^{\prime}$ and it is understood that $\alpha \in \mathcal{A}, \beta \in \mathcal{B}$ and $\eta \in \mathcal{H}$. In what follows we will let $V_{\theta, i}:=Y_{i}-B\left(b, X_{i}\right)$ be the reduced form error so that $A\left(\alpha, \sigma, X_{i}\right) V_{\theta, i}=\epsilon_{i}$. Each $P_{\theta}$ is absolutely continuous with respect to Lebesgue measure on $\mathbb{R}^{K+d-1}$, with (Lebesgue) density given by

$$
\begin{equation*}
p_{\theta}\left(W_{i}\right)=\left|\operatorname{det} A\left(\alpha, \sigma, X_{i}\right)\right| \prod_{k=1}^{K} \eta_{k}\left(e_{k}^{\prime} A\left(\alpha, \sigma, X_{i}\right) V_{\theta, i}\right) \times \eta_{0}\left(\tilde{X}_{i}\right), \tag{S3}
\end{equation*}
$$

and hence log-density

$$
\begin{equation*}
\ell_{\theta}\left(W_{i}\right)=\log \left|\operatorname{det} A\left(\alpha, \sigma, X_{i}\right)\right|+\sum_{k=1}^{K} \log \eta_{k}\left(e_{k}^{\prime} A\left(\alpha, \sigma, X_{i}\right) V_{\theta, i}\right)+\log \eta_{0}\left(\tilde{X}_{i}\right) . \tag{S4}
\end{equation*}
$$

The differentiable paths we consider have the following form.
Let $H=H_{0} \times \prod_{k=1}^{K} H_{k}$, where each $H_{k}$ is as defined following (6). Given a direction $(g, h) \in \mathbb{R}^{L} \times H$, the measures $P_{t}$ are those corresponding to the density with form as in (S3) evaluated at $\theta_{t}:=\left(\gamma+t g, \eta_{t}\right)$ where the $k$-th coordinate of $\eta_{t}$ is $\eta_{k, t}^{h_{k}}:=\eta_{k}\left(1+t h_{k}\right)$ $(k=0, \ldots, K)$.

[^1]We have the following analogues of Lemmas 1, 2 and 3.
Lemma S1. Suppose Assumptions $S 1$ and S2 hold and that $(\alpha, \beta)$ is an interior point of $\mathcal{A} \times \mathcal{B}$. For each $(g, h) \in \mathbb{R}^{L} \times H:=\mathcal{V}$, the map $t \mapsto P_{\theta_{t}}$ is a differentiable path, with score function $g^{\prime} \dot{\ell}_{\theta}+\tilde{h}_{0}+\sum_{k=1}^{K} \tilde{h}_{k}$, where $\dot{\ell}_{\theta}:=\nabla_{\gamma} \log p_{\theta}, \tilde{h}_{0}(W):=h_{0}(\tilde{X})$ and $\tilde{h}_{k}(W):=$ $h_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) . \dot{\ell}_{\theta}$ has the form $\dot{\ell}_{\theta}=\left(\dot{\ell}_{\theta, \alpha}^{\prime}, \dot{\ell}_{\theta, b}^{\prime}, \dot{\ell}_{\theta, \sigma}^{\prime}\right)^{\prime}$, with

$$
\begin{aligned}
\dot{\ell}_{\theta, \alpha, l}(W):= & \sum_{k=1}^{K} \sum_{j=1, j \neq k}^{K} \zeta_{l, k, j}^{\alpha}(\alpha, \sigma, X) \phi_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) e_{j}^{\prime} A(\alpha, \sigma, X) V_{\theta} \\
& +\sum_{k=1}^{K} \zeta_{l, k, k}^{\alpha}(\alpha, \sigma, X)\left[\phi_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}+1\right] ; \\
\dot{\ell}_{\theta, b, l}(W):= & -\sum_{k=1}^{K} \phi_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) e_{k}^{\prime} A(\alpha, \sigma, X) D_{b, l}(b, X) ; \\
\dot{\ell}_{\theta, \sigma, l}(W):= & \sum_{k=1}^{K} \sum_{j=1, j \neq k}^{K} \zeta_{l, k, j}^{\sigma}(\alpha, \sigma, X) \phi_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) e_{j}^{\prime} A(\alpha, \sigma, X) V_{\theta} \\
& +\sum_{k=1}^{K} \zeta_{l, k, k}^{\sigma}(\alpha, \sigma, X)\left[\phi_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}+1\right]
\end{aligned}
$$

Proof. Let $g=(a, \varrho, s) \in \mathbb{R}^{L_{\alpha}} \times \mathbb{R}^{L_{b}} \times \mathbb{R}^{L_{\sigma}}$. The $\log$ density of $W$ under $\theta_{t}$ is then

$$
\begin{aligned}
\ell_{\theta_{t}}(W)= & \log p_{\theta_{t}}(W) \\
= & \log \eta_{0}(\tilde{X})+\log \left(1+t h_{0}(\tilde{X})\right)+\log |\operatorname{det}(A(\alpha+t a, \sigma+t s, X))| \\
& +\sum_{k=1}^{K} \log \eta_{k}\left(e_{k}^{\prime} A(\alpha+t a, \sigma+t s, X)(Y-B(b+t \varrho, X))\right) \\
& +\sum_{k=1}^{K} \log \left(1+t h_{k}\left(e_{k}^{\prime} A(\alpha+t a, \sigma+t s, X)(Y-B(b+t \varrho, X))\right)\right) .
\end{aligned}
$$

By Lemma $\mathrm{S} 6, t \mapsto \sqrt{p_{\theta_{t}}}$ is continuously differentiable (pointwise) in a neighbourhood $\mathcal{V}$ of 0 . Moreover, if we define $q_{t}(W):=\left.\frac{\partial \log p_{\theta_{x}}(W)}{\partial x}\right|_{x=t}$ and $Q_{t}:=P_{\theta_{t}} q_{t}(W)^{2}, Q_{t}$ is finite and continuous in a neighbourhood of 0 by the uniformly integrability of $\left\{q_{t}(W)^{2}: t \in \mathcal{V}\right\}$ along with the pointwise continuity of $t \mapsto q_{t}(W)$, both of which follow from Lemma S6. Hence, by Lemma 1.8 in van der Vaart (2002), $t \mapsto P_{\theta_{t}}$ is a differentiable path with score function
given by the derivative of $\ell_{\theta_{t}}(W)$ at $t=0$, which is:

$$
\begin{align*}
& \sum_{k=1}^{K} \phi_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) e_{k}^{\prime} \sum_{l=1}^{L_{\alpha}} a_{l} D_{\alpha, l}(\alpha, \sigma, X) V_{\theta}+\sum_{l=1}^{L_{\alpha}} a_{l} \operatorname{tr}\left(A(\alpha, \sigma, X)^{-1} D_{\alpha, l}(\alpha, \sigma, X)\right) \\
& +\sum_{k=1}^{K} \phi_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) e_{k}^{\prime} \sum_{l=1}^{L_{\sigma}} s_{l} D_{\sigma, l}(\alpha, \sigma, X) V_{\theta}+\sum_{l=1}^{L_{\sigma}} s_{l} \operatorname{tr}\left(A(\alpha, \sigma, X)^{-1} D_{\sigma, l}(\alpha, \sigma, X)\right) \\
& -\sum_{k=1}^{K} \phi_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) e_{k}^{\prime} A(\alpha, \sigma, X) \sum_{l=1}^{L_{b}} \varrho_{l} D_{b, l}(b, X) \\
& +h_{0}(\tilde{X})+\sum_{k=1}^{K} h_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) \tag{S5}
\end{align*}
$$

We can re-write the two expressions involving the trace as follows: for any $x \in\{\alpha, \sigma\}$ and appropriate index $l$ we have

$$
\begin{aligned}
\sum_{k=1}^{K} & \phi_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) e_{k}^{\prime} D_{x, l}(\alpha, \sigma, X) V_{\theta}+\operatorname{tr}\left(A(\alpha, \sigma, X)^{-1} D_{x, l}(\alpha, \sigma, X)\right) \\
= & \sum_{k=1}^{K} \phi_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) e_{k}^{\prime} D_{x, l}(\alpha, \sigma, X) A(\alpha, \sigma, X)^{-1} \epsilon+\operatorname{tr}\left(D_{x, l}(\alpha, \sigma, X) A(\alpha, \sigma, X)^{-1}\right) \\
= & \sum_{k=1}^{K} \sum_{j=1, j \neq k}^{K} \zeta_{l, k, j}^{x}(\alpha, \sigma, X) \phi_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) e_{j}^{\prime} A(\alpha, \sigma, X) V_{\theta} \\
& +\sum_{k=1}^{K} \zeta_{l, k, k}^{x}(\alpha, \sigma, X)\left[\phi_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}+1\right]
\end{aligned}
$$

for $\zeta_{l, k, j}^{x}(\alpha, \sigma, X):=e_{k}^{\prime} D_{x, l}(\alpha, \sigma, X) A(\alpha, \sigma, X)^{-1} e_{j}$. We may therefore write the derivative (S5) as $a^{\prime} \dot{\ell}_{\theta, \alpha}+\varrho^{\prime} \dot{\ell}_{\theta, b}+s^{\prime} \dot{\ell}_{\theta, \sigma}+\dot{\ell}_{\theta, \eta, h}$ where

$$
\dot{\ell}_{\theta, \eta, h}:=h_{0}(\tilde{X})+\sum_{k=1}^{K} h_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right)=\tilde{h}_{0}(W)+\sum_{k=1}^{K} \tilde{h}_{k}(W) .
$$

An elementary calculation reveals that $g^{\prime} \dot{\ell}_{\theta}=a^{\prime} \dot{\ell}_{\theta, \alpha}+\varrho^{\prime} \dot{\ell}_{\theta, b}+s^{\prime} \dot{\ell}_{\theta, \sigma}$.
Lemma S2. Suppose that Assumptions S1 and S2 hold and that $(\alpha, \beta)$ is an interior point of $\mathcal{A} \times \mathcal{B}$. For $(g, h) \in \mathcal{V}$ let

$$
\theta_{n}(g, h):=\theta+n^{-1 / 2}\left(g, \eta_{0} h_{0}, \ldots, \eta_{K} h_{K}\right)
$$

For any convergent sequence $\left(g_{n}, h_{n}\right) \rightarrow(g, h)$ (all in $\left.\mathcal{V}\right)$, define $R_{n}$ as

$$
R_{n}:=\log \prod_{i=1}^{n} \frac{p_{\theta_{n}\left(g_{n}, h_{n}\right)}\left(W_{i}\right)}{p_{\theta}\left(W_{i}\right)}-\frac{1}{\sqrt{n}} \sum_{i=1}^{n}\left[g^{\prime} \dot{\ell}_{\theta}\left(W_{i}\right)+\sum_{k=0}^{K} \tilde{h}_{k}\left(W_{i}\right)\right]+\frac{1}{2} \mathbb{E}\left[g^{\prime} \dot{\ell}_{\theta}\left(W_{i}\right)+\sum_{k=0}^{K} \tilde{h}_{k}\left(W_{i}\right)\right]^{2} .
$$

Then,

1. $R_{n} \xrightarrow{P_{\theta}} 0$,
2. Under $P_{\theta}$,

$$
\frac{1}{\sqrt{n}} \sum_{i=1}^{n}\left[g^{\prime} \dot{\ell}_{\theta}\left(W_{i}\right)+\sum_{k=0}^{K} \tilde{h}_{k}\left(W_{i}\right)\right] \rightsquigarrow \mathcal{N}\left(0, \mathbb{E}\left[g^{\prime} \dot{\ell}_{\theta}\left(W_{i}\right)+\sum_{k=0}^{K} \tilde{h}_{k}\left(W_{i}\right)\right]^{2}\right)
$$

3. The (product) measures $P_{\theta_{n}}^{n}$ and $P_{\theta}^{n}$ are mutually contiguous.

Proof. The proof proceeds verbatim as that of Lemma 2 on replacing Lemma 1 with Lemma S1.

Lemma S3. Suppose Assumptions S1 and S2 hold. Then the components of $\tilde{\ell}_{\theta}$ are as follows.
For $x=\alpha$ or $x=\sigma$,

$$
\begin{aligned}
& \tilde{\ell}_{\theta, x, l}(W)=\sum_{k=1}^{K} \sum_{j=1, j \neq k}^{K} \zeta_{l, k, j}^{x}(\alpha, \sigma, X) \phi_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) e_{j}^{\prime} A(\alpha, \sigma, X) V_{\theta} \\
& +\sum_{k=1}^{K}\left(\zeta_{l, k, k}^{x}(\alpha, \sigma, X)-\mathbb{E}\left[\zeta_{l, k, k}^{x}(\alpha, \sigma, X)\right]\right)\left[\phi_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}+1\right] \\
& +\sum_{k=1}^{K} \mathbb{E}\left[\zeta_{l, k, k}^{x}(\alpha, \sigma, X)\right]\left(\tau_{k, 1} e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}+\tau_{k, 2} \kappa\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right)\right),
\end{aligned}
$$

with $l$ in $\left\{1, \ldots, L_{\alpha}\right\}$ or $\left\{1, \ldots, L_{\sigma}\right\}$ (respectively); for $l=1, \ldots, L_{b}$,

$$
\begin{aligned}
\tilde{\ell}_{\theta, b, l}(W)=- & \sum_{k=1}^{K} \phi_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) e_{k}^{\prime}\left(A(\alpha, \sigma, X) D_{b, l}(b, X)-\mathbb{E}\left[A(\alpha, \sigma, X) D_{b, l}(b, X)\right]\right) \\
& +\sum_{k=1}^{K} e_{k}^{\prime} \mathbb{E}\left[A(\alpha, \sigma, X) D_{b, l}(b, X)\right]\left(\varsigma_{k, 1} e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}+\varsigma_{k, 2} \kappa\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right)\right)
\end{aligned}
$$

where the expectations are taken under $P_{\theta}$ and

$$
\tau_{k}:=M_{k}^{-1}\binom{0}{-2}, \varsigma_{k}:=M_{k}^{-1}\binom{1}{0}, \text { for } M_{k}:=\left(\begin{array}{cc}
1 & \mathbb{E}\left[\epsilon_{k}^{3}\right] \\
\mathbb{E}\left[\epsilon_{k}^{3}\right] & \mathbb{E}\left[\epsilon_{k}^{4}\right]-1
\end{array}\right)
$$

Proof. For each $h_{k} \in H_{k}$, define the corresponding $\tilde{h}_{k}$ as in the statement of Lemma S1 and let $\tilde{H}_{k}$ collect all such $\tilde{h}_{k}$ formed with $h_{k}$ ranging over $H_{k} .{ }^{\mathrm{S} 2}$ By the definition of $\tilde{\ell}_{\theta}$ in equation (31) and Theorem 4.11 in Rudin (1987) it suffices to show that each such component is (a) in $\left(\tilde{H}_{0}+\cdots+\tilde{H}_{K}\right)^{\perp}$ and (b) $\dot{\ell}_{\theta, x}-\tilde{\ell}_{\theta, x} \in \operatorname{cl}\left(\tilde{H}_{0}+\cdots+\tilde{H}_{K}\right)$, the form of which is given in Lemma S 8.

Case 1: $x=\alpha$ or $x=\sigma$. For (a) note that if $j \neq k$, then

$$
\begin{aligned}
\mathbb{E}\left[\zeta_{l, k, j}^{x}(\alpha, \sigma, X) \phi_{k}\left(\epsilon_{k}\right) \epsilon_{j} h_{0}(\tilde{X})\right] & =\mathbb{E}\left[\zeta_{l, k, j}^{x}(\alpha, \sigma, X) \phi_{k}\left(\epsilon_{k}\right) h_{0}(\tilde{X})\right] \mathbb{E}\left[\epsilon_{j}\right]=0 \\
\mathbb{E}\left[\zeta_{l, k, j}^{x}(\alpha, \sigma, X) \phi_{k}\left(\epsilon_{k}\right) \epsilon_{j} h_{m}\left(\epsilon_{m}\right)\right] & =\mathbb{E}\left[\zeta_{l, k, j}^{x}(\alpha, \sigma, X)\right] \mathbb{E}\left[\phi_{k}\left(\epsilon_{k}\right) \epsilon_{j} h_{m}\left(\epsilon_{m}\right)\right]=0
\end{aligned}
$$

where the last equality follows from independence and the fact that $m$ must differ from one of $k, j$. Additionally, writing $\tilde{\zeta}_{l, k, j}^{x}(X):=\zeta_{l, k, j}^{x}(\alpha, \sigma, X)-\mathbb{E}\left[\zeta_{l, k, j}^{x}(\alpha, \sigma, X)\right]$ and $\bar{\zeta}_{l, k, j}^{x}:=$ $\mathbb{E}\left[\zeta_{l, k, j}^{x}(\alpha, \sigma, X)\right]$, by independence and our moment assumptions (i.e. Assumption S2)

$$
\begin{aligned}
& \mathbb{E}\left[\left(\tilde{\zeta}_{l, k, j}^{x}(X)\left[\phi_{k}\left(\epsilon_{k}\right) \epsilon_{k}+1\right]+\bar{\zeta}_{l, k, j}^{x}\left[\tau_{k, 1} \epsilon_{k}+\tau_{k, 2} \kappa\left(\epsilon_{k}\right)\right]\right) h_{0}(\tilde{X})\right] \\
& \quad=\mathbb{E}\left[\tilde{\zeta}_{l, k, j}^{x}(X) h_{0}(\tilde{X})\right] \mathbb{E}\left[\phi_{k}\left(\epsilon_{k}\right) \epsilon_{k}+1\right]+\bar{\zeta}_{l, k, j}^{x} \mathbb{E}\left[\tau_{k, 1} \epsilon_{k}+\tau_{k, 2} \kappa\left(\epsilon_{k}\right)\right] \mathbb{E}\left[h_{0}(\tilde{X})\right] \\
& \quad=0
\end{aligned}
$$

and again using independence and the definition of $H_{k}$,

$$
\begin{aligned}
\mathbb{E} & {\left[\left(\tilde{\zeta}_{l, k, j}^{x}(X)\left[\phi_{k}\left(\epsilon_{k}\right) \epsilon_{k}+1\right]+\bar{\zeta}_{l, k, j}^{x}\left[\tau_{k, 1} \epsilon_{k}+\tau_{k, 2} \kappa\left(\epsilon_{k}\right)\right]\right) h_{j}\left(\epsilon_{j}\right)\right] } \\
& =\mathbb{E}\left[\tilde{\zeta}_{l, k, j}^{x}(X)\right] \mathbb{E}\left[\left(\phi_{k}\left(\epsilon_{k}\right) \epsilon_{k}+1\right) h_{j}\left(\epsilon_{j}\right)\right]+\bar{\zeta}_{l, k, j}^{x} \mathbb{E}\left[\left(\tau_{k, 1} \epsilon_{k}+\tau_{k, 2} \kappa\left(\epsilon_{k}\right)\right) h_{j}\left(\epsilon_{j}\right)\right] \\
& =0
\end{aligned}
$$

Since $\epsilon_{k}=e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}$, these observations and the form of $\tilde{\ell}_{\theta, x}$ establish (a). For (b), it suffices to show that

$$
f_{k}\left(\epsilon_{k}\right):=\phi_{k}\left(\epsilon_{k}\right) \epsilon_{k}+1-\tau_{k, 1} \epsilon_{k}-\tau_{k, 2} \kappa\left(\epsilon_{k}\right) \in H_{k}
$$

That $\mathbb{E}\left[f_{k}\left(\epsilon_{k}\right)\right]=0$ and $\mathbb{E}\left[f_{k}\left(\epsilon_{k}\right)^{2}\right]<\infty$ follows immediately from Assumption S2. That additionally $\mathbb{E}\left[f_{k}\left(\epsilon_{k}\right) \epsilon_{k}\right]=\mathbb{E}\left[f_{k}\left(\epsilon_{k}\right) \kappa\left(\epsilon_{k}\right)\right]=0$ is ensured by the choice of $\tau_{k}$.

[^2]Case 2: $x=b$. For (a) let $m(X):=A(\alpha, \sigma, X) D_{b, l}(b, X)$ and $\mu=\mathbb{E}[m(X)]$. Then,

$$
\begin{aligned}
\mathbb{E}\left[\phi_{k}\left(\epsilon_{k}\right) e_{k}^{\prime}(m(X)-\mu) h_{0}(\tilde{X})\right] & =\mathbb{E}\left[\phi_{k}\left(\epsilon_{k}\right)\right] \mathbb{E}\left[e_{k}^{\prime}(m(X)-\mu) h_{0}(\tilde{X})\right]=0 \\
\mathbb{E}\left[\phi_{k}\left(\epsilon_{k}\right) e_{k}^{\prime}(m(X)-\mu) h_{j}\left(\epsilon_{j}\right)\right] & =\mathbb{E}\left[\phi_{k}\left(\epsilon_{k}\right) h_{j}\left(\epsilon_{j}\right)\right] \mathbb{E}\left[e_{k}^{\prime}(m(X)-\mu)\right]=0 \\
\mathbb{E}\left[e_{k}^{\prime} \mu\left(\varsigma_{k, 1} \epsilon_{k}+\varsigma_{k, 2} \kappa\left(\epsilon_{k}\right)\right) h_{0}(\tilde{X})\right] & =e_{k}^{\prime} \mu \mathbb{E}\left[\varsigma_{k, 1} \epsilon_{k}+\varsigma_{k, 2} \kappa\left(\epsilon_{k}\right)\right] \mathbb{E}\left[h_{0}(\tilde{X})\right]=0 ;
\end{aligned}
$$

for $k \neq j$ by independence

$$
\mathbb{E}\left[e_{k}^{\prime} \mu\left(\varsigma_{k, 1} \epsilon_{k}+\varsigma_{k, 2} \kappa\left(\epsilon_{k}\right)\right) h_{j}\left(\epsilon_{j}\right)\right]=e_{k}^{\prime} \mu \mathbb{E}\left[\varsigma_{k, 1} \epsilon_{k}+\varsigma_{k, 2} \kappa\left(\epsilon_{k}\right)\right] \mathbb{E}\left[h_{j}\left(\epsilon_{j}\right)\right]=0
$$

whilst for $k=j$, the definition of $H_{k}$ ensures that

$$
\mathbb{E}\left[e_{k}^{\prime} \mu\left(\varsigma_{k, 1} \epsilon_{k}+\varsigma_{k, 2} \kappa\left(\epsilon_{k}\right)\right) h_{k}\left(\epsilon_{k}\right)\right]=e_{k}^{\prime} \mu \mathbb{E}\left[\varsigma_{k, 1} \epsilon_{k} h_{k}\left(\epsilon_{k}\right)+\varsigma_{k, 2} \kappa\left(\epsilon_{k}\right) h_{k}\left(\epsilon_{k}\right)\right]=0
$$

Since $\epsilon_{k}=e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}$, these observations and the form of $\tilde{\ell}_{\theta, b}$ establish (a). For (b) it suffices to show that

$$
q_{k}\left(\epsilon_{k}\right):=\left(\phi_{k}\left(\epsilon_{k}\right)+\varsigma_{k, 1} \epsilon_{k}+\varsigma_{k, 2} \kappa\left(\epsilon_{k}\right)\right)\left(-e_{k}^{\prime} \mu\right) \in H_{k} .
$$

That $\mathbb{E}\left[q_{k}\left(\epsilon_{k}\right)\right]=0$ and $\mathbb{E}\left[q_{k}\left(\epsilon_{k}\right)^{2}\right]<\infty$ follows immediately from Assumption S2. That additionally $\mathbb{E}\left[q_{k}\left(\epsilon_{k}\right) \epsilon_{k}\right]=\mathbb{E}\left[q_{k}\left(\epsilon_{k}\right) \kappa\left(\epsilon_{k}\right)\right]=0$ is ensured by the choice of $\varsigma_{k}$.

## S1.2 Log density score estimation

We work with a high level condition analagous to Assumption 4, adapted to the more general setting of equation (S1).

Assumption S3. Let $\nu_{n}$ be as in Assumption 3. We have estimators $\hat{\phi}_{k, n, \gamma}$ such that for (a) any sequence with elements $\theta_{n}=\left(\alpha_{0}, \beta_{n}, \eta\right) \in \Theta$ where $\left(\beta_{n}\right)_{n \in \mathbb{N}}$ is a deterministic sequence with $\sqrt{n}\left\|\beta_{n}-\beta\right\|=O(1)$ and (b) any array $\left(Z_{n, i}\right)_{n \in \mathbb{N}, i \leq n}$ with i.i.d. rows and such that $\mathbb{E} Z_{n, i}=0, \sup _{n \in \mathbb{N}} \mathbb{E} Z_{n, i}^{2}<\infty$ and $Z_{n, i} \Perp \epsilon_{i, k}$ for each $n, i$, and $k$,

$$
\begin{align*}
& \frac{1}{n} \sum_{i=1}^{n}\left[\hat{\phi}_{k, n, \gamma_{n}}\left(A_{k, \gamma_{n}, i} V_{\theta_{n}, i}\right)-\phi_{k}\left(A_{k, \gamma_{n}, i} V_{\theta_{n}, i}\right)\right] Z_{n, i}=o_{P_{\theta_{n}}^{n}}\left(n^{-1 / 2}\right),  \tag{S6}\\
& \frac{1}{n} \sum_{i=1}^{n}\left(\left[\hat{\phi}_{k, n, \gamma_{n}}\left(A_{k, \gamma_{n}, i} V_{\theta_{n}, i}\right)-\phi_{k}\left(A_{k, \gamma_{n}, i} V_{\theta_{n}, i}\right)\right] Z_{n, i}\right)^{2}=o_{P_{\theta_{n}}^{n}}\left(\nu_{n}\right) . \tag{S7}
\end{align*}
$$

where $A_{k, \gamma_{n}, i}:=e_{k}^{\prime} A\left(\alpha_{0}, \sigma_{n}, X_{i}\right), V_{\theta_{n}, i}:=Y_{i}-B\left(b_{n}, X_{i}\right)$.

We additionally impose the following condition, which is necessary in this more general setup, due to the term

$$
\sum_{k=1}^{K}\left(\zeta_{l, k, k}^{x}(\alpha, \sigma, X)-\mathbb{E}\left[\zeta_{l, k, k}^{x}(\alpha, \sigma, X)\right]\right)\left[\phi_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right) e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}+1\right]
$$

which appears in $\tilde{\ell}_{\theta, x}$ for $x \in\{\alpha, \sigma\}$ when $A(\alpha, \sigma, X)$ depends on $X .{ }^{\mathrm{S} 3}$
Assumption S4. In the context of Assumption S3, additionally

$$
\begin{equation*}
\frac{1}{n} \sum_{i=1}^{n}\left(\left[\hat{\phi}_{k, n, \gamma_{n}}\left(A_{k, \gamma_{n}, i} V_{\theta_{n}, i}\right)-\phi_{k}\left(A_{k, \gamma_{n}, i} V_{\theta_{n}, i}\right)\right] A_{k, \gamma_{n}, i} V_{\theta_{n}, i}\right)^{2}=o_{P_{\theta_{n}}^{n}}\left(\nu_{n}\right) \tag{S8}
\end{equation*}
$$

Lemmas S4 and S5 below demonstrate that the estimator defined in (11) satisfies the high-level conditions in Assumptions S3 and S4 provided Assumption S2 holds along with Assumption 3 and some additional conditions given in the statement of Lemma S5. The proofs of these Lemmas are given in Section $S 5$ below.

Lemma S4. Suppose Assumptions S2 and 3 hold. Then, $\hat{\phi}_{k, n, \gamma}$ as defined in (11) satisfies Assumption S3.

Lemma S5. Suppose Assumptions S2 and 3 hold. Additionally suppose that for some $\mathrm{M}_{k, n} \geq$ $\max \left\{\left|\Xi_{k, n}^{L}\right|,\left|\Xi_{k, n}^{U}\right|\right\}$,

1. $\delta_{k, n}^{-3} \Delta_{k, n} \mathbb{E}\left[\epsilon_{i, k}^{2} \mathbf{1}\left\{\left|\epsilon_{i, k}\right|>\mathrm{M}_{k, n}\right\}\right]=o\left(\nu_{n}\right)$;
2. $\mathbb{E}\left[\epsilon_{i, k}^{4} \boldsymbol{1}\left\{\left|\epsilon_{i, k}\right|>\mathbf{M}_{k, n}\right\}\right]=o\left(\nu_{n}^{2}\right)$;
3. $\mathrm{M}_{k, n}^{2}\left\|\phi_{k, n}^{(3)}\right\|_{\infty}^{2} \delta_{k, n}^{6}=o\left(\nu_{n}\right)$.

Then, $\hat{\phi}_{k, n, \gamma}$ as defined in (11) satisfies Assumption S4.
Remark 2. For $\varrho<\rho$ where $\mathbb{E}\left|\epsilon_{k}\right|^{\rho}<\infty$, one has

$$
\mathbb{E}\left[\left|\epsilon_{k}\right|^{\varrho} \mathbf{1}\left\{\left|\epsilon_{k}\right|>\mathrm{M}_{k, n}\right\}\right]=\mathbb{E}\left[\left|\epsilon_{k}\right|^{\rho}\left|\epsilon_{k}\right|^{\varrho-\rho} \mathbf{1}\left\{\left|\epsilon_{k}\right|>\mathrm{M}_{k, n}\right\}\right] \leq \mathbb{E}\left|\epsilon_{k}\right|^{\rho} \mathrm{M}_{k, n}^{\varrho-\rho},
$$

and thus the speed at which $\mathrm{M}_{k, n}$ is required to increase to satisfy conditions 1, 2 in Lemma $S 5$ decreases with the number of finite moments of $\epsilon_{k}$.

[^3]
## S1.3 The test and its asymptotic properties

Since $\tilde{\ell}_{\theta}$ has a slightly different form in the setting considered in this section (compared to that considered in the main text; compare Lemmas 3 and S3), we amend our estimator $\hat{\ell}_{n, \gamma}$ accordingly. First let $\hat{\tau}_{k, n, \gamma}$ and $\hat{\varsigma}_{k, n, \gamma}$ be given by

$$
\hat{\tau}_{k, n, \gamma}=\hat{M}_{k, n, \gamma}^{-1}\binom{0}{-2}, \hat{\varsigma}_{k, n, \gamma}=\hat{M}_{k, n, \gamma}^{-1}\binom{1}{0}, \hat{M}_{k, n, \gamma}=\frac{1}{n} \sum_{i=1}^{n}\left(\begin{array}{cc}
1 & \left(A_{k, \gamma, i} V_{\gamma, i}\right)^{3} \\
\left(A_{k, \gamma, i} V_{\gamma, i}\right)^{3} & \left(A_{k, \gamma, i} V_{\gamma, i}\right)^{4}-1
\end{array}\right) .
$$

The estimators for the components corresponding to $\alpha$ and $\sigma$ are:

$$
\begin{align*}
\hat{\ell}_{n, \gamma, \alpha, l}\left(W_{i}\right):= & \sum_{k=1}^{K} \sum_{j=1, j \neq k}^{K} \zeta_{l, k, j, \gamma, i}^{\alpha} \hat{\phi}_{k, n, \gamma}\left(A_{k, \gamma, i} V_{\gamma, i}\right) A_{j, \gamma, i} V_{\gamma, i} \\
& +\sum_{k=1}^{K}\left(\zeta_{l, k, k, \gamma, i}^{\alpha}-\bar{\zeta}_{l, k, k, n, \gamma}^{\alpha}\right)\left(\hat{\phi}_{k, n, \gamma}\left(A_{k, \gamma, i} V_{\gamma, i}\right) A_{k, \gamma, 2} V_{\gamma, i}+1\right) \\
& +\sum_{k=1}^{K} \bar{\zeta}_{l, k, k, n, \gamma}^{\alpha}\left(\hat{\tau}_{k, n, \gamma, 1} A_{k, \gamma, i} V_{\gamma, i}+\hat{\tau}_{k, n, \gamma, 2} \kappa\left(A_{k, \gamma, i} V_{\gamma, i}\right)\right) ; \\
\hat{\ell}_{n, \gamma, \sigma, l}\left(W_{i}\right):= & \sum_{k=1}^{K} \sum_{j=1, j \neq k}^{K} \zeta_{l, k, j, \gamma, i}^{\sigma} \hat{\phi}_{k, n, \gamma}\left(A_{k, \gamma, i} V_{\gamma, i}\right) A_{j, \gamma, i} V_{\gamma, i}  \tag{S9}\\
& +\sum_{k=1}^{K}\left(\zeta_{l, k, k, \gamma, i}^{\sigma}-\bar{\zeta}_{l, k, k, n, \gamma}^{\sigma}\right)\left(\hat{\phi}_{k, n, \gamma}\left(A_{k, \gamma, i} V_{\gamma, i}\right) A_{k, \gamma, i} V_{\gamma, i}+1\right) \\
& +\sum_{k=1}^{K} \bar{\zeta}_{l, k, k, n, \gamma}^{\sigma}\left(\hat{\tau}_{k, n, \gamma, 1} A_{k, \gamma, i} V_{\gamma, i}+\hat{\tau}_{k, n, \gamma, 2} \kappa\left(A_{k, \gamma, i} V_{\gamma, i}\right)\right) ;
\end{align*}
$$

with $\zeta_{l, k, j, \gamma, i}^{\alpha}:=\zeta_{l, k, j}^{\alpha}\left(\alpha, \sigma, X_{i}\right), \bar{\zeta}_{l, k, j, n, \gamma}^{\alpha}:=\frac{1}{n} \sum_{i=1}^{n} \zeta_{l, k, j, \gamma, i}^{\alpha}, A_{k, \gamma, i}:=e_{k}^{\prime} A\left(\alpha, \sigma, X_{i}\right), V_{\gamma, i}:=V_{\theta, i}:=$ $Y_{i}-B X_{i}, \bar{X}_{n}:=\frac{1}{n} \sum_{i=1}^{n} X_{i}$. For the components corresponding to $b$,

$$
\begin{align*}
\hat{\ell}_{n, \gamma, b}\left(W_{i}\right):= & -\sum_{k=1}^{K} \hat{\phi}_{k, n, \gamma}\left(A_{k, \gamma, i} V_{\gamma, i}\right)\left(A_{k, \gamma, i}\left(X_{i}^{\prime} \otimes I_{K}\right)-\frac{1}{n} \sum_{i=1}^{n}\left[A_{k, \gamma, i}\left(X_{i}^{\prime} \otimes I_{K}\right)\right]\right) \\
& +\sum_{k=1}^{K}\left(\frac{1}{n} \sum_{i=1}^{n}\left[A_{k, \gamma, i}\left(X_{i}^{\prime} \otimes I_{K}\right)\right]\right)\left(\hat{\varsigma}_{k, n, \gamma, 1} A_{k, \gamma, i} V_{\gamma, i}+\hat{\varsigma}_{k, n, \gamma, 2} \kappa\left(A_{k, \gamma, i} V_{\gamma, i}\right)\right) . \tag{S10}
\end{align*}
$$

The estimator $\hat{I}_{n, \gamma}$ is given by

$$
\hat{I}_{n, \gamma}:=\frac{1}{n} \sum_{i=1}^{n} \hat{\ell}_{n, \gamma}\left(W_{i}\right) \hat{\ell}_{n, \gamma}\left(W_{i}\right)^{\prime}
$$

Remark 3. If $A(\alpha, \sigma, X)=A(\alpha, \sigma)$ and $B(b, X)=\operatorname{vec}^{-1}(b) X$ (as considered in the main text), the estimators given in (S9) and (S10) are numerically identical to those in (9).
$\hat{S}_{\gamma}$ is then defined as in (14) and we have the following Theorem (cf. Theorem 1), the proof of which is analogous to that of Theorem 1.

Theorem S1. Suppose that Assumptions S1, S2, S3 and S4 hold and suppose that $\beta$ is an interior point of $\mathcal{B}$. Let $r_{n}=\operatorname{rank}\left(\hat{\mathcal{I}}_{\bar{\gamma}}^{t}\right)$ and denote by $c_{n}$ the $1-a$ quantile of the $\chi_{r_{n}}^{2}$ distribution, for any $a \in(0,1)$. Then

$$
\limsup _{n \rightarrow \infty} \sup _{\theta \in \Theta_{0, n}} P_{\theta}\left(\hat{S}_{\bar{\gamma}}>c_{n}\right) \leq a
$$

with inequality only if $\operatorname{rank}\left(\tilde{\mathcal{I}}_{\theta_{0}}\right)=0$ where $\theta_{0}=\left(\alpha_{0}, \beta, \eta\right)$.
Proof. It suffices to show the conditions of Corollary 1 hold. There are 5 conditions which we verify in order: items $1,2,3 \&$ equation (38) of the statement of Theorem 2.

Condition 1: This follows verbatim as the demonstration of Condition 1 in the proof of Theorem 1 on replacing Lemma 1 with Lemma S1.

Condition 2: This follows by repeated addition and subtraction along with the convergence in probability and stochastic boundedness results of Lemma S11, the moment conditions in Assumption S2 and the local boundedness given by Assumption S1 Part 4.

Condition 3: This follows verbatim as the demonstration of Condition 3 in the proof of Theorem 1 on replacing "the local Lipschitz continuity of each $\beta \mapsto \zeta_{l, j, k}^{x}(\alpha, \sigma)$ and $\beta \mapsto$ $A(\alpha, \sigma)$ " with "the local Lipschitz continuity of each $\beta \mapsto \zeta_{l, j, k}^{x}(\alpha, \sigma, X)$ and $\beta \mapsto A(\alpha, \sigma, X)$ " and removing the reference to Lemma 4. ${ }^{\text {S4 }}$

Condition 4: This follows verbatim as the demonstration of Condition 4 in the proof of Theorem 1 on replacing Lemmas 1 and 2 with Lemmas S1 and S2.

## S2 Supporting results for the main Theorems

The following supporting results apply to the model introduced in Section S1. The model considered in the main text is a special case of this model with $A(\alpha, \sigma, X)=A(\alpha, \sigma)$ and $B(b, X)=\operatorname{vec}^{-1}(b) X$, for which Assumptions 1, 2 and 4 imply S1, S2 and S3 respectively. In consequence the results in this section apply a fortiori to the case considered in the main text.

[^4]Lemma S6. Suppose that Assumptions S1 and S2 hold and that $(\alpha, \beta)$ is an interior point of $\mathcal{A} \times \mathcal{B}$. Let $\varphi(g, h)=\left(g, \eta_{0} h_{0}, \ldots, \eta_{K} h_{K}\right)$. Then

1. $t \mapsto \sqrt{p_{\theta+t \varphi(g, h)}(w)}$ is (pointwise) continuously differentiable in a neighbourhood $\mathcal{U} \subset$ $[0, \infty)$ of zero. ${ }^{\text {S5 }}$

Moreover, if we define $q_{\theta,(g, h), u}(w):=\left.\frac{\partial \log p_{\theta+t \varphi(g, h)}(w)}{\partial t}\right|_{t=u}$, then
2. $\left\{q_{\theta,(g, h), u}(W)^{2}: u \in \mathcal{V}\right\}$ is uniformly $P_{\theta+u \varphi(g, h)}$ - integrable for some neighbourhood of zero $\mathcal{V} \subset \mathcal{U}$.

Proof. For all sufficiently small $t, \theta+t \varphi(g, h) \in \Theta$; in such an interval, the continuous differentiability follows directly from Assumptions S1 and S2 along with the definition of $H$.

Under $P_{\theta+u \varphi(g, h)}, q_{\theta,(g, h), u}(W)$ has the same law as

$$
\begin{align*}
Z_{u}:= & \frac{h_{0}(\tilde{X})}{1+u h_{0}(\tilde{X})}+\sum_{k=1}^{K} \frac{h_{k}\left(\epsilon_{k}\right)+u h_{k}^{\prime}\left(\epsilon_{k}\right) e_{k}^{\prime}\left[\mathrm{D}_{1, u} V_{\theta+u \varphi(g, h)}+\mathrm{D}_{2, u}\right]}{1+u h_{k}\left(\epsilon_{k}\right)}  \tag{S11}\\
& +\operatorname{tr}\left(A(\alpha+u a, \sigma+u s, X)^{-1} \mathrm{D}_{1, \mathrm{u}}\right)+\sum_{k=1}^{K} \phi_{k}\left(\epsilon_{k}\right) e_{k}^{\prime}\left[\mathrm{D}_{1, u} V_{\theta+u \varphi(g, h)}+\mathrm{D}_{2, u}\right] .
\end{align*}
$$

where

$$
\mathrm{D}_{1, u}:=\sum_{l=1}^{L_{\alpha}} a_{l} D_{\alpha, l}(\alpha+u a, \sigma+u s, X)+\sum_{l=1}^{L_{\sigma}} s_{l} D_{\sigma, l}(\alpha+u a, \sigma+u s, X)
$$

and

$$
\mathrm{D}_{2, u}:=A(\alpha+u a, \sigma+u s, X) \sum_{l=1}^{L_{b}} \varrho_{l} D_{b, l}(b+u \varrho, X)
$$

The definition of $H$ ensures that for all sufficiently small $u$ (i.e. $u \in \mathcal{V}$ ), the denominators $1+u h_{0}(\tilde{X})$ and $1+u h_{k}\left(\epsilon_{k}\right)$ are bounded, as are $h_{0}(\tilde{X}), h_{k}\left(\epsilon_{k}\right)$ and $u h_{k}^{\prime}\left(\epsilon_{k}\right)$. Assumption S1 ensures the same is true of $\mathrm{D}_{1, u}$, the trace term, $A(\alpha+t a, \sigma+t s, X)$ and its inverse. These bounds, along with the finite moments given by Assumption S2 allow the application of Jensen's and Hölder's inequalties to obtain that $\sup _{u \in \mathcal{V}} \mathbb{E}\left|Z_{u}\right|^{2+\delta / 2}<\infty$, implying the claimed uniform integrability.

Lemma S7. Suppose that Assumptions S1 and S2 hold and let $\mathcal{V}=\mathbb{R}^{L} \times H$ be equipped with the norm ${ }^{\mathrm{S} 6}$

$$
\|(g, h)\|:=\sqrt{\|g\|^{2}+\sum_{k=0}^{K}\left\|\tilde{h}_{k}\right\|_{L_{2}\left(P_{\theta}\right)}^{2}} .
$$

[^5]Then, the functions $(g, h) \mapsto \frac{1}{\sqrt{n}} \sum_{i=1}^{n}\left[g^{\prime} \dot{\ell}_{\theta}+\sum_{k=0}^{K} \tilde{h}_{k}\right]$ (i.e. indexed by $n$ ) are equicontinuous on compacts in $L_{2}\left(P_{\theta}\right)$ and the functions $(g, h) \mapsto P_{\theta_{n}(g, h)}^{n}$ (i.e. indexed by $n$ ) are equicontinuous on compacts in the total variation metric.

Proof. For any $(g, h),\left(g^{\star}, h^{\star}\right) \in \mathcal{V}$, by the fact the observations are i.i.d. and any $h \in H$ is mean zero, as is $\dot{\ell}_{\theta}$,

$$
\left\|\frac{1}{\sqrt{n}} \sum_{i=1}^{n}\left[\left(g^{\star}-g\right)^{\prime} \dot{\ell}_{\theta}+\sum_{k=0}^{K}\left(\tilde{h}_{k}^{\star}-\tilde{h}_{k}\right)\right]\right\|_{L_{2}\left(P_{\theta}^{n}\right)}^{2}=\left\|\left(g^{\star}-g\right)^{\prime} \dot{\varphi}_{\theta}+\sum_{k=0}^{K}\left(\tilde{h}_{k}^{\star}-\tilde{h}_{k}\right)\right\|_{L_{2}\left(P_{\theta}\right)}^{2} .
$$

Therefore, left hand side in the display above can be made arbitrarily small, uniformly in $n$, by taking $\left\|\left(g^{\star}, h^{\star}\right)-(g, h)\right\|$ sufficiently small and hence the first claim holds. For the second claim we note that each $(g, h) \mapsto P_{\theta_{n}(g, h)}^{n}$ is continuous by the pointwise continuity of the densities and Scheffés Lemma. Then, let $K \subset \mathcal{V}=\mathbb{R}^{L} \times H$ be compact. We will now show that for any convergent sequence $\left(g_{n}, h_{n}\right) \rightarrow(g, h)$ in $K, d_{T V}\left(P_{\theta_{n}\left(g_{n}, h_{n}\right)}^{n}, P_{\theta_{n}(g, h)}^{n}\right) \rightarrow 0$ as $n \rightarrow \infty .{ }^{\text {S7 }}$ For this, by Lemma S 17 and the triangle inequality, it is sufficient to show that

$$
\begin{equation*}
\log \frac{p_{\theta_{n}\left(g_{n}, h_{n}\right)}^{n}}{p_{\theta_{n}\left(g_{n}, h\right)}}=o_{P_{\theta_{n}\left(g_{n}, h\right)}^{n}}(1), \quad \log \frac{p_{\theta_{n}\left(g_{n}, h\right)}^{n}}{p_{\theta_{n}(g, h)}}=o_{P_{\theta_{n}(g, h)}^{n}}^{n}(1) . \tag{S12}
\end{equation*}
$$

For these we first note that since $h_{k}$ is bounded,

$$
\begin{align*}
\left\|\tilde{h}_{k, n}-\tilde{h}_{k}\right\|_{L_{2}\left(P_{\theta_{n}\left(g_{n}, h\right)}^{n}\right)}^{2} & =\int\left[h_{n, k}(x)-h_{k}(x)\right]^{2} \eta_{k}(x)\left(1+h_{k}(x) / \sqrt{n}\right) \mathrm{d} x  \tag{S13}\\
& \leq\left\|h_{n, k}-h_{k}\right\|_{L_{2}\left(P_{\theta}^{n}\right)}+\left\|h_{n, k}-h_{k}\right\|_{L_{2}\left(P_{\theta}^{n}\right)}\left\|h_{k}\right\|_{L_{\infty}\left(P_{\theta}^{n}\right)} / \sqrt{n}
\end{align*}
$$

Next introduce the notation: ${ }^{\text {S8 }}$

$$
u_{k, n, i}:= \begin{cases}e_{k}^{\prime} A\left(\theta_{n}\left(g_{n}, h\right), X\right) V_{\theta_{n}\left(g_{n}, h\right), i}=e_{k}^{\prime} A\left(\theta_{n}\left(g_{n}, h_{n}\right), X\right) V_{\theta_{n}\left(g_{n}, h_{n}\right), i} & \text { if } k=1, \ldots, K \\ \tilde{X}_{i} & \text { if } k=0\end{cases}
$$

Equation (S13) implies that $\left(\tilde{h}_{k, n}\right)_{n \in \mathbb{N}}$ is uniformly square $P_{\theta_{n}\left(g_{n}, h\right)}^{n}$ integrable, and hence the

[^6]Lindeberg condition holds for $h_{k, n}\left(u_{k, n, i}\right) / \sqrt{n}$. In particular, under $P_{\theta_{n}\left(g_{n}, h\right)}^{n}$,

$$
\begin{aligned}
\lim _{n \rightarrow \infty} & \sum_{i=1}^{n} \mathbb{E}\left[\frac{h_{k, n}\left(u_{k, n, i}\right)^{2}}{n} \mathbf{1}\left\{\left|h_{k, n}\left(u_{k, n, i}\right)\right|>\delta \sqrt{n}\right\}\right] \\
& =\lim _{n \rightarrow \infty} \frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left[h_{k, n}\left(u_{k, n, i}\right)^{2} \mathbf{1}\left\{\left|h_{k, n}\left(u_{k, n, i}\right)\right|>\delta \sqrt{n}\right\}\right] \\
& =\lim _{n \rightarrow \infty} \mathbb{E}\left[h_{k, n}\left(u_{k, n, i}\right)^{2} \mathbf{1}\left\{\left|h_{k, n}\left(u_{k, n, i}\right)\right|>\delta \sqrt{n}\right\}\right] \\
& =0
\end{aligned}
$$

for any $\delta>0$. This implies uniform asymptotic negligability (e.g. Gut, 2005, Remark 7.2.4):

$$
\begin{equation*}
\max _{1 \leq i \leq n} \frac{\left|h_{k, n}\left(u_{k, n, i}\right)\right|}{\sqrt{n}} \xrightarrow{P_{\theta_{n}\left(g_{n}, h\right)}^{n}} 0 . \tag{S14}
\end{equation*}
$$

Then, to prove the first claim in (S12) observe

$$
\log \frac{p_{\theta_{n}\left(g_{n}, h_{n}\right)}^{n}}{p_{\theta_{n}\left(g_{n}, h\right)}^{n}}=\sum_{k=0}^{K} \sum_{t=1}^{n} \log \left(1+h_{k, n}\left(u_{k, n, i}\right) / \sqrt{n}\right)-\log \left(1+h_{k}\left(u_{k, n, i}\right) / \sqrt{n}\right)
$$

hence it suffices to show that each

$$
l_{n, k}:=\sum_{t=1}^{n} \log \left(1+h_{k, n}\left(u_{k, n, i}\right) / \sqrt{n}\right)-\log \left(1+h_{k}\left(u_{k, n, i}\right) / \sqrt{n}\right) \xrightarrow{P_{\theta_{n}\left(g_{n}, h\right)}^{n}} 0 .
$$

Let $\varepsilon \in(0,1)$ be fixed and define

$$
\begin{aligned}
E_{n} & :=\left\{\max _{1 \leq i \leq n}\left|h_{k, n}\left(u_{k, n, i}\right)\right| / \sqrt{n} \leq \varepsilon\right\} ; \\
F_{n} & :=\left\{\max _{1 \leq i \leq n}\left|h_{k}\left(u_{k, n, i}\right)\right| / \sqrt{n} \leq \varepsilon\right\} .
\end{aligned}
$$

Since $h_{k}$ is bounded, $P_{\theta_{n}\left(g_{n}, h\right)}^{n} F_{n} \rightarrow 1 ; P_{\theta_{n}\left(g_{n}, h\right)}^{n} E_{n} \rightarrow 1$ follows from equation S14. Hence $P_{\theta_{n}\left(g_{n}, h\right)}^{n} F_{n} \cap E_{n} \rightarrow 1$. On $E_{n} \cap F_{n}$ we can perform a two-term Taylor expansion of $\log (1+x)$
to obtain

$$
\begin{aligned}
& \log \left(1+h_{k, n}\left(u_{k, n, i}\right) / \sqrt{n}\right)-\log \left(1+h_{k}\left(u_{k, n, i}\right) / \sqrt{n}\right) \\
& =\frac{h_{k, n}\left(u_{k, n, i}\right)}{\sqrt{n}}-\frac{1}{2} \frac{h_{k, n}\left(u_{k, n, i}\right)^{2}}{n}-\frac{h_{k}\left(u_{k, n, i}\right)}{\sqrt{n}}+\frac{1}{2} \frac{h_{k}\left(u_{k, n, i}\right)^{2}}{n} \\
& \quad+R\left(\frac{h_{k, n}\left(u_{k, n, i}\right)}{\sqrt{n}}\right)-R\left(\frac{h_{k}\left(u_{k, n, i}\right)}{\sqrt{n}}\right),
\end{aligned}
$$

where $|R(x)| \leq|x|^{3}$. It follows that

$$
\begin{aligned}
l_{n, k}=\frac{1}{\sqrt{n}} & \sum_{i=1}^{n} h_{k, n}\left(u_{k, n, i}\right)-h_{k}\left(u_{k, n, i}\right)-\frac{1}{2} \frac{1}{n} \sum_{i=1}^{n}\left[h_{k, n}\left(u_{k, n, i}\right)^{2}-h_{k}\left(u_{k, n, i}\right)^{2}\right] \\
& +\sum_{i=1}^{n} R\left(\frac{h_{k, n}\left(u_{k, n, i}\right)}{\sqrt{n}}\right)-R\left(\frac{h_{k}\left(u_{k, n, i}\right)}{\sqrt{n}}\right) .
\end{aligned}
$$

We will show that the remainder terms vanish. In particular, one has

$$
\sum_{i=1}^{n}\left|R\left(\frac{h_{k, n}\left(u_{k, n, i}\right)}{\sqrt{n}}\right)\right| \leq \sum_{i=1}^{n}\left|\frac{h_{k, n}\left(u_{k, n, i}\right)}{\sqrt{n}}\right|\left|\frac{h_{k, n}\left(u_{k, n, i}\right)^{2}}{n}\right| \leq \max _{1 \leq i \leq n} \frac{\left|h_{k, n}\left(u_{k, n, i}\right)\right|}{\sqrt{n}} \frac{1}{n} \sum_{i=1}^{n} h_{k, n}\left(u_{k, n, i}\right)^{2} .
$$

By Markov's inequality and equations (S13), (S14), this converges to zero in $P_{\theta_{n}\left(g_{n}, h\right)}^{n}$ probability. The same evidently holds for the case where $h_{k, n}=h_{k}$ for each $n \in \mathbb{N}$. Thus,

$$
l_{n, k}=\frac{1}{\sqrt{n}} \sum_{i=1}^{n} h_{k, n}\left(u_{k, n, i}\right)-h_{k}\left(u_{k, n, i}\right)-\frac{1}{2} \frac{1}{n} \sum_{i=1}^{n}\left[h_{k, n}\left(u_{k, n, i}\right)^{2}-h_{k}\left(u_{k, n, i}\right)^{2}\right]+o_{P_{\theta_{n}\left(g_{n}, h\right)}^{n}}(1),
$$

and it remains to show that $\frac{1}{\sqrt{n}} \sum_{i=1}^{n} h_{k, n}\left(u_{k, n, i}\right)-h_{k}\left(u_{k, n, i}\right)$ and $\frac{1}{n} \sum_{i=1}^{n}\left[h_{k, n}\left(u_{k, n, i}\right)^{2}-h_{k}\left(u_{k, n, i}\right)^{2}\right]$ also converge to zero in probability under $P_{\theta_{n}\left(g_{n}, h\right)}^{n}$. The second of these follows directly from (S13), Markov's inequality and the reverse triangle inequality since

$$
\begin{aligned}
P_{\theta_{n}\left(g_{n}, h\right)}^{n}\left(\left|\frac{1}{n} \sum_{i=1}^{n}\left[h_{k, n}\left(u_{k, n, i}\right)^{2}-h_{k}\left(u_{k, n, i}\right)^{2}\right]\right|>\varepsilon\right) & \leq \varepsilon^{-1} \frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left[h_{k, n}\left(u_{k, n, i}\right)^{2}-h_{k}\left(u_{k, n, i}\right)^{2}\right] \\
& =\varepsilon^{-1} \mathbb{E}\left[h_{k, n}\left(u_{k, n, i}\right)^{2}-h_{k}\left(u_{k, n, i}\right)^{2}\right] \\
& \rightarrow 0
\end{aligned}
$$

For the remaining term, we start by noting that

$$
\mathbb{E}\left[h_{k, n}\left(u_{k, n, i}\right)-h_{k}\left(u_{k, n, i}\right)\right]=\frac{\mathbb{E}\left[\left(h_{k, n}\left(\epsilon_{k}\right)-h_{k}\left(\epsilon_{k}\right)\right) h_{k}\left(\epsilon_{k}\right)\right]}{\sqrt{n}}
$$

$$
\left|\frac{1}{\sqrt{n}} \sum_{i=1}^{n} \mathbb{E}\left[h_{k, n}\left(u_{k, n, i}\right)\right]-\mathbb{E}\left[h_{k}\left(u_{k, n, i}\right)\right]\right| \leq \frac{1}{n} \sum_{i=1}^{n}\left\|h_{k, n}-h_{k}\right\|_{L_{2}\left(P_{\theta}^{n}\right)}\left\|h_{k}\right\|_{L_{2}\left(P_{\theta}^{n}\right)} \rightarrow 0 .
$$

Thus it suffices to show that

$$
\frac{1}{\sqrt{n}} \sum_{i=1}^{n} h_{k, n}\left(u_{k, n, i}\right)-h_{k}\left(u_{k, n, i}\right) \xrightarrow{P_{\theta_{n}\left(g_{n}, h\right)}^{n}} 0
$$

for $h_{k, n}\left(u_{k, n, i}\right):=h_{k, n}\left(u_{k, n, i}\right)-\mathbb{E}\left[h_{k, n}\left(u_{k, n, i}\right)\right]$ and $h_{k}\left(u_{k, n, i}\right):=h_{k, n}\left(u_{k, n, i}\right)-\mathbb{E}\left[h_{k}\left(u_{k, n, i}\right)\right]$. By the reverse triangle inequality and (S13),

$$
\mathbb{E}\left[\left(h_{k, n}\left(u_{k, n, i}\right)-h_{k}\left(u_{k, n, i}\right)\right)^{2}\right] \rightarrow 0, \quad \text { uniformly in } i .
$$

Using this, the independence of the $W_{i}$ and Markov's inequality:

$$
P_{\theta_{n}\left(g_{n}, h\right)}^{n}\left(\left|\frac{1}{\sqrt{n}} \sum_{i=1}^{n} h_{k, n}\left(u_{k, n, i}\right)-h_{k}\left(u_{k, n, i}\right)\right|>\varepsilon\right) \leq \frac{1}{\varepsilon^{2}} \frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left[\left(h_{k, n}\left(u_{k, n, i}\right)-h_{k}\left(u_{k, n, i}\right)\right)^{2}\right] \rightarrow 0
$$

This establishes that $\sum_{k=1}^{K} l_{n, k} \xrightarrow{P_{\theta_{n}\left(g_{n}, h\right)}^{n}} 0$, as required.
For the second condition in (S12), by Lemma S2 part $3 P_{\theta_{n}(g, h)}^{n} \triangleleft \triangleright P_{\theta}^{n} .{ }^{\text {S9 }}$ Hence it suffices to show that $\log \frac{p_{\theta_{n}\left(g_{n}, h\right)}^{n}}{p_{\theta_{n}(g, h)}}=o_{P_{\theta}^{n}}(1)$. We first show that,

$$
\begin{aligned}
\log \frac{p_{\theta_{n}\left(g_{n}, 0\right)}^{n}}{p_{\theta}^{n}} & =\frac{1}{\sqrt{n}} \sum_{t=1}^{n} g^{\prime} \dot{\ell}_{\theta}\left(W_{i}\right)-\mathbb{E}\left(\frac{1}{\sqrt{n}} \sum_{t=1}^{n} g^{\prime} \dot{\ell}_{\theta}\left(W_{i}\right)\right)^{2}+o_{P_{\theta}^{n}}(1) \\
\log \frac{p_{\theta_{n}(g, 0)}^{n}}{p_{\theta}^{n}} & =\frac{1}{\sqrt{n}} \sum_{t=1}^{n} g^{\prime} \dot{\ell}_{\theta}\left(W_{i}\right)-\mathbb{E}\left(\frac{1}{\sqrt{n}} \sum_{t=1}^{n} g^{\prime} \dot{\ell}_{\theta}\left(W_{i}\right)\right)^{2}+o_{P_{\theta}^{n}}(1)
\end{aligned}
$$

where the expectations are taken under $P_{\theta}^{n}$. Here we may proceed analogously to Lemma S1. In particular, by an argument analogous to that showing condition 1 in Lemma S 6 , $g \mapsto \sqrt{p_{\theta_{n}(g, 0)}}$ is continuously differentiable, whilst an argument analogous to that showing condition 2 in Lemma 56 yields that $\left\{q_{\theta,(g, 0)}(W)^{2}: g \in \mathcal{U}\right\}$ is uniformly $P_{\theta+\varphi(g, 0)}$ - integrable for some neighbourhood $\mathcal{U} \subset \mathbb{R}^{L}$ of 0 . Application of Lemma 7.6 and Theorem 7.2 in van der Vaart (1998) then yields the two likelihood expansions in the display above. To complete

[^7]the proof set
$$
\tilde{u}_{k, n, i}:=e_{k}^{\prime} A\left(\theta_{n}\left(g_{n}, h\right), X\right) V_{\theta_{n}\left(g_{n}, h\right), i}, \quad u_{k, n, i}:=e_{k}^{\prime} A\left(\theta_{n}(g, h), X\right) V_{\theta_{n}(g, h), i},
$$
and observe that
\[

$$
\begin{aligned}
& \log \frac{p_{\theta_{n}\left(g_{n}, h\right)}^{n}}{p_{\theta_{n}(g, h)}^{n}}-\left[\log \frac{p_{\theta_{n}\left(g_{n}, 0\right)}^{n}}{p_{\theta}^{n}}-\log \frac{p_{\theta_{n}(g, 0)}^{n}}{p_{\theta}^{n}}\right] \\
& \quad=\sum_{k=1}^{K} \sum_{i=1}^{n} \log \left(1+\frac{h_{k}\left(\tilde{u}_{k, n, i}\right)}{\sqrt{n}}\right)-\log \left(1+\frac{h_{k}\left(u_{k, n, i}\right)}{\sqrt{n}}\right),
\end{aligned}
$$
\]

where the bracketed term is $o_{P_{\theta}^{n}}(1)$ by the preceding argument. Hence it suffices to show that an arbitrary $k$-th element of the outer sum on the right hand side is also $o_{P_{\theta}^{n}}(1)$. Let $\varepsilon \in(0,1)$ be fixed and define

$$
E_{n}:=\left\{\max _{1 \leq i \leq n}\left|h_{k}\left(\tilde{u}_{k, n, i}\right)\right| / \sqrt{n} \leq \varepsilon\right\}, \quad F_{n}:=\left\{\max _{1 \leq i \leq n}\left|h_{k}\left(u_{k, n, i}\right)\right| / \sqrt{n} \leq \varepsilon\right\}
$$

Since $h_{k}$ is bounded $P_{\theta}^{n}\left(E_{n} \cap F_{n}\right) \rightarrow 1$. On this set we may perform a two-term Taylor expansion of $\log (1+x)$ to obtain

$$
\begin{aligned}
& \log \left(1+\frac{h_{k}\left(\tilde{u}_{k, n, i}\right)}{\sqrt{n}}\right)-\log \left(1+\frac{h_{k}\left(u_{k, n, i}\right)}{\sqrt{n}}\right) \\
& \quad=\frac{h_{k}\left(\tilde{u}_{k, n, i}\right)-h_{k}\left(u_{k, n, i}\right)}{\sqrt{n}}-\frac{1}{2} \frac{h_{k}\left(\tilde{u}_{k, n, i}\right)^{2}-h_{k}\left(u_{k, n, i}\right)^{2}}{n}+R\left(\frac{h_{k}\left(\tilde{u}_{k, n, i}\right)}{\sqrt{n}}\right)-R\left(\frac{h_{k}\left(u_{k, n, i}\right)}{\sqrt{n}}\right),
\end{aligned}
$$

where $|R(x)| \leq|x|^{3}$. For the remainder terms one has for any $u_{i}$,

$$
\sum_{i=1}^{n}\left|R\left(\frac{h_{k}\left(u_{i}\right)}{\sqrt{n}}\right)\right| \leq \max _{1 \leq i \leq n} \frac{h_{k}\left(u_{i}\right)}{\sqrt{n}} \frac{1}{n} \sum_{i=1}^{n} h_{k}\left(u_{i}\right)^{2} \lesssim \frac{1}{\sqrt{n}},
$$

since $h_{k}$ is bounded. For the first term in Taylor expansion, note that the derivative (in $\theta, \sigma$ ) of $A(\theta, \sigma, X)$ is bounded on a neighbourhood of $(\theta, \sigma)$ (by Assumption S1). Combine this with the boundedness of $h_{k}^{\prime}$ and the mean value theorem to conclude that

$$
\left|h_{k}\left(\tilde{u}_{k, n, i}\right)-h_{k}\left(u_{k, n, i}\right)\right| \lesssim n^{-1 / 2}\left\|g_{n}-g\right\|\left[\left\|\epsilon_{i}\right\|+\sqrt{\sum_{l=1}^{L_{b}} D_{b, l}\left(b+\varrho_{l, n}, X_{i}\right)^{2}}\right]
$$

for some $\varrho_{l, n}$ with $\left\|\varrho_{l, n}\right\| \leq\left\|g_{n}-g\right\|$. Since $h_{k}$ is bounded,

$$
\left|h_{k}\left(\tilde{u}_{k, n, i}\right)^{2}-h_{k}\left(u_{k, n, i}\right)^{2}\right| \lesssim n^{-1 / 2}\left\|g_{n}-g\right\|\left[\left\|\epsilon_{i}\right\|+\sqrt{\sum_{l=1}^{L_{b}} D_{b, l}\left(b+\varrho_{l, n}, X_{i}\right)^{2}}\right]
$$

Therefore, using the moment bounds in Assumption S2 parts 1 and 4

$$
\begin{aligned}
& \sum_{i=1}^{n}\left|\frac{h_{k}\left(\tilde{u}_{k, n, i}\right)-h_{k}\left(u_{k, n, i}\right)}{\sqrt{n}}-\frac{1}{2} \frac{h_{k}\left(\tilde{u}_{k, n, i}\right)^{2}-h_{k}\left(u_{k, n, i}\right)^{2}}{n}\right| \\
& \quad \lesssim\left\|g_{n}-g\right\|\left(1+\frac{1}{\sqrt{n}}\right) \frac{1}{n} \sum_{i=1}^{n}\left[\left\|\epsilon_{i}\right\|+\sqrt{\sum_{l=1}^{L_{b}} D_{b, l}\left(b+\varrho_{l, n}, X_{i}\right)^{2}}\right]=o_{P_{\theta}^{n}}(1) .
\end{aligned}
$$

This completes the demonstration of (S12) and hence the proof.
Lemma S8. Suppose that Assumptions S1 and S2 hold. Then,

1. cl $H_{0}$ is the space of functions $h_{0}: \mathbb{R}^{d-1} \rightarrow \mathbb{R}$ such that $\mathbb{E} h_{0}\left(\tilde{X}_{i}\right)^{2}<\infty, \mathbb{E} h_{0}(\tilde{X})=0$;
2. For $k=1, \ldots, K, \operatorname{cl} H_{k}$ is the space of functions $h_{k}: \mathbb{R} \rightarrow \mathbb{R}$ such that $\mathbb{E} h_{k}\left(\epsilon_{k}\right)^{2}<\infty$,

$$
\mathbb{E}\left[h_{k}\left(\epsilon_{k}\right)\right]=\mathbb{E}\left[\epsilon_{k} h_{k}\left(\epsilon_{k}\right)\right]=\mathbb{E}\left[\kappa\left(\epsilon_{k}\right) h_{k}\left(\epsilon_{i, k}\right)\right]=0
$$

Additionally, define $H_{0}^{\star}$ as the space of functions $\tilde{h}_{0}(W):=h_{0}(\tilde{X})$ for $h_{0} \in \operatorname{cl} \tilde{H}_{0}$ and $H_{k}^{\star}$ as the space of functions $\tilde{h}_{k}(W):=h_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right)$ for $h_{k} \in \operatorname{cl} \tilde{H}_{k}(k=1, \ldots, K)$. Then

$$
H^{\star}:=H_{0}^{\star}+\cdots+H_{K}^{\star} \subset L_{2}\left(P_{\theta}\right) \quad \text { and } \quad H^{\star}=\operatorname{cl}\left(\tilde{H}_{0}+\cdots+\tilde{H}_{K}\right)
$$

Proof. For $1 \& 2$ let $H_{k}^{*}$ denote the set of functions described in the statement (for $k=$ $0, \ldots, K)$. Clearly any convergent sequence in this space has a limit also in this space and hence $H_{k}^{*}$ is closed. For any $h_{k} \in \tilde{H}_{k}^{*}$ there is a sequence $\left(h_{k, n}\right)_{n \in \mathbb{N}}$ such that each $h_{k, n} \in H_{k}$ and $h_{k, n} \rightarrow h_{k}$ in squared mean (e.g. Newey, 1991, Lemma C.7) and hence cl $H_{k}=H_{k}^{*}$. . 10

For the second part, the first claim follows since $e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}$ has the same law as $\epsilon_{k}$ under $P_{\theta}$ and hence each $P_{\theta}\left[\tilde{h}_{k}(W)^{2}\right]<\infty$. For the second claim, as $\tilde{X}, \epsilon_{1}, \ldots, \epsilon_{K}$ are independent, $\tilde{H}_{0}^{\star}, \ldots, \tilde{H}_{K}^{\star}$ are pairwise orthogonal. As the (finite) sum of closed pairwise orthogonal subspaces is closed (e.g. Conway, 1985, p. 39) we have that $\operatorname{cl}\left(\tilde{H}_{0}+\ldots+\tilde{H}_{K}\right) \subset$ $H^{\star}$. For the reverse inclusion let $\tilde{h}=\sum_{k=0}^{K} \tilde{h}_{k} \in H^{\star}$. By the definition of $H^{\star}$ there are

[^8] $1>\mathbb{E}\left(\epsilon_{k}^{3}\right)^{2}$ imposed in Assumption S2.
$\tilde{h}_{0, n}(W):=h_{0, n}(\tilde{X})$ such that $\tilde{h}_{0, n} \in \tilde{H}_{0}$ and $P_{\theta}\left[\tilde{h}_{0, n}(W)-\tilde{h}_{0}(W)\right]^{2} \rightarrow 0$ and $\tilde{h}_{k, n}(W):=$ $h_{k, n}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right)$ such that $\tilde{h}_{k, n} \in H_{k}$ and $P_{\theta}\left[\tilde{h}_{0, k}(W)-\tilde{h}_{k}(W)\right]^{2} \rightarrow 0$. Hence $\tilde{h}_{n}:=$ $\sum_{k=0}^{K} \tilde{h}_{k, n} \in \tilde{H}_{0}+\ldots+\tilde{H}_{K}$ and converges to $\tilde{h}$, implying that $\tilde{h} \in \operatorname{cl}\left(\tilde{H}_{0}+\cdots+\tilde{H}_{K}\right)$.

Lemma S9. Suppose that Assumptions S1 and S2 hold. Then $\sup _{n \in \mathbb{N}} P_{\tilde{\theta}_{n}}\left\|\tilde{\ell}_{\tilde{\theta}_{n}}\right\|^{2+\delta / 2}<\infty$ and hence $\left(\left\|\tilde{\ell}_{\tilde{\theta}_{n}}\right\|^{2}\right)_{n \in \mathbb{N}}$ is uniformly $P_{\tilde{\theta}_{n}}$-integrable.

Proof. As each component of $\tilde{\ell}_{\tilde{\theta}_{n}}$ lies in $L_{2}\left(P_{\tilde{\theta}_{n}}\right)$ by its definition as an orthogonal projection, it suffices to show that $\lim \sup _{n \in \mathbb{N}} P_{\tilde{\theta}_{n}}\left[\left\|\tilde{\tilde{\theta}}_{\tilde{\theta}_{n}}\right\|^{2+\delta / 2}\right]<\infty$. Let $d_{n}:=\left(b_{n}, s_{n}\right):=\sqrt{n}\left(\beta_{n}-\beta\right)$, with $b_{n} \in \mathbb{R}^{L_{b}}$ and $s_{n} \in \mathbb{R}^{L_{\sigma}}$, so that $\tilde{\theta}_{n}=\theta_{n}\left(g_{n}, 0\right)$ with $g_{n}=\left(0, b_{n}, s_{n}\right)$. Then, under $P_{\tilde{\theta}_{n}}, e_{k}^{\prime} A\left(\alpha, \sigma+s_{n} / \sqrt{n}, X\right) V_{\tilde{\theta}_{n}}$ has the same law as $\epsilon_{k}$. This, along with the observations that $\mathbb{E}\left[\left|\phi_{k}\left(\epsilon_{k}\right)\right|^{4+\delta}\right]<\infty, \mathbb{E}\left|\epsilon_{k}\right|^{4+\delta}$ (both for $\left.k=1, \ldots, K\right), \mathbb{E}\left[\left\|D_{b, l}(b, X)\right\|^{4+\delta}\right]<\infty$ and the local boundedness conditions in Assumption S1 part 4 allow the application of Jensen's and Hölder's inequalities to conclude that $\lim \sup _{n \in \mathbb{N}} P_{\tilde{\theta}_{n}}\left[\left\|\tilde{\ell}_{\tilde{\theta}_{n}}\right\|^{2+\delta / 2}\right]<\infty$ as desired.

Lemma S10. Suppose that Assumptions S1 and S2 hold. Then,

$$
\lim _{n \rightarrow \infty} \int\left\|\tilde{\ell}_{\tilde{\theta}_{n}} \sqrt{p_{\tilde{\theta}_{n}}}-\tilde{\ell}_{\theta} \sqrt{p_{\theta}}\right\|^{2} \mathrm{~d} \lambda=0
$$

Proof. Re-write the integral as

$$
\begin{equation*}
\int\left\|\tilde{\ell}_{\tilde{\theta}_{n}} \sqrt{p_{\tilde{\theta}_{n}}}-\tilde{\ell}_{\theta} \sqrt{p_{\theta}}\right\|^{2} \mathrm{~d} \lambda=\sum_{l=1}^{L} \int\left[\tilde{\ell}_{\tilde{\theta}_{n}, l} \sqrt{p_{\tilde{\theta}_{n}}}-\tilde{\ell}_{\theta, l} \sqrt{p_{\theta}}\right]^{2} \mathrm{~d} \lambda . \tag{S15}
\end{equation*}
$$

It is evidently sufficient to show that each of the integrals in the sum on the rhs converges to zero. For this note that inspection of the forms of $\tilde{\ell}_{\theta}$ and $p_{\theta}$ reveals that $\tilde{\ell}_{\tilde{\theta}_{n}} \rightarrow \tilde{\ell}_{\theta}$ and $p_{\tilde{\theta}_{n}} \rightarrow p_{\theta}$ pointwise. Hence each $\tilde{\ell}_{\tilde{\theta}_{n}, l} \sqrt{p_{\tilde{\theta}_{n}}} \rightarrow \tilde{\ell}_{\theta, l} \sqrt{p_{\theta}}$ pointwise and, by Scheffé's Lemma, $P_{\tilde{\theta}_{n}} \xrightarrow{T V} P_{\theta}$. Combine this observation with Lemma $S 9$ and Corollary 2.9 in Feinberg, Kasyanov and Zgurovsky (2016) to obtain $\lim _{n \rightarrow \infty} \int\left|\tilde{\ell}_{\tilde{\theta}_{n}, l} \sqrt{p_{\tilde{\theta}_{n}}}\right|^{2} \mathrm{~d} \lambda=\int\left|\tilde{\ell}_{\theta, l} \sqrt{p_{\theta}}\right|^{2}<\infty$. Apply Proposition 2.29 in van der Vaart (1998) to conclude.

Lemma S11. Suppose that Assumptions S1, S2 and S3 hold. Then, for each $(k, j)$ with $k \neq j$, each $l$, each $x \in\{\alpha, \sigma\}$ and each $\varrho \in\{\tau, \varsigma\}$, the following terms are $o_{P_{\bar{\theta}_{n}}^{n}}(1)$ :

1. $\bar{\zeta}_{l, k, k, n, \gamma_{n}}^{x}-P_{\tilde{\theta}_{n}}\left[\zeta_{l, k, k, \gamma_{n}, i}^{x}\right]$;
2. $\frac{1}{\sqrt{n}} \sum_{i=1}^{n}\left(\phi_{k}\left(A_{k, \gamma_{n}} V_{\gamma_{n}, i}\right)-\hat{\phi}_{k, n, \gamma_{n}}\left(A_{k, \gamma_{n} i} V_{\gamma_{n}, i}\right)\right) \zeta_{l, k, j, i}^{x} A_{j, \gamma_{n}, i} V_{\gamma_{n}, i}$;
3. $\frac{1}{\sqrt{n}} \sum_{i=1}^{n}\left(\phi_{k}\left(A_{k, \gamma_{n} i} V_{\gamma_{n}, i}\right)-\hat{\phi}_{k, n, \gamma_{n}}\left(A_{k, \gamma_{n} i} V_{\gamma_{n}, i}\right)\right) A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\left(\zeta_{l, k, j, \gamma_{n}, i}^{x}-\bar{\zeta}_{l, k, j, n, \gamma_{n}}^{x}\right)$;
4. $\frac{1}{\sqrt{n}} \sum_{i=1}^{n}\left(\left[\varrho_{k, n, \gamma_{n}, 1}-\varrho_{k, 1}\right] A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}+\left[\varrho_{k, n, \gamma_{n}, 2}-\varrho_{k, 2}\right] \kappa\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)\right)$;
5. $\frac{1}{n} \sum_{i=1}^{n}\left[A_{k, \gamma_{n}, i} D_{b, l}\left(b_{n}, X_{i}\right)\right]-P_{\tilde{\theta}_{n}}\left[A_{k, \gamma_{n}, i} D_{b, l}\left(b_{n}, X_{i}\right)\right]$;
6. $\frac{1}{\sqrt{n}} \sum_{i=1}^{n}\left(\phi_{k}\left(A_{k, \gamma_{n} i} V_{\gamma_{n}, i}\right)-\hat{\phi}_{k, n, \gamma_{n}}\left(A_{k, \gamma_{n} i} V_{\gamma_{n}, i}\right)\right)\left(\left[A_{k, \gamma_{n}, i} D_{b, l}\left(b_{n}, X_{i}\right)\right]-\frac{1}{n} \sum_{i=1}^{n}\left[A_{k, \gamma_{n}, i} D_{b, l}\left(b_{n}, X_{i}\right)\right]\right)$;
and the following terms are $O_{P_{\bar{\theta}_{n}}^{n}}(1)$ :
7. $\frac{1}{\sqrt{n}} \sum_{i=1}^{n}\left(\phi_{k}\left(A_{k, \gamma_{n} i} V_{\gamma_{n}, i}\right) A_{k, \gamma_{n} i} V_{\gamma_{n}, i}+1\right)$;
8. $\frac{1}{\sqrt{n}} \sum_{i=1}^{n} \varrho_{k, 1} A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}+\varrho_{k, 2} \kappa\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)$;
9. $\frac{1}{\sqrt{n}} \sum_{i=1}^{n} \phi_{k}\left(A_{k, \gamma_{n} i} V_{\gamma_{n}, i}\right)$.

Proof. Under $P_{\tilde{\theta}_{n}}, \tilde{X}$ is distributed according to the density $\eta_{0}$ whilst $A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}$ has the same law as $\epsilon_{k}$. We will use these facts without explicit reference in the rest of the proof.

1. The triangular array $\left(\zeta_{l, k, k, \gamma_{n}, i}^{x}\right)_{n \in \mathbb{N}, i=1, \ldots, n}$ has i.i.d. rows and the variance of $\zeta_{l, k, k, \gamma_{n}, i}^{x}$ is bounded above uniformly in $n$ by Assumption S1. The claim then follows from a WLLN for triangular arrays (e.g. Durrett, 2019, Theorem 2.2.6).
2. Let $Z_{n, i}:=\zeta_{l, k, j, i}^{x} A_{j, \gamma_{n}, i} V_{\gamma_{n}, i}$. The triangular array $\left(Z_{n, i}\right)_{n \in \mathbb{N}, i=1, \ldots, n}$ has i.i.d. rows, $Z_{n, i} \Perp \epsilon_{i, k}, Z_{n, i}$ is mean zero and the variance of $Z_{n, i}$ is bounded above uniformly in $n$ by Assumptions S1 and S2. The claim then follows by Assumption S3.
3. By Cauchy - Schwarz one has

$$
\begin{aligned}
& \frac{1}{\sqrt{n}} \sum_{i=1}^{n}\left(\phi_{k}\left(A_{k, \gamma_{n} i} V_{\gamma_{n}, i}\right)-\hat{\phi}_{k, n, \gamma_{n}}\left(A_{k, \gamma_{n} i} V_{\gamma_{n}, i}\right)\right) A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\left(\zeta_{l, k, j, \gamma_{n}, i}^{x}-\bar{\zeta}_{l, k, j, n, \gamma_{n}}^{x}\right) \\
& \leq\left[\frac{1}{n} \sum_{i=1}^{n}\left(\phi_{k}\left(A_{k, \gamma_{n} i} V_{\gamma_{n}, i}\right)-\hat{\phi}_{k, n, \gamma_{n}}\left(A_{k, \gamma_{n} i} V_{\gamma_{n}, i}\right)\right)^{2}\left(\zeta_{l, k, j, \gamma_{n}, i}^{x}-\bar{\zeta}_{l, k, j, n, \gamma_{n}}^{x}\right)^{2}\right]^{1 / 2} \\
& \quad \times\left[\frac{1}{n} \sum_{i=1}^{n}\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)^{2}\right]^{1 / 2}
\end{aligned}
$$

Take $Z_{n, i}:=\zeta_{l, k, j, \gamma_{n}, i}^{x}-\bar{\zeta}_{l, k, j, n, \gamma_{n}}^{x}$. The triangular array $\left(Z_{n, i}\right)_{n \in \mathbb{N}, i=1, \ldots, n}$ has i.i.d. rows, $Z_{n, i} \Perp \epsilon_{i, k}, Z_{n, i}$ is mean zero and the variance of $Z_{n, i}$ is bounded above uniformly in $n$ by Assumption S1. Therefore, the first factor on the right hand side is $o_{P_{\bar{\theta}_{n}}^{n}}(1)$ by Assumption S3. The second right hand side factor is $O_{P_{\bar{\theta}_{n}}^{n}}(1)$ by Assumption S2.
4. $\varrho_{k, n, \gamma_{n}} \xrightarrow{P_{\theta_{n}}^{n}} \varrho_{k}$ by Lemma S12. Assumption $S 2$ and the central limit theorem imply that $\frac{1}{\sqrt{n}} \sum_{i=1}^{n} A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}$ and $\frac{1}{\sqrt{n}} \sum_{i=1}^{n} \kappa\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)$ are $O_{P_{\hat{\theta}_{n}}}(1)$.
5. Let $U_{n, i}:=\operatorname{vec}\left(A_{k, \gamma_{n}, i} D_{b, l}\left(b_{n}, X\right)\right)$. Then for each component $U_{n, i, l},\left(U_{n, i, l}\right)_{n \in \mathbb{N}, i=1, \ldots, n}$ is a triangular array with i.i.d. rows and the variance of $U_{n, i, l}$ is bounded above uniformly in $n$ by Assumptions S1 and S2. The claim then follows from a WLLN for triangular arrays (e.g. Durrett, 2019, Theorem 2.2.6).
6. Put $Z_{n, i}:=\left[A_{k, \gamma_{n}, i} D_{b, l}\left(b_{n}, X\right)\right]-\frac{1}{n} \sum_{i=1}^{n}\left[A_{k, \gamma_{n}, i} D_{b, l}\left(b_{n}, X\right)\right]$. Then, the triangular array $\left(Z_{n, i}\right)_{n \in \mathbb{N}, i=1, \ldots, n}$ has i.i.d. rows, $Z_{n, i} \Perp \epsilon_{i, k}, Z_{n, i}$ is mean zero and the variance of $Z_{n, i}$ is bounded above uniformly in $n$ by Assumptions S1 and S2. The claim follows by Assumption S3.

Each of the remaining items follow from the central limit theorem given Assumption S2.
Lemma S12. If Assumption S2 holds, $\left\|\varrho_{k, n, \gamma_{n}}-\varrho_{k}\right\|=o_{P_{\bar{\theta}_{n}}^{n}}\left(\nu_{n, p}\right)$ for $\varrho \in\{\tau, \varsigma\}$. S11
Proof. Under $P_{\tilde{\theta}_{n}}, \hat{M}_{k, n, \gamma_{n}}$ has the same law as $M_{k, n}:=\frac{1}{n} \sum_{i=1}^{n}\left(\begin{array}{cc}1 & \epsilon_{i, k}^{3} \\ \epsilon_{i, k}^{3} & \epsilon_{i, k}^{4}-1\end{array}\right)$. Therefore, it suffices to show that $\left\|M_{k, n}^{-1} w-M_{k}^{-1} w\right\|=o_{P_{\bar{\theta}_{n}}}\left(\nu_{n, p}\right)$ for any fixed $w \in \mathbb{R}^{2}$. Since the map $M \mapsto M^{-1}$ is Lipschitz continuous at a positive definite matrix,

$$
\left\|M_{k, n}^{-1} w-M_{k}^{-1} w\right\|_{2} \leq\|w\|\left\|M_{k, n}^{-1}-M_{k}^{-1}\right\|_{2} \lesssim\left\|M_{k, n}-M_{k}\right\|_{2},
$$

and thus it suffices to show that $\left\|M_{k, n}-M_{k}\right\|_{2}=o_{P_{\theta_{n}}^{n}}\left(\nu_{n, p}\right)$. If $v:=\delta / 4 \geq 1$, we have that by Theorem 2.5.11 in Durrett (2019)

$$
\begin{aligned}
& \frac{1}{n} \sum_{i=1}^{n}\left[\epsilon_{i, k}^{3}-\mathbb{E}\left(\epsilon_{i, k}^{3}\right)\right]=o_{P_{\bar{\theta}_{n}}^{n}}\left(n^{-1 / 2} \log (n)^{1 / 2+\rho}\right) \\
& \frac{1}{n} \sum_{i=1}^{n}\left[\epsilon_{i, k}^{4}-\mathbb{E}\left(\epsilon_{i, k}^{4}\right)\right]=o_{P_{\bar{\theta}_{n}}^{n}}\left(n^{-1 / 2} \log (n)^{1 / 2+\rho}\right)
\end{aligned}
$$

for any $\rho>0$, which implies that

$$
\left\|M_{k, n}-M_{k}\right\|_{2} \leq\left\|M_{k, n}-M_{k}\right\|_{F}=o_{P_{\bar{\theta}_{n}}^{n}}\left(n^{-1 / 2} \log (n)^{1 / 2+\rho}\right) .
$$

${ }^{\text {S11 }} \nu_{n, p}$ is as defined in Assumption 3: $p:=\min \{1+\delta / 4,2\}$ and $\nu_{n, p}:=\left\{\begin{array}{ll}n^{(1-p) / p} & \text { for } p \in(1,2) \\ n^{-1 / 2} \log (n)^{1 / 2+\rho} & \text { for } p=2\end{array}\right.$, for some $\rho>0$.

If $0<v<1$, by Theorems 2.5.11 \& 2.5.12 in Durrett (2019), for any $\rho>0$,

$$
\begin{aligned}
& \frac{1}{n} \sum_{i=1}^{n}\left[\left(\epsilon_{i, k}\right)^{3}-\mathbb{E}\left(\epsilon_{i, k}\right)^{3}\right]=\left\{\begin{array}{lc}
o_{P_{\bar{\theta}_{n}}^{n}}\left(n^{-1 / 2} \log (n)^{1 / 2+\rho}\right) & \text { if } v \in[1 / 2,1) \\
o_{P_{\bar{\theta}_{n}}^{n}}\left(n^{\frac{1-p}{p}}\right) & \text { if } v \in(0,1 / 2)
\end{array}\right. \\
& \frac{1}{n} \sum_{i=1}^{n}\left[\left(\epsilon_{i, k}\right)^{4}-\mathbb{E}\left(\epsilon_{i, k}\right)^{4}\right]=o_{P_{\bar{\theta}_{n}}^{n}}\left(n^{\frac{1-p}{p}}\right)
\end{aligned}
$$

which together imply that

$$
\left\|M_{k, n}-M_{k}\right\|_{2} \leq\left\|M_{k, n}-M_{k}\right\|_{F}=o_{P_{\theta_{n}}^{n}}\left(n^{\frac{1-p}{p}}\right) .
$$

Lemma S13. Suppose that Assumptions S1, S2 and S3 hold. Then, for each $(k, j)$ with $k \neq j$, each $l$, each $x \in\{\alpha, \sigma\}$ and each $\varrho \in\{\tau, \varsigma\}$, the following terms are $o_{P_{\bar{\theta}_{n}}^{n}}\left(\nu_{n}\right)$ :

$$
\begin{aligned}
& \text { 1. } \frac{1}{n} \sum_{i=1}^{n}\left(\phi_{k}\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)-\hat{\phi}_{k, n, \gamma_{n}}\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)\right)^{2}\left(A_{j, \gamma_{n}, i} V_{\gamma_{n}, i} \zeta_{l, k, j, \gamma_{n}, i}^{x}\right)^{2} ; \\
& \text { 2. }\left(P_{\tilde{\theta}_{n}}\left[\zeta_{l, k, k, \gamma_{n}, i}^{x}\right]-\bar{\zeta}_{l, k, k, n, \gamma_{n}}^{x}\right)^{2} ; \\
& \text { 3. } \frac{1}{n} \sum_{i=1}^{n}\left(\phi_{k}\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right) A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}+1\right)^{2}\left(P_{\tilde{\theta}_{n}}\left[\zeta_{l, k, k, \gamma_{n}, i}^{x}\right]-\bar{\zeta}_{l, k, k, n, \gamma_{n}}^{x}\right)^{2} ; \\
& \text { 4. } \frac{1}{n} \sum_{i=1}^{n}\left(\varrho_{k, 1} A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}+\varrho_{k, 2} \kappa\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)\right)^{2}\left(P_{\tilde{\theta}_{n}}\left[\zeta_{l, k, k, \gamma_{n}, i}^{x}\right]-\bar{\zeta}_{l, k, k, n, \gamma_{n}}^{x}\right)^{2} ; \\
& \text { 5. } \frac{1}{n} \sum_{i=1}^{n}\left(\bar{\zeta}_{l, k, k, n, \gamma_{n}}^{x}\left(\left[\varrho_{k, n, \gamma_{n}, 1}-\varrho_{k, 1}\right] A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}+\left[\hat{\varrho}_{k, n, \gamma_{n}, 2}-\varrho_{k, 2}\right] \kappa\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)\right)\right)^{2} ; \\
& \text { 6. }\left(P_{\tilde{\theta}_{n}}\left[A_{k, \gamma_{n}, i} D_{b, l}\left(b_{n}, X_{i}\right)\right]-[\mathrm{ADbX}]_{n}\right)^{2} ; \\
& \text { 7. } \frac{1}{n} \sum_{i=1}^{n}\left(\varrho_{k, 1} A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}+\varrho_{k, 2} \kappa\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)\right)^{2}\left(P_{\tilde{\theta}_{n}}\left[A_{k, \gamma_{n}, i} D_{b, l}\left(b_{n}, X_{i}\right)\right]-[\mathrm{ADbX}]_{n}\right)^{2} ; \\
& \text { 8. } \frac{1}{n} \sum_{i=1}^{n}\left([\mathrm{ADbX}]_{n}\left(\left[\varrho_{k, n, \gamma_{n}, 1}-\varrho_{k, 1}\right] A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}+\left[\hat{\varrho}_{k, n, \gamma_{n}, 2}-\varrho_{k, 2}\right] \kappa\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)\right)\right)^{2} \\
& \text { 9. } \frac{1}{n} \sum_{i=1}^{n}\left(\phi_{k}\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)-\hat{\phi}_{k, n, \gamma_{n}}\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)\right)^{2}\left(A_{k, \gamma_{n}, i} D_{b, l}\left(b_{n}, X_{i}\right)-[\mathrm{ADbX}]_{n}\right)^{2} ; \\
& \text { 10. } \frac{1}{n} \sum_{i=1}^{n}\left(\phi_{k}\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)\right)^{2}\left(P_{\tilde{\theta}_{n}}\left[A_{k, \gamma_{n}, i} D_{b, l}\left(b_{n}, X_{i}\right)\right]-[\mathrm{ADbX}]_{n}\right)^{2},
\end{aligned}
$$

where $[\mathrm{ADbX}]_{n}:=\frac{1}{n} \sum_{i=1}^{n} A_{k, \gamma_{n}, i} D_{b, l}\left(b_{n}, X_{i}\right)$.
Proof. Under $P_{\tilde{\theta}_{n}}, \tilde{X}$ is distributed according to the density $\eta_{0}$ whilst $A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}$ has the same law as $\epsilon_{k}$. We will use these facts without explicit reference in the rest of the proof.

1. Let $Z_{n, i}:=A_{j, \gamma_{n}, i} V_{\gamma_{n}, i} \zeta_{l, k, j, \gamma_{n}, i}^{x}$. This is independent of $\epsilon_{i, k}$, is mean-zero and has variance bounded above uniformly in $n$ by Assumptions S1 and S2. The claim then follows by Assumption S3.
2. Let $Z_{n, i}:=\left(\zeta_{l, k, k, \gamma_{n}, i}^{x}-P_{\tilde{\theta}_{n}}\left[\zeta_{l, k, k, \gamma_{n}, i}^{x}\right]\right)$ and note that $\sup _{n \in \mathbb{N}} \mathbb{E} Z_{n, i}^{2+\varepsilon}<\infty$ for a $\varepsilon>0$ (by Assumption S1). By the Lindeberg CLT one then has that $\sum_{i=1}^{n} Z_{n, i}=O_{P_{\bar{\theta}_{n}}^{n}}(\sqrt{n})$ and hence $\left(P_{\tilde{\theta}_{n}}\left[\zeta_{l, k, k, \gamma_{n}, i}^{x}\right]-\bar{\zeta}_{l, k, k, n, \gamma_{n}}^{x}\right)^{2}=o_{P_{\bar{\theta}_{n}}^{n}}\left(\nu_{n}\right)$.
3. By Assumption S2, $\frac{1}{n} \sum_{i=1}^{n}\left(\phi_{k}\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right) A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}+1\right)^{2}=O_{P_{\widehat{\theta}_{n}}}(1)$. Use 2 .
4. By Assumption S2, $\frac{1}{n} \sum_{i=1}^{n}\left(\varrho_{k, 1} A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}+\varrho_{k, 2} \kappa\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)\right)^{2}=O_{P_{\vec{\theta}_{n}}}(1)$. Use 2.
5. By Assumption $\mathrm{S} 1, \bar{\zeta}_{l, k, k, n, \gamma_{n}}^{x}$ is bounded uniformly for all sufficiently large $n$. By Assumption S2, $\frac{1}{n} \sum_{i=1}^{n}\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)^{2}$ and $\frac{1}{n} \sum_{i=1}^{n} \kappa\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)^{2}$ are $O_{P_{\vec{\theta}_{n}}^{n}}(1)$. Combine with Lemma S12.
6. Let $Z_{n, i}:=\left(A_{k, \gamma_{n}, i} D_{b, l}\left(b_{n}, X_{i}\right)-P_{\tilde{\theta}_{n}}\left[A_{k, \gamma_{n}, i} D_{b, l}\left(b_{n}, X_{i}\right)\right]\right)$ and note that $\sup _{n \in \mathbb{N}} \mathbb{E} Z_{n, i}^{2+\varepsilon}<$ $\infty$ for a $\varepsilon>0$ (by Assumptions S1 and S2). By the Lindeberg CLT one then has that $\sum_{i=1}^{n} Z_{n, i}=O_{P_{\hat{\theta}_{n}}^{n}}(\sqrt{n})$ and hence $\left(P_{\tilde{\theta}_{n}}\left[A_{k, \gamma_{n}, i} D_{b_{n}, l}\left(b_{n}, X_{i}\right)\right]-[\mathrm{ADbX}]_{n}\right)^{2}=o_{P_{\tilde{\theta}_{n}}^{n}}\left(\nu_{n}\right)$.
7. By Assumption S2, $\frac{1}{n} \sum_{i=1}^{n}\left(\varrho_{k, 1} A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}+\varrho_{k, 2} \kappa\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)\right)^{2}=O_{P_{\widehat{\theta}_{n}}}(1)$. Use 6.
8. Take $\left[\mathrm{ADbX}_{n}\right]$ out of the summation. By Assumption S 2 and 6. this is $O_{P_{\bar{\theta}_{n}}^{n}}(1)$. By Assumption S2, $\frac{1}{n} \sum_{i=1}^{n}\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)^{2}$ and $\frac{1}{n} \sum_{i=1}^{n} \kappa\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)^{2}$ are $O_{P_{\theta_{n}}^{n}}(1)$. Combine with Lemma S12.
9. For $Z_{n, i}:=A_{k, \gamma_{n}, i} D_{b, l}\left(b_{n}, X_{i}\right)-[\mathrm{ADbX}]_{n}, Z_{n, i}$ is independent of $\epsilon_{i, k}$, mean-zero and has variance bounded uniformly in $n$ by Assumptions S1 and S2. The claim follows from Assumption S3.
10. $\frac{1}{n} \sum_{i=1}^{n}\left(\phi_{k}\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)\right)^{2}=O_{P_{\theta_{n}}^{n}}(1)$ by Assumption S2. Use 6 .

Lemma S14. Suppose that Assumptions S1 and S4 hold. Then, for each $k$, each $l$, each $x \in\{\alpha, \sigma\}$,
$\frac{1}{n} \sum_{i=1}^{n}\left(\phi_{k}\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)-\hat{\phi}_{k, n, \gamma_{n}}\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)\right)^{2}\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\left[\zeta_{l, k, k, \gamma_{n}, i}^{x}-\bar{\zeta}_{l, k, k, n, \gamma_{n}}^{x}\right]\right)^{2}=o_{P_{\bar{\theta}_{n}}^{n}}\left(\nu_{n}\right)$.
Proof. By Assumption S1, $\left[\zeta_{l, k, k, \gamma_{n}, i}^{x}-\bar{\zeta}_{l, k, k, n, \gamma_{n}}^{x}\right]^{2}$ is uniformly bounded for all large enough $n$. Hence it suffices that by Assumption S4,

$$
\frac{1}{n} \sum_{i=1}^{n}\left(\phi_{k}\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)-\hat{\phi}_{k, n, \gamma_{n}}\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)\right)^{2}\left(A_{k, \gamma_{n}, i} V_{\gamma_{n}, i}\right)^{2}=o_{P_{\theta_{n}}^{n}}\left(\nu_{n}\right) .
$$

## S3 Additional auxillary results

We present a few additional results that explicitly prove some claims made in the main text. First, we show that if two errors $\epsilon_{i, k}$ and $\epsilon_{i, j}$ are Gaussian $\tilde{I}_{\theta, \alpha \alpha}$ becomes singular, which implies the singularity of $\tilde{\mathcal{I}}_{\theta}$ if $\tilde{I}_{\theta, \beta \beta}$ is non-singular (cf. Propositions 8.2.4 and 8.2.8 in Bernstein (2009)). Second, we provide an explicit example of a density which satisfies the first part of the Assumption 2 but not the second. Third we prove that if Assumption 2 part 1 holds then a sufficient condition for part 2 is that $\eta_{k}$ has tails that decay to zero at a polynomial rate.

Lemma S15. Consider the LSEM model (3) and suppose that Assumptions 1 and 2 hold. Define the random vector $Q$ in $\mathbb{R}^{K^{2}}$ as

$$
Q=\left(Q_{1}^{\prime}, \ldots, Q_{K}^{\prime}\right)^{\prime},
$$

where the $j$-th element of $Q_{k}$ for $j \in[K]$ is given by

$$
Q_{k, j}= \begin{cases}\phi_{k}\left(\epsilon_{k}\right) \epsilon_{j} & \text { if } k \neq j \\ \tau_{k, 1} \epsilon_{k}+\tau_{k, 2} \kappa\left(\epsilon_{k}\right) & \text { if } k=j\end{cases}
$$

Next define the matrix $\zeta \in \mathbb{R}^{K^{2} \times L_{\alpha}}$ according to

$$
\zeta=\left(\operatorname{vec}\left(\left[D_{\alpha, 1}(\alpha, \sigma) A(\alpha, \sigma)^{-1}\right]^{\prime}\right), \ldots, \operatorname{vec}\left(\left[D_{\alpha, L_{\alpha}}(\alpha, \sigma) A(\alpha, \sigma)^{-1}\right]^{\prime}\right)\right) .
$$

Then where $\tilde{\ell}_{\theta}$ is the effective score function as defined in lemma 3, the law of $\tilde{\ell}_{\theta, 1}$ under $P_{\theta}$ is equal to that of $\zeta^{\prime} Q$. Moreover,
(i) $\mathbb{E} Q Q^{\prime}$ is non-singular if and only if for each pair $(k, j)$ with $k \neq j$ and each $k, j \in[K]$ we have that $\left[\mathbb{E} \phi_{k}^{2}\left(\epsilon_{k}\right)\right]\left[\mathbb{E} \phi_{j}^{2}\left(\epsilon_{j}\right)\right] \neq 1$.
(ii) $\tilde{I}_{\theta, \alpha \alpha}$ is non-singular if $\operatorname{rank}(\zeta)=L_{\alpha}$ and $\mathbb{E} Q Q^{\prime}$ is non-singular.
(iii) If $\operatorname{rank}(\zeta)<L_{\alpha}$ then $\tilde{I}_{\theta, \alpha \alpha}$ is singular.
(iv) If $L_{\alpha}=K^{2}$ and $\mathbb{E} Q Q^{\prime}$ is singular then $\tilde{I}_{\theta, \alpha \alpha}$ is singular.
(v) If $\mathbb{E} Q Q^{\prime}$ is singular, $\tilde{I}_{\theta, \alpha \alpha}$ may be singular when $\operatorname{rank}(\zeta)=L_{\alpha}<K^{2}$.

In particular, if both $\epsilon_{k}$ and $\epsilon_{j}(k \neq j)$ have a Gaussian distribution and $L_{\alpha}=K^{2}, \tilde{I}_{\theta, \alpha \alpha}$ is singular.

Proof. For (i), let $j, k, m, i$ all be in $[K]$. We will consider the entries of the matrix $\mathbb{E} Q Q^{\prime}$, which are of the form $\mathbb{E}\left[Q_{k, j} Q_{m, i}\right]$. In particular, the $s, t$-th element of the matrix is given by the form $\mathbb{E}\left[Q_{k, j} Q_{m, i}\right]$ where $(k-1) K+j=s$ and $(m-1) K+i=t$. If $k=j=m=i$ we have $s=t$ and $\mathbb{E}\left[Q_{k, j} Q_{m, i}\right]=\mathbb{E}\left[\tau_{k, 1} \epsilon_{k}+\tau_{k, 2} \kappa\left(\epsilon_{k}\right)\right]^{2}$. The other diagonal entries occur when $k=m \neq j=i$, and have the form $\mathbb{E}\left[Q_{k, j} Q_{m, i}\right]=\mathbb{E}\left[\phi_{k}^{2}\left(\epsilon_{k}\right)\right]$. Inspection of the other possible cases reveals that the only other case with non-zero entries is $k=i \neq m=j$ which has value $\mathbb{E}\left[Q_{k, j} Q_{m, i}\right]=\mathbb{E}\left[\phi_{k}\left(\epsilon_{k}\right) \epsilon_{k}\right] \mathbb{E}\left[\phi_{k}\left(\epsilon_{m}\right) \epsilon_{m}\right]=1$ by assumption 2.

Therefore for any $k, j \in[K]$, column $(k-1) K+j$ has non-zero entries in row $(k-1) K+j$ only if $k=j$ and otherwise in rows $(k-1) K+j$ and $(j-1) K+k$, with values $\mathbb{E} \phi_{k}^{2}\left(\epsilon_{k}\right)$ and 1 respectively. There are therefore no columns that can be linearly related to column $(k-1) K+j$ if $k=j$. If $k \neq j$, then column $(k-1) K+j$ has zeros everywhere except row $(k-1) K+j$ where it has $\mathbb{E} \phi_{k}^{2}\left(\epsilon_{k}\right)$ and row $(j-1)+k$ where it has 1 . Column $(j-1) K+k$ has zeros everywhere except row $(j-1) K+k$ where it has $\mathbb{E} \phi_{j}^{2}\left(\epsilon_{j}\right)$ and row $(k-1) K+j$ where it has 1 . Since no other columns have entries in these rows, it follows that column $(k-1) K+j$ is linearly independent of all the other columns if and only if it is linearly independent of column $(j-1) K+k$, which occurs if and only if $\left[\mathbb{E} \phi_{k}^{2}\left(\epsilon_{k}\right)\right]\left[\mathbb{E} \phi_{j}^{2}\left(\epsilon_{j}\right)\right] \neq 1$.

For (ii), suppose that $\operatorname{rank}(\zeta)=L_{\alpha}$ and $\mathbb{E} Q Q^{\prime}$ is non-singular. Then there is a (unique) positive definite $\left[\mathbb{E} Q Q^{\prime}\right]^{1 / 2}$ and we have $\tilde{I}_{\theta, \alpha \alpha}=\left(\left[\mathbb{E} Q Q^{\prime}\right]^{1 / 2} \zeta\right)^{\prime}\left(\left[\mathbb{E} Q Q^{\prime}\right]^{1 / 2} \zeta\right)$ which has full rank, since $\left(\left[\mathbb{E} Q Q^{\prime}\right]^{1 / 2} \zeta\right)$ has full column rank.

For the remaining parts note first that

$$
\tilde{I}_{\theta, \alpha \alpha}=\mathbb{E} \tilde{\ell}_{\theta, 1} \tilde{\ell}_{\theta, 1}^{\prime}=\zeta^{\prime}\left[\mathbb{E} Q Q^{\prime}\right] \zeta
$$

and so $\operatorname{rank}\left(\tilde{I}_{\theta, \alpha \alpha}\right) \leq \min \left\{\operatorname{rank}\left(\zeta^{\prime} \mathbb{E} Q Q^{\prime}\right), \operatorname{rank}(\zeta)\right\}$. Hence if $\operatorname{rank}(\zeta)<L_{\alpha}, \operatorname{rank}\left(\tilde{I}_{\theta, \alpha \alpha}\right)<L_{\alpha}$ implying (iii).

For (iv), suppose that $\operatorname{rank}\left(\mathbb{E} Q Q^{\prime}\right)<K^{2}=L_{\alpha}$. Then, there is a non-zero $x \in \mathbb{R}^{L_{\alpha}}$ such that $\mathbb{E} Q Q^{\prime} x=0$ and hence $\zeta^{\prime} \mathbb{E} Q Q^{\prime} x=0$. Hence $\operatorname{dim}\left(\operatorname{ker}\left(\zeta^{\prime} \mathbb{E} Q Q^{\prime}\right)\right) \geq 1$. It follows that $\operatorname{rank}\left(\zeta^{\prime} \mathbb{E} Q Q^{\prime}\right) \leq L_{\alpha}-1<L_{\alpha}$ and hence $\operatorname{rank}\left(\tilde{I}_{\theta, \alpha \alpha}\right) \leq \min \left\{\operatorname{rank}\left(\zeta^{\prime} \mathbb{E} Q Q^{\prime}\right), \operatorname{rank}(\zeta)\right\}<L_{\alpha}$.

For (v) suppose that $K=2, \epsilon_{1}$ and $\epsilon_{2}$ are both Gaussian and $A(\alpha)=\left[\begin{array}{c}\cos (\alpha)-\sin (\alpha) \\ \sin (\alpha) \\ \cos (\alpha)\end{array}\right]$. We have for $l \in\{1,2\}, \phi_{l}(z)=-z$, hence $\phi_{l}^{2}(z)=z^{2}$ and so $\mathbb{E} \phi_{l}^{2}\left(\epsilon_{l}\right)=1$. $D_{\alpha, 1}(\gamma)=$ $\left[\begin{array}{cc}-\sin (\alpha) & -\cos (\alpha) \\ \cos (\alpha) & -\sin (\alpha)\end{array}\right]$ and hence

$$
D_{\alpha, 1}(\alpha) A(\alpha)^{-1}=D_{\alpha, 1}(\alpha) A(\alpha)^{\prime}=\left[\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right]
$$

which implies $\zeta=(0,-1,1,0)^{\prime}$ and hence $\operatorname{rank}(\zeta)=1=L_{\alpha}<K^{2}=4$. Explicit calculation
reveals that

$$
\mathbb{E} Q Q^{\prime}=\left[\begin{array}{cccc}
8 / 9 & 0 & 0 & 0 \\
0 & 1 & 1 & 0 \\
0 & 1 & 1 & 0 \\
0 & 0 & 0 & 8 / 9
\end{array}\right]
$$

which is clearly singular with rank 3 . We have

$$
\tilde{I}_{\theta, \alpha \alpha}=\zeta^{\prime}\left[\mathbb{E} Q Q^{\prime}\right] \zeta=\zeta^{\prime}\left[\begin{array}{cccc}
8 / 9 & 0 & 0 & 0 \\
0 & 1 & 1 & 0 \\
0 & 1 & 1 & 0 \\
0 & 0 & 0 & 8 / 9
\end{array}\right]\left[\begin{array}{c}
0 \\
-1 \\
1 \\
0
\end{array}\right]=\zeta^{\prime}\left[\begin{array}{l}
0 \\
0 \\
0 \\
0
\end{array}\right]=0 .
$$

For the last part, suppose that $k \neq j$ and $\epsilon_{k}$ and $\epsilon_{j}$ are both Gaussian. Since both have zero mean and unit variance, we have for $l \in\{k, j\}, \phi_{l}(z)=-z$, hence $\phi_{l}^{2}(z)=z^{2}$ and so $\mathbb{E} \phi_{l}^{2}\left(\epsilon_{l}\right)=1 . E_{\theta} Q Q^{\prime}$ is singular by (i) and hence by (iv) $\tilde{I}_{\theta, \alpha \alpha}$ is singular.

Example S1 (Necessity of part 2 of assumption 2). Suppose that $\tilde{\epsilon}_{k} \sim \chi_{2}^{2}$ and let $\epsilon_{k}=$ $\left(\tilde{\epsilon}_{k}-2\right) / 2$. Then $\epsilon_{k}$ has mean zero, variance one and density function $\eta_{k}(z)=\exp (-z-1)$ on its support $[-1, \infty)$ on which we also have that $\phi_{k}(z)=-1$. Explicit calculation reveals that part 1 of assumption 2 is satisfied. However, $\mathbb{E} \phi_{k}(z)=-1 \neq 0$ as would be required by part 2 of assumption 2.

Note also that this example does not satisfy the requirements of lemma S16: we have $a_{k}=-1, b_{k}=\infty$ and

$$
\lim _{z \downarrow a_{k}} \eta_{k}(x)=\lim _{z \downarrow-1} \exp (-z-1)=1 \neq 0,
$$

and hence the required condition is violated for $r=0$.
Lemma S16. Let $a_{k}=\inf \left\{x \in \mathbb{R} \cup\{-\infty\}: \eta_{k}(x)>0\right\}$ and $b_{k}=\sup \{x \in \mathbb{R} \cup\{\infty\}$ : $\left.\eta_{k}(x)>0\right\}$. Suppose that, for $r=0,1,2,3$ : (i) if $a_{k}=-\infty$ then $\eta_{k}(x)=o\left(x^{-3}\right)$ as $x \rightarrow-\infty$, else $a_{k}^{r} \lim _{x \downarrow a_{k}} \eta_{k}(x)=0$, and (ii) if $b_{k}=\infty$ then $\eta_{k}(x)=o\left(x^{-3}\right)$ as $x \rightarrow \infty$, else $b_{k}^{r} \lim _{x \uparrow b_{k}} \eta_{k}(x)=0$. Then, if part 1 of assumption 2 holds, part 2 is also satisfied.

Proof. Let $r \in\{0,1,2,3\}, b_{k}=\sup \left\{x \in \mathbb{R}: \eta_{k}(x)>0\right\}$ and $a_{k}=\inf \left\{x \in \mathbb{R}: \eta_{k}(x)>0\right\}$. We have, by integration by parts, with $G_{k}$ denoting the measure on $\mathbb{R}$ corresponding to $\eta_{k}$,

$$
\int \phi_{k}(z) z^{r} \mathrm{~d} G_{k}=\int \frac{\eta_{k}^{\prime}(z)}{\eta_{k}(z)} \eta_{k}(z) z^{r} \mathrm{~d} z=\int \eta_{k}^{\prime}(z) z^{r} \mathrm{~d} z=\left.\eta_{k}(z) z^{r}\right|_{a_{k}} ^{b_{k}}-\int \eta_{k}(z) \frac{\mathrm{d} z^{r}}{\mathrm{~d} z} \mathrm{~d} z
$$

Our hypothesis ensures that $\left.z^{r} \eta_{k}(z)\right|_{a_{k}} ^{b_{k}}=0$. Therefore we have $G_{k} \phi_{k}(z) z^{r}=-G_{k} \frac{\mathrm{~d}}{\mathrm{~d} z} z^{r}$. For
$r=0$ this equals zero as $\frac{\mathrm{d}}{\mathrm{d} z} z^{0}=\frac{\mathrm{d}}{\mathrm{d} z} 1=0$. For $r \in\{1,2,3\}$ we have $\frac{\mathrm{d} z^{r}}{\mathrm{~d} z}=r z^{r-1}$ and hence $G_{k} \phi_{k}(z) z^{r}=-r G_{k} z^{r-1}$. Since $G_{k} 1=1, G_{k} z=0$, and $G_{k} z^{2}=1$, the result follows.

Lemma S17. Suppose that $P_{n}$ and $Q_{n}$ are probability measures (with each pair $\left(P_{n}, Q_{n}\right)$ defined on a common measurable space) with corresponding densities $p_{n}$ and $q_{n}$ (with respect to some $\sigma$-finite measure $\nu_{n}$ ). Let $l_{n}=\log q_{n} / p_{n}$ be the log-likelihood ratio. ${ }^{\text {S12 }}$ If

$$
l_{n}=o_{P_{n}}(1),
$$

then $d_{T V}\left(P_{n}, Q_{n}\right) \rightarrow 0$.
Proof. By the continuous mapping theorem

$$
\frac{q_{n}}{p_{n}}=\exp \left(l_{n}\right) \xrightarrow{P_{n}} 1 .
$$

Le Cam's first lemma (e.g. van der Vaart, 1998, Lemma 6.4) then implies that $Q_{n} \triangleleft P_{n}$. Let $\phi_{n}$ be arbitrary measurable functions valued in $[0,1]$. Since the $\phi_{n}$ are uniformly tight, Prohorov's theorem ensures that for any arbitrary subsequence $\left(n_{j}\right)_{j \in \mathbb{N}}$ there exists a further subsequence $\left(n_{m}\right)_{m \in \mathbb{N}}$ such that $\phi_{n_{m}} \rightsquigarrow \phi \in[0,1]$ under $P_{n_{m}}$. Therefore by Slutsky's Theorem

$$
\left(\phi_{n_{m}}, \exp \left(l_{n_{m}}\right)\right) \rightsquigarrow(\phi, 1) \quad \text { under } P_{n_{m}} .
$$

By Le Cam's third Lemma (e.g. van der Vaart, 1998, Theorem 6.6), under $Q_{m_{n}}$ the law of $\phi_{n_{m}}$ converges weakly to the law of $\phi$. Since each $\phi_{n} \in[0,1]$

$$
\lim _{m \rightarrow \infty}\left[Q_{n_{m}} \phi_{n_{m}}-P_{n_{m}} \phi_{n_{m}}\right]=0
$$

As $\left(n_{j}\right)_{j \in \mathbb{N}}$ was arbitrary, the preceding display holds also along the original sequence.

## S4 A consistent estimator of the Moore - Penrose psuedoinverse

As is well known, the Moore - Penrose psuedoinverse of a matrix is not a continuous function on the space of positive semi-definite matrices (see e.g. Ben-Israel and Greville, 2003, Section 6.6). In consequence, if one has a consistent estimator $\bar{M}_{n}$ of some matrix $M$, it need not follow that $\check{M}^{\dagger}$ is consistent for $M^{\dagger}$. A necessary and sufficient condition for this convergence

[^9]in probability to occur is that $\operatorname{rank}\left(\check{M}_{n}\right)=\operatorname{rank}(M)$ with probability approaching one as $n \rightarrow \infty$ (Andrews, 1987, Theorem 2).

Here we provide a simple construction, based on the knowledge of the speed of convergence of $\check{M}_{n}$ to $M$, which results in an estimator $\hat{M}_{n}$ which is consistent for $M$ and satisfies rank $\hat{M}_{n}=\operatorname{rank} M$ with probability approaching one as $n \rightarrow \infty$ and, in consequence, $\hat{M}_{n}^{\dagger}$ is consistent for $M^{\dagger}$.

The construction proposed here is very similar to a special case of that considered by Dufour and Valéry (2016). We provide a direct proof for this construction rather than relying on Proposition 9.1 in Dufour and Valéry (2016) as the latter would require the introduction of an additional rate term ( $b_{n}$ in their notation) which satisfies a given condition (their Assumption 2.2). For our purposes we need only a single rate term (essentially the equivalent of $c_{n}$ in their notation) and thus there are fewer conditions to verify.

In particular, suppose that the sequence of (random) positive semi-definite (symmetric) matrices $\left(\check{M}_{n}\right)_{n \in \mathbb{N}}$ (of fixed dimension $\left.L \times L\right)$ satisfy

$$
\begin{equation*}
P_{n}\left(\left\|\check{M}_{n}-M_{n}\right\|_{2}<v_{n}\right) \rightarrow 1, \tag{S16}
\end{equation*}
$$

for a sequence $\left(P_{n}\right)_{n \in \mathbb{N}}$ of probability measures, a known non-negative sequence $\boldsymbol{v}_{n} \rightarrow 0$ and a sequence of deterministic matrices $M_{n} \rightarrow M$ with $\operatorname{rank}\left(M_{n}\right)=\operatorname{rank}(M)$ for all sufficiently large $n .{ }^{\text {S13 }}$ Let $\check{M}_{n}=\check{U}_{n} \check{\Lambda}_{n} \check{U}_{n}^{\prime}$ be the corresponding eigendecompositions and define

$$
\begin{equation*}
\hat{M}_{n}:=\check{U}_{n} \Lambda_{n}\left(v_{n}\right) \check{U}_{n}^{\prime} \tag{S17}
\end{equation*}
$$

where $\Lambda_{n}\left(\nu_{n}\right)$ is a diagonal matrix with the $\nu_{n}$-truncated eigenvalues of $\check{M}_{n}$ on the main diagonal and $\check{U}_{n}$ is the matrix of corresponding orthonormal eigenvectors. That is, if $\left(\check{\lambda}_{n, i}\right)_{i=1}^{L}$ denote the non-increasing eigenvalues of $\check{M}_{n}$, then the $(i, i)$-th element of $\Lambda_{n}\left(\boldsymbol{v}_{n}\right)$ is $\check{\lambda}_{n, i} \boldsymbol{1}\left(\check{\lambda}_{n, i} \geq v_{n}\right)$.

Proposition S1. If (S16) holds, $M_{n} \rightarrow M$ and for all $n$ greater than some $N \in \mathbb{N}$ $\operatorname{rank}\left(M_{n}\right)=\operatorname{rank}(M)$, then $\hat{M}_{n} \xrightarrow{P_{n}} M$ and

$$
P_{n}\left(\operatorname{rank}\left(\hat{M}_{n}\right)=\operatorname{rank}(M)\right) \rightarrow 1
$$

${ }^{\text {S13 }}$ (S16) is implied by $\left\|\check{M}_{n}-M_{n}\right\|=o_{P_{n}}\left(v_{n}\right)$ for any matrix norm. Moreover, the existence of such a sequence $\left(v_{n}\right)_{n \in \mathbb{N}}$ is guaranteed if $\left\|\check{M}_{n}-M_{n}\right\|_{2} \rightarrow 0$ in $P_{n}$-probability, however its explicit knowledge is necessary to perform the subsequent construction. In most cases $M_{n}=M$ for all $n \in \mathbb{N}$.
where $\hat{M}_{n}$ is defined as in (S17). In consequence,

$$
\hat{M}_{n}^{\dagger} \xrightarrow{P_{n}} M^{\dagger} .
$$

Proof. Throughout let $\hat{r}_{n}:=\operatorname{rank}\left(\hat{M}_{n}\right), r:=\operatorname{rank}(M), R_{n}:=\left\{\hat{r}_{n}=r\right\}$ and $\lambda_{l}, \lambda_{n, l}, \check{\lambda}_{n, l}$ and $\hat{\lambda}_{n, l}$ respectively the $l$-th largest eigenvalue of $M, M_{n}, \check{M}_{n}$ and $\hat{M}_{n}$.

Start with the case $r=0$. By Weyl's perturbation theorem (e.g. Bhatia, 1997, Corollary III.2.6) and the fact that $M_{n}=0$ for all $n$ larger than some $N \in \mathbb{N}$,

$$
P_{n}\left(R_{n}\right)=P_{n}\left(\max _{l=1, \ldots, L}\left|\check{\lambda}_{n, l}\right|<v_{n}\right) \geq P_{n}\left(\left\|\check{M}_{n}-M_{n}\right\|_{2}<v_{n}\right) \rightarrow 1
$$

On the sets $R_{n}$ we have that $\hat{M}_{n}=0=M$ and so $\hat{M}_{n} \xrightarrow{P_{n}} M$ as $P\left(R_{n}\right) \rightarrow 1$.
Now suppose that $r>0$. let $\underline{v}:=\lambda_{r} / 2>0$ and note that (S16) implies that $\| \check{M}_{n}-$ $M_{n} \|_{2}=o_{P_{n}}(1)$ and so, by Weyl's perturbation theorem, $\max _{l=1, \ldots, L}\left|\check{\lambda}_{n, l}-\lambda_{n, l}\right| \leq \| \check{M}_{n}-$ $M_{n} \|_{2}=o_{P_{n}}(1)$. Hence, defining $E_{n}:=\left\{\check{\lambda}_{n, r} \geq \nu_{n}\right\}$, for $n$ large enough such that $v_{n}<\underline{v}$ and $\left\|M_{n}-M\right\|_{2}<\underline{v} / 2$ we have

$$
P_{n}\left(E_{n}\right)=P_{n}\left(\check{\lambda}_{n, r} \geq v_{n}\right) \geq P_{n}\left(\check{\lambda}_{n, r} \geq \underline{v}\right) \geq P_{n}\left(\left|\check{\lambda}_{n, r}-\lambda_{n, r}\right|<\underline{v} / 2\right) \rightarrow 1 .
$$

If $r=L$ we have that $R_{n} \supset E_{n}$ and therefore $P_{n}\left(R_{n}\right) \rightarrow 1$. Additionally, if $\check{\lambda}_{n, L} \geq v_{n}$ then $\hat{\lambda}_{n, l}=\check{\lambda}_{n, l}$ for each $l=1, \ldots, L$ and hence $\hat{M}_{n}=\check{M}_{n}$, implying $\left\|\hat{M}_{n}-M\right\|_{2} \leq \| \check{M}_{n}-$ $M_{n}\left\|_{2}+\right\| M_{n}-M \|_{2}=o_{P_{n}}(1)$.

Now suppose instead that $r<L$ and define $F_{n}:=\left\{\check{\lambda}_{n, r+1}<\nu_{n}\right\}$. It follows by Weyl's perturbation theorem and the fact that $\lambda_{n, l}=0$ for $l>r$ and $n \geq N$ that as $n \rightarrow \infty$

$$
P_{n}\left(F_{n}\right)=P_{n}\left(\check{\lambda}_{n, r+1}<v_{n}\right) \geq P_{n}\left(\left\|\check{M}_{n}-M_{n}\right\|_{2}<v_{n}\right) \rightarrow 1 .
$$

Since $R_{n} \supset E_{n} \cap F_{n}$, this implies that $P_{n}\left(R_{n}\right) \rightarrow 1$ as $n \rightarrow \infty$. Additionally, if $\check{\lambda}_{n, r} \geq v_{n}$, $\check{\lambda}_{n, r+1}<\nu_{n}$ and $\left\|\check{M}_{n}-M\right\|_{2} \leq v$, we have that $\hat{\lambda}_{n, k}=\check{\lambda}_{n, k}$ for $k \leq r$ and $\hat{\lambda}_{n, l}=0=\lambda_{l}$ for $l>r$ and so

$$
\left\|\Lambda_{n}\left(v_{n}\right)-\Lambda\right\|_{2}=\max _{l=1, \ldots, r}\left|\hat{\lambda}_{n, l}-\lambda_{l}\right|=\max _{l=1, \ldots, r}\left|\check{\lambda}_{n, l}-\lambda_{l}\right| \leq\left\|\check{\Lambda}_{n}-\Lambda\right\|_{2} \leq\left\|\check{M}_{n}-M\right\|_{2} \leq v
$$

and hence $\left\{\left\|\check{M}_{n}-M\right\|_{2} \leq v\right\} \cap E_{n} \cap F_{n} \subset\left\{\left\|\Lambda_{n}\left(v_{n}\right)-\Lambda\right\|_{2} \leq v\right\}$, from which it follows that $\Lambda_{n}\left(v_{n}\right) \xrightarrow{P_{n}} \Lambda$ as $\left\|\check{M}_{n}-M\right\|_{2} \leq\left\|\check{M}_{n}-M_{n}\right\|_{2}+\left\|M_{n}-M\right\|_{2} \xrightarrow{P_{n}} 0$. Suppose that $\left(\lambda_{1}, \ldots, \lambda_{r}\right)$ consists of $s$ distinct eigenvalues with values $\lambda^{1}>\lambda^{2}>\cdots>\lambda^{s}$ and multiplicities $\mathfrak{m}_{1}, \ldots, \mathfrak{m}_{s}$
(each at least one). ${ }^{\text {S14 }} \lambda^{s+1}=0$ is an eigenvalue with multiplicity $\mathfrak{m}_{s+1}=L-r$. Let $l_{i}^{k}$ for $k=1, \ldots, s+1$ and $i=1, \ldots, \mathfrak{m}_{k}$ denote the column indices of the eigenvectors in $U$ corresponding to each $\lambda^{k}$. For each $\lambda^{k}$, the total eigenprojection is $\Pi_{k}:=\sum_{i=1}^{\mathfrak{m}_{k}} u_{l_{i}^{k}} u_{l_{i}^{k}}^{\prime}$. ${ }^{\text {S15 }}$ Total eigenprojections are continuous. ${ }^{\text {S16 }}$ Therefore, if we construct $\Pi_{n, k}$ in in an analogous fashion to $\Pi_{k}$ but replace columns of $U$ with columns of $\check{U}_{n}$, we have $\Pi_{n, k} \xrightarrow{P_{n}} \Pi_{k}$ for each $k=1, \ldots, s+1$ since $\check{M}_{n} \xrightarrow{P_{n}} M$. Spectrally decompose $M$ as $M=\sum_{k=1}^{s} \lambda^{k} \Pi_{k}$, where the sum runs to $s$ rather than $s+1$ since $\lambda^{s+1}=0$. Then,

$$
\hat{M}_{n}=\sum_{k=1}^{s+1} \sum_{i=1}^{\mathfrak{m}_{k}} \hat{\lambda}_{n, l_{i}^{k}} u_{n, l_{i}^{k}} u_{n, l_{i}^{k}}^{\prime}=\sum_{k=1}^{s+1} \sum_{i=1}^{\mathfrak{m}_{k}}\left(\hat{\lambda}_{n, l_{i}^{k}}-\lambda^{k}\right) u_{n, l_{i}^{k}} u_{n, l_{i}^{k}}^{\prime}+\sum_{k=1}^{s} \lambda^{k} \Pi_{n, k},
$$

whence

$$
\left\|\hat{M}_{n}-M\right\|_{2} \leq \sum_{k=1}^{s+1} \sum_{i=1}^{\mathfrak{m}_{k}}\left|\hat{\lambda}_{n, l_{i}^{k}}-\lambda^{k}\right|\left\|u_{n, l_{i}^{k}} u_{n, l_{i}^{k}}^{\prime}\right\|_{2}+\sum_{k=1}^{s}\left|\lambda^{k}\right|\left\|\Pi_{n, k}-\Pi_{k}\right\|_{2} \xrightarrow{P_{n}} 0
$$

by $\hat{\Pi}_{n, k} \xrightarrow{P_{n}} \Pi_{k}, \hat{\Lambda}_{n}\left(v_{n}\right) \xrightarrow{P_{n}} \Lambda$ and since we have $\left\|u_{n, l_{i}^{k}} u_{n, l_{i}^{\prime}}^{\prime}\right\|_{2}=1$ for any $i, k, n$. Combine this with $P_{n}\left(R_{n}\right) \rightarrow 1$ and Lemma 1 in Andrews (1987) to conclude.

## S5 Log density score estimation

In this section we discuss the details for the estimation of the log density scores $\phi_{k}$. We first provide a detailed description of the construction of the estimator (11). Secondly we provide a proofs of Lemma 4, i.e. we show that this estimate satisfies Assumption 4. Thirdly we provide proofs of Lemmas S4 and S5. The analysis here (in addition to the proposed estimator) is based on Chen and Bickel (2006) and Jin (1992), with small tweaks to fit the setup of the present paper.

## S5.1 B-spline based log density score estimation

For $\xi_{1}<\cdots<\xi_{N}$ a knot sequence, the first order B-splines are defined according to $b_{i}^{(1)}(x):=$ $\mathbf{1}_{\left[\xi_{i}, \xi_{i+1}\right)}(x)$. Subsequent order B-splines can be computed according to the recurrence relation

$$
\begin{equation*}
b_{i}^{(l)}(x)=\frac{x-\xi_{i}}{\xi_{i+l-1}-\xi_{i}} b_{i}^{(l-1)}(x)+\frac{\xi_{i+l}-x}{\xi_{i+l}-\xi_{i+1}} b_{i+1}^{(l-1)}(x) \tag{S18}
\end{equation*}
$$

[^10]for $l>1$ and $i=1, \ldots, N-l$. A $l$-th order B-spline is $l-2$ times differentiable in $x$ with first derivative
\[

$$
\begin{equation*}
c_{i}^{(l)}(x)=\frac{l-1}{\xi_{i+l-1}-\xi_{i}} b_{i}^{(l-1)}(x)-\frac{l-1}{\xi_{i+l}-\xi_{i+1}} b_{i+1}^{(l-1)}(x) . \tag{S19}
\end{equation*}
$$

\]

See de Boor (2001) for more details on B-splines.
Let $b_{k, n}=\left(b_{k, n, 1}, \ldots, b_{k, n, \mathrm{~B}_{k, n}}\right)^{\prime}$ be a collection of $\mathrm{B}_{k, n}$ cubic (i.e. 4 -th order) B-splines and let $c_{k, n}=\left(c_{k, n, 1}, \ldots, c_{k, n, \mathrm{~B}_{k, n}}\right)^{\prime}$ be their derivatives: $c_{k, n, i}(x):=\frac{\mathrm{d} b_{k, n, i}(x)}{\mathrm{d} x}$ for each $i \in$ $\left\{1, \ldots, \mathrm{~B}_{k, n}\right\}$. The knots of the splines, $\xi_{k, n}=\left(\xi_{k, n, i}\right)_{i=1}^{K_{k, n}}$ are equally spaced in $\left[\Xi_{k, n}^{L}, \Xi_{k, n}^{U}\right]$ with $\delta_{k, n}:=\xi_{k, n, i+1}-\xi_{k, n, i}>0$. ${ }^{\text {S17 }}$ For each $(k, n)$ pair the relationships between the number of knots $\left(K_{k, n}\right)$, the number of spline functions $\left(\mathrm{B}_{k, n}\right)$ and $\delta_{k, n}$ are given by $\mathrm{B}_{k, n}=K_{k, n}-4$ and $K_{k, n}=1+\left(\Xi_{k, n}^{U}-\Xi_{k, n}^{L}\right) / \delta_{k, n}$. S18

Since the B-splines vanish at infinity for any $n \in \mathbb{N}$, integration by parts gives that

$$
\begin{align*}
\int & \left(\phi_{k}(z)-\psi_{k, n}^{\prime} b_{k, n}(z)\right)^{2} \eta_{k}(z) \mathrm{d} z \\
& =\int \phi_{k}(z)^{2} \eta_{k}(z) \mathrm{d} z+\int\left(\psi_{k, n}^{\prime} b_{k, n}\right)^{2} \eta_{k}(z) \mathrm{d} z+2 \int \psi_{k, n}^{\prime} c_{k, n}(z) \eta_{k}(z) \mathrm{d} z  \tag{S20}\\
& =\mathbb{E} \phi_{k}\left(\epsilon_{k}\right)^{2}+\psi_{k, n}^{\prime} \mathbb{E}\left[b_{k, n}\left(\epsilon_{k}\right) b_{k, n}\left(\epsilon_{k}\right)^{\prime}\right] \psi_{k, n}+2 \psi_{k, n}^{\prime} \mathbb{E} c_{k, n}\left(\epsilon_{k}\right),
\end{align*}
$$

where we integrate over the support of $\phi_{k, n}$ (which is also the support of $b_{k, n}$ and $c_{k, n}$ ). This mean-squared error is minimsed by: ${ }^{\text {S19 }}$

$$
\begin{equation*}
\psi_{k, n}:=-\mathbb{E}\left[b_{k, n}\left(\epsilon_{k}\right) b_{k, n}\left(\epsilon_{k}\right)^{\prime}\right]^{-1} \mathbb{E}\left[c_{k, n}\left(\epsilon_{k}\right)\right] . \tag{S21}
\end{equation*}
$$

Replace the population expectations with sample counterparts to define the estimate of $\psi_{k, n}$

$$
\begin{equation*}
\hat{\psi}_{k, n, \gamma}:=-\left[\frac{1}{n} \sum_{i=1}^{n} b_{k, n}\left(A_{n, k, i} V_{\theta_{n}, i}\right) b_{k, n}\left(A_{n, k, i} V_{\theta_{n}, i}\right)^{\prime}\right]^{-1} \frac{1}{n} \sum_{i=1}^{n} c_{k, n}\left(A_{n, k, i} V_{\theta_{n}, i}\right) \tag{S22}
\end{equation*}
$$

where $A_{n, k, i}$ and $V_{\theta_{n}, i}$ are defined as in Assumption 4. The estimate for $\phi_{k}$ is

$$
\begin{equation*}
\hat{\phi}_{k, n, \gamma}(z):=\hat{\psi}_{k, n, \gamma}^{\prime} b_{k, n}(z) . \tag{S23}
\end{equation*}
$$

We note that computing (S23) effective only requires computing the B-spline regression coefficients $\hat{\psi}_{k, n, \gamma}$ in (S22). To implement the score test we need to estimate $K$ density scores, hence the computational cost is quite modest.

[^11]
## S5.2 Proof of Lemmas 4, S4 \& S5

Proof of Lemma 4. Under $P_{\theta_{n}}, A_{k, \gamma_{n}, i} V_{\theta_{n}, i} \simeq \epsilon_{k} \sim \eta_{k}$. We start by showing that $\hat{\phi}_{k, n}:=$ $\hat{\phi}_{k, n, \gamma_{n}}$ (where $\left.\gamma_{n}=\left(\alpha_{0}, \beta_{n}\right)\right)$ satisfies equation (35). We have

$$
\begin{align*}
\left|\frac{1}{n} \sum_{i=1}^{n} \hat{\phi}_{k, n}\left(\epsilon_{i, k}\right) Z_{n, i}-\phi_{k}\left(\epsilon_{i, k}\right) Z_{n, i}\right| \leq & \left|\frac{1}{n} \sum_{i=1}^{n}\left[\hat{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)\right] Z_{n, i}\right| \\
& +\left|\frac{1}{n} \sum_{i=1}^{n}\left[\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k, n}\left(\epsilon_{i, k}\right)\right] Z_{n, i}\right|  \tag{S24}\\
& +\left|\frac{1}{n} \sum_{i=1}^{n}\left[\phi_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k}\left(\epsilon_{i, k}\right)\right] Z_{n, i}\right|,
\end{align*}
$$

where $\phi_{k, n}:=\phi_{k} \mathbf{1}_{\left[\Xi_{k, n}^{L}, \Xi_{k, n}^{U}\right]}$ as in Assumption 3, $\tilde{\phi}_{k, n}(z):=\psi_{k, n}^{\prime} b_{k, n}(z)$ and $\hat{\phi}_{k, n}(z):=\hat{\psi}_{k, n, \gamma_{n}}^{\prime} b_{k, n}(z)$. To establish (35) it suffices to show that each of these three terms on the right hand side are $o_{\mathrm{P}}\left(n^{-1 / 2}\right) .{ }^{\mathrm{S} 20}$

For the last term in (S24), by assumption $\mathbb{E} \mathbb{1}\left\{\epsilon_{i, k} \notin\left[\Xi_{k, n}^{L}, \Xi_{k, n}^{U}\right]\right\} \downarrow 0$ and hence by independence, Cauchy-Schwarz and $\sup _{n \in \mathbb{N}} \mathbb{E} Z_{n, i}^{2}<\infty$,

$$
\begin{align*}
\mathbb{E}\left(\left[\phi_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k}\left(\epsilon_{i, k}\right)\right]^{2} Z_{n, i}^{2}\right) & =\mathbb{E}\left[\phi_{k}\left(\epsilon_{i, k}\right)^{2} \mathbf{1}\left\{\epsilon_{i, k} \notin\left[\Xi_{k, n}^{L}, \Xi_{k, n}^{U}\right]\right\}\right] \mathbb{E} Z_{n, i}^{2} \\
& \leq\left[\mathbb{E} \phi_{k}\left(\epsilon_{i, k}\right)^{4}\right]^{1 / 2}\left[\mathbb{E} \mathbf{1}\left\{\epsilon_{i, k} \notin\left[\Xi_{k, n}^{L}, \Xi_{k, n}^{U}\right]\right\}\right]^{1 / 2} \mathbb{E} Z_{n, i}^{2}  \tag{S25}\\
& \rightarrow 0 .
\end{align*}
$$

By Markov's inequality it follows that for any $v>0$,

$$
\mathrm{P}\left(\left|\frac{1}{\sqrt{n}} \sum_{i=1}^{n}\left[\phi_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k}\left(\epsilon_{i, k}\right)\right] Z_{n, i}\right|>v\right) \leq \frac{n \mathbb{E}\left(\left[\phi_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k}\left(\epsilon_{i, k}\right)\right]^{2} Z_{n, i}^{2}\right)}{n v} \rightarrow 0
$$

For the second term, we note that by our hypotheses and lemma S18 we have

$$
\begin{align*}
\mathbb{E}\left(\left[\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k, n}\left(\epsilon_{i, k}\right)\right]^{2} Z_{n, i}^{2}\right) & =\mathbb{E}\left(\left[\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k, n}\left(\epsilon_{i, k}\right)\right]^{2}\right) \mathbb{E} Z_{n, i}^{2}  \tag{S26}\\
& \leq C^{2} \delta_{k, n}^{6}\left\|\phi_{k, n}^{(3)}\right\|_{\infty}^{2} \mathbb{E} Z_{n, i}^{2} \rightarrow 0
\end{align*}
$$

as $n \rightarrow \infty$, and hence again by Markov's inequality for any $v>0$,

$$
\mathrm{P}\left(\left|\frac{1}{\sqrt{n}} \sum_{i=1}^{n}\left[\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k, n}\left(\epsilon_{i, k}\right)\right] Z_{n, i}\right|>v\right) \leq \frac{n \mathbb{E}\left(\left[\tilde{\phi}_{k}\left(\epsilon_{i, k}\right)-\phi_{k, n}\left(\epsilon_{i, k}\right)\right]^{2} Z_{n, i}^{2}\right)}{n v} \rightarrow 0
$$

[^12]For the first term, by Cauchy-Schwarz

$$
\left|\frac{1}{n} \sum_{i=1}^{n}\left[\hat{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)\right] Z_{n, i}\right| \leq\left\|\hat{\psi}_{k, n}-\psi_{k, n}\right\|_{2}\left\|\frac{1}{n} \sum_{i=1}^{n} b_{k, n}\left(\epsilon_{i, k}\right) Z_{n, i}\right\|_{2}=o_{\mathrm{P}}\left(n^{-1 / 2}\right),
$$

by lemmas S19 and S20.
Next, we show that $\hat{\phi}_{k, n}$ satisfies equation (36). We have:

$$
\begin{align*}
\frac{1}{n} \sum_{i=1}^{n}\left(\left[\hat{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k}\left(\epsilon_{i, k}\right)\right] Z_{n, i}\right)^{2} \leq \frac{4}{n} & \sum_{i=1}^{n}\left[\hat{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)\right]^{2} Z_{n, i}^{2} \\
& +\frac{4}{n} \sum_{i=1}^{n}\left[\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k, n}\left(\epsilon_{i, k}\right)\right]^{2} Z_{n, i}^{2}  \tag{S27}\\
& +\frac{4}{n} \sum_{i=1}^{n}\left[\phi_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k}\left(\epsilon_{i, k}\right)\right]^{2} Z_{n, i}^{2}
\end{align*}
$$

We will show that ( $1 / 4$ of) each of the right hand side terms is $o_{\mathrm{P}}\left(\nu_{n}\right)$ under our assumptions, which is sufficient for equation (36). For the last term, for any $v>0$, by Markov's inequality, independence and Cauchy-Schwarz we have

$$
\mathrm{P}\left(\left|\frac{1}{n} \sum_{i=1}^{n}\left[\phi_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k}\left(\epsilon_{i, k}\right)\right]^{2} Z_{n, i}^{2}\right|>v \nu_{n}\right) \lesssim \frac{\left[\mathbb{E} \mathbf{1}\left\{\epsilon_{i, k} \notin\left[\Xi_{k, n}^{L}, \Xi_{k, n}^{U}\right]\right\}\right]^{1 / 2} \mathbb{E} Z_{n, i}^{2}}{v \nu_{n}}=o(1) .
$$

For the second term, for any $v>0$, by Markov's inequality, independence and lemma S18:

$$
\begin{aligned}
\mathrm{P}\left(\left|\frac{1}{n} \sum_{i=1}^{n}\left[\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k, n}\left(\epsilon_{i, k}\right)\right]^{2} Z_{n, i}^{2}\right|>v \nu_{n}\right) & \leq \frac{\mathbb{E}\left(\left[\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k, n}\left(\epsilon_{i, k}\right)\right]^{2}\right) \mathbb{E} Z_{n, i}^{2}}{v \nu_{n}} \\
& \leq \frac{C \delta_{k, n}^{6}\left\|\phi_{k, n}^{(3)}\right\|_{\infty}^{2} \mathbb{E} Z_{n, i}^{2}}{v \nu_{n}} \\
& =o(1) .
\end{aligned}
$$

Finally, for the first term in the decomposition, by lemma S20 and Assumption 3-part (ii):

$$
\frac{1}{n} \sum_{i=1}^{n}\left[\hat{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)\right]^{2} Z_{n, i}^{2} \leq\left\|\hat{\psi}_{k, n}-\psi_{k, n}\right\|_{2}^{2}\left[\frac{1}{n} \sum_{i=1}^{n}\left\|b_{k, n}\left(\epsilon_{i, k}\right)\right\|_{2}^{2} Z_{n, i}^{2}\right]=o_{\mathrm{P}}\left(\nu_{n}\right)
$$

Proof of Lemma S4. The proof proceeds verbatim as that of Lemma 4 once references to equations (35), (36) are replaced by equations (S6), (S7) since under the conditions of the present Lemma, one still has $A_{n, \gamma_{n}, i} V_{\theta_{n}, i} \simeq \epsilon_{k} \sim \eta_{k}$ under $P_{\theta_{n}}$.

Proof of Lemma S5. We use a similar decomposition to as in the Proof of Lemma 4:

$$
\begin{align*}
\frac{1}{n} \sum_{i=1}^{n}\left(\left[\hat{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k}\left(\epsilon_{i, k}\right)\right] \epsilon_{k, i}\right)^{2} \leq \frac{4}{n} & \sum_{i=1}^{n}\left[\hat{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)\right]^{2} \epsilon_{k, i}^{2} \\
& +\frac{4}{n} \sum_{i=1}^{n}\left[\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k, n}\left(\epsilon_{i, k}\right)\right]^{2} \epsilon_{k, i}^{2}  \tag{S28}\\
& +\frac{4}{n} \sum_{i=1}^{n}\left[\phi_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k}\left(\epsilon_{i, k}\right)\right]^{2} \epsilon_{k, i}^{2} .
\end{align*}
$$

We will show that ( $1 / 4$ of) each of the right hand side terms is $o_{\mathrm{P}}\left(\nu_{n}\right)$ under our assumptions, which is sufficient for equation (S8), since under $P_{\theta_{n}}, A_{k, \gamma_{n}, i} V_{\theta_{n}, i} \simeq \epsilon_{k} \sim \eta_{k}$. For the last term, for any $v>0$, by Markov's inequality, Cauchy - Schwarz and the first additional condition in Lemma $S 5$ we have

$$
\mathrm{P}\left(\left|\frac{1}{n} \sum_{i=1}^{n}\left[\phi_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k}\left(\epsilon_{i, k}\right)\right]^{2} \epsilon_{k, i}^{2}\right|>v \nu_{n}\right) \lesssim \frac{\left(\mathbb{E}\left[\epsilon_{k, i}^{4} \mathbf{1}\left\{\epsilon_{i, k} \notin\left[\Xi_{k, n}^{L}, \Xi_{k, n}^{U}\right]\right\}\right]\right)^{1 / 2}}{v \nu_{n}}=o(1) .
$$

For the second term, first note that by Lemma S20

$$
\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)^{2} \leq\left\|\psi_{k, n}\right\|_{2}^{2}\left\|b_{k, n}\left(\epsilon_{i, k}\right)\right\|_{2}^{2} \leq\left\|\psi_{k, n}\right\|_{2}^{2} \leq\left\|\Gamma_{k, n}^{-1}\right\|_{2}^{2}\left\|C_{k, n}\right\|_{2}^{2}=O\left(\delta_{k, n}^{-3} \Delta_{k, n}\right) .
$$

Thus, for any $v>0$, by Markov's inequality, Cauchy - Schwarz, the additional conditions in Lemma S5 and Lemma S18:

$$
\begin{aligned}
& \mathrm{P}\left(\left|\frac{1}{n} \sum_{i=1}^{n}\left[\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k, n}\left(\epsilon_{i, k}\right)\right]^{2} \epsilon_{k, i}^{2}\right|>v \nu_{n}\right) \\
& \quad \leq \frac{\mathbb{E}\left(\left[\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k, n}\left(\epsilon_{i, k}\right)\right]^{2} \epsilon_{k, i}^{2}\right)}{v \nu_{n}} \\
& \quad \leq \frac{\mathrm{M}_{k, n}^{2} \mathbb{E}\left(\left[\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k, n}\left(\epsilon_{i, k}\right)\right]^{2}\right)}{v \nu_{n}}+\frac{\mathbb{E}\left[\left(\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)^{2}+\phi_{k}\left(\epsilon_{i, k}\right)^{2}\right) \epsilon_{i, k}^{2} \mathbf{1}\left\{\left|\epsilon_{i, k}\right|>\mathrm{M}_{k, n}\right\}\right]}{v \nu_{n}} \\
& \quad \lesssim \frac{\mathrm{M}_{k, n}^{2} C \delta_{k, n}^{6}\left\|\phi_{k, n}^{(3)}\right\|_{\infty}^{2}}{v \nu_{n}}+\frac{\delta_{k, n}^{-3} \Delta_{k, n} \mathbb{E}\left[\epsilon_{i, k}^{2} \mathbf{1}\left\{\left|\epsilon_{i, k}\right|>\mathrm{M}_{k, n}\right\}\right]}{v \nu_{n}}+\frac{\left[\mathbb{E}\left(\epsilon_{i, k}^{4} \mathbf{1}\left\{\left|\epsilon_{i, k}\right|>\mathrm{M}_{k, n}\right\}\right)\right]^{1 / 2}}{v \nu_{n}} \\
& \quad=o(1) .
\end{aligned}
$$

Finally, for the first term in the decomposition, by lemma $\mathrm{S} 20,\left\|b_{k, n}\left(\epsilon_{i, k}\right)\right\|_{2}^{2} \leq 1$ (e.g. de Boor,

2001, equation (36), p. 96), Assumption S2, the WLLN and Assumption 3-part (ii)

$$
\frac{1}{n} \sum_{i=1}^{n}\left[\hat{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)\right]^{2} \epsilon_{k, i}^{2} \leq\left\|\hat{\psi}_{k, n, \gamma_{n}}-\psi_{k, n}\right\|_{2}^{2}\left[\frac{1}{n} \sum_{i=1}^{n} \epsilon_{k, i}^{2}\right]=o_{\mathrm{P}}\left(\nu_{n}\right)
$$

## S5.3 Technical lemmas

Lemma S18 (Cf. Lemma A.5, Chen and Bickel, 2006). Let $\phi_{k, n}$ be defined as in Assumption 3 and $\tilde{\phi}_{k, n}:=\psi_{k, n}^{\prime} b_{k, n}$. If part (iv) of Assumption 3 holds,

$$
\mathbb{E}\left(\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k, n}\left(\epsilon_{i, k}\right)\right)^{2} \leq C^{2} \delta_{k, n}^{6}\left\|\phi_{k, n}^{(3)}\right\|_{\infty}^{2}
$$

Proof. By the definition of $\tilde{\phi}_{k, n}$ and lemma S22 we have

$$
\mathbb{E}\left(\tilde{\phi}_{k, n}\left(\epsilon_{i, k}\right)-\phi_{k, n}\left(\epsilon_{i, k}\right)\right)^{2}=\inf _{g \in \mathscr{G}_{4}\left(\xi_{k, n}\right)} \mathbb{E}\left(g\left(\epsilon_{i, k}\right)-\phi_{k, n}\left(\epsilon_{i, k}\right)\right)^{2} \leq C^{2} \delta_{k, n}^{6}\left\|\phi_{k, n}^{(3)}\right\|_{\infty}^{2}
$$

where the equality follows since $\psi_{k, n}$ is the minimiser of (S20) where we integrate over the support of $\phi_{k, n}$ (which is also the support of $b_{k, n}$ and $c_{k, n}$ ).

Lemma S19 (Cf. Lemma A.3, Chen and Bickel, 2006). Suppose assumptions 2 (or S2) and 3 hold. If $Z_{n, i}$ is independent of $\epsilon_{i, k}$ and $\sup _{n \in \mathbb{N}, i \leq 1, \ldots, n} \mathbb{E} Z_{n, i}^{2}<\infty$, then

$$
\left\|\frac{1}{n} \sum_{i=1}^{n} b_{k, n}\left(\epsilon_{i, k}\right) Z_{n, i}\right\|_{2}=O_{\mathrm{P}}\left(n^{-1 / 2}\right) .
$$

Proof. By $\sum_{m=1}^{\mathrm{B}_{k, n}} b_{k, n, m}(x)^{2} \leq 1$ (e.g. de Boor, 2001, equation (36), p. 96) and our hypotheses

$$
\mathbb{E}\left(\left\|\frac{1}{n} \sum_{i=1}^{n} b_{k, n}\left(\epsilon_{i, k}\right) Z_{n, i}\right\|_{2}^{2}\right)=\frac{1}{n} \mathbb{E}\left(\sum_{m=1}^{\mathrm{B}_{k, n}} b_{n, k, m}\left(\epsilon_{i, k}\right)^{2}\right) \mathbb{E} Z_{n, i}^{2} \leq \frac{\mathbb{E} Z_{n, i}^{2}}{n} .
$$

Fix $\epsilon>0$ and take $M>0$ large enough such that $\sup _{n \in \mathbb{N}, i \leq 1, \ldots, n} \mathbb{E} Z_{n, i}^{2} / M^{2}<\epsilon$. Markov's inequality yields

$$
\mathrm{P}\left(\sqrt{n}\left\|\frac{1}{n} \sum_{i=1}^{n} b_{k, n}\left(\epsilon_{i, k}\right) Z_{n, i}\right\|_{2}>M\right) \leq \frac{\mathbb{E}\left(n\left\|\frac{1}{n} \sum_{i=1}^{n} b_{k, n}\left(\epsilon_{i, k}\right) Z_{n, i}\right\|_{2}^{2}\right)}{M^{2}} \leq \frac{\mathbb{E} Z_{n, i}^{2}}{M^{2}}<\epsilon
$$

Lemma S20 (Cf. Lemma A.2, Chen and Bickel, 2006). Suppose that Assumptions 2 (or S2) and 3 hold. Then, for

$$
\hat{\Gamma}_{k, n}:=\frac{1}{n} \sum_{i=1}^{n} b_{k, n}\left(\epsilon_{i, k}\right) b_{k, n}\left(\epsilon_{i, k}\right)^{\prime}, \quad \Gamma_{k, n}:=\mathbb{E}\left[b_{k, n}\left(\epsilon_{k}\right) b_{k, n}\left(\epsilon_{k}\right)^{\prime}\right]
$$

and

$$
\hat{C}_{k, n}:=\frac{1}{n} \sum_{i=1}^{n} c_{k, n}\left(\epsilon_{i, k}\right), \quad C_{k, n}:=\mathbb{E}\left[c_{k, n}\left(\epsilon_{k}\right)\right],
$$

we have that
(i) $\left\|C_{k, n}\right\|_{2}=O\left(\delta_{k, n} \mathrm{~B}_{k, n}^{1 / 2}\right)$,
(ii) $\left\|\hat{C}_{k, n}-C_{k, n}\right\|_{2}=O_{\mathrm{P}}\left(\sqrt{\frac{\mathrm{B}_{k, n} \log \mathrm{~B}_{k, n}}{n \delta_{k, n}^{2}}}\right)$,
(iii) $\left\|\hat{\Gamma}_{k, n}-\Gamma_{k, n}\right\|_{2}=O_{\mathrm{P}}\left(\sqrt{\frac{\mathrm{B}_{k, n} \log \mathrm{~B}_{k, n}}{n}}\right)$,
(iv) $\left\|\Gamma_{k, n}\right\|_{2}=O\left(\delta_{k, n}\right)$
(v) $\left\|\Gamma_{k, n}^{-1}\right\|_{2}=O\left(\delta_{k, n}^{-2}\right)$.

In particular, $\left\|\hat{\Gamma}_{k, n}^{-1} \hat{C}_{k, n}-\psi_{k, n}\right\|_{2}=O_{\mathrm{P}}\left(n^{-1 / 2} \Delta_{k, n} \delta_{k, n}^{-4}\left(\Delta_{k, n} \delta_{k, n}^{-1}\right)^{\iota}\right)=o_{\mathrm{P}}(1)$ and $\left\|\hat{\Gamma}_{k, n}\right\|_{2}=o_{\mathrm{P}}(1)$.
Proof. The proof follows the relevant parts of the proof of lemma A. 2 in Chen and Bickel (2006). Firstly, from the representation of the derivative of the cubic spline (e.g. de Boor, 2001) $c_{k, n, i}=\left(b_{k, n, i}^{(3)}-b_{k, n, i+1}^{(3)}\right) / \delta_{k, n}$. We have, for large enough $n \in \mathbb{N}$,

$$
\begin{aligned}
\left|C_{k, n, i}\right|=\left|\mathbb{E}\left[c_{k, n, i}\left(\epsilon_{k}\right)\right]\right| & =\delta_{k, n}^{-1}\left|\int b_{k, n, i}^{(3)}(t) \eta_{k}(t) \mathrm{d} t-\int b_{k, n, i+1}^{(3)}(t) \eta_{k}(t) \mathrm{d} t\right| \\
& =\delta_{k, n}^{-1}\left|\int b_{k, n, i}^{(3)}(t) \eta_{k}(t) \mathrm{d} t-\int b_{k, n, i}^{(3)}(t) \eta_{k}\left(t+\delta_{k, n}\right) \mathrm{d} t\right| \\
& \leq\left|\int b_{k, n, i}^{(3)}(t) \frac{\eta_{k}\left(t+\delta_{k, n}\right)-\eta_{k}(t)}{\delta_{k, n}} \mathrm{~d} t\right| \\
& \leq 2\left\|\eta_{k}^{\prime}\right\|_{\infty} \int b_{k, n, i}^{(3)}(t) \mathrm{d} t \\
& \leq 6\left\|\eta_{k}^{\prime}\right\|_{\infty} \delta_{k, n},
\end{aligned}
$$

where the last inequality is due to (20) on p. 91 in de Boor (2001) and the fact that splines
(of any order) take values in $[0,1] .{ }^{\text {S21 }}$ It follows immediately that for large enough $n \in \mathbb{N}$,

$$
\sum_{i=1}^{\mathrm{B}_{k, n}} C_{k, n, i}^{2} \leq \sum_{i=1}^{\mathrm{B}_{k, n}} 6^{2}\left\|\eta_{k}^{\prime}\right\|_{\infty}^{2} \delta_{k, n}^{2}=\mathrm{B}_{k, n} 6^{2}\left\|\eta_{k}^{\prime}\right\|_{\infty}^{2} \delta_{k, n}^{2}
$$

from which (i) follows.
As noted above $c_{k, n, i}=\left(b_{k, n, i}^{(3)}-b_{k, n, i+1}^{(3)}\right) / \delta_{k, n}$. Since splines (of any order) take values in $[0,1]$, it follows that $c_{k, n, i} \in\left[-\delta_{k, n}^{-1}, \delta_{k, n}^{-1}\right]$. Hence, by Hoeffdings's inequality for $t \geq 0$ we have

$$
\mathrm{P}\left(\left|\frac{1}{n} \sum_{i=1}^{n} c_{k, n, m}\left(\epsilon_{i, k}\right)-\mathbb{E} c_{n, k, m}\left(\epsilon_{i, k}\right)\right| \geq t\right) \leq 2 \exp \left(\frac{-n^{2} t^{2}}{2 n \delta_{k, n}^{-2}}\right)=2 \exp \left(-n t^{2} \delta_{k, n}^{2} / 2\right)
$$

Therefore,

$$
\begin{aligned}
\mathrm{P}\left(\left\|\hat{C}_{k, n}-C_{k, n}\right\|_{2} \geq t\right) & \leq \sum_{m=1}^{\mathrm{B}_{k, n}} \mathrm{P}\left(\left|\frac{1}{n} \sum_{i=1}^{n} c_{k, n, m}\left(\epsilon_{i, k}\right)-\mathbb{E} c_{n, k, m}\left(\epsilon_{i, k}\right)\right| \geq \frac{t}{\sqrt{\mathrm{~B}_{k, n}}}\right) \\
& \leq 2 \mathrm{~B}_{k, n} \exp \left(-n t^{2} \mathrm{~B}_{k, n}^{-1} \delta_{k, n}^{2} / 2\right),
\end{aligned}
$$

and so for any fixed $\epsilon>0$ we can take $t=\sqrt{\frac{4 \mathrm{~B}_{k, n} \log \mathrm{~B}_{k, n}}{n \delta_{k, n}^{2}}}$ to obtain (ii) as then

$$
\mathrm{P}\left(\left\|\hat{C}_{k, n}-C_{k, n}\right\|_{2} \geq t\right) \leq 2 \mathrm{~B}_{k, n}^{-1} \rightarrow 0
$$

Since for any $m, s \in\left\{1, \ldots, \mathrm{~B}_{k, n}\right\}$ we have $b_{k, n, m} b_{k, n, s} \in[0,1]$ it follows by Hoeffding's inequality that for any $t \geq 0$

$$
\mathrm{P}\left(\left|\frac{1}{n} \sum_{i=1}^{n} b_{k, n, m}\left(\epsilon_{i, k}\right) b_{k, n, s}\left(\epsilon_{i, k}\right)-\mathbb{E}\left[b_{k, n, m}\left(\epsilon_{i, k}\right) b_{k, n, s}\left(\epsilon_{i, k}\right)\right]\right| \geq t\right) \leq 2 \exp \left(-2 n t^{2}\right)
$$

Therefore, since $\left\|\hat{\Gamma}_{k, n}-\Gamma_{k, n}\right\|_{2} \leq\left\|\hat{\Gamma}_{k, n}-\Gamma_{k, n}\right\|_{F}$ and both $\hat{\Gamma}_{k, n}$ and $\Gamma_{k, n}$ are zero for all ( $m, s$ )

[^13]entries where $|m-s|>3$ (de Boor, 2001, (20), p. 91) we have that
\[

$$
\begin{aligned}
& \mathrm{P}\left(\left\|\hat{\Gamma}_{k, n}-\Gamma_{k, n}\right\|_{2} \geq t\right) \\
& \quad \leq \mathrm{P}\left(\left\|\hat{\Gamma}_{k, n}-\Gamma_{k, n}\right\|_{F} \geq t\right) \\
& \quad \leq \sum_{m=1}^{\mathrm{B}_{k, n}} \sum_{s=\max (m-3,1)}^{\min \left(\mathrm{B}_{k, n}, m+3\right)} \mathrm{P}\left(\left|\frac{1}{n} \sum_{i=1}^{n} b_{k, n, m}\left(\epsilon_{i, k}\right) b_{k, n, s}\left(\epsilon_{i, k}\right)-\mathbb{E}\left[b_{k, n, m}\left(\epsilon_{i, k}\right) b_{k, n, s}\left(\epsilon_{i, k}\right)\right]\right| \geq \frac{t}{\sqrt{7 \mathrm{~B}_{k, n}}}\right) \\
& \quad \leq 14 \mathrm{~B}_{k, n} \exp \left(\frac{-2 n t^{2}}{7 \mathrm{~B}_{k, n}}\right) .
\end{aligned}
$$
\]

Putting $t=\sqrt{\frac{7 \mathrm{~B}_{k, n} \log \mathrm{~B}_{k, n}}{n}}$ we obtain (iii) as

$$
\mathrm{P}\left(\left\|\hat{\Gamma}_{k, n}-\Gamma_{k, n}\right\|_{2} \geq t\right) \leq 14 \mathrm{~B}_{k, n}^{-1} \rightarrow 0
$$

Since $\Gamma_{k, n}$ is symmetric and positive (semi-)definite we have that: ${ }^{\text {S22 }}$

$$
\left\|\Gamma_{k, n}\right\|_{2} \leq\left\|\Gamma_{k, n}\right\|_{\infty}=\max _{m=1, \ldots, \mathrm{~B}_{k, n}} \sum_{s=1}^{\mathrm{B}_{k, n}} \mathbb{E} b_{n, k, m}\left(\epsilon_{k}\right) b_{k, n, s}\left(\epsilon_{k}\right) .
$$

Then, since for any $z \in \mathbb{R}$, each row of $b_{k, n}(z) b_{k, n}(z)^{\prime}$ has at most 7 non-zero entries, ${ }^{\text {S23 }}$ all of which are bounded above by 1 we have

$$
\begin{aligned}
\left\|\Gamma_{k, n}\right\|_{2} & \leq \max _{m=1, \ldots, \mathrm{~B}_{k, n}} \sum_{s=1}^{\mathrm{B}_{k, n}} \mathbb{E} b_{n, k, m}\left(\epsilon_{k}\right) b_{k, n, s}\left(\epsilon_{k}\right) \\
& =\max _{m=1, \ldots, \mathrm{~B}_{k, n}} \sum_{s=1}^{\mathrm{B}_{k, n}} \int_{\xi_{k, n, m}}^{\xi_{k, n, m+4}} b_{k, n, m}(z) b_{k, n, s}(z) \eta_{k}(z) \mathrm{d} z \\
& \leq \max _{m=1, \ldots, B_{k, n}} 7\left\|\eta_{k}\right\|_{\infty} 4 \delta_{k, n} \\
& =28\left\|\eta_{k}\right\|_{\infty} \delta_{k, n}
\end{aligned}
$$

which yields (iv) in conjunction with requirement (iii) of Assumption 3.

By Assumption 3 part (v), on $\left[\Xi_{k, n}^{L}, \Xi_{k, n}^{U}\right]$ we have $\eta(x) \geq c \delta_{k, n}$. Hence $\eta(x)-c \delta_{k, n} \geq 0$ and so $\int b_{k, n} b_{k, n}^{\prime}\left(\eta-c \delta_{k, n}\right) \lambda=\int\left(b_{k, n} \sqrt{\eta-c \delta_{k, n}}\right)\left(b_{k, n} \sqrt{\eta-c \delta_{k, n}}\right)^{\prime} \lambda$. Note that the functions $b_{k, i} \sqrt{\eta-c \delta_{k, n}}$ satisfy $\int\left(b_{k, i} \sqrt{\eta-c \delta_{k, n}}\right)^{2} \mathrm{~d} \lambda<\infty$ and hence belong to $L_{2}(\lambda)$. It follows that

[^14]the matrix $\int b_{k, n} b_{k, n}^{\prime}\left(\eta-c \delta_{k}\right) \mathrm{d} \lambda$ is a Gram matrix and hence positive semi-definite. This implies that $\Gamma_{k, n} \succeq c \delta_{k, n} \tilde{\Gamma}_{k, n}$ where $\tilde{\Gamma}_{k, n}$ is defined as in lemma S21. Hence, by the Rayleigh quotient theorem (see e.g. Theorem 4.2.2 in Horn and Johnson, 2013) and lemma S21
$$
\lambda_{\min }\left(\Gamma_{k, n}\right) \geq \lambda_{\min }\left(c \delta_{k, n} \tilde{\Gamma}_{k, n}\right)=c \delta_{k, n} \lambda_{\min }\left(\tilde{\Gamma}_{k, n}\right) \geq c v \delta_{k, n}^{2}
$$
for a $v>0$, which may be used to conclude that (v) holds via
$$
\left\|\Gamma_{k, n}^{-1}\right\|_{2}=\frac{1}{\lambda_{\min }\left(\Gamma_{k, n}\right)} \leq(c v)^{-1} \delta_{k, n}^{-2}
$$

To demonstrate the last claim, note that with the results just derived, under our assumptions we have,

$$
\left\|\hat{C}_{k, n}\right\|_{2} \leq\left\|\hat{C}_{k, n}-C_{k, n}\right\|_{2}+\left\|C_{k, n}\right\|_{2}=O_{\mathrm{P}}\left(\sqrt{\frac{\mathrm{~B}_{k, n} \log \mathrm{~B}_{k, n}}{n \delta_{k, n}^{2}}}\right)+O\left(\delta_{k, n} \sqrt{\mathrm{~B}_{k, n}}\right)=O_{\mathrm{P}}\left(\delta_{k, n} \sqrt{\mathrm{~B}_{k, n}}\right)
$$

and, using inequality (5.8.2) from Horn and Johnson (2013),

$$
\begin{align*}
\left\|\hat{\Gamma}_{k, n}^{-1}\right\|_{2} & \leq\left\|\Gamma_{k, n}^{-1}\left(I+\left[\hat{\Gamma}_{k, n}-\Gamma_{k, n}\right] \Gamma_{k, n}^{-1}\right)^{-1}\right\|_{2} \\
& \leq\left\|\Gamma_{k, n}^{-1}\right\|_{2}\left\|\left(I+\left[\hat{\Gamma}_{k, n}-\Gamma_{k, n}\right] \Gamma_{k, n}^{-1}\right)^{-1}\right\|_{2} \\
& \leq\left\|\Gamma_{k, n}^{-1}\right\|_{2}\left(1-\left\|\left[\hat{\Gamma}_{k, n}-\Gamma_{k, n}\right] \Gamma_{k, n}^{-1}\right\|_{2}\right)^{-1}  \tag{S29}\\
& \leq\left\|\Gamma_{k, n}^{-1}\right\|_{2}\left(1-\left\|\hat{\Gamma}_{k, n}-\Gamma_{k, n}\right\|_{2}\left\|\Gamma_{k, n}^{-1}\right\|_{2}\right)^{-1} \\
& =O_{\mathrm{P}}\left(\delta_{k, n}^{-2}\right) .
\end{align*}
$$

Using these intermediate results along with (ii) - (v) and our hypotheses we obtain that

$$
\begin{aligned}
\left\|\hat{\psi}_{k, n}-\psi_{k, n}\right\|_{2} & =\left\|\hat{\Gamma}_{k, n}^{-1} \hat{C}_{k, n}-\Gamma_{k, n}^{-1} C_{k, n}\right\|_{2} \\
& \leq\left\|\left(\hat{\Gamma}_{k, n}^{-1}-\Gamma_{k, n}^{-1}\right) \hat{C}_{k, n}\right\|_{2}+\left\|\Gamma_{k, n}^{-1}\left(\hat{C}_{k, n}-C_{k, n}\right)\right\|_{2} \\
& \leq\left\|\Gamma_{k, n}^{-1}\right\|_{2}\left\|\Gamma_{k, n}-\hat{\Gamma}_{k, n}\right\|_{2}\left\|\hat{\Gamma}_{k, n}^{-1}\right\|_{2}\left\|\hat{C}_{k, n}\right\|_{2}+\left\|\Gamma_{k, n}^{-1}\right\|_{2}\left\|\hat{C}_{k, n}-C_{k, n}\right\|_{2} \\
& =O_{\mathrm{P}}\left(\sqrt{\frac{\mathrm{~B}_{k, n}^{2} \log \mathrm{~B}_{k, n}}{\delta_{k, n}^{6} n}}\right)+O_{\mathrm{P}}\left(\sqrt{\frac{\mathrm{~B}_{k, n} \log \mathrm{~B}_{k, n}}{\delta_{k, n}^{6} n}}\right) \\
& =o_{\mathrm{P}}(1)
\end{aligned}
$$

by Assumption 3 part (ii), since we have $\mathrm{B}_{k, n} \leq \Delta_{k, n} \delta_{k, n}^{-1}$ and hence the dominant term above
vanishes since for all large enough $n$,

$$
\sqrt{\frac{\mathrm{B}_{k, n}^{2} \log \mathrm{~B}_{k, n}}{\delta_{k, n}^{6} n}} \leq n^{-1 / 2} \Delta_{k, n} \delta_{k, n}^{-4} \log \left(\Delta_{k, n} \delta_{k, n}^{-1}\right) \leq n^{-1 / 2} \Delta_{k, n} \delta_{k, n}^{-4}\left(\Delta_{k, n} \delta_{k, n}^{-1}\right)^{\iota}=o(1)
$$

Finally, by (iii) and (iv) and Assumption 3 part (ii) we have

$$
\left\|\hat{\Gamma}_{k, n}\right\|_{2} \leq\left\|\hat{\Gamma}_{k, n}-\Gamma_{k, n}\right\|_{2}+\left\|\Gamma_{k, n}\right\|_{2}=O_{\mathrm{P}}\left(\sqrt{\frac{\mathrm{~B}_{k, n} \log \mathrm{~B}_{k, n}}{n}}\right)+O\left(\delta_{k, n}\right)=o_{\mathrm{P}}(1),
$$

since $\delta_{k, n} \rightarrow 0$ and for all large enough $n$,

$$
\sqrt{\frac{\mathrm{B}_{k, n} \log \mathrm{~B}_{k, n}}{n}} \leq n^{-1 / 2} \Delta_{k, n} \delta_{k, n}^{-1} \log \left(\Delta_{k, n} \delta_{k, n}^{-1}\right) \leq \delta_{k, n}^{3} n^{-1 / 2} \Delta_{k, n} \delta_{k, n}^{-4}\left(\Delta_{k, n} \delta_{k, n}^{-1}\right)^{\iota}=o(1)
$$

Lemma S21. The smallest eigenvalue of the $\mathrm{B}_{k, n} \times \mathrm{B}_{k, n}$ Gram matrix $\tilde{\Gamma}_{k, n}:=\int b_{k, n} b_{k, n}^{\prime} \mathrm{d} \lambda$ satisfies

$$
\lambda_{\min }\left(\tilde{\Gamma}_{k, n}\right) \geq v \delta_{k, n}>0
$$

for a $v>0$.
Proof. Since $b_{k, n, m}(x) b_{k, n, s}(x)$ is non-zero only for $|m-s| \leq 3$ and each $b_{k, n, m}$ is non-zero only on $\left[\xi_{m, k, n}, \xi_{m+4, k, n}\right)$ (e.g. (20) p. 91 of de Boor, 2001), $\tilde{\Gamma}_{k, n}$ is a symmetric banded Toeplitz matrix. ${ }^{\text {S24 }}$ Its entries can be computed by direct integration:

$$
\left[\tilde{\Gamma}_{k, n}\right]_{m, s}=\delta_{k, n} \times\left\{\begin{array}{ll}
\frac{151}{315} & \text { if } m=s \\
\frac{397}{1680} & \text { if }|m-s|=1 \\
\frac{1}{42} & \text { if }|m-s|=2 . \\
\frac{1}{5040} & \text { if }|m-s|=3 \\
0 & \text { if }|m-s|>3
\end{array} .\right.
$$

Let $f_{0}:=\frac{151}{315}, f_{1}:=f_{-1}:=\frac{397}{1680}, f_{2}:=f_{-2}:=\frac{1}{42}$ and $f_{3}:=f_{-3}:=\frac{1}{5040}$ and let $f_{s}:=0$ for $|s|>3$. Now, let $f(\theta):=\sum_{s=-3}^{3} f_{s} e^{i(s \theta)}$. Then, $\tilde{\Gamma}_{k, n} / \delta_{k, n}$ is then the matrix generated by $f$ in the sense that $\tilde{\Gamma}_{k, n} / \delta_{k, n}=\mathscr{T}_{n}(f):=\sum_{s=-\min \left(\mathrm{B}_{k, n}-1,3\right)}^{\min \left(\mathrm{B}_{k, n}-1,3\right.} f_{k} J_{n}^{s}$ where each $J_{n}^{s}$ is the $\mathrm{B}_{k, n} \times \mathrm{B}_{k, n}$ matrix which is zero everywhere except for the $(i, j)$-th entries where $i-j=s$, where it has a value of $1 .{ }^{\text {S25 }}$ Since $f \in L_{1}([-\pi, \pi])$ and is real on $[-\pi, \pi]$ by Theorem 6.1 in Garoni and Serra-

[^15]Capizzano (2017) we have that $\lambda_{\min }\left(\tilde{\Gamma}_{k, n}\right)=\delta_{k, n} \lambda_{\min }\left(\tilde{\Gamma}_{k, n} / \delta_{k, n}\right) \geq \delta_{k, n} \inf _{\theta \in[-\pi, \pi]} f(\theta)=\delta_{k, n} v$, where $v:=\inf _{\theta \in[-\pi, \pi]} f(\theta)>0$.

Lemma S22. Suppose $\xi \in \mathbb{R}^{N+1}$ such that $a=\xi_{0}<\xi_{1}<\cdots<\xi_{N}=b$, $h:=\max _{i \in[N]} \xi_{i}-$ $\xi_{i-1}$, and let $\mathscr{G}_{l}(\xi)$ be the linear space formed by degree $l$ splines with knots $\xi$. Then, if $f \in C^{l-1}[a, b]$ we have that

$$
\inf _{g \in \mathscr{G}_{l}(\xi)}\|g-f\|_{\infty} \leq \frac{(l+1)!}{2^{l}} h^{l-1}\left\|f^{(l-1)}\right\|_{\infty}=c_{l} h^{l-1}\left\|f^{(l-1)}\right\|_{\infty}
$$

where $c_{l}$ depends only on $l$.
Proof. This is a special case of Theorem 20.3 in Powell (1981).

## S6 Power optimality under strong identification

In either the setting considered in the main text or that introduced in Section S1, consider local alternatives of the type given in (17). We now prove the limiting power statements claimed in equations (18), (19) and (20).

Proposition S2. Suppose that Assumptions 1, 2 and 3 (or S1, S2, S3 and S4) hold, $\alpha \in \mathbb{R}$ and $\tilde{\mathcal{I}}_{\theta}>0$. Then, (18) holds.

Proof. Apply Proposition S3 in the case where $L_{\alpha}=1$ to obtain

$$
\lim _{n \rightarrow \infty} P_{\theta_{n}(q, d, h)}^{n} \varphi_{n}=1-\mathrm{P}\left(\chi_{1}^{2}\left(\tilde{\mathcal{I}}_{\theta} q^{2}\right) \leq c_{a}\right) .
$$

The right hand side is the power function of the test $\psi(Z):=\mathbf{1}\left\{Z^{2}>c_{a}\right\}$ for $Z \sim \mathcal{N}\left(\tilde{\mathcal{I}}_{\theta}^{1 / 2} q, 1\right)$. If $X=Z-\tilde{\mathcal{I}}_{\theta}^{1 / 2} q$, then

$$
\psi(Z)=\mathbf{1}\left\{\left(X-\tilde{\mathcal{I}}_{\theta}^{1 / 2} q\right)^{2}>c_{a}\right\}=\mathbf{1}\left\{\left|X-\tilde{\mathcal{I}}_{\theta}^{1 / 2} q\right|>z_{a / 2}\right\}, \quad X \sim \mathcal{N}(0,1)
$$

hence $\mathbb{E} \psi(Z)$ is (18).
Proposition S3. Suppose that Assumptions 1, 2 and 3 (or S1, S2, S3 and S4) hold and $\tilde{\mathcal{I}}_{\theta}$ is positive definite. Then, (19) holds.

Proof. The proof of Theorem 1 (or Theorem S1) showed that the conditions of Theorem 2 hold. Therefore, by (40), $c_{n}$ is equal to the $1-a$ quantile of a $\chi_{L_{\alpha}}^{2}$ distribution with
probability 1 for all large enough $n$. By (38), (39), Le Cam's third lemma (e.g. Example 6.7 in van der Vaart (1998)) and Theorem 12.14 in Rudin (1991),

$$
\sqrt{n} \mathbb{P}_{n} \hat{\kappa}_{n, \bar{\gamma}_{n}} \rightsquigarrow \mathcal{N}\left(\tilde{\mathcal{I}}_{\theta} q, \tilde{\mathcal{I}}_{\theta}\right) \quad \text { under } P_{\theta_{n}(q, d, h)}^{n}
$$

By condition 3, the mutual contiguity which follows from (38) and Example 6.5 in van der Vaart (1998), Proposition S1 and Theorem 9.2.3 in Rao and Mitra (1971)

$$
\hat{S}_{n, \bar{\gamma}_{n}} \rightsquigarrow \chi_{L_{\alpha}}^{2}\left(q^{\prime} \tilde{\mathcal{I}}_{\theta} q\right) \quad \text { under } P_{\theta_{n}(q, d, h)}^{n}
$$

from which the result follows.
Proposition S4. Suppose that Assumptions 1, 2 and 3 (or S1, S2, S3 and S4) hold and $\tilde{\mathcal{I}}_{\theta}$ is positive definite. Then, (20) holds.

Proof. By arguing exactly as in Proposition S 3 with convergent sequences $\left(q_{n}, g_{n}, h_{n}\right) \rightarrow$ $(q, d, h)$ replacing the fixed $(q, d, h)$ in that Proposition one obtains that

$$
\hat{S}_{n, \bar{\gamma}_{n}} \rightsquigarrow \chi_{L_{\alpha}}^{2}\left(q^{\prime} \tilde{\mathcal{I}}_{\theta} q\right) \quad \text { under } P_{\theta_{n}\left(q_{n}, d_{n}, h_{n}\right)}^{n}
$$

and hence

$$
\begin{equation*}
\lim _{n \rightarrow \infty} P_{\theta_{n}\left(q_{n}, g_{n}, h_{n}\right)}^{n} \varphi_{n}=1-\mathrm{P}\left(\chi_{L_{\alpha}}^{2}\left(q^{\prime} \tilde{\mathcal{I}}_{\theta} q\right) \leq c_{a}\right) \tag{S30}
\end{equation*}
$$

with $c_{a}$ the $1-a$ quantile of a $\chi_{L_{\alpha}}^{2}$ distribution. The proof is completed by a standard subsequence argument. Note first that the map $(q, d, h) \mapsto q^{\prime} \tilde{\mathcal{I}}_{\theta} q$ from $\mathcal{V} \rightarrow \mathbb{R}$ is continuous. As $K_{u}^{\star}$ is compact this function attains its infimum, hence

$$
u=\inf \left\{q^{\prime} \tilde{\mathcal{I}}_{\theta} q:(q, d, h) \in K_{u}^{\star}\right\}=\min \left\{q^{\prime} \tilde{\mathcal{I}}_{\theta} q:(q, d, h) \in K_{u}^{\star}\right\}
$$

Taking $\left(q_{\star}, d_{\star}, h_{\star}\right) \in K_{u}^{\star}$ such that $q_{\star} \tilde{\mathcal{I}}_{\theta} q_{\star}=u$, we have by (S30)

$$
\begin{equation*}
\limsup _{n \rightarrow \infty} \inf _{(q, d, h) \in K_{u}^{\star}} P_{\theta_{n}(q, d, h)}^{n} \varphi_{n} \leq \lim _{n \rightarrow \infty} P_{\theta_{n}\left(q_{\star}, d_{\star}, h_{\star}\right)}^{n} \varphi_{n}=1-\mathrm{P}\left(\chi_{L_{\alpha}}^{2}(u) \leq c_{a}\right)=: \mathcal{R} . \tag{S31}
\end{equation*}
$$

There is a sequence $\left(v_{n}\right)_{n \in \mathbb{N}} \subset K_{u}^{\star}$ and a subsequence $\left(n_{j}\right)_{j \in \mathbb{N}}$ such that

$$
\lim _{j \rightarrow \infty} v_{n_{j}}=v_{\star}=\left(q_{\star}, d_{\star}, h_{\star}\right) \in K_{u}^{\star}
$$

and

$$
\begin{equation*}
\mathcal{S}:=\liminf _{n \rightarrow \infty} \inf _{(q, d, h) \in K_{u}^{*}} P_{\theta_{n}(q, d, h)}^{n} \varphi_{n}=\lim _{j \rightarrow \infty} P_{\theta_{n_{j}}\left(q_{n_{j}}, d_{n_{j}}, h_{n_{j}}\right)}^{n_{j}} \varphi_{n_{j}} . \tag{S32}
\end{equation*}
$$

Construct a new sequence $\left(v_{m}^{*}\right)_{m \in \mathbb{N}}$ as follows. For all $m \in\left[n_{j}, n_{j+1}\right) \cap \mathbb{N}$ for some $j \in \mathbb{N}$ put $v_{m}^{*}=v_{n_{j}}$ and for $m=1, \ldots, n_{1}$ put $v_{m}^{*}=v_{n_{1}}$. By construction $\lim _{m \rightarrow \infty} v_{m}^{*}=v_{\star}$. By (S30)

$$
\lim _{m \rightarrow \infty} P_{\theta_{m}\left(v_{m}^{*}\right)}^{m} \varphi_{m}=1-\mathrm{P}\left(\chi_{r}^{2}\left(u^{\star}\right) \leq c_{a}\right) \geq \mathcal{R}, \quad \text { with } \quad u^{\star}=\left(q_{\star}\right)^{\prime} \tilde{\mathcal{I}}_{\theta} q_{\star} \geq u
$$

For any $\varepsilon>0$, there is a $M \in \mathbb{N}$ such that if $m \geq M, P_{\theta_{m}\left(v_{m}^{*}\right)}^{m} \varphi_{m} \geq \mathcal{R}-\varepsilon$ by the preceding display. Taking a subsequence $n_{j_{k}}$ such that for all $k \in \mathbb{N}$ we have $m_{k}=n_{j_{k}} \geq M$ gives

$$
\mathcal{S}=\mathcal{S}-P_{\theta_{n_{j_{k}}}\left(v_{n_{j_{k}}}^{*}\right)}^{n_{j_{2}}} \varphi_{n_{j_{k}}}+P_{\theta_{m_{k}}\left(v_{m_{k}}^{*}\right)}^{m_{k}} \varphi_{m_{k}} \geq \mathcal{S}-P_{\theta_{n_{j_{k}}}\left(v_{n_{j_{k}}}^{*}\right.}^{n_{j_{k}}} \varphi_{n_{j_{k}}}+\mathcal{R}-\varepsilon .
$$

Take $k \rightarrow \infty$ to conclude (via (S32)) that $\mathcal{S} \geq \mathcal{R}-\varepsilon$. Since $\varepsilon>0$ was arbitrary, it follows that $\mathcal{S} \geq \mathcal{R}$. Combine with equations (S31) and (S32) to obtain (20).

## S7 Additional simulation results

In this section we provide a number of additional simulation results.

## S7.1 Truncation in the baseline model

In our main simulations we truncated the effective information matrix estimate at machine precision, i.e. $\nu_{n}^{1 / 2}=10^{-308}$. Here we investigate the sensitivity of the rejection frequencies to this choice. Specifically, we replicate Table 2 from the main text, fixing $B=6$, but allowing for different truncation rates $\nu_{n}^{1 / 2}=10^{-308}, 10^{-5}, 10^{-1}$. S26 The value $10^{-1}$ is a high truncation value which implies that we end up truncating often when all densities are Gaussian. The results are shown in Table S1.

We find that the results are not sensitive to the truncation parameter choice. Comparing machine precision to $\nu_{n}^{1 / 2}=10^{-5}$ yields no differences at all, whereas $\nu_{n}^{1 / 2}=10^{-1}$ makes the test slightly conservative. Closer inspection reveals that the under rejection is due to cases where all eigenvalues are truncated and hence $\operatorname{rank}\left(\hat{\mathcal{I}}_{\hat{\gamma}}^{t}\right)=0$. In Theorem 1 this corresponds to the conservative case.

## S7.2 Additional power results for the baseline model

Figure 4 in the main text compared the power of different tests for the baseline model $Y_{i}=A^{-1} \epsilon_{i}$ for the case where $n=1000$. Here we show the results for $n=200$ and $n=500$.

[^16]Specifically, Figures S1 and S2 show the results.
Overall, the patterns that we find are similar as in the main text. One thing that stands out is that the $S^{\text {gmm }}$ test over-rejects for these smaller sample sizes, essentially confirming the results in Table 3. It is possible that a more careful selection of the relevant higher order moments will improve this finding.

Besides this our two main findings from the main text hold. First, the standard LM test is the preferred approach whenever the true density is known, but the semi-parametric score test comes close in terms of power. Second, for all other densities the semi-parametric score test shows the highest power.

## S7.3 Additional power results for the LSEM

Figure 5 in the main text compared the power of different tests for the LSEM model for the case where $n=1000$. Here we show the results for $n=200$ and $n=500$. Specifically, Figures S3 and S4 show the results.

We find that for $n=200$ the power of tests is generally quite low, indicating that for small sample sizes little can be learned by exploiting deviations from the Gaussian density. This holds most notably for the Student's $t$ densities, the skewed unimodal density and the bimodal density. Intuitively, given a small sample these densities are hard to distinguish from the normal density and little can be learned about the parameter $\alpha$. A reassuring finding is that the null rejection frequency of the test remains well controlled. These findings persist when we increase to $n=500$, though the power does improve as one would expect.

Overall, the implementing the test with one-step efficient estimates leads to higher power, but the null rejection frequency of the test is controlled less well. Therefore we recommend using OLS estimates for $\beta$ when the sample size is small.

## S7.4 Heteroskedastic LSEM model

In this section we study the empirical rejection frequency (under the null) of the semiparametric score test for the heteroskedastic baseline model. Specifically we consider

$$
\begin{equation*}
Y_{i}=A\left(\alpha, \sigma, X_{i}\right)^{-1} \epsilon_{i} \quad A\left(\alpha, \sigma, X_{i}\right)^{-1}=L(\sigma) D\left(\sigma, \tilde{X}_{i}\right)^{1 / 2} R(\alpha)^{\prime}, \tag{S33}
\end{equation*}
$$

where $R(\alpha)$ is a rotation matrix parametrized by the Cayley transformation of a skewsymmetric matrix (e.g. Gouriéroux, Monfort and Renne, 2017), $L(\sigma)$ is lower triangular with positive diagonal elements and $D\left(\sigma, \tilde{X}_{i}\right)$ is a diagonal matrix with diagonal elements
given by

$$
\left[D\left(\sigma, \tilde{X}_{i}\right)\right]_{j j}=\exp \left(\sigma_{j 1}^{\prime} \tilde{X}_{i}\right), \quad j=1, \ldots, K
$$

where $\sigma_{j 1}$ is a $(d-1) \times 1$ parameter vector. Note that the average scaling of the errors is captured by $L(\sigma)$ and $D\left(\sigma, \tilde{X}_{i}\right)$ is the only heteroskedastic part. More elaborate specifications that allow off-diagonal elements of $L$ to depend on $X_{i}$ are also possible.

The results for different sample sizes, dimensions $K$ and number of explanatory variables are shown in Table S 2 . Overall, we find a similar pattern as for the LSEM model from the main text (cf Table 4). When $K=5$ and the sample size is small, i.e. $n=200$, the test tends to over-reject. The over-rejection vanishes for larger sample sizes. A slight difference is observed for heavy tailed densities (e.g. $t(5)$ ) where even with $n=1000$ there is still some over-rejection.

## S8 Additional empirical results

In this section we present some additional results for the returns to schooling application of section 6. Specifically, we consider the more flexible model from Section S1 which allows for conditional heteroskedasticity.

Starting from the baseline linear IV model with a possibly scalar endogenous instrument:

$$
\begin{align*}
y_{i} & =\alpha_{1} w_{i}+b_{y}^{\prime} X_{i}+u_{i} \\
w_{i} & =\pi z_{i}+b_{w}^{\prime} X_{i}+v_{i}  \tag{S34}\\
z_{i} & =B_{z} X_{i}+\left(\alpha_{2} / \sigma_{u}\right) u_{i}+e_{i}
\end{align*}
$$

We now allow the scaling of the errors $\sigma_{u}, \sigma_{v}$ and $\sigma_{e}$ to be a flexible functions of $X_{i}$. Specifically, we follow Wooldridge (2012, Chapter 8) and model the scales using flexible functions, i.e.

$$
\sigma_{j}\left(X_{i}\right)=\sigma_{j, 0} \exp \left(\sigma_{j 1} \tilde{X}_{i, 1}+\ldots+\sigma_{j d} \tilde{X}_{i, d-1}\right), \quad j=u, v, e
$$

see also Romano and Wolf (2017) for more elaborate specifications. The coefficients $\sigma_{i k}$ are estimated along with the other well identified parameters. Following (23) we write the model in our general form

$$
\begin{align*}
Y_{i} & =B X_{i}+A^{-1}\left(\alpha, \sigma, X_{i}\right) \epsilon_{i},  \tag{S35}\\
A^{-1}\left(\alpha, \sigma, X_{i}\right) & =\left[\begin{array}{ccc}
\sigma_{u}\left(X_{i}\right)+\alpha_{1} \sigma_{v}\left(X_{i}\right) \rho+\alpha_{1} \pi \alpha_{2} & \alpha_{1} \sqrt{1-\rho^{2}} \sigma_{v}\left(X_{i}\right) & \alpha_{1} \pi \sigma_{e}\left(X_{i}\right) \\
\rho \sigma_{v}\left(X_{i}\right)+\pi^{\prime} \alpha_{2} & \sqrt{1-\rho^{2}} \sigma_{v}\left(X_{i}\right) & \pi^{\prime} \sigma_{e}\left(X_{i}\right) \\
\alpha_{2} & 0 & \sigma_{e}\left(X_{i}\right)
\end{array}\right],
\end{align*}
$$

which shows that the model is a special case of (S1). For this specification we reconstruct the confidence set for $\alpha=\left(\alpha_{1}, \alpha_{2}\right)$. The result is shown in Figure S5.

We find that the confidence region is quite similar when compared to the homoskedastic one. The volume is slightly smaller and there is more mass on the probability that $\alpha_{2}$ is positive. Importantly however, the main conclusion remains the same. Even when relaxing the instrument validity assumption the effect of education is positive and quite precisely identified.

An obvious caveat is that this result is obtained under the additional assumption that the model for heteroskedasticity is correctly specified. An open question is how to handle model mis-specification in the class semi-parametric LSEM models. We leave this for future research.

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Figure S1: Power Comparison Baseline model $n=200$


Notes: Empirical power curves for the baseline model with $k=2$ and $n=200$. Each plot corresponds to the choice for densities $\epsilon_{k}$, for $k \geq 2$, where the numbers correspond to the different densities listed in Figure 3. The solid red line corresponds to $S_{\hat{\gamma}}$, the dashed blue line to $\mathrm{LM}^{\mathrm{mle}}$, the dotted pink line to $\mathrm{LM}^{\mathrm{pmle}}$ and the dot-dashed green line to $S^{g m m}$.

Figure S2: Power Comparison Baseline model $n=500$


Notes: Empirical power curves for the baseline model with $k=2$ and $n=500$. Each plot corresponds to the choice for densities $\epsilon_{k}$, for $k \geq 2$, where the numbers correspond to the different densities listed in Figure 3 . The solid red line corresponds to $S_{\hat{\gamma}}$, the dashed blue line to $\mathrm{LM}^{\mathrm{mle}}$, the dotted pink line to $\mathrm{LM}^{\mathrm{pmle}}$ and the dot-dashed green line to $S^{g m m}$.

Figure S3: Power LSEM $n=200$


Notes: Empirical power curves for the LSEM model with $k=2, d=2$ and $n=200$. Each plot corresponds to the choice for densities $\epsilon_{i, k}$, for $k \geq 2$, where the numbers correspond to the different densities shown in Figure 3. The solid red line corresponds to the empirical rejection frequency of the $\hat{S}_{\hat{\gamma}}$ test where $\hat{\gamma}=\left(\alpha_{0}, \hat{\beta}\right)$, with $\hat{\beta}$ the OLS estimator. The dashed blue line corresponds to the empirical rejection frequency of the $\hat{S}_{\hat{\gamma}}$ test where $\hat{\gamma}=\left(\alpha_{0}, \hat{\beta}\right)$, with $\hat{\beta}$ the one-step efficient MLE estimator.

Figure S4: Power LSEM $n=500$


Notes: Empirical power curves for the LSEM model with $k=2, d=2$ and $n=500$. Each plot corresponds to the choice for densities $\epsilon_{i, k}$, for $k \geq 2$, where the numbers correspond to the different densities shown in Figure 3. The solid red line corresponds to the empirical rejection frequency of the $\hat{S}_{\hat{\gamma}}$ test where $\hat{\gamma}=\left(\alpha_{0}, \hat{\beta}\right)$, with $\hat{\beta}$ the OLS estimator. The dashed blue line corresponds to the empirical rejection frequency of the $\hat{S}_{\hat{\gamma}}$ test where $\hat{\gamma}=\left(\alpha_{0}, \hat{\beta}\right)$, with $\hat{\beta}$ the one-step efficient MLE estimator.

## Figure S5: Confidence sets: Returns to schooling



Notes: We show $95 \%$ (light gray) and $67 \%$ (dark gray) confidence sets for $\alpha=\left(\alpha_{1}, \alpha_{2}\right)$, where $\alpha_{1}$ captures the effect of education on $\log$ wages and $\alpha_{2}$ capture the correlation between the instrument (proximity to schooling interacted with parental education) and the error of the log wage equation. The red line indicates the confidence interval under the restriction of instrument exogeneity, i.e. $\alpha_{2}=0$. Figure (a) shows the result after inverting the $\hat{S}_{\hat{\gamma}}$ test statistic with heteroskedastic errors. Figure (b) shows the result after inverting the same test statistic but with homoskedastic errors.

Table S1: Rejection Frequencies $\hat{S}_{\hat{\gamma}}$ Test for Baseline model: truncation

| $n$ | $K$ | $\nu_{n}^{1 / 2}$ | 1 | 2 | 3 | 4 | 5 | 6 | 7 | 8 | 9 | 10 |
| :--- | :--- | :--- | :--- | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: |
| 200 | 2 | $10^{-308}$ | 0.051 | 0.047 | 0.048 | 0.041 | 0.050 | 0.049 | 0.047 | 0.049 | 0.050 | 0.044 |
| 200 | 2 | $10^{-5}$ | 0.051 | 0.047 | 0.048 | 0.041 | 0.050 | 0.049 | 0.047 | 0.049 | 0.050 | 0.044 |
| 200 | 2 | $10^{-1}$ | 0.051 | 0.047 | 0.048 | 0.041 | 0.050 | 0.049 | 0.047 | 0.049 | 0.050 | 0.044 |
| 200 | 3 | $10^{-308}$ | 0.046 | 0.041 | 0.049 | 0.036 | 0.045 | 0.052 | 0.046 | 0.048 | 0.049 | 0.047 |
| 200 | 3 | $10^{-5}$ | 0.046 | 0.041 | 0.049 | 0.036 | 0.045 | 0.052 | 0.046 | 0.048 | 0.049 | 0.047 |
| 200 | 3 | $10^{-1}$ | 0.046 | 0.043 | 0.049 | 0.036 | 0.044 | 0.052 | 0.046 | 0.049 | 0.049 | 0.045 |
| 200 | 5 | $10^{-308}$ | 0.034 | 0.040 | 0.037 | 0.037 | 0.034 | 0.044 | 0.041 | 0.048 | 0.044 | 0.042 |
| 200 | 5 | $10^{-5}$ | 0.034 | 0.040 | 0.037 | 0.037 | 0.034 | 0.044 | 0.041 | 0.048 | 0.044 | 0.042 |
| 200 | 5 | $10^{-1}$ | 0.041 | 0.039 | 0.040 | 0.036 | 0.037 | 0.047 | 0.042 | 0.050 | 0.044 | 0.040 |
| 500 | 2 | $10^{-308}$ | 0.050 | 0.044 | 0.052 | 0.045 | 0.051 | 0.052 | 0.052 | 0.043 | 0.049 | 0.049 |
| 500 | 2 | $10^{-5}$ | 0.050 | 0.044 | 0.052 | 0.045 | 0.051 | 0.052 | 0.052 | 0.043 | 0.049 | 0.049 |
| 500 | 2 | $10^{-1}$ | 0.050 | 0.044 | 0.052 | 0.045 | 0.031 | 0.052 | 0.052 | 0.043 | 0.049 | 0.049 |
| 500 | 3 | $10^{-308}$ | 0.048 | 0.046 | 0.040 | 0.047 | 0.050 | 0.055 | 0.054 | 0.047 | 0.051 | 0.048 |
| 500 | 3 | $10^{-5}$ | 0.048 | 0.046 | 0.040 | 0.047 | 0.050 | 0.055 | 0.054 | 0.047 | 0.051 | 0.048 |
| 500 | 3 | $10^{-1}$ | 0.038 | 0.048 | 0.042 | 0.045 | 0.050 | 0.055 | 0.054 | 0.047 | 0.051 | 0.051 |
| 500 | 5 | $10^{-308}$ | 0.042 | 0.038 | 0.041 | 0.039 | 0.045 | 0.050 | 0.040 | 0.050 | 0.052 | 0.043 |
| 500 | 5 | $10^{-5}$ | 0.042 | 0.038 | 0.041 | 0.039 | 0.045 | 0.050 | 0.040 | 0.050 | 0.052 | 0.043 |
| 500 | 5 | $10^{-1}$ | 0.043 | 0.034 | 0.050 | 0.040 | 0.047 | 0.051 | 0.041 | 0.050 | 0.052 | 0.042 |
| 1000 | 2 | $10^{-308}$ | 0.056 | 0.048 | 0.045 | 0.047 | 0.050 | 0.053 | 0.049 | 0.049 | 0.045 | 0.050 |
| 1000 | 2 | $10^{-5}$ | 0.056 | 0.048 | 0.045 | 0.047 | 0.050 | 0.053 | 0.049 | 0.049 | 0.045 | 0.050 |
| 1000 | 2 | $10^{-1}$ | 0.010 | 0.048 | 0.041 | 0.047 | 0.050 | 0.053 | 0.049 | 0.049 | 0.045 | 0.050 |
| 1000 | 3 | $10^{-308}$ | 0.046 | 0.044 | 0.046 | 0.042 | 0.049 | 0.050 | 0.046 | 0.051 | 0.049 | 0.047 |
| 1000 | 3 | $10^{-5}$ | 0.046 | 0.040 | 0.046 | 0.042 | 0.049 | 0.050 | 0.046 | 0.051 | 0.049 | 0.047 |
| 1000 | 3 | $10^{-1}$ | 0.039 | 0.044 | 0.035 | 0.043 | 0.049 | 0.050 | 0.046 | 0.050 | 0.049 | 0.047 |
| 1000 | 5 | $10^{-308}$ | 0.044 | 0.042 | 0.043 | 0.038 | 0.045 | 0.050 | 0.043 | 0.050 | 0.049 | 0.046 |
| 1000 | 5 | $10^{-5}$ | 0.044 | 0.042 | 0.043 | 0.038 | 0.045 | 0.050 | 0.043 | 0.050 | 0.049 | 0.046 |
| 1000 | 5 | $10^{-1}$ | 0.043 | 0.050 | 0.044 | 0.036 | 0.050 | 0.053 | 0.042 | 0.053 | 0.049 | 0.047 |
|  |  |  |  |  |  |  |  |  |  |  |  |  |

Notes: The table shows the empirical rejection frequencies for the $S_{\hat{\gamma}}$ test based on $S=5,000$ Monte Carlo replications for the baseline model $Y_{i}=A^{-1} \epsilon_{i}$. The test has nominal level $a=0.05$. The columns denote the sample size $n$, the dimension of the model $K$, the truncation rate $\nu_{n}^{1 / 2}$ and the choice for densities $\epsilon_{i k}$, for $k \geq 2$, where the numbers correspond to the different densities shown in Figure 3.

Table S2: Rejection Frequencies $\hat{S}_{\hat{\gamma}}$ Test for Heteroskedastic model

| $n$ | $K$ | $d$ | 1 | 2 | 3 | 4 | 5 | 6 | 7 | 8 | 9 | 10 |
| :--- | :--- | :--- | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: |
| 200 | 2 | 2 | 0.061 | 0.061 | 0.065 | 0.072 | 0.054 | 0.053 | 0.054 | 0.040 | 0.056 | 0.045 |
| 200 | 2 | 3 | 0.063 | 0.069 | 0.070 | 0.085 | 0.067 | 0.061 | 0.058 | 0.047 | 0.062 | 0.051 |
| 200 | 3 | 2 | 0.074 | 0.088 | 0.092 | 0.127 | 0.076 | 0.071 | 0.081 | 0.047 | 0.081 | 0.056 |
| 200 | 3 | 3 | 0.079 | 0.093 | 0.103 | 0.145 | 0.080 | 0.078 | 0.082 | 0.044 | 0.081 | 0.065 |
| 200 | 5 | 2 | 0.126 | 0.167 | 0.197 | 0.279 | 0.132 | 0.097 | 0.068 | 0.056 | 0.057 | 0.080 |
| 200 | 5 | 3 | 0.151 | 0.180 | 0.209 | 0.307 | 0.151 | 0.107 | 0.065 | 0.062 | 0.059 | 0.080 |
| 500 | 2 | 2 | 0.050 | 0.060 | 0.057 | 0.075 | 0.058 | 0.054 | 0.035 | 0.045 | 0.061 | 0.051 |
| 500 | 2 | 3 | 0.054 | 0.060 | 0.062 | 0.079 | 0.063 | 0.055 | 0.040 | 0.048 | 0.052 | 0.050 |
| 500 | 3 | 2 | 0.061 | 0.074 | 0.079 | 0.110 | 0.060 | 0.063 | 0.044 | 0.046 | 0.078 | 0.051 |
| 500 | 3 | 3 | 0.070 | 0.079 | 0.084 | 0.115 | 0.064 | 0.058 | 0.052 | 0.048 | 0.074 | 0.050 |
| 500 | 5 | 2 | 0.084 | 0.113 | 0.139 | 0.201 | 0.091 | 0.075 | 0.050 | 0.060 | 0.097 | 0.069 |
| 500 | 5 | 3 | 0.094 | 0.132 | 0.158 | 0.229 | 0.095 | 0.090 | 0.047 | 0.053 | 0.091 | 0.061 |
| 1000 | 2 | 2 | 0.059 | 0.060 | 0.057 | 0.066 | 0.053 | 0.050 | 0.026 | 0.040 | 0.057 | 0.045 |
| 1000 | 2 | 3 | 0.055 | 0.055 | 0.062 | 0.072 | 0.049 | 0.053 | 0.027 | 0.046 | 0.054 | 0.053 |
| 1000 | 3 | 2 | 0.056 | 0.062 | 0.069 | 0.087 | 0.056 | 0.056 | 0.030 | 0.047 | 0.072 | 0.050 |
| 1000 | 3 | 3 | 0.053 | 0.067 | 0.076 | 0.102 | 0.054 | 0.055 | 0.035 | 0.045 | 0.065 | 0.057 |
| 1000 | 5 | 2 | 0.071 | 0.092 | 0.101 | 0.150 | 0.074 | 0.051 | 0.048 | 0.042 | 0.051 | 0.051 |
| 1000 | 5 | 3 | 0.072 | 0.092 | 0.100 | 0.145 | 0.071 | 0.052 | 0.049 | 0.046 | 0.052 | 0.050 |

Notes: The table shows the empirical rejection frequencies for the $S_{\hat{\gamma}}$ test based on $S=5,000$ Monte Carlo replications for the heteroskedastic model $Y_{i}=A\left(\alpha, \sigma, X_{i}\right)^{-1} \epsilon_{i}$. The test has nominal level $a=0.05$. The columns denote the sample size $n$, the dimension of the model $K$, the number explanatory variables $d$ and the choice for densities $\epsilon_{i k}$, for $k \geq 2$, where the numbers correspond to the different densities shown in Figure 3.


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[^1]:    ${ }^{\text {S1 Part }} 2$ of Assumption S2 serves to simplify the form of the effective score function derived in Lemma S3 and is not necessary to set up the model.

[^2]:    ${ }^{\text {S2 }}$ That is, for each $h_{0} \in H_{0}$ define $\tilde{h}_{0}: \mathcal{W} \rightarrow \mathbb{R}$ acccording to $\tilde{h}_{0}(W):=h_{0}(\tilde{X})$ and let $\tilde{H}_{0}$ collect the $\tilde{h}_{0}$ functions so formed. Similarly, for each $h_{k} \in H_{k}(k=1, \ldots, K)$, define $\widetilde{h}_{k}: \mathcal{W} \rightarrow \mathbb{R}$ according to $\tilde{h}_{k}(W):=h_{k}\left(e_{k}^{\prime} A(\alpha, \sigma, X) V_{\theta}\right)$ and let let $\tilde{H}_{k}$ collect the $\tilde{h}_{k}$ functions so formed.

[^3]:    ${ }^{\text {S3 }}$ Compare the forms of the effective scores given in Lemmas 3 and S3.

[^4]:    ${ }^{\text {S4 }}$ Lemmas S4 and S5 are not necessary here since the high level Assumptions S3 and S4 are directly assumed.

[^5]:    ${ }^{\text {S5 }}$ If $\theta+t \varphi(g, h) \in \Theta$ for all $t \in[0,1], \mathcal{U}$ may be taken to include $[0,1]$.
    ${ }^{\text {S6 }}$ Each $\tilde{h}_{k}$ is as defined in the statement of Lemma S1.

[^6]:    ${ }^{\text {S7 }}$ That this convergence holds for any convergent sequence in a compact subset $K$ is equivalent to equicontinuity on $K$, given the continuity of $(g, h) \mapsto P_{\theta_{n}(g, h)}^{n}$ already noted.
    ${ }^{\mathrm{S} 8} A(\theta, X):=A(\alpha, \sigma, X)$.

[^7]:    ${ }^{\text {S9 }}$ The present Lemma is used in the proof of Lemma S2, but is used only to handle the case where $\left(g_{n}, h_{n}\right)$ are not constant in $n$, which is the relevant case here.

[^8]:    

[^9]:    ${ }^{\mathrm{S} 12} l_{n}$ may be defined arbitrarily when $p_{n}=0$.

[^10]:    S14 The superscripts on the $\lambda \mathrm{s}$ are indices, not exponents.
    ${ }^{\text {S15 }}$ See e.g Chapter 8.8 of Magnus and Neudecker (2019).
    ${ }^{\text {S16 }}$ E.g. Theorem 8.7 of Magnus and Neudecker (2019).

[^11]:    ${ }^{\text {S17 }}$ For each $k=1, \ldots, K$ the sequences $\left(\Xi_{k, n}^{L}\right)_{n \in \mathbb{N}},\left(\Xi_{k, n}^{U}\right)_{n \in \mathbb{N}},\left(\mathrm{~B}_{k, n}\right)_{n \in \mathbb{N}}$ and $\left(\delta_{k, n}\right)_{n \in \mathbb{N}}$ are deterministic.
    ${ }^{\text {S18 }}$ Implicitly we choose $K_{k, n}$ and the endpoints and $\delta_{k, n}$ adjusts such that these formulae hold; this way we do not need to adjust anything to ensure these are integers.
    ${ }^{\text {S } 19}$ This differs from the expression in Chen and Bickel (2006) by a factor of -1 as they estimate $-\phi_{k}$.

[^12]:    ${ }^{\text {S20 }}$ Here we implicitly assume (without loss of generality) that all the $\epsilon_{i}$ and $Z_{n, i}$ are defined on a common probability space $(\Omega, \mathcal{F}, \mathrm{P})$.

[^13]:    ${ }^{\text {S21 }}$ This is evident from their definition. See also property (36) (p. 96) of de Boor (2001).

[^14]:    ${ }^{\text {S22 }}$ See e.g. Theorem 5.6.9 in Horn and Johnson (2013).
    ${ }^{\mathrm{S} 23} b_{k, n, m}(z)=0$ outside $\left[\xi_{k, n, m}, \xi_{k, n, m+4}\right)$. See (20) on p. 91 in de Boor (2001).

[^15]:    ${ }^{\text {S24 }}$ As can be easily verified, unlike in the case of linear $(\kappa=2)$ or quadratic splines $(\kappa=3)$, this matrix is not diagonally dominant. In the case of $\kappa \in\{2,3\}$ this argument could be completed in a simpler fashion by using the Gershgorin circle theorem.
    ${ }^{\text {S25 }}$ See section 6.1 in Garoni and Serra-Capizzano (2017), noting that it is clear that $f \in L_{1}([-\pi, \pi])$.

[^16]:    ${ }^{\text {S26 }}$ Recall that the specification corresponds to the baseline model $Y_{i}=A^{-1} \epsilon_{i}$, with $A$ a rotation matrix parametrized by the Cayley transform. The first shock is always drawn from a Gaussian distribution whereas the remaining $k=2, \ldots, K$ are from different distributions whose densities are shown in Figure 3.

