# Online supplementary materials for "Narrowest-Over-Threshold Detection of Multiple Change-points and Change-point-like Features" 

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This document contains the following parts:
A. Simulation models
B. More details on the contrast functions and their construction
C. More details on the computational aspects of NOT and its solution path
D. Further extension of NOT: noise with slow-varying variance
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H. Additional real data example: oil price
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## A. Simulation models

(M1) teeth: piecewise-constant $f_{t}$ (in Scenario (S1)), $T=512, q=7$ change-points at $\tau=64,128, \ldots, 448$, with the corresponding jump sizes $-2,2,-2, \ldots,-2$, starting intercept $f_{1}=1, \sigma_{t}=1$ for $t=1, \ldots, T$.
(M2) blocks: piecewise-constant $f_{t}$ (in Scenario (S1)), $T=2024, q=11$ change-points at $\tau=205,267,308,472,512,820,902,1332,1557,1598,1659$, with the corresponding jump sizes $1.464,-1.830,1.098,-1.464,1.830,-1.537,0.768,1.574,-1.135,0.769,-1.537$, starting intercept $f_{1}=0, \sigma_{t}=1$ for $t=1, \ldots, T$. This signal is widely analysed in the literature, see e.g. Fryzlewicz (2014).
(M3) wave1: piecewise-linear $f_{t}$ without jumps in the intercept (in Scenario (S2)), $T=$ 1408, $q=7$ change-points at $\tau=256,512,768,1024,1152,1280,1344$, with the corresponding changes in slopes $1 \cdot 2^{-6},-2 \cdot 2^{-6}, 3 \cdot 2^{-6} \ldots,-7 \cdot 2^{-6}$, starting intercept $f_{1}=1$ and slope $f_{2}-f_{1}=2^{-8}, \sigma_{t}=1$ for $t=1, \ldots, T$.
(M4) wave2: piecewise-linear $f_{t}$ without jumps in the intercept (in Scenario (S2)), $T=$ 1500, $q=9$ change-points at $\tau=150,300, \ldots, 1350$, with the corresponding changes in slopes $2^{-5},-2^{-5}, 2^{-5}, \ldots,-2^{-5}$, starting intercept $f_{1}=2^{-1}$ and slope $f_{2}-f_{1}=$ $2^{-6}, \sigma_{t}=1$ for $t=1, \ldots, T$.
(M5) mix: piecewise-linear $f_{t}$ with possible jumps at change-points (in Scenario (S3)), length $T=2048, q=7$ change-points at $\tau=256,512, \ldots, 1792$, with the corresponding sizes of jump $0,-1,0,0,2,-1,0$ and changes in the slope $2^{-6},-2^{-6},-2^{-6}$, $2^{-6}, 0,2^{-6},-2^{-5}$, starting value for the intercept $f_{1}=0$ and slope $f_{2}-f_{1}=0, \sigma_{t}=1$ for $t=1, \ldots, T$.
(M6) vol: piecewise-constant $f_{t}$ and $\sigma_{t}$ (in Scenario (S4)), $T=2048, q=7$ changepoints at $\tau=256,512, \ldots, 1792$ with the corresponding jumps in $f_{t}$ and $\sigma_{t}$ being $1,0,-2,0,2,-1,0$ and $0,1,0,1,0,-1,1$, respectively, initial values $f_{1}=\sigma_{1}=1$.
(M7) quad: piecewise-quadratic $f_{t}$ (in Scenario (S5)), $T=1000, q=3$ change-points at $\tau=100,250,500$, with the corresponding changes in the intercept $2,-2,0$, in the slope $0,-10^{-1}, 10^{-1}$ and in the quadratic coefficient $0,0,2 \times 10^{-5}$, the initial values $f_{1}=f_{2}-f_{1}=f_{3}-2 f_{2}+f_{1}=0, \sigma_{t}=1$ for all $t=1, \ldots, T$.
(M8) smile: piecewise-linear $f_{t}$ with possible jumps at change-points (designed to test NOT under misspecification), $T=2048, q=6$ change-points at $\tau=256,512,768$, $1280,1536,1792$, with the corresponding sizes of jump $0,-4,0,0,4,0$ and changes in the slope $-2^{-5}, 0,2^{-6}, 2^{-6}, 0,-2^{-5}$, starting value for the intercept $f_{1}=0$ and slope $f_{2}-f_{1}=2^{-6}, \sigma_{t}=1$ for $t=1, \ldots, T$.

## B. More details on the contrast functions and their construction

## B.1. Scenario (S1)

Here $f_{t}$ is piecewise-constant. For any integer triple $(s, e, b)$ with $0 \leq s<e \leq T$ and $s<b<e$, recalling that we have defined the contrast vector $\boldsymbol{\psi}_{s, e}^{b}=\left(\psi_{s, e}^{b}(1), \ldots, \psi_{s, e}^{b}(T)\right)^{\prime}$ as

$$
\psi_{(s, e]}^{b}(t)= \begin{cases}\sqrt{\frac{e-b}{(e-s)(b-s)}}, & t=s+1, \ldots, b \\ -\sqrt{\frac{b-s}{(e-s)(e-b)}}, & t=b+1, \ldots, e \\ 0, & \text { otherwise }\end{cases}
$$

Also, if $b \notin\{s+1, \ldots, e-1\}$, then we set $\psi_{(s, e]}^{b}(t)=0$ for all $t$.
For any vector $\mathbf{v}=\left(v_{1}, \ldots, v_{T}\right)^{\prime}$, we define the contrast function as $\mathcal{C}_{(s, e]}^{b}(\mathbf{v})=\left|\left\langle\mathbf{v}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle\right|$. Therefore, if $s<b<e$, then

$$
\mathcal{C}_{(s, e]}^{b}(\mathbf{v})=\left|\sqrt{\frac{e-b}{(e-s)(b-s)}} \sum_{t=s+1}^{b} v_{t}-\sqrt{\frac{b-s}{(e-s)(e-b)}} \sum_{t=b+1}^{e} v_{t}\right| .
$$

Otherwise, $\mathcal{C}_{(s, e]}^{b}(\mathbf{v})=0$. This recovers the well-known CUSUM statistic in the changepoint detection literature. It can be shown that $\left[\mathcal{C}_{(s, e]}^{b}(\mathbf{Y})\right]^{2}=\sigma_{0}^{2} \mathcal{R}_{(s, e]}^{b}(\mathbf{Y})$ for every $(s, e, b)$ with $0 \leq s<b<e \leq T$, thus $\mathcal{C}_{(s, e]}^{b}(\cdot)$ fulfills the requirements for the contrast function listed in Section 2.3.

In addition, for any $0 \leq s<e \leq T$, we define the constant vector for the interval $(s, e]$ as

$$
\mathbf{1}_{(s, e]}(t)= \begin{cases}(e-s)^{-1 / 2}, & t=s+1, \ldots, e \\ 0, & \text { otherwise }\end{cases}
$$

and write $\mathbf{1}_{(s, e]}=\left(\mathbf{1}_{(s, e]}(1), \ldots, \mathbf{1}_{(s, e]}(T)\right)^{\prime}$. Then it is easy to check that $\mathbf{1}_{(s, e]}$ and $\boldsymbol{\psi}_{(s, e]}^{b}$ are orthonormal. This explains why the CUSUM is invariant to shifts in the mean.

## B.2. Scenario (S2)

Here $f_{t}$ is piecewise-linear and continuous. For any triple $(s, e, b)$ with $0 \leq s<e \leq T$ and $s+1<b<e$, consider the contrast vector $\phi_{(s, e]}^{b}=\left(\phi_{(s, e]}^{b}(1), \ldots, \phi_{(s, e]}^{b}(T)\right)^{\prime}$ with
$\phi_{(s, e]}^{b}(t)= \begin{cases}\alpha_{(s, e]}^{b} \beta_{(s, e]}^{b}[\{3(b-s)+(e-b)-1\} t-\{b(e-s-1)+2(s+1)(b-s)\}], & t=s+1, \ldots, b \\ -\frac{\alpha_{s(s, e]}}{\beta_{(s, e]}}[\{3(e-b)+(b-s)+1\} t-\{b(e-s-1)+2 e(e-b+1)\}], & t=b+1, \ldots, e, \\ 0, & \text { otherwise. }\end{cases}$
where $\alpha_{s, e}^{b}=\left(\frac{6}{l\left(l^{2}-1\right)(1+(e-b+1)(b-s)+(e-b)(b-s-1))}\right)^{1 / 2}, \beta_{s, e}^{b}=\left(\frac{(e-b+1)(e-b)}{(b-s-1)(b-s)}\right)^{1 / 2}$ and $l=$ $e-s$. If $b \notin\{s+2, \ldots, e-1\}$, then we set $\phi_{(s, e]}^{b}(t)=0$ for all $t$. The contrast function is then defined as

$$
\mathcal{C}_{(s, e]}^{b}(\mathbf{v})=\left|\left\langle\mathbf{v}, \phi_{(s, e]}^{b}\right\rangle\right| .
$$

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To explain the rationale behind $\phi_{(s, e]}^{b}$, we first define the "linear" vector for the interval $(s, e]$ with $e-s>1, \gamma_{(s, e]}=\left(\gamma_{(s, e]}(1), \ldots, \gamma_{(s, e]}(T)\right)^{\prime}$, as

$$
\gamma_{(s, e]}(t)= \begin{cases}\left\{\frac{1}{12}(e-s-1)(e-s)(e-s+1)\right\}^{-1 / 2}\left(t-\frac{e+s+1}{2}\right), & t=s+1, \ldots, e ; \\ 0, & \text { otherwise }\end{cases}
$$

Then we have that $\phi_{(s, e]}^{b}$ is orthonormal to both $\mathbf{1}_{(s, e]}$ and $\gamma_{(s, e]}$ (note that $\gamma_{(s, e]}$ itself is orthonormal to $\mathbf{1}_{(s, e]}$. The orthonormality of the vectors $\mathbf{1}_{(s, e]}, \boldsymbol{\gamma}_{(s, e]}$ and $\boldsymbol{\phi}_{(s, e]}^{b}$ is important in deriving the identity $\sigma_{0}^{2} \mathcal{R}_{(s, e]}^{b}(\mathbf{Y})=\mathcal{C}_{(s, e]}^{b}(\mathbf{Y})^{2}$ below, and helps improve the numerical efficiency and stability in our implementation of NOT. In particular, it means that the contrast function is invariant to both mean shifts and slope shifts on a given interval. In fact, $\phi_{(s, e]}^{b}$ can be derived by (i) applying the Gram-Schmidt process on the following vector (linear with a kink at $b$ on ( $s, e]$ )

$$
\tilde{\phi}_{(s, e]}^{b}(t)= \begin{cases}t-b, & t=b+1, \ldots, e \\ 0, & \text { otherwise }\end{cases}
$$

with respect to $\mathbf{1}_{(s, e]}$ and $\boldsymbol{\gamma}_{(s, e]}$, and (ii) normalisation such that $\|\cdot\|_{2}=1$. Now write the restriction of $\mathbf{v}$ on the interval $(s, e]$ as $\left.\mathbf{v}\right|_{(s, e]}=\left(0, \ldots, 0, v_{s+1}, \ldots, v_{e}, 0, \ldots, 0\right)^{\prime}$. Fix any $(s, e, b)$, given the restriction imposed on $\boldsymbol{\Theta}$ in (S2), the best approximation of $\left.\mathbf{Y}\right|_{(s, e]}$ (in the $\ell_{2}$ distance) with a single kink at $b$ is a linear combination of $\mathbf{1}_{(s, e]}, \boldsymbol{\gamma}_{(s, e]}$ and $\phi_{(s, e]}^{b}$ (all mutually orthonormal). Therefore,

$$
\begin{aligned}
& \sigma_{0}^{2} \mathcal{R}_{(s, e]}^{b}(\mathbf{Y}) \\
& =\min _{a_{0}, a_{1} \in \mathbb{R}}\left\|\left.\mathbf{Y}\right|_{(s, e]}-a_{0} \mathbf{1}_{(s, e]}-a_{1} \boldsymbol{\gamma}_{(s, e]}\right\|_{2}^{2}-\min _{a_{0}, a_{1}, a_{2} \in \mathbb{R}}\left\|\left.\mathbf{Y}\right|_{(s, e]}-a_{0} \mathbf{1}_{(s, e]}-a_{1} \boldsymbol{\gamma}_{(s, e]}-a_{2} \boldsymbol{\phi}_{(s, e]}^{b}\right\|_{2}^{2} \\
& =\left\||\mathbf{Y}|_{(s, e]}-\left\langle\mathbf{Y}, \gamma_{(s, e]}\right\rangle \boldsymbol{\gamma}_{(s, e]}-\left\langle\mathbf{Y}, \mathbf{1}_{(s, e]}\right\rangle \mathbf{1}_{(s, e]}\right\|^{2} \\
& \quad-\left\|\left.\mathbf{Y}\right|_{(s, e]}-\left\langle\mathbf{Y}, \boldsymbol{\phi}_{(s, e]}^{b}\right\rangle \boldsymbol{\phi}_{(s, e]}^{b}-\left\langle\mathbf{Y}, \boldsymbol{\gamma}_{(s, e]}\right\rangle \boldsymbol{\gamma}_{(s, e]}-\left\langle\mathbf{Y}, \mathbf{1}_{(s, e]}\right\rangle \mathbf{1}_{(s, e]}\right\|^{2} \\
& =\left\langle\mathbf{Y}, \boldsymbol{\phi}_{(s, e]}^{b}\right\rangle^{2}=\mathcal{C}_{(s, e]}^{b}(\mathbf{Y})^{2} .
\end{aligned}
$$

Thus the aforementioned requirements for the contrast function are satisfied.

## B.3. Scenario (S3)

Here $f_{t}$ is a piecewise-linear but not necessarily continuous function. We use the following contrast function for any $s+1<b<e-1$ :

$$
\begin{equation*}
\mathcal{C}_{(s, e]}^{b}(\mathbf{v})=\left(\left\langle\mathbf{v}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle^{2}+\left\langle\mathbf{v}, \boldsymbol{\gamma}_{(s, b]}\right\rangle^{2}+\left\langle\mathbf{v}, \boldsymbol{\gamma}_{(b, e]}\right\rangle^{2}-\left\langle\mathbf{v}, \boldsymbol{\gamma}_{(s, e]}\right\rangle^{2}\right)^{1 / 2} . \tag{13}
\end{equation*}
$$

Otherwise, for $b \notin\{s+2, \ldots, e-2\}$, we set $\mathcal{C}_{(s, e]}^{b}(\mathbf{v})=0$.

This construction is justified by noting that

$$
\begin{aligned}
\sigma_{0}^{2} \mathcal{R}_{(s, e]}^{b}(\mathbf{Y})= & \min _{a_{0}, a_{1} \in \mathbb{R}}\left\|\left.\mathbf{Y}\right|_{(s, e]}-a_{0} \mathbf{1}_{(s, e]}-a_{1} \boldsymbol{\gamma}_{(s, e]}\right\|_{2}^{2} \\
& -\left(\min _{a_{0}, a_{1} \in \mathbb{R}}\left\|\left.\mathbf{Y}\right|_{(s, b]}-a_{0} \mathbf{1}_{(s, b]}-a_{1} \boldsymbol{\gamma}_{(s, b]}\right\|_{2}^{2}+\min _{a_{0}, a_{1} \in \mathbb{R}}\left\|\left.\mathbf{Y}\right|_{(b, e]}-a_{0} \mathbf{1}_{(b, e]}-a_{1} \boldsymbol{\gamma}_{(b, e]}\right\|_{2}^{2}\right) \\
= & \min _{a_{0}, a_{1} \in \mathbb{R}}\left\|\left.\mathbf{Y}\right|_{(s, e]}-a_{0} \mathbf{1}_{(s, e]}-a_{1} \boldsymbol{\gamma}_{(s, e]}\right\|_{2}^{2} \\
& -\min _{a_{0}, a_{1}, a_{2}, a_{3} \in \mathbb{R}}\left\|\left.\mathbf{Y}\right|_{(s, e]}-a_{0} \mathbf{1}_{(s, e]}-a_{1} \boldsymbol{\gamma}_{(s, b]}-a_{2} \boldsymbol{\gamma}_{(b, e]}-a_{3} \boldsymbol{\psi}_{(s, e]}^{b}\right\|_{2}^{2} \\
= & \mathcal{C}_{(s, e]}^{b}(\mathbf{Y})^{2},
\end{aligned}
$$

where we also used the orthonormality among $\mathbf{1}_{(s, e]}, \boldsymbol{\psi}_{(s, e]}^{b}, \boldsymbol{\gamma}_{(s, b]}$ and $\boldsymbol{\gamma}_{(b, e]}$ in the above derivation.

## B.4. Scenario (S4)

Here both $f_{t}$ and $\sigma_{t}$ are piecewise-constant. For any $s+1<b<e-1$, we propose

$$
\begin{equation*}
\mathcal{C}_{(s, e]}^{b}(\mathbf{v})=(e-s) \log \left(\hat{\sigma}_{(s, e]}(\mathbf{v})\right)-(b-s) \log \left(\hat{\sigma}_{(s, b]}(\mathbf{v})\right)-(e-b) \log \left(\hat{\sigma}_{(b, e]}(\mathbf{v})\right), \tag{14}
\end{equation*}
$$

where

$$
\hat{\sigma}_{(s, e]}^{2}(\mathbf{v})=\frac{1}{e-s} \sum_{t=s+1}^{e}\left(v_{t}-\frac{1}{e-s} \sum_{t=s+1}^{e} v_{t}\right)^{2}=\left\langle\mathbf{v}^{2}, \mathbf{1}_{(s, e]}^{2}\right\rangle-\left\langle\mathbf{v}, \mathbf{1}_{(s, e]}^{2}\right\rangle^{2} .
$$

Otherwise, for $b \notin\{s+2, \ldots, e-2\}$, we set $\mathcal{C}_{(s, e]}^{b}(\mathbf{v})=0$. In this Scenario, it is straightforward to verify that $\mathcal{C}_{(s, e]}^{b}(\mathbf{Y})=\mathcal{R}_{(s, e]}^{b}(\mathbf{Y})$. (N.B. $\mathbf{1}_{(s, e]}^{2} \neq \mathbf{1}_{(s, e]}$ due to the normalising constant.) In practice, for numerical stability, we use $\log _{\epsilon}(\cdot):=\log \{\max (\cdot, \epsilon)\}$ instead of $\log (\cdot)$ in (14) with a pre-given small $\epsilon>0$.

## C. More details on the compututational aspects of NOT and its solution path

## C.1. Computing contrast functions in a linear time

The practical performance (in terms of computational cost) of Algorithm 1 relies on the fast computation of the contrast functions discussed in Section 2.3 on any given interval $(s, e]$. Here we show that in all scenarios listed in Section 2.3, the cost of computing $\left\{\mathcal{C}_{(s, e]}^{b}(\mathbf{Y})\right\}_{b=s+1}^{e-1}$ is $O(e-s)$.

Note that the key ingredients in $\mathcal{C}_{(s, e]}^{b}(\mathbf{Y})$ under the different scenarios are functions of the inner products, i.e. $\left\langle\mathbf{Y}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle,\left\langle\mathbf{Y}, \boldsymbol{\phi}_{(s, e]}^{b}\right\rangle,\left\langle\mathbf{Y}, \boldsymbol{\gamma}_{(s, b]}\right\rangle,\left\langle\mathbf{Y}, \boldsymbol{\gamma}_{(b, e]}\right\rangle,\left\langle\mathbf{Y}, \mathbf{1}_{(s, b]}^{2}\right\rangle$, $\left\langle\mathbf{Y}, \mathbf{1}_{(b, e]}^{2}\right\rangle,\left\langle\mathbf{Y}^{2}, \mathbf{1}_{(s, b]}^{2}\right\rangle$ and $\left.\left\langle\mathbf{Y}^{2}, \mathbf{1}_{(b, e]}^{2}\right\rangle\right)$ for $b=s+1, \ldots, e-1$. For a fixed interval $(s, e]$,
by simple algebra, we observe that $\left\langle\mathbf{Y}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle$ and $\left\langle\mathbf{Y}, \phi_{(s, e]}^{b}\right\rangle$ can be decomposed as

$$
\begin{aligned}
\left\langle\mathbf{Y}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle & =\overleftarrow{a}_{\boldsymbol{\psi}, b} \sum_{t=s+1}^{b} Y_{t}-\vec{a}_{\boldsymbol{\psi}, b} \sum_{t=b+1}^{e} Y_{t} \\
& :=\overleftarrow{a}_{\boldsymbol{\psi}, b} \overleftarrow{\pi}_{b}^{(0)}(\mathbf{Y})-\vec{a}_{\boldsymbol{\psi}, b} \vec{\pi}_{b}^{(0)}(\mathbf{Y}), \\
\left\langle\mathbf{Y}, \boldsymbol{\phi}_{(s, e]}^{b}\right\rangle & =\overleftarrow{a}_{\phi, b}^{(1)} \sum_{t=s+1}^{b} t Y_{t}-\vec{a}_{\phi, b}^{(1)} \sum_{t=b+1}^{e} t Y_{t}+\overleftarrow{a}_{\phi, b}^{(0)} \sum_{t=s+1}^{b} Y_{t}-\vec{a}_{\phi, b}^{(0)} \sum_{t=b+1}^{e} Y_{t} \\
& :=\overleftarrow{a}_{\phi, b}^{(1)} \overleftarrow{\pi}_{b}^{(1)}(\mathbf{Y})-\vec{a}_{\phi, b}^{(1)} \vec{\pi}_{b}^{(1)}(\mathbf{Y})+\overleftarrow{a}_{\phi, b}^{(0)} \overleftarrow{\pi}_{b}^{(0)}(\mathbf{Y})-\vec{a}_{\phi, b}^{(0)} \vec{\pi}_{b}^{(0)}(\mathbf{Y}),
\end{aligned}
$$

where $\overleftarrow{a}_{\boldsymbol{\psi}, b}, \vec{a}_{\boldsymbol{\psi}, b}, \overleftarrow{a}_{\phi, b}^{(1)}, \vec{a}_{\phi, b}^{(1)}, \overleftarrow{a}_{\phi, b}^{(0)}$ and $\vec{a}_{\phi, b}^{(0)}$ are scalars that do not depend on $\mathbf{Y}$, and can all be computed at the cost of $O(1)$ using equations given in Section 2.3. Here for notational convenience, we use overhead arrows to indicate whether a scalar or a function is associated with observations with indices $\leq b$ (i.e. over ( $s, b]$, using $\leftarrow$ ) or with indices $>b$ (i.e. over $(b, e]$, using $\vec{\cdot}$ ). We also suppress their dependence on $s$ and $e$ in the notation. In addition, the following recursive formulae hold

$$
\begin{aligned}
\overleftarrow{\pi}_{b+1}^{(k)}(\mathbf{Y}) & =\overleftarrow{\pi}_{b}^{(k)}(\mathbf{Y})+(b+1)^{k} Y_{b+1} \\
\vec{\pi}_{b}^{(k)}(\mathbf{Y}) & =\vec{\pi}_{b+1}^{(k)}(\mathbf{Y})+(b+1)^{k} Y_{b+1}
\end{aligned}
$$

with $\overleftarrow{\pi}_{s}^{(k)}(\mathbf{Y})=\vec{\pi}_{e}^{(k)}(\mathbf{Y})=0$ for $k=0,1$. Consequently, $\overleftarrow{\pi}_{b}^{(k)}(\mathbf{Y})$ and $\vec{\pi}_{b}^{(k)}(\mathbf{Y})$ for all $b \in\{s+1, \ldots, e-1\}$ and $k=0,1$ (thereby $\left\langle\mathbf{Y}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle$ and $\left.\left\langle\mathbf{Y}, \boldsymbol{\phi}_{(s, e]}^{b}\right\rangle\right)$ can be computed in a single pass through $Y_{s+1}, \ldots, Y_{e}$. Similar approach can be applied to the remaining inner products involved in the definitions of the contrast functions given in Section 2.3, which demonstrates that in all these cases the computation of $\left\{\mathcal{C}_{(s, e]}^{b}(\mathbf{Y})\right\}_{b=s+1}^{e-1}$ scales linearly with the length of the sub-interval.

## C.2. Details of the NOT solution path algorithm

As mentioned in Section 3.2 of the main paper, we have developed Algorithm 2 that computes the entire threshold-indexed solution path $\left\{\mathcal{T}\left(\zeta_{T}\right)\right\}_{\zeta_{T} \geq 0}$ quickly, and have implemented it in our R package not. Detailed pseudo-code is provided on the next page.

The construction of Algorithm 2 stems from two observations. First, for any fixed threshold $\zeta_{T}$, Algorithm 1 implies a binary tree data structure that is constructed according to the order of the detection of each change-point. More specifically, in our implementation, each tree node N contains the following information.
(a) The current interval of interest is (N.s, N.e].
(b) From all elements in $F_{T}^{M}$ that are also subsets of (N.s, N.e], we find the narrowest-overthreshold sub-interval. Within that sub-interval, let N.c be the maximum achieved value of the contrast function over all possible locations of the feature, and N.b be the corresponding location (i.e. the detected change-point location over (N.s, N.e]).
(c) N.Left and N.Right point to the nodes of the next detected change-points in (N.s, N.b] and (N.b, N.e], respectively.

```
Algorithm 2 NOT solution path
Input: Data vector \(\mathbf{Y}\), all sub-intervals \(\left(s_{m}, e_{m}\right] \in F_{T}^{M}\) together with
    \(b_{m}:=\operatorname{argmax}_{s_{m}<b \leq e_{m}} \mathcal{C}_{\left(s_{m}, e_{m}\right]}^{b}(\mathbf{Y}), \quad c_{m}:=\mathcal{C}_{\left(s_{m}, e_{m}\right]}^{b_{m}}(\mathbf{Y}) \quad\) and \(\quad l_{m}:=e_{m}-s_{m}\).
Output: Thresholds \(0=\zeta_{T}^{(1)}<\ldots<\zeta_{T}^{(N)}\) and sets of estimated change-points
    \(\mathcal{T}\left(\zeta_{T}^{(1)}\right), \ldots, \mathcal{T}\left(\zeta_{T}^{(N)}\right)\).
To start the algorithm: Call SolutionPath()
procedure BuildBinaryTree \(\left((s, e], \zeta_{T}, \mathrm{~N}\right)\)
        \(\mathcal{M}_{(s, e]}:=\) set of those \(m \in\{1, \ldots, M\}\) such that \(\left(s_{m}, e_{m}\right] \subset(s, e]\)
        \(\mathcal{O}_{(s, e]}:=\) set of \(m \in \mathcal{M}_{(s, e]}\) such that \(c_{m}>\zeta_{T}\)
        if \(\mathcal{O}_{(s, e]}=\emptyset\) then \(\mathrm{N}=\mathrm{NULL}\)
        else
            \(k:=\) any element of \(\operatorname{argmin}_{m \in \mathcal{O}_{(s, e]}} l_{m}\)
            N.b \(:=b_{k}\), N.c \(:=c_{k}\), N.Left \(:=\) NULL, N.Right \(:=\) NULL
            BuildBinaryTree ( \(\left(s\right.\), N.b], \(\zeta_{T}\), N.Left)
            BuildBinaryTree((N.b, e], \(\zeta_{T}\), N.Right)
        end if
end procedure
procedure \(\operatorname{UpdateBinaryTree}\left((s, e], \zeta_{T}, \mathrm{~N}\right)\)
        if N.c \(\leq \zeta_{T}\) then
            BuildBinaryTree \(\left((s, e], \zeta_{T}, \mathrm{~N}\right)\)
        else
            if \(N . L e f t \neq\) NULL then
                UpdateBinary Tree ( \(s\), N.b], \(\zeta_{T}\), N.Left)
            end if
            if N.Right \(\neq\) NULL then
                    UpdateBinaryTree((n.b, e], \(\zeta_{T}\), N.Right)
            end if
        end if
    end procedure
    procedure SolutionPath()
        Set \(\mathrm{N}_{\mathrm{r}}:=\) NULL, \(i:=1, \zeta_{T}^{(1)}:=0\)
        BuildBinaryTree ( \(\left.(0, T], \zeta_{T}^{(1)}, N_{\mathrm{r}}\right)\)
        while \(N_{r} \neq\) NULL do
            \(\mathcal{D}:=\left\{\mathrm{N}_{\mathrm{r}}\right.\) and all its children nodes \(\}\)
            \(\mathcal{T}\left(\zeta_{T}^{(i)}\right):=\{\mathrm{N} . \mathrm{b} \mid \mathrm{N} \in \mathcal{D}\}\)
            \(\zeta_{T}^{(i+1)}:=\min _{\mathrm{N} \in \mathcal{D}}\{\mathrm{N} . \mathrm{c}\}\)
            UpdateBinaryTree \(\left((0, T], \zeta_{T}^{(i+1)}, \mathrm{N}_{\mathrm{r}}\right)\)
            \(i:=i+1\)
        end while
end procedure
```

Table 2. Intervals considered in Figure 7a and corresponding maxima of the contrast function $\mathcal{C}_{(s, e]}^{b}(\cdot)$ given by (8), all calculated for a sample path of $Y_{t}, t=1, \ldots, 1000$ generated from model (1) with the signal $f_{t}$ given by $(2)$ and the noise $\varepsilon_{t} \sim \mathcal{N}\left(0,0.05^{2}\right)$.

| $s$ | $e$ | $e-s$ | $\operatorname{argmax}_{s<b \leq e} \mathcal{C}_{(s, e]}^{b}(\mathbf{Y})$ | $\max _{s<b \leq e} \mathcal{C}_{(s, e]}^{b}(\mathbf{Y})$ |
| :---: | :---: | :---: | :---: | :---: |
| 0 | 1000 | 1000 | 490 | 10.19 |
| 9 | 245 | 236 | 43 | 0.08 |
| 224 | 450 | 226 | 344 | 0.76 |
| 499 | 750 | 251 | 651 | 0.83 |
| 749 | 950 | 211 | 746 | 0.03 |
| 449 | 550 | 101 | 471 | 0.07 |

We then treat the first detected change-point over $(0, T]$ as the root of the tree and construct its branches in a recursive fashion afterwards. Second, suppose that we have already constructed the tree for $\zeta_{T}$ with root $\mathrm{N}_{\mathrm{r}}$. For $\zeta_{T}^{\prime}>\zeta_{T}$, the new tree's root is unchanged if $N_{r} . c>\zeta_{T}^{\prime}$. This observation remains valid for $N_{r}$.Left and $N_{r}$. Right and all subsequent nodes. Therefore, a branch of the tree has to be reconstructed only if N.c $\leq \zeta_{T}^{\prime}$ for some node N . In this way, the tree constructed for $\zeta_{T}$ can be used as a starting point to finding the tree corresponding to $\zeta_{T}^{\prime}$, thus significantly reducing the computational time in comparison to constructing the tree from scratch.

Next, we elaborate on the complexity of Algorithm 2. As explained previously, finding solutions of Algorithm 1 for a single threshold $\zeta_{T}$ is equivalent to the construction of a binary tree, which can be performed with the BuildBinaryTree routine given in Algorithm 2. Computational cost of this operation is no larger than $O\left(M K_{\zeta_{T}}\right)$, where $K_{\zeta_{T}}$ denotes the height of the constructed binary tree with the threshold $\zeta_{T}$. The computational complexity of finding the entire solution path using Algorithm 2 is therefore (in the worst case) $O(M K N)$, where $N$ and $K$ are, respectively, the number of solutions and the maximum tree depth over the entire solution path. However, this is a rough estimate which assumes that for each threshold on the path the binary tree has a different root node, which, from our empirical experience, is highly unlikely to occur in practice. Typically, the consecutive trees on the path differ just slightly (see e.g. our next Section C.3), which significantly reduces the amount of computation that Algorithm 2 requires. As such, we find that the computational complexity of Algorithm 2 is more like $O(M T)$ in practice.

## C.3. An illustrative example

In this part, we revisit the example shown in the Introduction of our paper, and provide a simple illustration of how Algorithm 1 and Algorithm 2 work on a simulated dataset. Figure 7 shows the generated data $\left\{Y_{t}\right\}_{t=1}^{1000}$ following Scenario (S2), where the signal $f_{t}$ is as in (2) and $\sigma_{t}=0.05$. The contrast function (8) is evaluated for 5 intervals. We observe that the contrast function corresponding to $(0,1000]$, being the longest interval here, attains its maximum at $b=490$, which is far from the true change-points located at $\tau=350$ and $\tau=650$. Furthermore, $\max _{b} \mathcal{C}_{(0,1000]}^{b}(\mathbf{Y})$ is much larger than the corresponding value for the other intervals considered in Table 2. However, thanks to the fact that we focus on the narrowest-over-threshold intervals, Algorithm 1 (for any $\zeta_{T} \in(0.08,0.83)$ ) picks at its first iteration an interval with exactly one change-point (depending on $\zeta_{T}$, it is either (224, 450] or $(499,750])$ and the maximum of the contrast function computed is close to one of the true change-points.


Fig. 7. An application of the NOT methodology to $Y_{t}$ generated from model (1) with the signal $f_{t}$ given by (2) and i.i.d. $\varepsilon_{t} \sim \mathcal{N}\left(0,0.05^{2}\right)$. Figure 7a: contrast function $\mathcal{C}_{(s, e]}^{b}(\mathbf{Y})$ given by (8) evaluated for all $b \in(s, e]$ with intervals $(s, e]$ specified in Table 2. For intervals containing one change-point, $\mathcal{C}_{(s, e]}^{b}(\mathbf{Y})$ attains its maximum at $b$ close to the actual change-point. When there are two changepoints (black solid line), the maximum is far from both change-points, despite $\max _{b} \mathcal{C}_{(s, e]}^{b}(\mathbf{Y})$ being large. Figure 7b: observed $Y_{t}$ (thin grey), true signal (thick dashed black), signal estimated picking the change-point candidate based on the interval corresponding to the largest contrast function (dotted-dashed navy) and the narrowest-over-threshold intervals (dashed red).


Fig. 8. First four segmentation trees obtained by Algorithm 2 applied to a realization of $\left(Y_{1}, \ldots, Y_{1000}\right)^{\prime}$ presented in Figure 7. The larger the node, the larger the corresponding value of $\max _{b} \mathcal{C}_{(s, e]}^{b}(\mathbf{Y})$. Here $\mathcal{C}_{(s, e]}^{b}(\cdot)$ is given by (8). The grey nodes correspond to the smallest contrast function for each tree that are updated as Algorithm 2 proceeds.

Figure 8 shows how Algorithm 2 proceeds in the example presented in Figure 7. At the initial stage that can be seen in Figure 8a, the threshold is set to $\zeta_{T}^{(1)}=0$ and $b=471$, the maximum of the contrast function computed for the shortest interval $(449,550]$ is taken as the root of the binary tree. Then we construct its left and right branches by considering only those intervals specified in Table 2 with $(s, e] \subset(0,471]$ and $(s, e] \subset(471,1000]$, respectively, and the procedure continues for the resulting nodes. Next, the node with the smallest value of the contrast function is determined $(b=746)$ and the threshold is set to the corresponding minimum $\zeta_{T}^{(2)}=0.03$. This guarantees that as Algorithm 2 proceeds, there will be at least one update in the binary tree. In our example, the $b=746$ node is removed and, as the maximum for $(499,750] \subset(471,1000]$ exceeds the threshold, the $b=651$ node is inserted its place. Subsequently, we identify the node with the smallest contrast again $(b=471)$, update the threshold to $\zeta_{T}^{(3)}=0.07$ and reconstruct the entire tree, as $b=471$ in Figure 8 b constitutes its root. Algorithm 2 keeps running until the resulting tree shrinks to NULL. In this example, the fourth solution on the path (Figure 8d) contains exactly two nodes being close to the true change-points.

## D. Further extension of NOT: noise with slow-varying variance

In all scenarios considered previously, we assumed a constant or piecewise constant $\sigma_{t}$. Now we discuss how NOT can be extended to handle $\sigma_{t}$ of a more general form. We model $\mathbf{Y}$ through (3) with $\varepsilon_{t} \stackrel{\text { i.i.d. }}{\sim} \mathcal{N}(0,1)$. To fix ideas, here we focus on the case of piecewise constant signal $f_{t}$ (i.e. similar to (S1)), but with a slowly-varying $\sigma_{t}$, i.e. $\sigma_{t}=\sigma(t / T)$ with $\sigma(\cdot)$ being an unknown smooth function from $[0,1] \rightarrow(0, \infty)$.

## D.1. Methodology

We propose to estimate the change-points in three steps:
(a) We estimate $\sigma(\cdot)$ using a standard nonparametric method, such as spline smoothing, on $\left\{\left(t / T, \sqrt{\pi}\left|Y_{t+1}-Y_{t}\right| / 2\right)\right\}_{t=2}^{T}$, which we denote as $\hat{\sigma}(\cdot)$. Also, we write $\hat{\sigma}_{t}=\hat{\sigma}(t / T)$ for $t=1, \ldots, T$.
(b) We perform our NOT solution path algorithm using the contrast function $\hat{\mathcal{C}}_{(s, e]}^{b}(\mathbf{Y})=$ $\left|\left\langle\mathbf{Y}, \hat{\boldsymbol{\psi}}_{(s, e]}^{b}\right\rangle\right|$ for any tuple $(s, b, e)$ with $0 \leq s<b<e \leq T$ in our NOT procedure, with $\hat{\boldsymbol{\psi}}_{(s, e]}^{b}=\left(\hat{\psi}_{(s, e]}^{b}(1), \ldots, \hat{\psi}_{(s, e]}^{b}(T)\right)^{\prime}$ and

$$
\hat{\psi}_{(s, e]}^{b}(t)= \begin{cases}\hat{\sigma}_{t}^{-2} \sqrt{\left(\Omega_{e}^{2}-\Omega_{b}^{2}\right)\left(\Omega_{e}^{2}-\Omega_{s}^{2}\right)^{-1}\left(\Omega_{b}^{2}-\Omega_{s}^{2}\right)^{-1}}, & t=s+1, \ldots, b \\ -\hat{\sigma}_{t}^{-2} \sqrt{\left(\Omega_{b}^{2}-\Omega_{s}^{2}\right)\left(\Omega_{e}^{2}-\Omega_{s}^{2}\right)^{-1}\left(\Omega_{e}^{2}-\Omega_{b}^{2}\right)^{-1}}, & t=b+1, \ldots, e \\ 0, & \text { otherwise }\end{cases}
$$

where $\Omega_{t}^{2}=\sum_{i=1}^{t} 1 / \hat{\sigma}_{i}^{2}$ (and by default $\Omega_{0}^{2}=0$ ). As before, if $b \notin\{s+1, \ldots, e-1\}$, then we set $\hat{\psi}_{(s, e]}^{b}(t)=0$ for all $t$. We remark that this contrast function originates from the generalised log-likelihood ratio, and can be viewed as a weighted and scaled version of (6) based on CUSUM statistic. Its derivation can be found in Section D.2. In addition, when $\hat{\sigma}_{t}$ is constant (say $\hat{\sigma}_{t}=1$ for all $t$ ), we would recover (6).
(c) We pick the best model along the solution path via the sSIC criterion, with the log-likelihood for each segment given by

$$
\log \ell\left(Y_{\hat{\tau}_{j-1}+1}, \ldots, Y_{\hat{\tau}_{j}} ; \hat{\boldsymbol{\Theta}}_{j}\right)=\sum_{t=\hat{\tau}_{j-1}+1}^{\hat{\tau}_{j}}\left\{-\frac{\left(Y_{t}-\hat{Y}_{\left(\hat{\tau}_{j-1}, \hat{\gamma}_{j}\right.}\right)^{2}}{2 \hat{\sigma}_{t}^{2}}-\frac{\log \left(2 \pi \hat{\sigma}_{t}^{2}\right)}{2}\right\},
$$

where $\hat{Y}_{\left(\hat{\tau}_{j-1}, \hat{\tau}_{j}\right]}=\left(\sum_{t=\hat{\tau}_{j-1}+1}^{\hat{\tau}_{j}} \hat{\sigma}_{t}^{-2} Y_{t}\right) /\left(\sum_{t=\hat{\tau}_{j-1}+1}^{\hat{\tau}_{j}} \hat{\sigma}_{t}^{-2}\right)$.
To make this solution complete, a suitable choice of the smoothing parameter would have to be considered in the first step. This is a standard problem in nonparametrics, and several solutions exist, e.g. those based on (leave-one-out) cross-validation. We leave a detailed study of this issue for future research.

## D.2. Detailed derivation of the corresponding contrast function

Here we derive the contrast function from the generalised log-likelihood ratio.

Given an interval $(s, e]$. Suppose that there is a change-point at $b$, then under the normality assumption, the log-likelihood is

$$
-\frac{1}{2} \sum_{t=s+1}^{b} \frac{\left(Y_{t}-\mu_{L}\right)^{2}}{\hat{\sigma}_{t}^{2}}-\frac{1}{2} \sum_{t=b+1}^{e} \frac{\left(Y_{t}-\mu_{R}\right)^{2}}{\hat{\sigma}_{t}^{2}}-\frac{1}{2} \sum_{t=s+1}^{e} \log \left(2 \pi \hat{\sigma}_{t}^{2}\right)
$$

which is maximised at

$$
\mu_{L}=\left(\sum_{t=s+1}^{b} \frac{1}{\hat{\sigma}_{t}^{2}}\right)^{-1}\left(\sum_{t=s+1}^{b} \frac{Y_{t}}{\hat{\sigma}_{t}^{2}}\right) \quad \text { and } \quad \mu_{R}=\left(\sum_{t=b+1}^{e} \frac{1}{\hat{\sigma}_{t}^{2}}\right)^{-1}\left(\sum_{t=b+1}^{e} \frac{Y_{t}}{\hat{\sigma}_{t}^{2}}\right) .
$$

Now suppose there is no change-point over $(s, e]$, then the log-likelihood is

$$
-\frac{1}{2} \sum_{t=s+1}^{e} \frac{\left(Y_{t}-\mu\right)^{2}}{\hat{\sigma}_{t}^{2}}-\frac{1}{2} \sum_{t=s+1}^{e} \log \left(2 \pi \hat{\sigma}_{t}^{2}\right)
$$

which is maximised at

$$
\mu=\left(\sum_{t=s+1}^{e} \frac{1}{\hat{\sigma}_{t}^{2}}\right)^{-1}\left(\sum_{t=s+1}^{e} \frac{Y_{t}}{\hat{\sigma}_{t}^{2}}\right)=\frac{\Omega_{b}^{2}-\Omega_{s}^{2}}{\Omega_{e}^{2}-\Omega_{s}^{2}} \mu_{L}+\frac{\Omega_{e}^{2}-\Omega_{b}^{2}}{\Omega_{e}^{2}-\Omega_{s}^{2}} \mu_{R}
$$

where $\Omega_{t}^{2}=\sum_{i=1}^{t} 1 / \hat{\sigma}_{i}^{2}$ (and by default $\Omega_{0}^{2}=0$ ). After some algebraic manupulation, we have that the generalised log-likelihood is

$$
\begin{aligned}
\mathcal{R}_{(s, e]}^{b}(\mathbf{Y}) & =\frac{1}{2}\left\{\left(\Omega_{b}^{2}-\Omega_{s}^{2}\right) \mu_{L}^{2}+\left(\Omega_{e}^{2}-\Omega_{b}^{2}\right) \mu_{R}^{2}-\left(\Omega_{e}^{2}-\Omega_{s}^{2}\right) \mu^{2}\right\} \\
& =\frac{1}{2} \frac{\left(\Omega_{b}^{2}-\Omega_{s}^{2}\right)\left(\Omega_{e}^{2}-\Omega_{s}^{2}\right)}{\Omega_{e}^{2}-\Omega_{b}^{2}}\left(\mu_{L}-\mu_{R}\right)^{2} \\
& =\frac{1}{2}\left(\sqrt{\frac{\Omega_{e}^{2}-\Omega_{b}^{2}}{\left(\Omega_{e}^{2}-\Omega_{s}^{2}\right)\left(\Omega_{b}^{2}-\Omega_{s}^{2}\right)}} \sum_{t=s+1}^{b} \frac{Y_{t}}{\hat{\sigma}_{t}^{2}}-\sqrt{\frac{\Omega_{b}^{2}-\Omega_{s}^{2}}{\left(\Omega_{e}^{2}-\Omega_{s}^{2}\right)\left(\Omega_{e}^{2}-\Omega_{b}^{2}\right)}} \sum_{t=b+1}^{e} \frac{Y_{t}}{\hat{\sigma}_{t}^{2}}\right)^{2} \\
& =\frac{1}{2}\left|\left\langle\mathbf{Y}, \hat{\boldsymbol{\psi}}_{(s, e]}^{b}\right\rangle\right|^{2} \\
& =\frac{1}{2}\left\{\hat{\mathcal{C}}_{(s, e]}^{b}(\mathbf{Y})\right\}^{2} .
\end{aligned}
$$

## E. Additional simulation results

In addition to the results presented in Section 5, here we present Tables 3-6 that summarise the results for three different distributions of the noise $\varepsilon_{t}$, where (b) $\varepsilon_{t} \stackrel{\text { i.i.d. }}{\sim} \mathcal{N}(0,2)$, (c) $\varepsilon_{t} \stackrel{\text { i.i.d. }}{\sim} \operatorname{Laplace}\left(0,2^{-1 / 2}\right),(\mathrm{d})$ i.i.d. scaled Student- $t_{5}$ in Table 5, and (e) $\varepsilon_{t}$ follows zero-mean unit-variance Gaussian $\operatorname{AR}(1)$ with $\varphi=0.3$.

Table 3. Distribution of $\hat{q}-q$ for data generated according to (3) with the noise term $\varepsilon_{t}$ being i.i.d. $\mathcal{N}(0,2)$ for various choices of $f_{t}$ and $\sigma_{t}$ given in Section A and competing methods listed in Section 5. Also, the average Mean-Square Error of the resulting estimate of the signal $f_{t}$, average inverse V-measure $d_{H}$, average V distance $d_{V}$ and average computation time in seconds using a single core of an Intel Xeon 3.6 GHz CPU with 16 GB of RAM, all calculated over 100 simulated data sets. Bold: methods with the largest empirical frequency of $\hat{q}-q=0$ or smallest average of $d_{H}$ or $d_{V}$, and those within $10 \%$ of the highest or lowest accordingly.

| Method | Model | $\hat{q}-q$ |  |  |  |  |  |  | MSE | $d_{H} \times 10^{2}$ | $d_{V}$ | time |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | $\leq-3$ | -2 | -1 | 0 | 1 | 2 | $\geq 3$ |  |  |  |  |
| B\&P |  | 0 | 0 | 0 | 100 | 0 | 0 | 0 | 0.107 | 0.93 | 0.037 | 1.337 |
| e-cp3o |  | 0 | 0 | 0 | 100 | 0 | 0 | 0 | 0.132 | 0.98 | 0.052 | 0.12 |
| FDRSeg |  | 0 | 0 | 0 | 83 | 14 | 1 | 2 | 0.134 | 1.51 | 0.054 | 0.093 |
| NMCD |  | 0 | 0 | 0 | 92 | 8 | 0 | 0 | 0.151 | 1.44 | 0.059 | 1.067 |
| NOT |  | 0 | 0 | 0 | 97 | 3 | 0 | 0 | 0.111 | 1.05 | 0.038 | 0.044 |
| NOT HT | (M1) | 0 | 0 | 0 | 97 | 3 | 0 | 0 | 0.126 | 1.25 | 0.044 | 0.058 |
| NP-PELT |  | 0 | 0 | 0 | 73 | 25 | 2 | 0 | 0.141 | 1.55 | 0.048 | 0.019 |
| PELT |  | 0 | 2 | 0 | 98 | 0 | 0 | 0 | 0.115 | 1.22 | 0.039 | 0.002 |
| S3IB |  | 0 | 0 | 0 | 90 | 8 | 1 | 1 | 0.114 | 1.13 | 0.038 | 0.076 |
| SMUCE |  | 0 | 2 | 14 | 84 | 0 | 0 | 0 | 0.185 | 2.14 | 0.064 | 0.056 |
| WBS |  | 0 | 0 | 0 | 93 | 7 | 0 | 0 | 0.113 | 1.15 | 0.038 | 0.074 |
| B\&P |  | 100 | 0 | 0 | 0 | 0 | 0 | 0 | 0.145 | 8.78 | 0.155 | 30.413 |
| e-cp3o |  | 100 | 0 | 0 | 0 | 0 | 0 | 0 | 0.223 | 7.74 | 0.15 | 2.425 |
| FDRSeg |  | 20 | 33 | 36 | 7 | 2 | 2 | 0 | 0.073 | 3.31 | 0.066 | 2.576 |
| NMCD |  | 46 | 27 | 21 | 6 | 0 | 0 | 0 | 0.076 | 4.29 | 0.074 | 4.324 |
| NOT |  | 28 | 30 | 27 | 13 | 2 | 0 | 0 | 0.066 | 3.4 | 0.059 | 0.077 |
| NOT HT | (M2) | 49 | 27 | 19 | 2 | 3 | 0 | 0 | 0.083 | 4.26 | 0.077 | 0.138 |
| NP-PELT |  | 4 | 9 | 30 | 21 | 21 | 10 | 5 | 0.068 | 3.74 | 0.062 | 0.239 |
| PELT |  | 91 | 7 | 2 | 0 | 0 | 0 | 0 | 0.114 | 8.21 | 0.122 | 0.004 |
| S3IB |  | 37 | 34 | 17 | 10 | 1 | 1 | 0 | 0.071 | 4.15 | 0.068 | 0.342 |
| SMUCE |  | 100 | 0 | 0 | 0 | 0 | 0 | 0 | 0.144 | 5.95 | 0.13 | 0.022 |
| WBS |  | 26 | 32 | 29 | 13 | 0 | 0 | 0 | 0.067 | 3.55 | 0.062 | 0.145 |
| B\&P |  | 0 | 0 | 100 | 0 | 0 | 0 | 0 | 0.258 | 4.25 | 0.155 | 54.381 |
| NOT | (M3) | 0 | 0 | 0 | 97 | 2 | 1 | 0 | 0.033 | 1.59 | 0.073 | 0.35 |
| TF |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.032 | 8.42 | 0.216 | 46.038 |
| B\&P |  | 13 | 53 | 28 | 6 | 0 | 0 | 0 | 0.322 | 6.11 | 0.204 | 62.421 |
| NOT | (M4) | 0 | 0 | 0 | 100 | 0 | 0 | 0 | 0.037 | 2.01 | 0.097 | 0.335 |
| TF |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.03 | 4.47 | 0.151 | 47.536 |
| B\&P |  | 0 | 0 | 9 | 91 | 0 | 0 | 0 | 0.046 | 3.52 | 0.115 | 119.454 |
| NOT | (M5) | 0 | 0 | 7 | 92 | 1 | 0 | 0 | 0.047 | 3.65 | 0.117 | 0.334 |
| TF |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.041 | 5.9 | 0.24 | 57.36 |
| e-cp3o |  | 11 | 12 | 12 | 33 | 20 | 5 | 7 | 0.145 | 6.91 | 0.164 | 1.738 |
| HSMUCE |  | 97 | 3 | 0 | 0 | 0 | 0 | 0 | 0.091 | 12.77 | 0.209 | 0.051 |
| NMCD |  | 0 | 0 | 18 | 70 | 12 | 0 | 0 | 0.06 | 4.04 | 0.068 | 4.135 |
| NOT | (M6) | 0 | 0 | 13 | 85 | 2 | 0 | 0 | 0.047 | 2.6 | 0.048 | 0.455 |
| NP-PELT |  | 0 | 0 | 1 | 19 | 26 | 24 | 30 | 0.126 | 3.17 | 0.068 | 0.279 |
| PELT |  | 9 | 18 | 31 | 37 | 5 | 0 | 0 | 0.069 | 8.17 | 0.087 | 0.011 |
| SegNeigh |  | 0 | 0 | 3 | 49 | 36 | 10 | 2 | 0.053 | 1.98 | 0.048 | 17.211 |
| B\&P |  | 0 | 0 | 42 | 58 | 0 | 0 | 0 | 0.073 | 6.21 | 0.132 | 29.222 |
| NOT | (M7) | 0 | 0 | 43 | 57 | 0 | 0 | 0 | 0.071 | 6.13 | 0.122 | 0.225 |
| TF |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.08 | 22.86 | 0.399 | 43.198 |

Table 4. Distribution of $\hat{q}-q$ for data generated according to (3) with the noise term $\varepsilon_{t}$ being i.i.d. Laplace $\left(0,(\sqrt{2})^{-1}\right)\left(\right.$ N.B. $\left.\operatorname{Var}\left(\varepsilon_{t}\right)=1\right)$ for various choices of $f_{t}$ and $\sigma_{t}$ given in Section A and competing methods listed in Section 5. Also, the average Mean-Square Error of the resulting estimate of the signal $f_{t}$, average Hausdorff distance $d_{H}$, average inverse V-measure $d_{V}$ and average computation time in seconds using a single core of an Intel Xeon 3.6 GHz CPU with 16 GB of RAM, all calculated over 100 simulated data sets. Bold: methods with the largest empirical frequency of $\hat{q}-q=0$ or smallest average of $d_{H}$ or $d_{V}$, and those within $10 \%$ of the highest or lowest accordingly.

| Method | $\hat{q}-q$ |  |  |  |  |  |  |  | MSE | $d_{H} \times 10^{2}$ | $d_{V}$ | time |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | Model | $\leq-3$ | -2 | -1 | 0 | 1 | 2 | $\geq 3$ |  |  |  |  |
| B\&P |  | 0 | 0 | 0 | 100 | 0 | 0 | 0 | 0.105 | 1.02 | 0.035 | 1.37 |
| e-cp3o |  | 0 | 0 | 0 | 100 | 0 | 0 | 0 | 0.123 | 0.91 | 0.049 | 0.116 |
| FDRSeg |  | 0 | 0 | 0 | 2 | 1 | 5 | 92 | 0.207 | 5.16 | 0.116 | 0.078 |
| NMCD |  | 0 | 0 | 0 | 92 | 7 | 1 | 0 | 0.15 | 1.55 | 0.059 | 1.053 |
| NOT |  | 0 | 0 | 0 | 93 | 4 | 3 | 0 | 0.114 | 1.3 | 0.038 | 0.043 |
| NOT HT | (M1) | 0 | 0 | 0 | 100 | 0 | 0 | 0 | 0.099 | 0.96 | 0.033 | 0.058 |
| NP-PELT |  | 0 | 0 | 0 | 55 | 31 | 12 | 2 | 0.154 | 2.07 | 0.05 | 0.018 |
| PELT |  | 0 | 0 | 0 | 58 | 17 | 16 | 9 | 0.17 | 1.89 | 0.042 | 0.001 |
| S3IB |  | 0 | 0 | 0 | 70 | 12 | 12 | 6 | 0.154 | 1.64 | 0.04 | 0.068 |
| SMUCE |  | 0 | 0 | 0 | 48 | 20 | 22 | 10 | 0.154 | 2.68 | 0.066 | 0.057 |
| WBS |  | 0 | 0 | 0 | 60 | 4 | 22 | 14 | 0.173 | 2.18 | 0.044 | 0.073 |
| B\&P |  | 100 | 0 | 0 | 0 | 0 | 0 | 0 | 0.144 | 8.7 | 0.151 | 32.221 |
| e-cp3o |  | 100 | 0 | 0 | 0 | 0 | 0 | 0 | 0.208 | 7.62 | 0.145 | 2.313 |
| FDRSeg |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.1 | 7.58 | 0.158 | 1.909 |
| NMCD |  | 24 | 36 | 35 | 4 | 1 | 0 | 0 | 0.06 | 3.57 | 0.056 | 4.354 |
| NOT |  | 63 | 14 | 12 | 6 | 4 | 1 | 0 | 0.085 | 5.09 | 0.082 | 0.078 |
| NOT HT | (M2) | 31 | 27 | 34 | 8 | 0 | 0 | 0 | 0.056 | 3.2 | 0.049 | 0.139 |
| NP-PELT |  | 1 | 0 | 12 | 22 | 24 | 16 | 25 | 0.084 | 4.17 | 0.06 | 0.223 |
| PELT |  | 25 | 23 | 15 | 15 | 10 | 8 | 4 | 0.112 | 5.35 | 0.081 | 0.004 |
| S3IB |  | 90 | 7 | 1 | 2 | 0 | 0 | 0 | 0.125 | 9.55 | 0.142 | 0.315 |
| SMUCE |  | 21 | 15 | 19 | 11 | 13 | 13 | 8 | 0.111 | 6.03 | 0.12 | 0.019 |
| WBS |  | 25 | 11 | 15 | 13 | 14 | 5 | 17 | 0.11 | 5.34 | 0.083 | 0.143 |
| B\&P |  | 0 | 0 | 100 | 0 | 0 | 0 | 0 | 0.255 | 4.1 | 0.153 | 54.071 |
| NOT | (M3) | 0 | 0 | 0 | 93 | 5 | 2 | 0 | 0.038 | 1.93 | 0.078 | 0.345 |
| TF |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.035 | 8.42 | 0.224 | 46.298 |
| B\&P |  | 10 | 49 | 35 | 6 | 0 | 0 | 0 | 0.311 | 6.27 | 0.204 | 61.911 |
| NOT | (M4) | 0 | 0 | 1 | 93 | 6 | 0 | 0 | 0.042 | 2.26 | 0.1 | 0.309 |
| TF |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.03 | 4.41 | 0.153 | 47.57 |
| B\&P |  | 0 | 0 | 10 | 90 | 0 | 0 | 0 | 0.044 | 3.47 | 0.112 | 118.603 |
| NOT | (M5) | 0 | 0 | 5 | 92 | 3 | 0 | 0 | 0.045 | 3.62 | 0.112 | 0.329 |
| TF |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.041 | 5.87 | 0.232 | 57.763 |
| e-cp3o |  | 34 | 19 | 9 | 11 | 6 | 7 | 14 | 0.304 | 9.94 | 0.225 | 1.693 |
| HSMUCE |  | 100 | 0 | 0 | 0 | 0 | 0 | 0 | 0.199 | 15.6 | 0.275 | 0.064 |
| NMCD |  | 4 | 18 | 44 | 31 | 3 | 0 | 0 | 0.114 | 9.25 | 0.116 | 4.085 |
| NOT | (M6) | 2 | 6 | 33 | 38 | 18 | 2 | 1 | 0.185 | 7.82 | 0.107 | 0.451 |
| NP-PELT |  | 0 | 1 | 0 | 14 | 24 | 23 | 38 | 0.364 | 5.44 | 0.109 | 0.32 |
| PELT |  | 26 | 13 | 35 | 22 | 4 | 0 | 0 | 0.226 | 13.41 | 0.148 | 0.013 |
| SegNeigh |  | 0 | 0 | 7 | 30 | 38 | 13 | 12 | 0.176 | 5.23 | 0.094 | 17.316 |
| B\&P |  | 0 | 0 | 39 | 60 | 1 | 0 | 0 | 0.067 | 6.23 | 0.123 | 28.903 |
| NOT | (M7) | 0 | 2 | 50 | 48 | 0 | 0 | 0 | 0.077 | 7.42 | 0.136 | 0.211 |
| TF |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.074 | 22.81 | 0.406 | 43.53 |

Table 5. Distribution of $\hat{q}-q$ for data generated according to (3) with the noise term $\varepsilon_{t}$ being i.i.d. $(3 / 5)^{1 / 2} t_{5}$ (N.B. $\operatorname{Var}\left(\varepsilon_{t}\right)=1$ ) for various choices of $f_{t}$ and $\sigma_{t}$ given in Section A of the online supplementary materials and competing methods listed in Section 5. Also, the average Mean-Square Error of the resulting estimate of the signal $f_{t}$, average Hausdorff distance $d_{H}$, average inverse V-measure $d_{V}$ and average computation time in seconds using a single core of an Intel Xeon 3.6 GHz CPU with 16 GB of RAM, all calculated over 100 simulated data sets. Bold: methods with the largest empirical frequency of $\hat{q}-q=0$ or smallest average of $d_{H}$ or $d_{V}$, and those within $10 \%$ of the highest or lowest accordingly.

| Method | Model | $\hat{q}-q$ |  |  |  |  |  |  | MSE | $d_{H} \times 10^{2}$ | $d_{V}$ | time |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | $\leq-3$ | -2 | -1 | 0 | 1 | 2 | $\geq 3$ |  |  |  |  |
| B\&P |  | 0 | 0 | 0 | 100 | 0 | 0 | 0 | 0.046 | 0.45 | 0.016 | 1.36 |
| e-cp3o |  | 0 | 0 | 0 | 100 | 0 | 0 | 0 | 0.087 | 0.58 | 0.04 | 0.119 |
| FDRSeg |  | 0 | 0 | 0 | 8 | 2 | 5 | 85 | 0.113 | 4.67 | 0.089 | 0.07 |
| NMCD |  | 0 | 0 | 0 | 97 | 3 | 0 | 0 | 0.089 | 0.67 | 0.041 | 1.07 |
| NOT |  | 0 | 0 | 0 | 96 | 4 | 0 | 0 | 0.049 | 0.53 | 0.017 | 0.047 |
| NOT HT | (M1) | 0 | 0 | 0 | 99 | 1 | 0 | 0 | 0.045 | 0.48 | 0.016 | 0.057 |
| NP-PELT |  | 0 | 0 | 0 | 75 | 12 | 12 | 1 | 0.081 | 1.35 | 0.031 | 0.015 |
| PELT |  | 0 | 0 | 0 | 53 | 7 | 25 | 15 | 0.106 | 1.89 | 0.026 | 0.002 |
| S3IB |  | 0 | 0 | 0 | 50 | 10 | 28 | 12 | 0.105 | 1.97 | 0.026 | 0.066 |
| SMUCE |  | 0 | 0 | 0 | 43 | 13 | 21 | 23 | 0.093 | 2.65 | 0.054 | 0.056 |
| WBS |  | 0 | 0 | 0 | 43 | 3 | 29 | 25 | 0.12 | 2.45 | 0.031 | 0.071 |
| B\&P |  | 100 | 0 | 0 | 0 | 0 | 0 | 0 | 0.126 | 5.71 | 0.128 | 33.68 |
| e-cp3o |  | 100 | 0 | 0 | 0 | 0 | 0 | 0 | 0.186 | 6.77 | 0.129 | 1.996 |
| FDRSeg |  | 0 | 0 | 0 | 0 | 0 | 2 | 98 | 0.042 | 7.02 | 0.11 | 1.56 |
| NMCD |  | 0 | 6 | 55 | 39 | 0 | 0 | 0 | 0.03 | 1.8 | 0.032 | 4.355 |
| NOT |  | 3 | 10 | 51 | 20 | 13 | 3 | 0 | 0.029 | 3.49 | 0.038 | 0.077 |
| NOT HT | (M2) | 0 | 3 | 52 | 44 | 1 | 0 | 0 | 0.023 | 1.48 | 0.022 | 0.136 |
| NP-PELT |  | 0 | 0 | 13 | 22 | 19 | 23 | 23 | 0.043 | 3.98 | 0.039 | 0.2 |
| PELT |  | 1 | 5 | 16 | 28 | 18 | 12 | 20 | 0.056 | 3.63 | 0.04 | 0.003 |
| S3IB |  | 26 | 18 | 23 | 21 | 9 | 3 | 0 | 0.058 | 4.21 | 0.054 | 0.299 |
| SMUCE |  | 1 | 9 | 10 | 22 | 24 | 6 | 28 | 0.05 | 5.49 | 0.074 | 0.016 |
| WBS |  | 2 | 3 | 24 | 7 | 22 | 11 | 31 | 0.058 | 4.49 | 0.046 | 0.143 |
| B\&P |  | 0 | 0 | 100 | 0 | 0 | 0 | 0 | 0.221 | 3.67 | 0.132 | 53.919 |
| NOT | (M3) | 0 | 0 | 0 | 97 | 3 | 0 | 0 | 0.016 | 1.05 | 0.054 | 0.395 |
| TF |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.019 | 8.36 | 0.221 | 46.891 |
| B\&P |  | 0 | 0 | 9 | 91 | 0 | 0 | 0 | 0.082 | 2.85 | 0.143 | 61.857 |
| NOT | (M4) | 0 | 0 | 0 | 98 | 1 | 1 | 0 | 0.017 | 1.29 | 0.07 | 0.371 |
| TF |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.018 | 4.41 | 0.151 | 48.119 |
| B\&P |  | 0 | 0 | 0 | 100 | 0 | 0 | 0 | 0.018 | 2.17 | 0.082 | 118.05 |
| NOT | (M5) | 0 | 0 | 2 | 90 | 7 | 1 | 0 | 0.021 | 2.53 | 0.086 | 0.368 |
| TF |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.026 | 5.98 | 0.26 | 59.006 |
| e-cp3o |  | 19 | 4 | 12 | 34 | 19 | 7 | 5 | 0.141 | 6.83 | 0.17 | 1.695 |
| HSMUCE |  | 100 | 0 | 0 | 0 | 0 | 0 | 0 | 0.098 | 12.68 | 0.212 | 0.052 |
| NMCD |  | 0 | 13 | 40 | 42 | 5 | 0 | 0 | 0.056 | 7.67 | 0.088 | 4.123 |
| NOT | (M6) | 0 | 3 | 11 | 51 | 23 | 9 | 3 | 0.08 | 5.09 | 0.084 | 0.463 |
| NP-PELT |  | 0 | 0 | 3 | 15 | 19 | 19 | 44 | 0.194 | 5.08 | 0.089 | 0.281 |
| PELT |  | 5 | 16 | 27 | 40 | 9 | 3 | 0 | 0.09 | 7.71 | 0.099 | 0.012 |
| SegNeigh |  | 0 | 0 | 7 | 26 | 28 | 20 | 19 | 0.094 | 4.33 | 0.077 | 17.3 |
| B\&P |  | 0 | 0 | 0 | 99 | 1 | 0 | 0 | 0.022 | 2.26 | 0.071 | 28.876 |
| NOT | (M7) | 0 | 0 | 6 | 86 | 8 | 0 | 0 | 0.027 | 3.03 | 0.078 | 0.226 |
| TF |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.049 | 23.29 | 0.442 | 42.538 |

Table 6. Distribution of $\hat{q}-q$ for data generated according to (3) with the noise term $\varepsilon_{t}$ being a zero-mean unit-variance Gaussian $\operatorname{AR}(1)$ process with $\varphi=0.3$ for various choices of $f_{t}$ and $\sigma_{t}$ given in Section A and competing methods listed in Section 5. Also, the average Mean-Square Error of the resulting estimate of the signal $f_{t}$, average Hausdorff distance $d_{H}$, average inverse V-measure $d_{V}$ and average computation time in seconds using a single core of an Intel Xeon 3.6 GHz CPU with 16 GB of RAM, all calculated over 100 simulated data sets. Bold: methods with the largest empirical frequency of $\hat{q}-q=0$ or smallest average of $d_{H}$ or $d_{V}$, and those within $10 \%$ of the highest or lowest accordingly.

| Method | Model | $\hat{q}-q$ |  |  |  |  |  |  | MSE | $d_{H} \times 10^{2}$ | $d_{V}$ | time |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | $\leq-3$ | -2 | -1 | 0 | 1 | 2 | $\geq 3$ |  |  |  |  |
| B\&P |  | 0 | 0 | 0 | 100 | 0 | 0 | 0 | 0.088 | 0.84 | 0.028 | 1.361 |
| e-cp3o |  | 0 | 0 | 0 | 100 | 0 | 0 | 0 | 0.126 | 0.99 | 0.05 | 0.116 |
| FDRSeg |  | 0 | 0 | 0 | 1 | 1 | 4 | 94 | 0.199 | 5.59 | 0.128 | 0.07 |
| NMCD |  | 0 | 0 | 0 | 63 | 29 | 6 | 2 | 0.145 | 2.36 | 0.06 | 1.048 |
| NOT |  | 0 | 0 | 0 | 64 | 18 | 7 | 11 | 0.113 | 2.13 | 0.04 | 0.046 |
| NOT HT | (M1) | 0 | 0 | 0 | 78 | 19 | 2 | 1 | 0.104 | 1.67 | 0.036 | 0.058 |
| NP-PELT |  | 0 | 0 | 0 | 39 | 31 | 20 | 10 | 0.134 | 2.63 | 0.05 | 0.017 |
| PELT |  | 0 | 0 | 0 | 73 | 21 | 3 | 3 | 0.106 | 1.88 | 0.036 | 0.001 |
| S3IB |  | 0 | 0 | 0 | 73 | 22 | 3 | 2 | 0.102 | 1.79 | 0.034 | 0.069 |
| SMUCE |  | 0 | 0 | 0 | 56 | 30 | 10 | 4 | 0.136 | 2.52 | 0.059 | 0.053 |
| WBS |  | 0 | 0 | 0 | 63 | 20 | 7 | 10 | 0.11 | 2.18 | 0.038 | 0.072 |
| B\&P |  | 100 | 0 | 0 | 0 | 0 | 0 | 0 | 0.136 | 6.67 | 0.14 | 30.394 |
| e-cp3o |  | 100 | 0 | 0 | 0 | 0 | 0 | 0 | 0.202 | 6.81 | 0.137 | 2.046 |
| FDRSeg |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.121 | 8.87 | 0.209 | 1.401 |
| NMCD |  | 1 | 9 | 37 | 35 | 13 | 4 | 1 | 0.056 | 2.92 | 0.055 | 4.316 |
| NOT |  | 1 | 8 | 34 | 25 | 9 | 12 | 11 | 0.053 | 3.63 | 0.053 | 0.082 |
| NOT HT | (M2) | 5 | 14 | 39 | 24 | 8 | 7 | 3 | 0.056 | 3.34 | 0.056 | 0.136 |
| NP-PELT |  | 0 | 1 | 1 | 10 | 17 | 14 | 57 | 0.067 | 4.98 | 0.074 | 0.192 |
| PELT |  | 1 | 11 | 30 | 38 | 10 | 9 | 1 | 0.048 | 2.73 | 0.045 | 0.003 |
| S3IB |  | 11 | 27 | 39 | 20 | 3 | 0 | 0 | 0.05 | 3.1 | 0.048 | 0.34 |
| SMUCE |  | 0 | 12 | 36 | 26 | 21 | 3 | 2 | 0.057 | 4.45 | 0.066 | 0.015 |
| WBS |  | 2 | 10 | 29 | 27 | 11 | 12 | 9 | 0.052 | 3.41 | 0.051 | 0.141 |
| B\&P |  | 0 | 0 | 91 | 9 | 0 | 0 | 0 | 0.245 | 4.37 | 0.147 | 53.676 |
| NOT | (M3) | 0 | 0 | 0 | 96 | 4 | 0 | 0 | 0.03 | 1.51 | 0.07 | 0.394 |
| TF |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.465 | 9.08 | 0.519 | 46.654 |
| B\&P |  | 0 | 1 | 25 | 74 | 0 | 0 | 0 | 0.136 | 3.74 | 0.159 | 61.576 |
| NOT | (M4) | 0 | 0 | 0 | 97 | 2 | 1 | 0 | 0.035 | 2.03 | 0.095 | 0.378 |
| TF |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.479 | 5 | 0.462 | 47.875 |
| B\&P |  | 0 | 0 | 0 | 98 | 2 | 0 | 0 | 0.04 | 3.28 | 0.113 | 117.832 |
| NOT | (M5) | 0 | 0 | 0 | 89 | 8 | 2 | 1 | 0.043 | 3.55 | 0.115 | 0.346 |
| TF |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.218 | 6.24 | 0.461 | 56.926 |
| e-cp3o |  | 19 | 9 | 16 | 23 | 13 | 10 | 10 | 0.224 | 8.25 | 0.19 | 1.659 |
| HSMUCE |  | 65 | 30 | 5 | 0 | 0 | 0 | 0 | 0.117 | 12.78 | 0.196 | 0.05 |
| NMCD |  | 1 | 0 | 5 | 28 | 29 | 18 | 19 | 0.178 | 5.36 | 0.093 | 4.097 |
| NOT | (M6) | 0 | 2 | 23 | 56 | 13 | 5 | 1 | 0.123 | 5.29 | 0.074 | 0.455 |
| NP-PELT |  | 0 | 0 | 0 | 0 | 1 | 1 | 98 | 0.482 | 5.49 | 0.127 | 0.219 |
| PELT |  | 9 | 17 | 28 | 40 | 6 | 0 | 0 | 0.126 | 8.78 | 0.1 | 0.011 |
| SegNeigh |  | 0 | 0 | 2 | 39 | 24 | 23 | 12 | 0.12 | 3.1 | 0.066 | 17.242 |
| B\&P |  | 0 | 0 | 2 | 86 | 11 | 1 | 0 | 0.045 | 4.13 | 0.101 | 28.884 |
| NOT | (M7) | 0 | 0 | 4 | 89 | 5 | 2 | 0 | 0.043 | 3.39 | 0.089 | 0.232 |
| TF |  | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.11 | 24.39 | 0.537 | 42.602 |

## F. Additional numerical experiments on the choice of $M$

## F.1. Setup

We now elaborate on the effect of the choice of $M$, the number of randomly drawn subintervals. We focus on Scenario (S1) and consider the models based on variations of (M1). All models listed below have piecewise-constant $f_{t}$ with equal-spaced change-points.
(M1-1) teeth-1: $T=512, q=1$ change-points at $\tau=256$, with the corresponding jump sizes $-2, f_{1}=1, \sigma_{t}=1$ for $t=1, \ldots, T$.
(M1-2) teeth-2: $T=512, q=3$ change-points at $\tau=128,256,384$, with the corresponding jump sizes $-2,2,-2, f_{1}=1, \sigma_{t}=1$ for $t=1, \ldots, T$.
(M1-3) teeth-3: $T=512, q=7$ change-points at $\tau=64,128, \ldots, 448$, with the corresponding jump sizes $-2,2,-2,2,-2, f_{1}=1, \sigma_{t}=1$ for $t=1, \ldots, T$. Note that this model is the same as (M1) teeth listed in Section A.
(M1-4) teeth-4: $T=512, q=15$ change-points at $\tau=32,64, \ldots, 480$, with the corresponding jump sizes $-2,2,-2, \ldots,-2, f_{1}=1, \sigma_{t}=1$ for $t=1, \ldots, T$.
(M1-5) teeth-5: $T=512, q=31$ change-points at $\tau=16,32, \ldots, 496$, with the corresponding jump sizes $-2,2,-2, \ldots,-2, f_{1}=1, \sigma_{t}=1$ for $t=1, \ldots, T$.
(M1-6) teeth-6: $T=512, q=63$ change-points at $\tau=8,16, \ldots, 504$, with the corresponding jump sizes $-2,2,-2, \ldots,-2, f_{1}=1, \sigma_{t}=1$ for $t=1, \ldots, T$.

We take $\varepsilon_{t} \stackrel{\text { i.i.d. }}{\sim} \mathcal{N}(0,1)$, run our NOT procedure using contrast function given by (6), and the threshold picked by the SIC, but with different $M=10,10^{2}, 10^{3}, 10^{4}, 10^{5}$.

## F.2. Results

In Table 7, we report the performance of NOT with the SIC in terms of estimates of $\mathbb{P}(\hat{q}=q), \mathbb{E}(\hat{q} / q), d_{H}$ and $d_{V}$ after 500 realisations. Here $d_{H}$ and $d_{V}$ are, respectively, the scaled Hausdorff distance measure and the inverse V-measure, both given in Section 5.3.

We see that when there is only a small number of change-points in the signal, a moderate $M$ would be sufficient for the purpose of identifying all the change-points. In this case, having a larger $M$ will be more computationally intensive, but will not necessarily improve the performance of the NOT procedure. As we increase the number of change-points (while fixing $T$, as well as the size of the jumps, etc), we see that a larger $M$ would be needed for satisfactory performance. For example, in Model (M1-4) (with $q=15$ ), our procedure with $M=10^{4}$ or $10^{5}$ estimates the number of the change-points correctly more than $95 \%$ of the time, while this proportion reduces to $87 \%$ for $M=10^{3}$, and virtually $0 \%$ for $M=10^{2}$ or smaller. In cases like that, increasing $M$ would be helpful.

Table 7. Performance of NOT with $\zeta_{T}$ picked via the SIC and different $M=10,10^{2}, 10^{3}, 10^{4}, 10^{5}$. Here data is generated according to (3) with the noise term $\varepsilon_{t}$ being i.i.d. $\mathcal{N}(0,1)$ for various models given in Section F.1. The estimated values of $\mathbb{P}(\hat{q}=q), \mathbb{E}(\hat{q} / q), d_{H}$ and $d_{V}$ are reported, all calculated over 500 simulated data sets.

| Model $(T=512)$ | Measure | $M$ |  |  |  |  |
| :--- | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | 10 | $10^{2}$ | $10^{3}$ | $10^{4}$ | $10^{5}$ |
| (M1-1): $q=1$ |  | 0.978 | 0.978 | 0.968 | 0.972 | 0.976 |
| (M1-2): $q=3$ |  | 0.418 | 0.988 | 0.982 | 0.976 | 0.978 |
| (M1-3): $q=7$ | $(\hat{q}=q)$ | 0.004 | 0.566 | 0.976 | 0.972 | 0.974 |
| (M1-4): $q=15$ |  | 0.006 | 0.872 | 0.958 | 0.956 |  |
| (M1-5): $q=31$ |  | 0 | 0 | 0.002 | 0.45 | 0.424 |
| (M1-6): $q=63$ |  | 0 | 0 | 0 | 0 | 0 |
| (M1-1): $q=1$ |  | 1.042 | 1.026 | 1.04 | 1.036 | 1.036 |
| (M1-2): $q=3$ |  | 0.925 | 1.004 | 1.006 | 1.009 | 1.008 |
| (M1-3): $q=7$ | $\mathbb{E}(\hat{q} / q)$ | 0.419 | 0.98 | 1.003 | 1.004 | 1.004 |
| (M1-4): $q=15$ |  | 0.06 | 0.565 | 1.003 | 1.003 | 1.003 |
| (M1-5): $q=31$ |  | 0.002 | 0.007 | 0.055 | 0.837 | 0.845 |
| (M1-6): $q=63$ |  | 0 | 0 | 0 | 0 | 0 |
| (M1-1): $q=1$ |  | 0.53 | 0.36 | 0.45 | 0.46 | 0.42 |
| (M1-2): $q=3$ |  | 14.23 | 0.33 | 0.36 | 0.39 | 0.38 |
| (M1-3): $q=7$ |  | $d_{H} \times 10^{2}$ | 22.71 | 5.03 | 0.53 | 0.53 |
| (M1-4): $q=15$ |  | 43.45 | 14.31 | 1.25 | 0.79 | 0.53 |
| (M1-5): $q=31$ |  | 49.82 | 49.04 | 45.49 | 8.4 | 7.92 |
| (M1-6): $q=63$ |  | 50 | 49.98 | 49.97 | 49.98 | 49.98 |
| (M1-1): $q=1$ |  | 0.019 | 0.014 | 0.015 | 0.015 | 0.015 |
| (M1-2): $q=3$ |  | 0.142 | 0.013 | 0.013 | 0.013 | 0.013 |
| (M1-3): $q=7$ |  | $d_{V}$ | 0.349 | 0.043 | 0.018 | 0.018 |
| (M1-4): $q=15$ |  | 0.848 | 0.242 | 0.029 | 0.026 | 0.027 |
| (M1-5): $q=31$ |  | 0.995 | 0.977 | 0.905 | 0.168 | 0.155 |
| (M1-6): $q=63$ |  | 1 | 0.999 | 0.999 | 0.999 | 0.999 |

Table 8. Performance of NOT with $\zeta_{T}$ picked via the AIC and different $M=10,10^{2}, 10^{3}, 10^{4}, 10^{5}$. Here data is generated according to (3) with the noise term $\varepsilon_{t}$ being i.i.d. $\mathcal{N}(0,1)$ for Model (M1-6) given in Section F.1. The estimated values of $\mathbb{P}(\hat{q}=q), \mathbb{E}(\hat{q} / q), d_{H}$ and $d_{V}$ are reported, all calculated over 500 simulated data sets.

|  | $M$ | $\mathbb{P}(\hat{q}=q)$ | $\mathbb{E}(\hat{q} / q)$ | $d_{H} \times 10^{2}$ | $d_{V}$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
| $(\mathrm{M} 1-6): q=63$ | 10 | 0 | 0.007 | 47.06 | 0.948 |
|  | $10^{2}$ | 0 | 0.041 | 37.29 | 0.776 |
|  | $10^{3}$ | 0 | 0.276 | 13.29 | 0.331 |
|  | $10^{4}$ | 0.08 | 0.937 | 1.76 | 0.063 |
|  | $10^{5}$ | 0.106 | 0.988 | 1.58 | 0.057 |

On the other hand, caution must be exercised for signals with an extremely large number of change-points, or the spacing of change-points to be highly non-homogeneous. For example, in Model (M1-6) (with $q=63$ ) where jumps in the signal occur at every 8 observations (which is itself a difficult problem), having $M=10^{5}$ or larger will not lead to any improvement of NOT with the SIC. However, we believe that this is partially due to the fact that the SIC penalty is no longer approperiate for this extreme scenario. One could alleviate the issue by using NOT with other less harsh penalty, or use methods designed to tackle frequent change-points such as Fryzlewicz (2018). For instance, by changing the SIC to the AIC in our procedure, we can see from Table 8 that with $M=10^{5}, \mathbb{P}(\hat{q}=q)$ increases from $0 \%$ to around $10 \%$. More importantly, there is huge improvement in terms of $\mathbb{E}(\hat{q} / q), d_{H}$ and $d_{V}$. In fact, a close inspection indicates that $\hat{q} / q \in[0.9,1.1]$ more than $90 \%$ of the time.

## G. More on model misspecification and model selection

We have demonstrated that NOT is relatively robust against the misspecification in the distribution of $\varepsilon_{t}$, when the truth is either correlated or heavy-tailed. Now we investigate the case where the signal $f_{t}$ is misspecified. In particular, we focus on the misspecification of the degree of the polynomials between consecutive change-points.

We simulate data according to (3) using the signal (M8) smile and noise of (a) i.i.d. $\mathcal{N}(0,1)$ and (b) i.i.d. $\mathcal{N}(0,2)$. Here the true signal is piecewise-linear but not necessarily continuous (i.e. from Scenario (S3)). We test NOT with the sSIC using contrast functions constructed from Scenarios (S1), (S3) and (S5), where the estimators are denoted by $\mathrm{NOT}_{0}$, $\mathrm{NOT}_{1}$ and $\mathrm{NOT}_{2}$, respectively. Again we take $\alpha=1$. Figure 9 shows a typical realisation of the estimates produced by NOT with different contrast functions, while Table 9 summarises the results.

For $\mathrm{NOT}_{0}$ (suitable for piecewise-constant signal), we see that unsurprisingly $\mathrm{NOT}_{0}$ significantly overestimates the number of change-points $q$. This is due to the bias-variance tradeoff in the sSIC, where the bias term only approaches zero as the estimated number of change-points $\hat{q} \rightarrow \infty$. Nevertheless, we observe that the set of change-point estimates from $\mathrm{NOT}_{0}$ typically includes the true change-points with jump, even though the construction of the contrast function (wrongly) assumes that the signal is piecewise-constant in the neighbourhood of these change-points. Furthermore, under the higher signal-to-noise ratio setting, $\mathrm{NOT}_{2}$, which is designed for piecewise-quadratic signal, is able to estimate the number of change-points $q$ correctly most of the time. However, since $\mathrm{NOT}_{2}$ is over-
parameterised in this setting of Scenario (S3), it tends to perform slightly worse than $\mathrm{NOT}_{1}$ in terms of both the MSE for the estimated signal, and the accuracy of the estimated locations of the change-points. Finally, under the lower signal-to-noise ratio setting, $\mathrm{NOT}_{2}$ tends to underestimate the number of change-points, thanks to the bias-variance tradeoff in the sSIC. Nevertheless, as is illustrated in Figures 9f, the estimated $f_{t}$ is quite close to the truth in terms of the $\ell_{2}$ distance. These findings suggest that NOT could still provide valuable insights in certain misspecified circumstances.


Fig. 9. Typical realisation of the estimates produced by different NOTs, with data generated from (M8) smile. Figure 9a-9f: data series $Y_{t}$ (thin grey), true signal $f_{t}$ (dashed black), $\hat{f}_{t}$ being the LS estimate of $f_{t}$ with the change-points estimated by NOT (thick red). Higher signal-to-noise ratio setting (with $\mathcal{N}(0,1)$ errors) in Figures 9a, 9c and 9e; lower signal-to-noise ratio setting (with $\mathcal{N}(0,2)$ errors) in Figures 9b, 9d and 9f. Here $\mathrm{NOT}_{0}, \mathrm{NOT}_{1}$ and $\mathrm{NOT}_{2}$ denote estimates from NOT with sSIC using contrast functions constructed from Scenarios (S1), (S3) and (S5), respectively.

In the same example, we also demonstrate that one could empirically select the degree of the polynomial for the NOT's contrast function via sSIC. Denote the sSIC scores corresponding to the estimates from $\mathrm{NOT}_{0}, \mathrm{NOT}_{1}$ and $\mathrm{NOT}_{2}$ by $\operatorname{sSIC}\left(\mathrm{NOT}_{0}\right), \mathrm{sSIC}\left(\mathrm{NOT}_{1}\right)$ and $\operatorname{sSIC}\left(\mathrm{NOT}_{2}\right)$ respectively. We propose to pick the estimator produced by $\mathrm{NOT}_{i^{*}}$ with

$$
i^{*}=\operatorname{argmin}_{i \in\{0,1,2\}} \mathrm{SSIC}\left(\mathrm{NOT}_{i}\right)
$$

As shown in Table 9, empirical results suggest that we are able to select the correct order of the polynomial for our NOT approach using the sSIC (with $\alpha=1$ ), especially when the signal-to-noise ratio is high.

Table 9. Distribution of $\hat{q}-q$ obtained by $\mathrm{NOT}_{0}, \mathrm{NOT}_{1}, \mathrm{NOT}_{2}$ for data generated according to (3) with the signal (M8) and the noise $\varepsilon_{t} \stackrel{\text { i.i.d. }}{\sim} \mathcal{N}(0,1)$ and $\mathcal{N}(0,2)$, the average MeanSquare Error of the resulting estimate of the signal over 100 simulations. The number of times each method selected by sSIC is also reported.

| Noise | Method | $\leq-3$ | -2 | -1 | 0 | 1 | 2 | $\geq 3$ | MSE | Number of times <br> selected by sSIC |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | NOT $_{0}$ | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.120 | 0 |
|  | NOT $_{1}$ | 0 | 0 | 0 | $\mathbf{9 9}$ | 1 | 0 | 0 | $\mathbf{0 . 0 1 5}$ | $\mathbf{1 0 0}$ |
|  | NOT $_{2}$ | 0 | 4 | 18 | 78 | 0 | 0 | 0 | 0.024 | 0 |
| $\mathcal{N}(0,2)$ | NOT $_{0}$ | 0 | 0 | 0 | 0 | 0 | 0 | 100 | 0.188 | 0 |
|  | NOT $_{1}$ | 0 | 0 | 0 | $\mathbf{1 0 0}$ | 0 | 0 | 0 | $\mathbf{0 . 0 3 2}$ | $\mathbf{9 4}$ |
|  | NOT $_{2}$ | 57 | 23 | 14 | 6 | 0 | 0 | 0 | 0.078 | 6 |

## H. Additional real data example: OPEC Reference Basket oil price

We perform change-point analysis on the daily Organisation of the Petroleum Exporting Countries (OPEC) Reference Basket oil price from 1 January, 2003 to 15 July, 2016. The data were obtained from the OPEC database through the R package Quandl (McTaggart et al., 2016). Instead of working with the raw price series, we analyse the log-returns series $Y_{t}=100 \log \left(P_{t} / P_{t-1}\right)$, where $P_{t}$ denotes the daily oil price. One of the stylised facts of the financial time series data is that the autocorrelation of assets returns are weak, while squared returns tend to exhibit strong autocorrelation, which is the case for the oil price time series (see Figure 10b). This phenomenon can be possibly explained by the existence of the structural breaks in the mean and variance structure of the data series (Mikosch and Stărică, 2004; Fryzlewicz et al., 2006). In this study, we apply NOT with the contrast function given by (14), which is designed to detect changes in both the mean and the volatility, as in Scenario (S4). For comparison, we also report change-points detected with the NMCD method of Zou et al. (2014).

We apply Algorithm 2 to compute the NOT solution path and choose the model achieving the lowest SIC given by (11), setting the number of intervals drawn $M=10000$ and the maximum number of change-points $q_{\max }=25$. Computations for the solution path and model selection are performed using the R package not (Baranowski et al., 2016). For the NMCD procedure, we use the nmcd routine from the R package nmcdr (Zou and Lancezhange, 2014), setting the maximum number of change-points to $q_{\max }=25$ as well.

Figure 10 illustrates the results of our analysis. The oil price time series and the locations of the change-points identified by NOT and NMCD can be seen in Figure 10a. Both methods discover 7 change-points, largely agreeing on their locations, in the sense that for 6 out of 7 features NOT detects, NMCD detects a change-point nearby. However, NMCD does not indicate any change-point around the first change-point identified by NOT on 29 April 2003. This date could potentially be related to the end of the 2003 invasion of Iraq, which initiated the upward trend in the oil price lasting almost ceaselessly until the beginning of the 2008-09 financial crisis. On the other hand, NMCD indicates two changepoints in the first quarter of 2016 , while NOT only finds one in that period. Table 10 lists the exact locations of the change-points detected by the two methods and the events that coincide with some of them. Figure 10 f shows the autocorrelation function for the squared residuals obtained by subtracting the sample mean and dividing by the standard deviations from the data in each segment. It appears that there is little autocorrelation in the squares of the residuals, suggesting that Scenario (S4) fits the data reasonably well.


Fig. 10. Change-point analysis on the daily OPEC Reference Basket oil price in USD from 1 January, 2003 to 15 July, 2016. Figure 10a: price series $P_{t}$ (thin grey), locations of the changepoints detected with NOT (vertical dotted lines) and NMCD (vertical dashed lines). Figure 10b: autocorrelation function of $Y_{t}^{2}$. Figure 10c: log-returns $Y_{t}=100 \log \left(P_{t} / P_{t-1}\right)$ (thin grey), the fitted piecewise-constant mean via NOT, $\hat{f}_{t}$ (thick red). Figure 10d: estimated residuals via NOT, $\hat{\varepsilon}_{t}=\left(Y_{t}-\hat{f}_{t}\right) / \hat{\sigma}_{t}$. Figure 10e: the centred log-returns $\left|Y_{t}-\hat{f}_{t}\right|$ (thin grey), fitted piecewise-constant volatility $\hat{\sigma}_{t}$ (thick red). Figure 10f: autocorrelation of $\hat{\varepsilon}_{t}^{2}$. The exact locations of the change-points detected via NOT are given in Table 10.

Table 10. Change-points detected using NOT and NMCD methods in the daily OPEC Reference Basket oil price data from 1 January 2003 to 15 July 2016, with some of them dated.

| NOT | NMCD | Event that coincides |
| :---: | :---: | :---: |
| 29 April 2003 | N/A | Invasion of Iraq |
| 1 September 2008 | 28 August 2008 | critical stage of the subprime mortgage crisis |
| 27 January 2009 | 22 January 2009 | tensions in the Gaza Strip |
| 1 October 2009 | 23 October 2009 |  |
| 12 November | 12 October 2012 | beginning of a period of low volatility |
| 2012 |  |  |
| $\begin{array}{ll} 30 & \text { September } \\ 2014 & \end{array}$ | 1 October 2014 |  |
| 5 January 2016 | 21 January 2016 | beginning of a sell-off leading the price to 12 -year low |
| N/A | $\begin{array}{ll} 22 & \text { February } \\ 2016 \end{array}$ |  |

## I. Proofs

## I.1. Some useful lemmas

## I.1.1. The piecewise-constant case

Lemma 1. Let $g(x, y)=\frac{x y}{x+y}$ and suppose that $\min (x, y)>0$. Then

$$
g(x, y) \geq \frac{1}{2} \min (x, y)
$$

Proof. Without loss of generality, assume that $x \geq y$. Then $g(x, y) \geq \frac{x y}{2 x} \geq y / 2=$ $\min (x, y) / 2$.

Lemma 2. Suppose $\mathbf{f}=\left(f_{1}, \ldots, f_{T}\right)^{\prime}$ is piecewise-constant vector as in Scenario (S1), and $\tau_{1}, \ldots, \tau_{q}$ are the locations of the change-points. Suppose $0 \leq s<e \leq T$, such that $\tau_{j-1} \leq s<\tau_{j}<e \leq \tau_{j+1}$ for some $j=1 \ldots, q$. Let $\eta=\min \left\{\tau_{j}-s, e-\tau_{j}\right\}$ and $\Delta_{j}^{\mathbf{f}}=\left|f_{\tau_{j}+1}-f_{\tau_{j}}\right|$. Then

$$
\mathcal{C}_{(s, e]}^{\tau_{j}}(\mathbf{f})=\max _{s<b<e} \mathcal{C}_{(s, e]}^{b}(\mathbf{f})\left\{\begin{array}{l}
\geq \frac{1}{\sqrt{2}} \eta^{1 / 2} \Delta_{j}^{\mathbf{f}} \\
\leq \eta^{1 / 2} \Delta_{j}^{\mathbf{f}}
\end{array}\right.
$$

Proof. For any $s<b<e$, by simple algebra, we have

$$
\mathcal{C}_{(s, e]}^{b}(\mathbf{f})= \begin{cases}\sqrt{\frac{b-s}{(e-s)(e-b)}}\left(e-\tau_{j}\right)\left|f_{\tau_{j}+1}-f_{\tau_{j}}\right|, & b \leq \tau_{j}  \tag{15}\\ \sqrt{\frac{\left(\tau_{j}-s\right)\left(e-\tau_{j}\right)}{e-s}}\left|f_{\tau_{j}+1}-f_{\tau_{j}}\right|, & b=\tau_{j} \\ \sqrt{\frac{e-b}{(e-s)(b-s)}}\left(\tau_{j}-s\right)\left|f_{\tau_{j}+1}-f_{\tau_{j}}\right|, & b \geq \tau_{j}\end{cases}
$$

$\operatorname{Now} \mathcal{C}_{(s, e]}^{\tau_{j}}(\mathbf{f})=\max _{s<b<e} \mathcal{C}_{(s, e]}^{b}(\mathbf{f})$ follows from the fact that $\mathcal{C}_{(s, e]}^{b}(\mathbf{f})$ is increasing (as a function of $b$ ) for $s<b \leq \tau_{j}$ and decreasing for $\tau_{j} \leq b<e$. To prove the lower bound, we set $\eta_{L}=\tau_{j}-s$ and $\eta_{R}=e-\tau_{j}$ and observe that $\eta_{L} \geq \eta$ and $\eta_{R} \geq \eta$. Therefore by Lemma 1, $\frac{\eta_{L} \eta_{R}}{\eta_{L}+\eta_{R}} \geq \frac{\eta}{2}$. Noting that $e-s=\eta_{L}+\eta_{R}$ we bound

$$
\mathcal{C}_{(s, e]}^{\tau_{j}}(\mathbf{f})=\sqrt{\frac{\left(\tau_{j}-s\right)\left(e-\tau_{j}\right)}{e-s}}\left|f_{\tau_{j}+1}-f_{\tau_{j}}\right|\left\{\begin{array}{l}
\geq(\eta / 2)^{1 / 2} \Delta_{j}^{\mathbf{f}} ; \\
\leq \eta^{1 / 2} \Delta_{j}^{\mathbf{f}}
\end{array}\right.
$$

which completes the proof.
Lemma 3. Suppose $\mathbf{f}=\left(f_{1}, \ldots, f_{T}\right)^{\prime}$ is piecewise-constant vector as in Scenario (S1), and $\tau_{1}, \ldots, \tau_{q}$ are the locations of the change-points. Suppose $0 \leq s<e \leq T$ such that $\tau_{j-1} \leq s \leq \tau_{j}$ and $\tau_{j+1} \leq e \leq \tau_{j+2}$ for some $j=1 \ldots, q-1$. Then

$$
\max _{s<b<e} \mathcal{C}_{(s, e]}^{b}(\mathbf{f}) \leq\left(\tau_{j}-s\right)^{1 / 2} \Delta_{j}^{\mathbf{f}}+\left(e-\tau_{j+1}\right)^{1 / 2} \Delta_{j+1}^{\mathbf{f}}
$$

where $\Delta_{j}^{\mathbf{f}}=\left|f_{\tau_{j}+1}-f_{\tau_{j}}\right|$.
Proof. Suppose that $b^{*}=\operatorname{argmax}_{s<b<e} \mathcal{C}_{(s, e]}^{b}(\mathbf{f})$. Then

$$
\begin{aligned}
0 & \leq\left\|\mathbf{f}-\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{b^{*}}\right\rangle \boldsymbol{\psi}_{(s, e]}^{b^{*}}-\left\langle\mathbf{f}, \mathbf{1}_{(s, e]}\right\rangle \mathbf{1}_{(s, e]}\right\|^{2}=\left\|\mathbf{f}-\left\langle\mathbf{f}, \mathbf{1}_{(s, e]}\right\rangle \mathbf{1}_{(s, e]}\right\|^{2}-\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{b^{*}}\right\rangle^{2} \\
& \leq\left\|\mathbf{f}-f_{\tau_{j}+1} \sqrt{e-s} \mathbf{1}_{(s, e]}\right\|^{2}-\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{b^{*}}\right\rangle^{2} \\
& =\left(\tau_{j}-s\right)\left(\Delta_{j}^{\mathbf{f}}\right)^{2}+\left(e-\tau_{j+1}\right)\left(\Delta_{j+1}^{\mathbf{f}}\right)^{2}-\left(\max _{s<b<e} \mathcal{C}_{(s, e]}^{b}(\mathbf{f})\right)^{2}
\end{aligned}
$$

It then follows that

$$
\max _{s<b<e} \mathcal{C}_{(s, e]}^{b}(\mathbf{f}) \leq \sqrt{\left(\tau_{j}-s\right)\left(\Delta_{j}^{\mathbf{f}}\right)^{2}+\left(e-\tau_{j+1}\right)\left(\Delta_{j+1}^{\mathbf{f}}\right)^{2}} \leq\left(\tau_{j}-s\right)^{1 / 2} \Delta_{j}^{\mathbf{f}}+\left(e-\tau_{j+1}\right)^{1 / 2} \Delta_{j+1}^{\mathbf{f}} .
$$

Lemma 4. Suppose $\mathbf{f}=\left(f_{1}, \ldots, f_{T}\right)^{\prime}$ is piecewise-constant vector as in Scenario (S1). Pick any interval $(s, e] \subset(0, T]$ such that $[s+1, e-1]$ contains exactly one change-point $\tau_{j}$. Let $\rho=\left|\tau_{j}-b\right|, \Delta_{j}^{\mathrm{f}}=\left|f_{\tau_{j}+1}-f_{\tau_{j}}\right|, \eta_{L}=\tau_{j}-s$ and $\eta_{R}=e-\tau_{j}$. Then,

$$
\left\|\boldsymbol{\psi}_{(s, e]}^{b}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle-\boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\right)\right\|_{2}^{2}=\left(\mathcal{C}_{(s, e]}^{\tau_{j}}(\mathbf{f})\right)^{2}-\left(\mathcal{C}_{(s, e]}^{b}(\mathbf{f})\right)^{2} .
$$

Moreover,
(a) for any $\tau_{j} \leq b<e,\left(\mathcal{C}_{(s, e]}^{\tau_{j}}(\mathbf{f})\right)^{2}-\left(\mathcal{C}_{(s, e]}^{b}(\mathbf{f})\right)^{2}=\frac{\rho \eta_{L}}{\rho+\eta_{L}}\left(\Delta_{j}^{\mathbf{f}}\right)^{2}$;
(b) for any $s<b \leq \tau_{j},\left(\mathcal{C}_{(s, e]}^{\tau_{j}}(\mathbf{f})\right)^{2}-\left(\mathcal{C}_{(s, e]}^{b}(\mathbf{f})\right)^{2}=\frac{\rho \eta_{R}}{\rho+\eta_{R}}\left(\Delta_{j}^{\mathbf{f}}\right)^{2}$.

Proof. First, we note that since there is only one change-point in [ $s+1, e-1$ ], the restriction of $\mathbf{f}$ on $(s, e]$, i.e. $\left.\mathbf{f}\right|_{(s, e]}=\left(0, \ldots, 0, f_{s+1}, \ldots, f_{e}, 0, \ldots, 0\right)^{\prime}$ can be decomposed into

$$
\left.\mathbf{f}\right|_{(s, e]}=\boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\right\rangle+\mathbf{1}_{(s, e]}\left\langle\mathbf{f}, \mathbf{1}_{(s, e]}\right\rangle,
$$

where we also used the fact that $\boldsymbol{\psi}_{(s, e]}^{\tau_{j}}$ and $\mathbf{1}_{(s, e]}$ are orthonormal. Note that $\boldsymbol{\psi}_{(s, e]}^{b}$ and $\mathbf{1}_{(s, e]}$ are also orthonormal, it follows that
$\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle=\left\langle\left.\mathbf{f}\right|_{(s, e]}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle=\left\langle\boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\right\rangle+\mathbf{1}_{(s, e]}\left\langle\mathbf{f}, \mathbf{1}_{(s, e]}\right\rangle, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle=\left\langle\boldsymbol{\psi}_{(s, e]}^{\tau_{j}}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\right\rangle$.
Therefore,

$$
\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{b}\right)^{b}=\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle\left\langle\boldsymbol{\psi}_{(s, e]}^{\tau_{j}}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\right\rangle,
$$

and thus

$$
\begin{aligned}
\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\right)^{2}-\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle^{2} & =\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\right)^{2}+\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{b}\right)^{2}-2\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle\left\langle\boldsymbol{\psi}_{(s, e]}^{\tau_{j}}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\right\rangle \\
& =\left\|\boldsymbol{\psi}_{(s, e]}^{b}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle-\boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\right\rangle\right\|_{2}^{2} .
\end{aligned}
$$

Here in the above final step, we used the fact that $\left\|\boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\right\|_{2}^{2}=\left\|\boldsymbol{\psi}_{(s, e]}^{b}\right\|_{2}^{2}=1$.
Second, for the sake of brevity, we only prove the case of $b \geq \tau_{j}$. Let $l=e-s, x=b-s$, and thus $\rho=x-\eta_{L}$. Using (15), we get

$$
\begin{aligned}
\left(\mathcal{C}_{(s, e]}^{\tau_{j}}(\mathbf{f})\right)^{2}-\left(\mathcal{C}_{(s, e]}^{b}(\mathbf{f})\right)^{2} & =\left(\frac{\eta_{L}\left(l-\eta_{L}\right)}{l}-\frac{\eta_{L}^{2}(l-x)}{l x}\right)\left|f_{\tau_{j}+1}-f_{\tau_{j}}\right|^{2} \\
& =\frac{\eta_{L}\left(x-\eta_{L}\right)}{x}\left(\Delta_{j}^{\mathrm{f}}\right)^{2}=\left(\frac{\rho \eta_{L}}{\eta_{L}+\rho}\right)\left(\Delta_{j}^{\mathrm{f}}\right)^{2} .
\end{aligned}
$$

## I.1.2. The piecewise-linear continuous case

Lemma 5. Suppose $\mathbf{f}=\left(f_{1}, \ldots, f_{T}\right)^{\prime}$ is piecewise-linear vector as in Scenario (S2), and $\tau_{1}, \ldots, \tau_{q}$ are the locations of the change-points. Suppose $0 \leq s<e \leq T$, such that $\tau_{j-1} \leq s+1<\tau_{j}<e \leq \tau_{j+1}$ for some $j=1 \ldots, q$. Let $\eta=\min \left\{\tau_{j}-s-1, e-\tau_{j}\right\}$ and $\Delta_{j}^{\mathbf{f}}=\left|2 f_{\tau_{j}}-f_{\tau_{j}-1}-f_{\tau_{j}+1}\right|$. Then

$$
\mathcal{C}_{(s, e]}^{\tau_{j}}(\mathbf{f})=\max _{s+1<b<e} \mathcal{C}_{(s, e]}^{b}(\mathbf{f})\left\{\begin{array}{l}
\geq \frac{1}{\sqrt{24}} \eta^{3 / 2} \Delta_{j}^{\mathbf{f}} \\
\leq \frac{1}{\sqrt{3}}(\eta+1)^{3 / 2} \Delta_{j}^{\mathbf{f}}
\end{array}\right.
$$

Proof. First, we show that $\mathcal{C}_{(s, e]}^{b}(\mathbf{f})$ is maximised at $b=\tau_{j}$. Using the notation from the proof of Lemma 4, we have that

$$
\left.\mathbf{f}\right|_{(s, e]}=\phi_{(s, e]}^{\tau_{j}}\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{\tau_{j}}\right\rangle+\gamma_{(s, e]}\left\langle\mathbf{f}, \mathbf{1}_{(s, e]}\right\rangle+\mathbf{1}_{(s, e]}\left\langle\mathbf{f}, \mathbf{1}_{(s, e]}\right\rangle .
$$

Therefore, it follows that

$$
\begin{equation*}
\left\|\left.\mathbf{f}\right|_{(s, e]}\right\|_{2}^{2}=\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{\tau_{j}}\right\rangle^{2}+\left\langle\mathbf{f}, \gamma_{(s, e]}\right\rangle^{2}+\left\langle\mathbf{f}, \mathbf{1}_{(s, e]}\right\rangle^{2} . \tag{16}
\end{equation*}
$$

For any $b \in\left\{s+2, \ldots, \tau_{j}-1, \tau_{j}+1, \ldots, e-1\right\}$, it is clear that $\left.\mathbf{f}\right|_{(s, e]}$ does not lie in the span of $\boldsymbol{\phi}_{(s, e]}^{b}, \gamma_{(s, e]}$ and $\mathbf{1}_{(s, e]}$. Consequently, by projecting $\left.\mathbf{f}\right|_{(s, e]}$ onto these three bases, we have that

$$
\begin{equation*}
\left\|\left.\mathbf{f}\right|_{(s, e]}\right\|^{2}>\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{b}\right\rangle^{2}+\left\langle\mathbf{f}, \boldsymbol{\gamma}_{(s, e]}\right\rangle^{2}+\left\langle\mathbf{f}, \mathbf{1}_{(s, e]}\right\rangle^{2} . \tag{17}
\end{equation*}
$$

Comparing (17) with (16) entails that $\left|\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{\tau_{j}}\right\rangle\right|>\left|\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{b}\right\rangle\right|$ for any $b \neq \tau_{j}$.
Secondly, set $\eta_{L}=\tau_{j}-s-1$ and $\eta_{R}=e-\tau_{j}$. After some calculation, we get that

$$
\mathcal{C}_{(s, e]}^{\tau_{j}}(\mathbf{f})=\left\{\frac{\eta_{L}\left(\eta_{L}+1\right) \eta_{R}\left(\eta_{R}+1\right)\left(2 \eta_{L} \eta_{R}+\eta_{L}+\eta_{R}+2\right)}{6 l\left(l^{2}-1\right)}\right\} \Delta_{j}^{\mathbf{f}}
$$

where $l=e-s$. Also, we have $\eta_{L} \geq \eta, \eta_{R} \geq \eta$ and $l=\eta_{L}+\eta_{R}+1$. To prove the lower bound, we observe that

$$
\begin{aligned}
& \left\{\frac{\eta_{L}\left(\eta_{L}+1\right) \eta_{R}\left(\eta_{R}+1\right)\left(2 \eta_{L} \eta_{R}+\eta_{L}+\eta_{R}+2\right)}{6 l\left(l^{2}-1\right)}\right\} \\
& \geq\left\{\frac{1}{6} \frac{\left(\eta_{L}+1\right) \eta_{R}}{l} \frac{\eta_{L}\left(\eta_{R}+1\right)}{l} \frac{2 \min \left(\eta_{L}, \eta_{R}\right)\left\{\max \left(\eta_{L}, \eta_{R}\right)+1\right\}}{l}\right\} \geq\left\{\frac{\eta^{3}}{24}\right\}
\end{aligned}
$$

where the last inequality is obtained applying Lemma 1 three times. For the upper bound, we notice that $2 \eta_{L} \eta_{R}+\eta_{L}+\eta_{R}+2 \leq 2\left(\eta_{L}+1\right)\left(\eta_{R}+1\right)$ which implies

$$
\left\{\frac{\eta_{L}\left(\eta_{L}+1\right) \eta_{R}\left(\eta_{R}+1\right)\left(2 \eta_{L} \eta_{R}+\eta_{L}+\eta_{R}+2\right)}{6 l\left(l^{2}-1\right)}\right\} \leq\left\{\frac{1}{3} \frac{\eta_{L} \eta_{R}\left(\eta_{L}+1\right)^{2}\left(\eta_{R}+1\right)^{2}}{(l-1) l^{2}}\right\} \leq\left\{\frac{(\eta+1)^{3}}{3}\right\}
$$

Lemma 6. Suppose $\mathbf{f}=\left(f_{1}, \ldots, f_{T}\right)^{\prime}$ is piecewise-linear vector as in Scenario (S2), and $\tau_{1}, \ldots, \tau_{q}$ are the locations of the change-points. Suppose $0 \leq s<e \leq T$ such that $\tau_{j-1} \leq$ $s+1 \leq \tau_{j}$ and $\tau_{j+1} \leq e \leq \tau_{j+2}$ for some $j=1 \ldots, q-1$. Then,

$$
\max _{s+1<b<e} \mathcal{C}_{(s, e]}^{b}(\mathbf{f}) \leq \frac{1}{\sqrt{3}}\left(\tau_{j}-s\right)^{3 / 2} \Delta_{j}^{\mathbf{f}}+\frac{1}{\sqrt{3}}\left(e-\tau_{j+1}+1\right)^{3 / 2} \Delta_{j+1}^{\mathbf{f}}
$$

and

$$
\max _{s+1<b<e} \mathcal{C}_{(s, e]}^{b}(\mathbf{f}) \leq\left(\tau_{j}-s-1\right)^{3 / 2} \Delta_{j}^{\mathrm{f}}+\left(e-\tau_{j+1}\right)^{3 / 2} \Delta_{j+1}^{\mathrm{f}},
$$

where $\Delta_{j}^{\mathbf{f}}=\left|2 f_{\tau_{j}}-f_{\tau_{j}-1}-f_{\tau_{j}+1}\right|$.
Proof. Suppose that $b^{*}=\operatorname{argmax}_{s \leq b \leq e} \mathcal{C}_{(s, e]}^{b}(\mathbf{f})$. Then

$$
\begin{aligned}
0 \leq & \left\|\left.\mathbf{f}\right|_{(s, e]}-\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{b^{*}}\right\rangle \boldsymbol{\phi}_{(s, e]}^{b^{*}}-\left\langle\mathbf{f}, \boldsymbol{\gamma}_{(s, e]}\right\rangle \gamma_{(s, e]}-\left\langle\mathbf{f}, \mathbf{1}_{(s, e]}\right\rangle \mathbf{1}_{(s, e]}\right\|^{2} \\
= & \left\|\left.\mathbf{f}\right|_{(s, e]}-\left\langle\mathbf{f}, \boldsymbol{\gamma}_{(s, e]}\right\rangle \gamma_{(s, e]}-\left\langle\mathbf{f}, \mathbf{1}_{(s, e]}\right\rangle \mathbf{1}_{(s, e]}\right\|^{2}-\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{b^{*}}\right\rangle^{2} \\
= & \frac{1}{6}\left(\tau_{j}-s-1\right)\left(\tau_{j}-s\right)\left(2 \tau_{j}-2 s-1\right)\left(\Delta_{j}^{\mathbf{f}}\right)^{2}+\frac{1}{6}\left(e-\tau_{j+1}\right)\left(e-\tau_{j+1}+1\right)\left(2 e-2 \tau_{j+1}+1\right)\left(\Delta_{j+1}^{\mathbf{f}}\right)^{2} \\
& -\left(\max _{s+1<b<e} \mathcal{C}_{(s, e]}^{b}(\mathbf{f})\right)^{2} .
\end{aligned}
$$

It then follows that

$$
\begin{aligned}
\max _{s+1<b<e} \mathcal{C}_{(s, e]}^{b}(\mathbf{f}) & \leq\left\{\left(\tau_{j}-s\right)^{3}\left(\Delta_{j}^{\mathrm{f}}\right)^{2} / 3+\left(e-\tau_{j+1}+1\right)^{3}\left(\Delta_{j+1}^{\mathrm{f}}\right)^{2} / 3\right\}^{1 / 2} \\
& \leq \frac{1}{\sqrt{3}}\left(\tau_{j}-s\right)^{3 / 2} \Delta_{j}^{\mathrm{f}}+\frac{1}{\sqrt{3}}\left(e-\tau_{j+1}+1\right)^{3 / 2} \Delta_{j+1}^{\mathrm{f}} .
\end{aligned}
$$

For the second claim, we note that $\left(\tau_{j}-s\right)\left(2 \tau_{j}-2 s-1\right) \leq 6\left(\tau_{j}-s-1\right)^{2}$ and $\left(e-\tau_{j+1}+\right.$ 1) $\left(2 e-2 \tau_{j+1}+1\right) \leq 6\left(e-\tau_{j+1}\right)^{2}$, so

$$
\begin{aligned}
\max _{s+1<b<e} \mathcal{C}_{(s, e]}^{b}(\mathbf{f}) & \leq\left\{\left(\tau_{j}-s-1\right)^{3}\left(\Delta_{j}^{\mathbf{f}}\right)^{2}+\left(e-\tau_{j+1}\right)^{3}\left(\Delta_{j+1}^{\mathbf{f}}\right)^{2}\right\}^{1 / 2} \\
& \leq\left(\tau_{j}-s-1\right)^{3 / 2} \Delta_{j}^{\mathbf{f}}+\left(e-\tau_{j+1}\right)^{3 / 2} \Delta_{j+1}^{\mathbf{f}} .
\end{aligned}
$$

Lemma 7. Suppose $\mathbf{f}=\left(f_{1}, \ldots, f_{T}\right)^{\prime}$ is piecewise-linear vector as in Scenario (S2), and $\tau_{1}, \ldots, \tau_{q}$ are the locations of the change-points. Suppose $0 \leq s<e \leq T$, such that $\tau_{j-1} \leq s+1<\tau_{j}<e \leq \tau_{j+1}$ for some $j=1 \ldots, q$. Let $\rho=\left|\tau_{j}-b\right|, \eta_{L}=\tau_{j}-s-1$, $\eta_{R}=e-\tau_{j}$ and $\Delta_{j}^{\mathrm{f}}=\left|2 f_{\tau_{j}}-f_{\tau_{j}-1}-f_{\tau_{j}+1}\right|$. Then,

$$
\begin{equation*}
\left\|\boldsymbol{\phi}_{(s, e]}^{b}\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{b}\right\rangle-\boldsymbol{\phi}_{(s, e]}^{\tau_{j}}\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{\tau_{j}}\right\rangle\right\|_{2}^{2}=\left(\mathcal{C}_{(s, e]}^{\tau_{j}}(\mathbf{f})\right)^{2}-\left(\mathcal{C}_{(s, e]}^{b}(\mathbf{f})\right)^{2} . \tag{18}
\end{equation*}
$$

Moreover,
(a) for any $\tau_{j} \leq b<e,\left(\mathcal{C}_{(s, e]}^{\tau_{j}}(\mathbf{f})\right)^{2}-\left(\mathcal{C}_{(s, e]}^{b}(\mathbf{f})\right)^{2} \geq \frac{1}{63} \min \left(\rho, \eta_{L}\right)^{3}\left(\Delta_{j}^{\mathbf{f}}\right)^{2}$;
(b) for any $s+1<b \leq \tau_{j}$, $\left(\mathcal{C}_{(s, e]}^{\tau_{j}}(\mathbf{f})\right)^{2}-\left(\mathcal{C}_{(s, e]}^{b}(\mathbf{f})\right)^{2} \geq \frac{1}{63} \min \left(\rho, \eta_{R}\right)^{3}\left(\Delta_{j}^{\mathbf{f}}\right)^{2}$.

Proof. The proof of (18) is very similar to that shown in Lemma 4, so is omitted for brevity. In the following, we only deal with the case of $\tau_{j} \leq b<e$. Note that

$$
\begin{aligned}
& \left\|\boldsymbol{\phi}_{(s, e]}^{b}\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{b}\right\rangle-\boldsymbol{\phi}_{(s, e]}^{\tau_{j}}\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{\tau_{j}}\right\rangle\right\|_{2}^{2} \\
& \left.=\| \boldsymbol{\phi}_{(s, e]}^{b}\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{b}\right\rangle+\gamma_{(s, e]} \mathbf{f}, \boldsymbol{\gamma}_{(s, e]}\right\rangle+\mathbf{1}_{(s, e]}\left\langle\mathbf{f}, \mathbf{1}_{(s, e]}\right\rangle-\left.\mathbf{f}\right|_{(s, e]} \|_{2}^{2} \\
& \geq \min _{a_{0}, a_{1} \in \mathbb{R}}\left\|\left.\mathbf{f}\right|_{(s, b]}-a_{0} \mathbf{1}_{(s, b]}-a_{1} \boldsymbol{\gamma}_{(s, b]}\right\|_{2}^{2}+\min _{a_{0}, a_{1} \in \mathbb{R}}\left\|\left.\mathbf{f}\right|_{(b, e]}-a_{0} \mathbf{1}_{(b, e]}-a_{1} \boldsymbol{\gamma}_{(b, e]}\right\|_{2}^{2} \\
& \geq \min _{a_{0}, a_{1} \in \mathbb{R}}\left\|\left.\mathbf{f}\right|_{(s, b]}-a_{0} \mathbf{1}_{(s, b]}-a_{1} \boldsymbol{\gamma}_{(s, b]}\right\|_{2}^{2} .
\end{aligned}
$$

Recalling the definitions of $\alpha_{(s, b]}^{\tau_{j}}$ and $\beta_{(s, b]}^{\tau_{j}}$ in (7), and writing $d=b-s$. After some calculations (similar to what has already been carried out in deriving $\phi_{(s, e]}^{b}$, as demonstrated in Section B), we obtain that

$$
\begin{aligned}
& \min _{a_{0}, a_{1} \in \mathbb{R}}\left\|\left.\mathbf{f}\right|_{(s, b]}-a_{0} \mathbf{1}_{(s, b]}-a_{1} \gamma_{(s, b]}\right\|_{2}^{2} \\
& =\left[\left(3 \eta_{L}+\rho+2\right) \alpha_{(s, b]}^{\tau_{j}} \beta_{(s, b]}^{\tau_{j}}+\left(3 \rho+\eta_{L}+2\right) \alpha_{(s, b]}^{\tau_{j}}\left(\beta_{(s, b]}^{\tau_{j}}\right)^{-1}\right]^{-2}\left(\Delta_{j}^{\mathbf{f}}\right)^{2} \\
& =\frac{1}{6}\left(\Delta_{j}^{\mathbf{f}}\right)^{2} d\left(d^{2}-1\right)\left[1+\rho \eta_{L}+(\rho+1)\left(\eta_{L}+1\right)\right] \times \\
& \quad\left[\left(d+2 \eta_{L}+1\right)^{2} \frac{\rho(\rho+1)}{\eta_{L}\left(\eta_{L}+1\right)}+(d+2 \rho+1)^{2} \frac{\eta_{L}\left(\eta_{L}+1\right)}{\rho(\rho+1)}+2\left(d+2 \eta_{L}+1\right)(d+2 \rho+1)\right]^{-1}
\end{aligned}
$$

Notice that the above equation is symmetric with respect to $\eta_{L}$ and $\rho$. Without loss of generality, here we proceed by assuming that $\eta_{L} \geq \rho$. Since $\left(d+2 \eta_{L}+1\right)+(d+2 \rho+1)=4 d$, it follows that $\left(d+2 \eta_{L}+1\right)(d+2 \rho+1) \leq 4 d^{2}$. Therefore,

$$
\begin{aligned}
& \left.\min _{a_{0}, a_{1} \in \mathbb{R}}| | \mathbf{f}\right|_{(s, b]}-a_{0} \mathbf{1}_{(s, b]}-a_{1} \gamma_{(s, b]} \|_{2}^{2} \\
& \geq \frac{1}{6}\left(\Delta_{j}^{\mathbf{f}}\right)^{2} d\left(d^{2}-1\right)\left[2\left(\eta_{L}+1\right) \rho\right]\left[(3 d)^{2}+(2 d)^{2} \frac{\left(\eta_{L}+1\right)^{2}}{\rho^{2}}+8 d^{2}\right]^{-1} \\
& \geq \frac{1}{6}\left(\Delta_{j}^{\mathbf{f}}\right)^{2} d^{2}(d-1)\left[2\left(\eta_{L}+1\right) \rho\right]\left[21 d^{2} \frac{\left(\eta_{L}+1\right)^{2}}{\rho^{2}}\right]^{-1} \geq \frac{1}{63} \rho^{3}\left(\Delta_{j}^{\mathbf{f}}\right)^{2}
\end{aligned}
$$

where in the last step, we used the fact that $\frac{d-1}{\eta_{L}+1} \geq 1$ for $\rho \geq 1$ (and note that the last above-displayed equation also holds if $\rho=0$ ).

Finally, we remark that the case of $s+1<b \leq \tau_{j}$ can be handled in a similar manner by symmetry.

Lemma 8. Suppose $\mathbf{f}=\left(f_{1}, \ldots, f_{T}\right)^{\prime}$ is piecewise-linear vector as in Scenario (S2), and $\tau_{1}, \ldots, \tau_{q}$ are the locations of the change-points. Suppose $0 \leq s<e \leq T$, such that $\tau_{j-1} \leq s+1<\tau_{j}<e \leq \tau_{j+1}$ for some $j=1 \ldots, q$. Let $\rho=\left|\tau_{j}-b\right|, \eta_{L}=\tau_{j}-s-1$, $\eta_{R}=e-\tau_{j}$ and $\Delta_{j}^{\mathrm{f}}=\left|2 f_{\tau_{j}}-f_{\tau_{j}-1}-f_{\tau_{j}+1}\right|$. Then, for any b satisfying $\tau_{j}-\min \left(\eta_{L}, \eta_{R}\right) / 2<$ $b<\tau_{j}+\min \left(\eta_{L}, \eta_{R}\right) / 2$, we have that

$$
\left(\mathcal{C}_{(s, e]}^{\tau_{j}}(\mathbf{f})\right)^{2}-\left(\mathcal{C}_{(s, e]}^{b}(\mathbf{f})\right)^{2} \geq \frac{\left(\Delta_{j}^{\mathbf{f}}\right)^{2}}{48}\left\{\min \left(\eta_{L}, \eta_{R}\right)-1\right\} \rho^{2}
$$

Proof. Here we focus on the scenario where $b>\tau_{j}$. By Lemma 7,

$$
\begin{aligned}
\left(\mathcal{C}_{(s, e]}^{\tau_{j}}(\mathbf{f})\right)^{2}-\left(\mathcal{C}_{(s, e]}^{b}(\mathbf{f})\right)^{2} & =\left\|\boldsymbol{\phi}_{s, e}^{b}\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{b}\right\rangle-\boldsymbol{\phi}_{(s, e]}^{\tau_{j}}\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{\tau_{j}}\right\rangle\right\|_{2}^{2} \\
& =\min _{a_{0}, a_{1}, a_{2} \in \mathbb{R}}\left\|\left.\mathbf{f}\right|_{(s, e]}-a_{0} \mathbf{1}_{(s, e]}-a_{1} \gamma_{(s, e]}-a_{2} \boldsymbol{\phi}_{(s, e]}^{b}\right\|_{2}^{2} \\
& =\left(\Delta_{j}^{\mathbf{f}}\right)^{2} \min _{a_{0}, a_{1}, a_{2} \in \mathbb{R}}\left\|\left.\tilde{\mathbf{f}}\right|_{(s, e]}-a_{0} \mathbf{1}_{(s, e]}-a_{1} \gamma_{(s, e]}-a_{2} \boldsymbol{\phi}_{(s, e]}^{b}\right\|_{2}^{2},
\end{aligned}
$$

where $\left.\tilde{\mathbf{f}}\right|_{(s, e]}:=\left(0, \ldots, 0,1, \ldots, e-\tau_{j}, 0, \ldots, 0\right)^{\prime}$, in which " 1 " appears at the $\left(\tau_{j}+1\right)$-th position. In the following, our aim is to bound the residual sum of squares resulted from fitting $\left.\tilde{\mathbf{f}}\right|_{(s, e]}$ via a piecewise-linear and continuous function with only one kink at $b$ over
$(s, e]$. Assuming that the fitted value of this vector at the $b$-th position is $m$, then, we have that

$$
\begin{aligned}
& \min _{a_{0}, a_{1}, a_{2} \in \mathbb{R}}\left\|\left.\tilde{\mathbf{f}}\right|_{(s, e]}-a_{0} \mathbf{1}_{(s, e]}-a_{1} \boldsymbol{\gamma}_{(s, e]}-a_{2} \boldsymbol{\phi}_{(s, e]}^{b}\right\|_{2}^{2} \\
& \quad \geq\left(\frac{2 m}{\eta_{L}+2 \rho}\right)^{2} \times \frac{1}{6}\left(\frac{\eta_{L}-1}{2}\right)\left(\frac{\eta_{L}+1}{2}\right) \eta_{L}+\left\{\frac{2(\rho-m)}{e-b}\right\}^{2} \times \frac{1}{6}\left(\frac{e-b-1}{2}\right)\left(\frac{e-b+1}{2}\right)(e-b) .
\end{aligned}
$$

Since $b<\tau_{j}+\eta_{R} / 2$, it follows that $e-b>\eta_{R} / 2$, and thus $e-b-1 \geq\left(\eta_{R}-1\right) / 2$. Moreover, the fact of $\rho<\min \left(\eta_{L}, \eta_{R}\right) / 2$ yields $\eta_{L}+2 \rho \leq 2 \eta_{L}$. Plugging these two inequalities into the previous equation, we have that

$$
\begin{aligned}
& \min _{a_{0}, a_{1}, a_{2} \in \mathbb{R}}|\tilde{\mathbf{f}}|_{(s, e]}-a_{0} \mathbf{1}_{(s, e]}-a_{1} \boldsymbol{\gamma}_{(s, e]}-a_{2} \phi_{(s, e]}^{b} \|_{2}^{2} \\
& \quad \geq m^{2} \frac{\eta_{L}-1}{24}+(\rho-m)^{2} \frac{\eta_{R}-1}{12} \geq \frac{1}{2} \min \left(\frac{\eta_{L}-1}{24}, \frac{\eta_{R}-1}{12}\right) \rho^{2}
\end{aligned}
$$

Consequently,

$$
\left(\mathcal{C}_{(s, e]}^{\tau_{j}}(\mathbf{f})\right)^{2}-\left(\mathcal{C}_{(s, e]}^{b}(\mathbf{f})\right)^{2} \geq \frac{\left(\Delta_{j}^{\mathbf{f}}\right)^{2}}{48}\left\{\min \left(\eta_{L}, \eta_{R}\right)-1\right\} \rho^{2} .
$$

By symmetry, the scenario of $b<\tau_{j}$ can be dealt with in a similar fashion. Finally, we remark that the constants here are not sharp, as we will only use this lemma to establish rate-type results later.

## I.2. Proof of Theorem 1

Here we informally discuss our proof strategy, which could be generalised to other scenarios.

- Intuitively speaking, lemmas from Appendix I. 1 deal with noiseless versions of the change-point estimation problems. In order to apply these results to show the consistency of estimated number of change-points, we need to control $\left\|\mathcal{C}_{(s, e]}^{b}(\mathbf{Y})-\mathcal{C}_{(s, e]}^{b}(\mathbf{f})\right\|$ for every tuple ( $s, e, b$ ), which can be achieved using Bonferroni in Step One.
- Note that for any fixed left-open and right-closed interval with start-point $s$ and endpoint $e$, to decide whether $b_{1}$ or $b_{2}$ is a more suitable change-point candidate inside this interval, we only need to look at the value of $\mathcal{C}_{(s, e]}^{b_{1}}(\mathbf{Y})-\mathcal{C}_{(s, e]}^{b_{2}}(\mathbf{Y})$. Therefore, when establishing the convergence rate of the estimated change-point location, we control the distance between $\mathcal{C}_{(s, e]}^{b_{1}}(\mathbf{Y})-\mathcal{C}_{(s, e]}^{b_{2}}(\mathbf{Y})$ and its noiseless analogue $\mathcal{C}_{(s, e]}^{b_{1}}(\mathbf{f})-\mathcal{C}_{(s, e]}^{b_{2}}(\mathbf{f})$ (after proper normalisation) for all tuples ( $s, e, b_{1}, b_{2}$ ) in Step Two.
- In Step Three, we show that given a properly chosen threshold and a large enough $M$, both bounds in Step One and Step Two hold, and for each change-point $\tau_{j}$, there exists an interval from $F_{T}^{M}$ that contains only this change-point and both its start- and endpoints are sufficiently far away from other change-points. Since we are dealing with the narrowest-over-threshold intervals, the actual intervals that our NOT algorithm pick must have length no longer than the ones we considered in Step Three, thus could only contain precisely one change-point.
- So in Step Four, it suffices to investigate a single change-point detection problem, where we can use lemmas from Appendix I. 1 and the bound in Step Two to establish the convergence rate for its location estimation.
- Finally, in Step Five, we show that after detecting all the change-points, the NOT algorithm stops with no further detection. This is because the remaining elements $(s, e] \in F_{T}^{M}$ to be considered either have no change-point inside, or have one/two change-points that are very close to its start- or/and end- points, thus their corresponding $\max _{b} \mathcal{C}_{(s, e]}^{b}(\mathbf{Y})$ cannot exceed the given threshold in views of the property of its noiseless analogue and the bound from Step One.

Now we proceed to the technical details.
Proof. We shall prove the following more specific result, which in turn implies (9).

$$
\begin{equation*}
\mathbb{P}\left(\hat{q}=q, \max _{j=1, \ldots, q}\left(\left|\hat{\tau}_{j}-\tau_{j}\right|\left(\Delta_{j}^{\mathbf{f}}\right)^{2}\right) \leq C_{3} \log T\right) \geq 1-T^{-1} /(6 \sqrt{\pi})-T \delta_{T}^{-1}\left(1-\delta_{T}^{2} T^{-2} / 36\right)^{M} \tag{19}
\end{equation*}
$$

Step One.
Let $\boldsymbol{\varepsilon}=\left(\varepsilon_{1}, \ldots, \varepsilon_{T}\right)^{\prime}$ and $\lambda_{T}=\sqrt{8 \log T}$. Define the set

$$
A_{T}=\left\{\max _{s, b, e: 0 \leq s<b<e \leq T}\left|\mathcal{C}_{(s, e]}^{b}(\varepsilon)\right| \leq \lambda_{T}\right\}
$$

Note that for any $0 \leq s<b<e \leq T, \mathcal{C}_{(s, e]}^{b}(\varepsilon)$ follows a standard normal distribution. Therefore, using the Bonferroni bound, we get

$$
\mathbb{P}\left(A_{T}^{c}\right) \leq \frac{T^{3}}{6} \frac{2 e^{-(\sqrt{8 \log T})^{2} / 2}}{\sqrt{8 \log T} \sqrt{2 \pi}} \leq \frac{T^{-1}}{12 \sqrt{\pi}}
$$

Moreover, because $\mathcal{C}_{(s, e]}^{b}(\mathbf{Y})-\mathcal{C}_{(s, e]}^{b}(\mathbf{f})=\mathcal{C}_{(s, e]}^{b}(\varepsilon)$, so $A_{T}$ also implies that

$$
\left\{\max _{s, b, e: 0 \leq s<b<e \leq T}\left|\mathcal{C}_{(s, e]}^{b}(\mathbf{Y})-\mathcal{C}_{(s, e]}^{b}(\mathbf{f})\right| \leq \lambda_{T}\right\}
$$

We remark that though the constant in $\lambda_{T}$ (i.e. $\sqrt{8}$ ) does not appear sharp (as it is rooted in the simple Bonferroni bound), it is sufficient for our purpose of establishing consistency and rate-type results later. We refer the readers to Dümbgen and Spokoiny (2001) and Rufibach and Walther (2010) for possible improvement over this constant.

Step Two.
Define the set

$$
B_{T}=\left\{\max _{j=1, \ldots, q} \max _{\substack{\tau_{j-1} \leq s<\tau_{j} \\ \tau_{j}<e \leq \tau_{j+1} \\ s<b<e}} \frac{\left|\left\langle\boldsymbol{\psi}_{(s, e]}^{b}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle-\boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\right\rangle, \boldsymbol{\varepsilon}\right\rangle\right|}{\left\|\boldsymbol{\psi}_{(s, e]}^{b}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle-\boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\right\rangle\right\|_{2}} \leq \lambda_{T}\right\}
$$

Again, for any $0 \leq s<b<e \leq T, \frac{\left|\left\langle\boldsymbol{\psi}_{(s, e]}^{b}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle-\boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\right\rangle, \boldsymbol{\varepsilon}\right\rangle\right|}{\| \boldsymbol{\psi}_{(s, e]}^{b}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle-\boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{\left.\boldsymbol{\tau}_{j}^{\tau_{j}}\right\rangle \|_{2}}\right.}$ follows a standard normal distribution, so using a similar argument, we get

$$
\mathbb{P}\left(B_{T}^{c}\right) \leq \frac{T^{-1}}{12 \sqrt{\pi}}
$$

Step Three.
To fix the ideas, for $j=1, \ldots, q$, we define intervals

$$
\begin{align*}
& \mathcal{I}_{j}^{L}=\left(\tau_{j}-\delta_{T} / 3-1, \tau_{j}-\delta_{T} / 6-1\right]  \tag{20}\\
& \mathcal{I}_{j}^{R}=\left(\tau_{j}+\delta_{T} / 6, \tau_{j}+\delta_{T} / 3\right] \tag{21}
\end{align*}
$$

Note that these intervals all contain at least one integer as long as $\delta_{T}>6$. This is always true for sufficiently large $T$, as it follows from Conditions 1 and 2 that $\delta_{T}>\underline{C} \log T / f$. Recall that $F_{T}^{M}$ is the set of $M$ randomly drawn intervals with pairs of endpoints in $\{0, \ldots, T-1\} \times\{1, \ldots, T\}$. Denote by $\left(s_{1}, e_{1}\right], \ldots,\left(s_{M}, e_{M}\right]$ the elements of $F_{T}^{M}$ and let

$$
\begin{equation*}
D_{T}^{M}=\left\{\forall j=1, \ldots, q, \exists k \in\{1, \ldots, M\}, \text { s.t. } s_{k} \in \mathcal{I}_{j}^{L} \text { and } e_{k} \in \mathcal{I}_{j}^{R}\right\} \tag{22}
\end{equation*}
$$

We have that

$$
\begin{aligned}
\mathbb{P}\left(\left(D_{T}^{M}\right)^{c}\right) & \leq \sum_{j=1}^{q} \Pi_{m=1}^{M}\left(1-\mathbb{P}\left(s_{m} \times e_{m} \in \mathcal{I}_{j}^{L} \times \mathcal{I}_{j}^{R}\right)\right) \\
& \leq q\left(1-\frac{\delta_{T}^{2}}{6^{2} T^{2}}\right)^{M} \leq \frac{T}{\delta_{T}}\left(1-\frac{\delta_{T}^{2}}{36 T^{2}}\right)^{M}
\end{aligned}
$$

Therefore, $\mathbb{P}\left(A_{T} \cap B_{T} \cap D_{T}^{M}\right) \geq 1-T^{-1} /(6 \sqrt{\pi})-T \delta_{T}^{-1}\left(1-\delta_{T}^{2} T^{-2} / 36\right)^{M}$.
In the rest of the proof, we assume that $A_{T}, B_{T}$ and $D_{T}^{M}$ all hold. We give the constants as follows:

$$
\underline{C}=\sqrt{6}\left(2 \sqrt{C_{3}}+4 \sqrt{2}\right)+1, \quad C_{1}=2 \sqrt{C_{3}}+2 \sqrt{2}, \quad C_{2}=\frac{1}{\sqrt{6}}-\frac{2 \sqrt{2}}{\underline{C}}, \quad C_{3}=32 \sqrt{2}+48 .
$$

These constants could be further refined by applying the Bonferroni bound more carefully. See also our remark at the end of Step One. But since our main aim is to establish the rate, we chose not to pursue this direction further. In addition, here we set $\underline{C}$ in such a way that $\underline{C} C_{2}>C_{1}$ (as well as $C_{2}>0$ ). This means that given $\delta_{T}^{1 / 2} \underline{f}_{T} \geq \underline{C} \sqrt{\log T}$, one have that $C_{2} \delta_{T}^{1 / 2} \underline{f}_{T}>C_{1} \sqrt{\log T}$, i.e. we can select $\zeta_{T} \in\left[C_{1} \sqrt{\log T}, C_{2} \delta_{T}^{1 / 2} \underline{f}_{T}\right)$.

## Step Four.

We focus on a generic interval $(s, e]$ such that

$$
\begin{equation*}
\exists j \in\{1, \ldots, q\}, \exists k \in\{1, \ldots, M\}, \text { s.t. }\left(s_{k}, e_{k}\right] \subset(s, e] \text { and } s_{k} \times e_{k} \in \mathcal{I}_{j}^{L} \times \mathcal{I}_{j}^{R} \tag{23}
\end{equation*}
$$

Fix such an interval $(s, e]$ and let $j \in\{1, \ldots, q\}$ and $k \in\{1, \ldots, M\}$ be such that (23) is satisfied. Let $b_{k}^{*}=\operatorname{argmax}_{s_{k}<b<e_{k}} \mathcal{C}_{\left(s_{k}, e_{k}\right]}^{b}(\mathbf{Y})$. By construction, $\left(s_{k}, e_{k}\right]$ satisfies $\tau_{j}-s_{k}>$
$\delta_{T} / 6$ and $e_{k}-\tau_{j}>\delta_{T} / 6$. Denote by

$$
\begin{aligned}
\mathcal{M}_{(s, e]} & =\left\{m:\left(s_{m}, e_{m}\right] \in F_{T}^{M},\left(s_{m}, e_{m}\right] \subset(s, e]\right\} \\
\mathcal{O}_{(s, e]} & =\left\{m \in \mathcal{M}_{(s, e]}: \max _{s_{m}<b<e_{m}} \mathcal{C}_{\left(s_{m}, e_{m}\right]}^{b}(\mathbf{Y})>\zeta_{T}\right\}
\end{aligned}
$$

Our first aim is to show that $\mathcal{O}_{(s, e]}$ is non-empty. This follows from Lemma 2 and the calculation below.

$$
\begin{aligned}
\mathcal{C}_{\left(s_{k}, e_{k}\right]}^{b_{k}^{*}}(\mathbf{Y}) & \geq \mathcal{C}_{\left(s_{k}, e_{k}\right]}^{\tau_{j}}(\mathbf{Y}) \\
& \geq \mathcal{C}_{\left(s_{k}, e_{k}\right]}^{b_{k}^{*}}(\mathbf{f})-\lambda_{T} \geq\left(\frac{\delta_{T}}{6}\right)^{1 / 2}\left|f_{\tau_{j}+1}-f_{\tau_{j}}\right|-\lambda_{T} \geq\left(\frac{\delta_{T}}{6}\right)^{1 / 2} \underline{f}_{T}-\lambda_{T} \\
& =\left(\frac{1}{\sqrt{6}}-\frac{\lambda_{T}}{\delta_{T}^{1 / 2} \underline{f}_{T}}\right) \delta_{T}^{1 / 2} \underline{f}_{T} \geq\left(\frac{1}{\sqrt{6}}-\frac{2 \sqrt{2}}{\underline{C}}\right) \delta_{T}^{1 / 2} \underline{f}_{T}=C_{2} \delta_{T}^{1 / 2} \underline{f}_{T}>\zeta_{T}
\end{aligned}
$$

Let $m^{*}=\operatorname{argmin}_{m \in \mathcal{O}_{s, e}}\left(e_{m}-s_{m}\right)$ and $b^{*}=\operatorname{argmax}_{s_{m^{*}}<b<e_{m^{*}}} \mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b}(\mathbf{Y})$. Observe that $\left(s_{m^{*}}, e_{m^{*}}\right)$ must contain at least one change-point. Indeed, if that was not the case, we would have $\mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b}(\mathbf{f})=0$ and

$$
\mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b^{*}}(\mathbf{Y})=\left|\mathcal{C}_{\left(s_{\left.m^{*}, e_{m^{*}}\right]}^{b^{*}}\right.}(\mathbf{Y})-\mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b^{*}}(\mathbf{f})\right| \leq \lambda_{T} \leq \zeta_{T}
$$

which contradicts $\mathcal{C}_{\left(s_{m^{*}, e_{m}}\right]}^{b^{*}}(\mathbf{Y})>\zeta_{T}$. On the other hand, $\left[s_{m^{*}}, e_{m^{*}}\right)$ cannot contain more than one change-points, because $e_{m^{*}}-s_{m^{*}} \leq e_{k}-s_{k} \leq \delta_{T}$, as we picked the narrowest-over-threshold interval.

Without loss of generality, assume $\tau_{j} \in\left(s_{m^{*}}, e_{m^{*}}\right)$. Denote by $\eta_{L}=\tau_{j}-s_{m^{*}}, \eta_{R}=$ $e_{m^{*}}-\tau_{j}$ and $\eta_{T}=\left(C_{1}-\sqrt{8}\right)^{2}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2} \log T$, where $\Delta_{j}^{\mathbf{f}}=\left|f_{\tau_{j}+1}-f_{\tau_{j}}\right|$. We claim that $\min \left(\eta_{L}, \eta_{R}\right)>\eta_{T}$, because otherwise $\min \left(\eta_{L}, \eta_{R}\right) \leq \eta_{T}$ and Lemma 2 would result in

$$
\begin{aligned}
\mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b^{*}}(\mathbf{Y}) & \leq \mathcal{C}_{\left(s_{\left.m^{*}, e_{m^{*}}\right]}^{b^{*}}\right.}(\mathbf{f})+\lambda_{T} \leq \mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{\tau_{j}}(\mathbf{f})+\lambda_{T} \leq \eta_{T}^{1 / 2} \Delta_{j}^{\mathbf{f}}+\lambda_{T} \\
& =\left(C_{1}-\sqrt{8}+\sqrt{8}\right) \sqrt{\log T}=C_{1} \sqrt{\log T} \leq \zeta_{T}
\end{aligned}
$$

which would contradict $\mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b^{*}}(\mathbf{Y})>\zeta_{T}$.
We are now in the position to prove $\left|b^{*}-\tau_{j}\right| \leq C_{3} \log T /\left(\Delta_{j}^{\mathbf{f}}\right)^{2}$. The arguments we use here are simpler and slightly more general than Lemma A. 3 of Fryzlewicz (2014). Our aim is to find $\epsilon_{T}$ such that for any $b \in\left\{s_{m^{*}}+1, \ldots, e_{m^{*}}-1\right\}$ with $\left|b-\tau_{j}\right|>\epsilon_{T}$, we always have

$$
\begin{equation*}
\left(\mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{\tau_{j}}(\mathbf{Y})\right)^{2}-\left(\mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b}(\mathbf{Y})\right)^{2}>0 \tag{24}
\end{equation*}
$$

This would then imply that $\left|b^{*}-\tau_{j}\right| \leq \epsilon_{T}$. By expansion and rearranging the terms (using the fact that $\left.f_{t}=Y_{t}+\varepsilon_{t}\right)$, we see that $(24)$ is equivalent to

$$
\begin{align*}
\left\langle\mathbf{f}, \boldsymbol{\psi}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{\tau_{j}}\right)^{2}-\left\langle\mathbf{f}, \boldsymbol{\psi}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b}\right\rangle^{2}> & \left\langle\boldsymbol{\varepsilon}, \boldsymbol{\psi}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b}\right\rangle^{2}-\left\langle\boldsymbol{\varepsilon}, \boldsymbol{\psi}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{\tau_{j}}\right\rangle^{2} \\
& +2\left\langle\varepsilon, \boldsymbol{\psi}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b}\left\langle\mathbf{f}, \boldsymbol{\psi}_{\left(s_{m^{*},}, e_{m^{*}}\right]}^{b}\right\rangle-\boldsymbol{\psi}_{\left(s_{\left.m^{*}, e_{m^{*}}\right]}^{\tau_{j}}\right.}^{\tau_{j}}\left\langle\mathbf{f}, \boldsymbol{\psi}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{\tau_{j}}\right\rangle\right\rangle . \tag{25}
\end{align*}
$$

In the following, we assume that $b \geq \tau_{j}$. The case that $b<\tau_{j}$ can be handled in a similar fashion. By Lemma 4, we have

$$
\left\langle\mathbf{f}, \boldsymbol{\psi}_{\left(s_{\left.m^{*}, e_{m^{*}}\right]}^{\tau_{j}}\right.}^{\tau_{j}}\right\rangle^{2}-\left\langle\mathbf{f}, \boldsymbol{\psi}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b}\right\rangle^{2}=\left(\mathcal{C}_{\left(s_{\left.m^{*}, e_{m^{*}}\right]}^{\tau_{j}}\right.}^{b}(\mathbf{f})\right)^{2}-\left(\mathcal{C}_{\left(s_{\left.m^{*}, e_{m^{*}}\right]}^{b}\right.}^{b}(\mathbf{f})\right)^{2}=\frac{\left|b-\tau_{j}\right| \eta_{L}}{\left|b-\tau_{j}\right|+\eta_{L}}\left(\Delta_{j}^{\mathbf{f}}\right)^{2}:=\kappa
$$

In addition, since $A_{T}$ and $B_{T}$ hold, we have that

$$
\begin{aligned}
& \left\langle\varepsilon, \boldsymbol{\psi}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b}\right\rangle^{2}-\left\langle\varepsilon, \boldsymbol{\psi}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{\tau_{j}}\right\rangle^{2} \leq \lambda_{T}^{2}, \\
& 2\left\langle\varepsilon, \boldsymbol{\psi}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b}\left\langle\mathbf{f}, \boldsymbol{\psi}_{\left(s_{\left.m^{*}, e_{m^{*}}\right]}^{b}\right.}^{b}\right\rangle-\boldsymbol{\psi}_{\left(s_{\left.m^{*}, e_{m^{*}}\right]}^{\tau_{j}}\right.}^{\tau_{j}}\left\langle\mathbf{f}, \boldsymbol{\psi}_{\left.\left(s_{\left.m^{*}, e_{m^{*}}\right]}^{\tau_{j}}\right\rangle\right\rangle}\right\rangle\right. \\
& \leq 2\left\|\boldsymbol{\psi}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b}\left\langle\mathbf{f}, \boldsymbol{\psi}_{\left(s_{\left.m^{*}, e_{m^{*}}\right]}^{b}\right.}^{b}\right\rangle-\boldsymbol{\psi}_{\left(s_{\left.m^{*}, e_{m^{*}}\right]}^{\tau_{j}}\right.}^{\tau_{j}}\left\langle\boldsymbol{f}, \boldsymbol{\psi}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{\tau_{j}}\right\rangle\right\|_{2} \lambda_{T}=2 \kappa^{1 / 2} \lambda_{T},
\end{aligned}
$$

where the last equality also comes from Lemma 4. Consequently, (25) can be deducted from the stronger inequality $\kappa-2 \lambda_{T} \kappa^{1 / 2}-\lambda_{T}^{2}>0$. This quadratic inequality is implied by $\kappa>(\sqrt{2}+1)^{2} \lambda_{T}^{2}$, and could be restricted further to

$$
\begin{equation*}
\frac{2\left|b-\tau_{j}\right| \eta_{L}}{\left|b-\tau_{j}\right|+\eta_{L}} \geq \min \left(\left|b-\tau_{j}\right|, \eta_{L}\right)>(32 \sqrt{2}+48)\left(\Delta_{j}^{\mathbf{f}}\right)^{-2} \log T=C_{3}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2} \log T \tag{26}
\end{equation*}
$$

But since

$$
\eta_{L} \geq \eta_{T}=\left(C_{1}-\sqrt{8}\right)^{2}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2} \log T=\left(2 \sqrt{C_{3}}\right)^{2}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2} \log T>C_{3}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2} \log T
$$

we see that the inequality in (26) is equivalent to $\left|b-\tau_{j}\right|>C_{3}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2} \log T$. To sum up, $\left|b^{*}-\tau_{j}\right|\left(\Delta_{j}^{\mathbf{f}}\right)^{2}>C_{3} \log T$ would result in (24), a contradiction. So we have proved that $\left|b^{*}-\tau_{j}\right|\left(\Delta_{j}^{\mathbf{f}}\right)^{2} \leq C_{3} \log T$.

## Step Five.

Using the arguments given above which are valid on the event $A_{T} \cap B_{T} \cap D_{T}^{M}$, we can now proceed with the proof of the theorem as follows. At the start of Algorithm 1 we have $s=0$ and $e=T$ and, provided that $q \geq 1$, condition (23) is satisfied. Therefore the algorithm detects a change-point $b^{*}$ in that interval such that $\left|b^{*}-\tau_{j}\right| \leq C_{3} \log T\left(\Delta_{j}^{\mathbf{f}}\right)^{-2}$ for some $j$. By construction, we also have that $\left|b^{*}-\tau_{j}\right|<2 / 3 \delta_{T}$. This in turn implies that for all $l=1, \ldots, q$ such that $\tau_{l} \in(s, e)$ and $l \neq j$ we have either $\mathcal{I}_{l}^{L}, \mathcal{I}_{l}^{R} \subset\left(s, b^{*}\right]$ or $\mathcal{I}_{l}^{L}, \mathcal{I}_{l}^{R} \subset\left(b^{*}, e\right]$. Therefore $(23)$ is satisfied within each segment containing at least one change-point. Note that before all $q$ change-points are detected, each change-point will not be detected twice. To see this, we suppose that $\tau_{j}$ has already been detected by $b^{*}$, then for all intervals $\left(s_{k}, e_{k}\right] \subset\left(\tau_{j}-C_{3} \log T\left(\Delta_{j}^{\mathbf{f}}\right)^{-2}, \tau_{j}+2 / 3 \delta_{T}\right] \cup\left(\tau_{j}-2 / 3 \delta_{T}, \tau_{j}+C_{3} \log T\left(\Delta_{j}^{\mathbf{f}}\right)^{-2}\right]$, Lemma 2, together with the event $A_{T}$, guarantees that

$$
\max _{s_{k}<b<e_{k}} \mathcal{C}_{\left(s_{k}, e_{k}\right]}^{b}(\mathbf{Y}) \leq \max _{s<b<e} \mathcal{C}_{\left(s_{k}, e_{k}\right]}^{b}(\mathbf{f})+\lambda_{T} \leq \sqrt{C_{3} \log T\left(\Delta_{j}^{\mathbf{f}}\right)^{-2}} \Delta_{j}^{\mathbf{f}}+\lambda_{T} \leq C_{1} \sqrt{\log T} \leq \zeta_{T}
$$

Once all the change-points are detected, we then only need to consider $\left(s_{k}, e_{k}\right]$ such that

$$
\left(s_{k}, e_{k}\right] \subset\left(\tau_{j}-C_{3} \log T\left(\Delta_{j}^{\mathbf{f}}\right)^{-2}, \tau_{j+1}+C_{3} \log T\left(\Delta_{j+1}^{f}\right)^{-2}\right]
$$

for $j=0, \ldots, q$, where we set $\Delta_{0}^{f}=\Delta_{q+1}^{f}=\infty$ for notational convenience. It follows from Lemma 3 (within the event $A_{T}$ ) that

$$
\begin{aligned}
\max _{s_{k}<b<e_{k}} \mathcal{C}_{\left(s_{k}, e_{k}\right]}^{b}(\mathbf{Y}) & \leq \max _{s_{k}<b<e_{k}} \mathcal{C}_{\left(s_{k}, e_{k}\right]}^{b}(\mathbf{f})+\lambda_{T} \\
& \leq \sqrt{C_{3} \log T\left(\Delta_{j}^{\mathbf{f}}\right)^{-2}} \Delta_{j}^{\mathbf{f}}+\sqrt{C_{3} \log T\left(\Delta_{j+1}^{f}\right)^{-2}} \Delta_{j+1}^{\mathrm{f}}+\lambda_{T} \\
& <\left(2 \sqrt{C_{3}}+\sqrt{8}\right) \sqrt{\log T}=C_{1} \sqrt{\log T} \leq \zeta_{T} .
\end{aligned}
$$

Hence the algorithm terminates with no further change-point detection.

## I.3. Proof of Theorem 2

Proof. The proof proceeds in analogy to the proof of Theorem 1. In five steps we shall establish the following result,
$\mathbb{P}\left(\hat{q}=q, \max _{j=1, \ldots, q}\left(\left|\hat{\tau}_{j}-\tau_{j}\right|\left(\Delta_{j}^{\mathbf{f}}\right)^{2 / 3}\right) \leq C_{3}(\log T)^{1 / 3}\right) \geq 1-T^{-1} /(6 \sqrt{\pi})-T \delta_{T}^{-1}\left(1-\delta_{T}^{2} T^{-2} / 36\right)^{M}$,
which in turn implies (10).

Step One and Step Two
We define the following two events

$$
\begin{aligned}
A_{T} & =\left\{\max _{s, b, e: 1 \leq s+1<b<e \leq T}\left|\mathcal{C}_{(s, e]}^{b}(\varepsilon)\right| \leq \lambda_{T}\right\}, \\
B_{T} & =\left\{\max _{j=1, \ldots, q} \max _{\substack{\tau_{j-1} \leq s+1<\tau_{j} \\
\tau_{j}<e \leq \tau_{j}+1 \\
s+1<b<e}} \frac{\left|\left\langle\phi_{(s, e]}^{b}\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{b}\right\rangle-\boldsymbol{\phi}_{(s, e]}^{\tau_{j}}\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{\tau_{j}}\right\rangle, \varepsilon\right\rangle\right|}{\left\|\boldsymbol{\phi}_{(s, e]}^{b}\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{b}\right\rangle-\boldsymbol{\phi}_{(s, e]}^{\tau_{s}}\left\langle\mathbf{f}, \boldsymbol{\phi}_{(s, e]}^{\tau_{j}}\right\rangle\right\|_{2}} \leq \lambda_{T}\right\},
\end{aligned}
$$

where $\lambda_{T}=\sqrt{8 \log T}$. Arguments as those used in Step One and Step Two of the proof of Theorem 1 show that $\mathbb{P}\left(A_{T}^{c}\right) \leq \frac{T^{-1}}{12 \sqrt{\pi}}$ and $\mathbb{P}\left(B_{T}^{c}\right) \leq \frac{T^{-1}}{12 \sqrt{\pi}}$.

Step Three
Here $\mathcal{I}_{j}^{L}, \mathcal{I}_{j}^{L}$ and $D_{T}^{M}$ are as defined in the proof of Theorem 1. In the rest of the proof, we assume that $A_{T}, B_{T}$ and $D_{T}^{M}$ all hold, where the last event is given by (22). Exactly as in the proof of Theorem 9 , we show that $\mathbb{P}\left(A_{T} \cap B_{T} \cap D_{T}^{M}\right) \geq 1-T^{-1} /(6 \sqrt{\pi})-T \delta_{T}^{-1}(1-$ $\left.\delta_{T}^{2} T^{-2} / 36\right)^{M}$.

We give the constants as follows:
$\underline{C}=72\left(4 \sqrt{2}+2 C_{3}^{3 / 2}\right)+1, \quad C_{1}=2 C_{3}^{3 / 2}+2 \sqrt{2}, \quad C_{2}=\frac{1}{72}-\frac{2 \sqrt{2}}{\underline{C}}, \quad C_{3}=2 \sqrt[3]{7}(3(1+\sqrt{2}))^{2 / 3}$.
Here we set $\underline{C}$ in such a way that $\underline{C C_{2}}>C_{1}$ (which also implies that $C_{2}>0$ ). Consequently, given $\delta_{T}^{3 / 2} \underline{f}_{T} \geq \underline{C} \sqrt{\log T}$ it is possible to select $\zeta_{T} \in\left[C_{1} \sqrt{\log T}, C_{2} \delta_{T}^{3 / 2} \underline{f}_{T}\right)$.

Again, these constants could be further refined. But since our main aim is to establish the rate, we chose not to pursue this direction here.

## Step Four

Consider a generic interval ( $s, e$ ] satisfying

$$
\begin{equation*}
\exists j \in\{1, \ldots, q\}, \exists k \in\{1, \ldots, M\} \text {, s.t. }\left(s_{k}, e_{k}\right] \subset(s, e] \text { and } s_{k} \times e_{k} \in \mathcal{I}_{j}^{L} \times \mathcal{I}_{j}^{R} \tag{28}
\end{equation*}
$$

and define events

$$
\begin{aligned}
\mathcal{M}_{(s, e]} & =\left\{m:\left(s_{m}, e_{m}\right] \in F_{T}^{M},\left(s_{m}, e_{m}\right] \subset(s, e]\right\}, \\
\mathcal{O}_{(s, e]} & =\left\{m \in \mathcal{M}_{(s, e]}: \max _{s_{m}+1<b<e_{m}} \mathcal{C}_{\left(s_{m}, e_{m}\right]}^{b}(\mathbf{Y})>\zeta_{T}\right\} .
\end{aligned}
$$

Let $b_{k}^{*}=\operatorname{argmax}_{s_{k}+1<b<e_{k}} \mathcal{C}_{\left(s_{k}, e_{k}\right]}^{b}(\mathbf{Y})$. We have

$$
\begin{aligned}
\mathcal{C}_{\left(s_{k}, e_{k}\right]}^{b_{k}^{*}}(\mathbf{Y}) & \geq \mathcal{C}_{\left(s_{k}, e_{k}\right]}^{\tau_{j}}(\mathbf{Y}) \\
& \geq \mathcal{C}_{\left(s_{k}, e_{k}\right]}^{b_{k}^{*}}(\mathbf{f})-\lambda_{T} \geq \frac{1}{\sqrt{24}}\left(\delta_{T} / 6\right)^{3 / 2} \Delta_{j}^{\mathrm{f}}-\lambda_{T} \geq \frac{1}{72} \delta_{T}^{3 / 2} \underline{f}_{T}-\lambda_{T} \\
& =\left(\frac{1}{72}-\frac{\lambda_{T}}{\delta_{T}^{3 / 2} \underline{f}_{T}}\right) \delta_{T}^{3 / 2} \underline{f}_{T} \geq\left(\frac{1}{72}-\frac{2 \sqrt{2}}{\underline{C}}\right) \delta_{T}^{3 / 2} \underline{f}_{T}=C_{2} \delta_{T}^{3 / 2} \underline{f}_{T}>\zeta_{T},
\end{aligned}
$$

where the third inequality above follows from Lemma 5 , therefore $\mathcal{O}_{s, e}$ is non-empty.
Let $m^{*}=\operatorname{argmin}_{m \in \mathcal{O}_{(s, e]}}\left(e_{m}-s_{m}\right)$ and $\left.b^{*}=\operatorname{argmax}_{s_{m^{*}}+1<b<e_{m^{*}}}\right)_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b}(\mathbf{Y})$. Arguing exactly as in Step Four in the proof of Theorem 1, we show that $\left(s_{m^{*}}+1, e_{m^{*}}\right)$ must contain exactly one change-point. Further, without loss of generality, assume that $\tau_{j} \in$ $\left(s_{m^{*}}+1, e_{m^{*}}\right)$. Let $\eta_{L}=\tau_{j}-s_{m^{*}}-1, \eta_{R}=e_{m^{*}}-\tau_{j}$ and

$$
\left.\eta_{T}=\left(\sqrt{3}\left(C_{1}-\sqrt{8}\right) \sqrt{\log T}\left(\Delta_{j}^{\mathbf{f}}\right)^{-1}\right)\right)^{2 / 3}-1
$$

We observe that $\min \left(\eta_{L}, \eta_{R}\right)>\eta_{T}$, as otherwise $\min \left(\eta_{L}, \eta_{R}\right) \leq \eta_{T}$ and Lemma 5 would imply

$$
\begin{aligned}
\mathcal{C}_{\left(s_{m^{*}}, e_{\left.m^{*}\right]}\right]}^{b^{*}}(\mathbf{Y}) & \leq \mathcal{C}_{\left(s_{m^{*}}, e_{\left.m^{*}\right]}\right]}^{b^{*}}(\mathbf{f})+\lambda_{T} \leq \mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{\tau_{j}}(\mathbf{f})+\lambda_{T} \leq \frac{1}{\sqrt{3}}\left(\eta_{T}+1\right)^{3 / 2} \Delta_{j}^{\mathbf{f}}+\lambda_{T} \\
& =\left(C_{1}-\sqrt{8}+\sqrt{8}\right) \sqrt{\log T}=C_{1} \sqrt{\log T} \leq \zeta_{T},
\end{aligned}
$$

contradicting $\mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b^{*}}(\mathbf{Y})>\zeta_{T}$.
We are now in the position to prove that $\left|b^{*}-\tau_{j}\right| \leq C_{3}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2 / 3}(\log T)^{1 / 3}:=\epsilon_{T}$. Let $b \in\left\{s_{m^{*}}+2, \ldots, e_{m^{*}}-1\right\}$. Our aim is to claim that when $\left|b-\tau_{j}\right|>\epsilon_{T}$,

$$
\begin{equation*}
\left(\mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{\tau_{j}}(\mathbf{Y})\right)^{2}-\left(\mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b}(\mathbf{Y})\right)^{2}>0 \tag{29}
\end{equation*}
$$

Since inequality (29) does not hold for $b=b^{*}$, proving this claim consequently demonstrates that $\left|b^{*}-\tau_{j}\right| \leq \epsilon_{T}$.

Without loss of generality, we consider the case of $b>\tau_{j}$. Using arguments as those in Step Four of the proof of Theorem 1 we can show that (29) is implied by $\kappa>(\sqrt{2}+1)^{2} \lambda_{T}^{2}$, where $\kappa=\left(\mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{\tau_{j}}(\mathbf{f})\right)^{2}-\left(\mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b}(\mathbf{f})\right)^{2}$. By Lemma $7, \kappa>(\sqrt{2}+1)^{2} \lambda_{T}^{2}$ is implied by

$$
\min \left(\left|b-\tau_{j}\right|, \eta_{L}\right)>\left(63\left(\Delta_{j}^{\mathbf{f}}\right)^{-2} \cdot 8(\sqrt{2}+1)^{2} \log T\right)^{1 / 3}=C_{3}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2 / 3}(\log T)^{1 / 3}
$$

However, for sufficiently large $T$,

$$
\begin{aligned}
\eta_{L}>\eta_{T} & =\left(\sqrt{3}\left(C_{1}-\sqrt{8}\right)\right)^{2 / 3}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2 / 3}(\log T)^{1 / 3}-1>\left(C_{1}-\sqrt{8}\right)^{2 / 3}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2 / 3}(\log T)^{1 / 3} \\
& >\left(C_{3}^{3 / 2}+\sqrt{8}-\sqrt{8}\right)^{2 / 3}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2 / 3}(\log T)^{1 / 3}=C_{3}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2 / 3}(\log T)^{1 / 3}=\epsilon_{T},
\end{aligned}
$$

hence $\left|b-\tau_{j}\right|>\epsilon_{T}$ implies (29), so it must hold that $\left|b^{*}-\tau_{j}\right| \leq \epsilon_{T}$.

Step Five
Using the arguments given above which are valid on the event $A_{T} \cap B_{T} \cap D_{T}^{M}$, we can now proceed with the proof of the theorem as follows. At the start of Algorithm 1 we have $s=0$ and $e=T$ and, provided that $q \geq 1$, condition (23) is satisfied. Therefore the algorithm detects a change-point $b^{*}$ in that interval such that $\left|b^{*}-\tau_{j}\right| \leq C_{3}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2 / 3}(\log T)^{1 / 3}$ for some $j$. By construction, we also have that $\left|b^{*}-\tau_{j}\right|<2 / 3 \delta_{T}$. This in turn implies that for all $l=1, \ldots, q$ such that $\tau_{l} \in(s+1, e)$ and $l \neq j$ we have either $\mathcal{I}_{l}^{L}, \mathcal{I}_{l}^{R} \subset\left(s, b^{*}\right]$ or $\mathcal{I}_{l}^{L}, \mathcal{I}_{l}^{R} \subset\left(b^{*}, e\right]$. Therefore (23) is satisfied within each segment containing at least one change-point. Note that before all $q$ change-points are detected, each change-point will not be detected twice. To see this, we suppose that $\tau_{j}$ has already been detected by $b^{*}$, then for all intervals $\left(s_{k}, e_{k}\right] \subset\left(\tau_{j}-\epsilon_{T}, \tau_{j}+2 / 3 \delta_{T}\right] \cup\left(\tau_{j}-2 / 3 \delta_{T}, \tau_{j}+\epsilon_{T}\right]$, Lemma 5, together with the event $A_{T}$, guarantees that for sufficiently large $T$

$$
\begin{aligned}
\max _{s_{k}+1<b<e_{k}} \mathcal{C}_{\left(s_{k}, e_{k}\right]}^{b}(\mathbf{Y}) & \leq \max _{s+1<b<e} \mathcal{C}_{\left(s_{k}, e_{k}\right]}^{b}(\mathbf{f})+\sqrt{8 \log T} \\
& \leq \frac{1}{\sqrt{3}}\left(C_{3}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2 / 3}(\log T)^{1 / 3}+1\right)^{3 / 2} \Delta_{j}^{\mathbf{f}}+\sqrt{8 \log T} \\
& \leq\left(2 C_{3}^{3 / 2}+\sqrt{8}\right) \sqrt{\log T}=C_{1} \sqrt{\log T} \leq \zeta_{T}
\end{aligned}
$$

Once all the change-points are detected, we then only need to consider $\left(s_{k}, e_{k}\right]$ such that

$$
\left(s_{k}, e_{k}\right] \subset\left(\tau_{j}-C_{3}\left(\Delta_{j}^{\mathrm{f}}\right)^{-2 / 3}(\log T)^{1 / 3}, \tau_{j+1}+C_{3}\left(\Delta_{j+1}^{\mathrm{f}}\right)^{-2 / 3}(\log T)^{1 / 3}\right]
$$

for $j=0, \ldots, q$, where we set $\Delta_{0}^{f}=\Delta_{q+1}^{f}=\infty$ for notational convenience. It follows from Lemma 6 (within the event $A_{T}$ ) that

$$
\begin{aligned}
\max _{s_{k}+1<b<e_{k}} \mathcal{C}_{\left(s_{k}, e_{k}\right]}^{b}(\mathbf{Y}) & \leq \max _{s_{k}+1<b<e_{k}} \mathcal{C}_{\left(s_{k}, e_{k}\right]}^{b}(\mathbf{f})+\sqrt{8 \log T} \\
& \leq\left(C_{3}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2 / 3}(\log T)^{1 / 3}\right)^{3 / 2} \Delta_{j}^{\mathrm{f}}+\left(C_{3}\left(\Delta_{j+1}^{\mathrm{f}}\right)^{-2 / 3}(\log T)^{1 / 3}\right)^{3 / 2} \Delta_{j+1}^{\mathrm{f}}+\sqrt{8 \log T} \\
& =\left(2 C_{3}^{3 / 2}+\sqrt{8}\right) \sqrt{\log T} \leq C_{1} \sqrt{\log T} \leq \zeta_{T} .
\end{aligned}
$$

Hence the algorithm terminates and no further change-points will be detected.

### 1.4. Proof of Theorem 3

Proof. Recall that $\left\{\varepsilon_{t}\right\}_{t=1}^{T}$ are i.i.d. $N\left(0, \sigma_{0}^{2}\right)$ with $\sigma_{0}=1$. For any candidate $\mathcal{T}\left(\zeta^{(k)}\right)$ on the NOT solution path, the sSIC criterion function in (S1) can be written as

$$
T \hat{\sigma}_{k}^{2}+\left(2 \hat{q}_{k}+1\right) \log ^{\alpha}(T)+\text { constant }
$$

where $\hat{\sigma}_{k}^{2}$ is the estimated variance of the noise (i.e. the residual sum of squares divided by $T$ ) based on $\mathcal{T}\left(\zeta^{(k)}\right)$, and $\hat{q}_{k}$ is the estimated number of change-points.

We now divide our proof into three parts.

Part I. About a particular model candidate on the NOT solution path
By Theorem 1, we know that with arbitrarily high probability for sufficiently large $T$, there exists $k^{*}$ such that $\mathcal{T}\left(\zeta^{\left(k^{*}\right)}\right)$ on the NOT solution path is a "good" candidate with $\hat{\tau}_{1}, \ldots, \hat{\tau}_{\hat{q}_{k^{*}}} \in \mathcal{T}\left(\zeta^{\left(k^{*}\right)}\right)$ satisfying $\hat{q}_{k^{*}}=q$ and $\max _{i=1}^{q}\left|\hat{\tau}_{i}-\tau_{i}\right| \leq C^{\prime} \log T$ for some $C^{\prime}>0$. In the rest of the proof, for presentational convenience, we condition on the event that such $k^{*}$ does exist throughout our analysis.

In addition, we recall that $\mathbf{1}_{(s, e]}=\left(\mathbf{1}_{(s, e]}(1), \ldots, \mathbf{1}_{(s, e]}(T)\right)^{\prime}$ with

$$
\mathbf{1}_{(s, e]}(t)= \begin{cases}(e-s)^{-1 / 2}, & t=s+1, \ldots, e  \tag{30}\\ 0, & \text { otherwise }\end{cases}
$$

and define the set

$$
E_{T}=\left\{\max _{s, e: 0 \leq s<e \leq T}\left|\left\langle\mathbf{1}_{(s, e]}, \varepsilon\right\rangle\right| \leq \sqrt{6 \log T}\right\}
$$

Using an argument similar to Step One of the proof of Theorem 1, we see that $\mathbb{P}\left(E_{T}^{c}\right)=$ $O\left(T^{-1}\right)$. Since we are only interested in proving a certain type of probabilistic statement for $T \rightarrow \infty$, here we could also assume that $E_{T}$ holds.

Let $\left\{\hat{f_{t}}\right\}_{t=1}^{T}$ be the fitted values using the candidate on the solution path with $\hat{\tau}_{1}, \ldots, \hat{\tau}_{\hat{q}_{k^{*}}} \in$ $\mathcal{T}\left(\zeta^{\left(k^{*}\right)}\right)$, and define $\tilde{f}_{t}=f_{\tau_{j}}$ for $t=\hat{\tau}_{j}, \ldots, \hat{\tau}_{j+1}-1$ for every $j=0,1, \ldots q$. Here for notational convenience, we suppressed the dependence of $\left\{\hat{f}_{t}\right\}_{t=1}^{T}$ and $\left\{\tilde{f}_{t}\right\}_{t=1}^{T}$ on $k^{*}$. It is easy to see that $f_{t}-\tilde{f}_{t}$ is piecewise-constant, only non-zero for $t$ between the true location of the change-point $\tau_{j}$ and its estimation $\hat{\tau}_{j}$, and exactly zero elsewhere. Write $\tilde{\mathbf{f}}=\left(\tilde{f}_{1}, \ldots, \tilde{f}_{T}\right)^{\prime}$. Then

$$
\begin{aligned}
T \hat{\sigma}_{k^{*}}^{2} & =\sum_{t=1}^{T}\left(\varepsilon_{t}+f_{t}-\hat{f}_{t}\right)^{2} \\
& \leq \sum_{t=1}^{T}\left(\varepsilon_{t}+f_{t}-\tilde{f}_{t}\right)^{2}=\sum_{t=1}^{T} \varepsilon_{t}^{2}+2\langle\varepsilon, \mathbf{f}-\tilde{\mathbf{f}}\rangle+\|\mathbf{f}-\tilde{\mathbf{f}}\|^{2} \\
& =\sum_{t=1}^{T} \varepsilon_{t}^{2}+4 q \bar{C} \sqrt{6 \log T} \sqrt{C^{\prime} \log T}+q(2 \bar{C})^{2} C^{\prime} \log T \\
& =\sum_{t=1}^{T} \varepsilon_{t}^{2}+\left(4 q \bar{C} \sqrt{6 C^{\prime}}+4 q C^{\prime} \bar{C}^{2}\right) \log T
\end{aligned}
$$

where the second last step follows from $E_{T}$, linearity of the inner product, and the fact that $\max _{i=1}^{q}\left|\hat{\tau}_{i}-\tau_{i}\right| \leq C^{\prime} \log T$.

## Part II. Estimation of the number of change-points

In this part, we prove that for NOT with the sSIC, $\mathbb{P}(\hat{q}=q) \rightarrow 1$ as $T \rightarrow \infty$. We accomplish this by showing separately that (i) $\mathbb{P}(\hat{q}>q) \rightarrow 0$ and (ii) $\mathbb{P}(\hat{q}<q) \rightarrow 0$.

First, for all $k$ with $\hat{q}_{k}>q$ and $\hat{\tau}_{1}, \ldots, \hat{\tau}_{\widehat{q}_{k}} \in \mathcal{T}\left(\zeta^{(k)}\right)$, we consider a "saturated oracle" candidate model with $\hat{q}_{k}+q$ change-points at $\hat{\tau}_{1}, \ldots, \hat{\tau}_{\hat{q}_{k}}, \tau_{1}, \ldots, \tau_{q}$ respectively. We reorder these $\hat{q}_{k}+q$ locations as $0=\stackrel{\circ}{\tau}_{0}<\stackrel{\circ}{\tau}_{1} \leq \ldots \leq{\stackrel{\circ}{\hat{q}_{k}+q}}^{\overbrace{\tau}}{\stackrel{\circ}{\hat{q}_{k}+q+1}}=T$, and denote the estimated variance of the errors corresponding this saturated oracle candidate by $\stackrel{\circ}{\sigma}_{k}^{2}$. Since for each
$j=0, \ldots, \hat{q}_{k}+q, f_{t}$ is constant over $\left\{1+\dot{\tau}_{j}, \ldots, \tau_{j+1}\right\}$, it then follows that

$$
\begin{aligned}
T \hat{\sigma}_{k}^{2} \geq T \stackrel{\circ}{\sigma}_{k}^{2} & =\sum_{j=0}^{\hat{q}_{k}+q} \sum_{t=1+\tilde{\tau}_{j}}^{\dot{\tau}_{j+1}}\left\{\varepsilon_{t}-\frac{1}{\tau_{j+1}-\dot{\tau}_{j}} \sum_{b=1+\tilde{\tau}_{j}}^{\mathfrak{\tau}_{j+1}} \varepsilon_{b}\right\}^{2} \\
& =\sum_{t=1}^{T} \varepsilon_{t}^{2}-\sum_{j=0}^{\hat{q}_{k}+q}\left\langle\varepsilon, \mathbf{1}_{1+\dot{\tau}_{j}, \tilde{\tau}_{j+1}}\right\rangle^{2} \geq \sum_{t=1}^{T} \varepsilon_{t}^{2}-6\left(q+\hat{q}_{k}+1\right) \log T,
\end{aligned}
$$

where the last line again follows from $E_{T}$. Note that in the above, for notational convenience, we have implicitly assumed that $\hat{\tau}_{1}, \ldots, \hat{\tau}_{\hat{q}_{k}}, \tau_{1}, \ldots, \tau_{q}$ are distinct. It is clear to see that the same argument holds even if $\left\{\hat{\tau}_{1}, \ldots, \hat{\tau}_{\hat{q}_{k}}\right\} \cap\left\{\tau_{1}, \ldots, \tau_{q}\right\} \neq \emptyset$. This means that for all $k$ with $\hat{q}_{k}>q$,

$$
\begin{aligned}
\operatorname{sSIC}(k)-\operatorname{sSIC}\left(k^{*}\right) \geq & T\left(\stackrel{\sigma}{\sigma}_{k}^{2}-\hat{\sigma}_{k^{*}}^{2}\right)+2\left(\hat{q}_{k}-q\right) \log ^{\alpha}(T) \\
\geq & \left\{\sum_{t=1}^{T} \varepsilon_{t}^{2}-6\left(q+\hat{q}_{k}+1\right) \log T\right\}-\left\{\sum_{t=1}^{T} \varepsilon_{t}^{2}+\left(4 q \bar{C} \sqrt{6 C^{\prime}}+4 q C^{\prime} \bar{C}^{2}\right) \log T\right\} \\
& \quad+2\left(\hat{q}_{k}-q\right) \log ^{\alpha}(T) \\
= & \left(\hat{q}_{k}-q\right)\left\{2 \log ^{\alpha}(T)-6 \log T\right\}-\left(12 q+4 q \bar{C} \sqrt{6 C^{\prime}}+4 q C^{\prime} \bar{C}^{2}+6\right) \log T \\
\geq & \left\{2 \log ^{\alpha}(T)-6 \log T\right\}-\left(12 q+4 q \bar{C} \sqrt{6 C^{\prime}}+4 q C^{\prime} \bar{C}^{2}+6\right) \log T>0
\end{aligned}
$$

for large enough $T$, which implies $\mathbb{P}(\hat{q}>q) \rightarrow 0$.
Second, for all $k$ with $\hat{q}_{k}<q$, it must be the case that one can find some $j^{*} \in\{1, \ldots, q\}$ such that the corresponding $\hat{f}_{t}$ is constant over $\left(\tau_{j^{*}}-\left\lfloor\delta_{T} / 2\right\rfloor, \tau_{j^{*}}+\left\lfloor\delta_{T} / 2\right\rfloor\right]$. Now consider the "intermediate" candidate model with $\hat{q}_{k}+3$ change-points at $\hat{\tau}_{1}, \ldots, \hat{\tau}_{\hat{q}_{k}}, \tau_{j^{*}}-$ $\left\lfloor\delta_{T} / 2\right\rfloor, \tau_{j^{*}}, \tau_{j^{*}}+\left\lfloor\delta_{T} / 2\right\rfloor$, and denote the corresponding estimated variance of errors by $\tilde{\sigma}_{k}^{2}$. Without loss of generality, assume that $f_{\tau_{j^{*}}+1}>f_{\tau_{j^{*}}}$. Then,

$$
\begin{aligned}
T \hat{\sigma}_{k}^{2}-T \tilde{\sigma}_{k}^{2} \geq & \sum_{t=\tau_{j^{*}}-\left\lfloor\delta_{T} / 2\right\rfloor+1}^{\tau_{j^{*}}}\left\{\varepsilon_{t}-\frac{\Delta_{j^{*}}^{\mathrm{f}}}{2}-\frac{1}{2\left\lfloor\delta_{T} / 2\right\rfloor} \sum_{b=\tau_{j^{*}}-\left\lfloor\delta_{T} / 2\right\rfloor+1}^{\tau_{j^{*}}+\left\lfloor\delta_{T} / 2\right\rfloor} \varepsilon_{b}\right\}^{2} \\
& +\sum_{t=\tau_{j^{*}+1}}^{\tau_{j^{*}}+\left\lfloor\delta_{T} / 2\right\rfloor}\left\{\varepsilon_{t}+\frac{\Delta_{j^{*}}^{\mathrm{f}}}{2}-\frac{1}{2\left\lfloor\delta_{T} / 2\right\rfloor} \sum_{b=\tau_{j^{*}}-\left\lfloor\delta_{T} / 2\right\rfloor+1}^{\tau_{j^{*}}+\left\lfloor\delta_{T} / 2\right\rfloor} \varepsilon_{b}\right\}^{2} \\
& -\sum_{t=\tau_{j^{*}}-\left\lfloor\delta_{T} / 2\right\rfloor+1}^{\tau_{j^{*}}} \varepsilon_{t}^{2}-\sum_{t=\tau_{j^{*}+1}}^{\tau_{j^{*}}+\left\lfloor\delta_{T} / 2\right\rfloor} \varepsilon_{t}^{2} \\
= & 2\left(\Delta_{j^{*}}^{\mathrm{f}} / 2\right)^{2}\left\lfloor\delta_{T} / 2\right\rfloor-\left\langle\varepsilon, \mathbf{1}_{\left(\tau_{j^{*}}-\left\lfloor\delta_{T} / 2\right\rfloor, \tau_{j^{*}}\left\lfloor\left\lfloor\delta_{T} / 2\right\rfloor\right]\right.}\right\rangle^{2} \\
& \quad-\Delta_{j^{*}}^{\mathbf{f}} \sqrt{\left\lfloor\delta_{T} / 2\right\rfloor}\left\{\left\langle\varepsilon, \mathbf{1}\left(\tau_{j^{*}}-\left\lfloor\delta_{T} / 2\right\rfloor, \tau_{j^{*}}\right\rfloor\right.\right. \\
\geq & -\left\langle\varepsilon, \mathbf{1}\left(\tau_{\left.j^{*}, \tau_{j^{*}}+\left\lfloor\delta_{T} / 2\right\rfloor\right\rfloor}\right\rangle\right\} \\
\geq & \frac{1}{2}\left(\Delta_{j^{*}}^{\mathrm{f}} \sqrt{\left\lfloor\delta_{T} / 2\right\rfloor}-2 \sqrt{6 \log T}\right)^{2}-12 \log T-6 \log T
\end{aligned}
$$

In the mean time, by adding $q-1$ more change-points, $\tau_{1}, \ldots, \tau_{j^{*}-1}, \tau_{j^{*}+1}, \ldots, \tau_{q}$, to the intermediate candidate model, we can show that using the same argument as in the first
half of Part II that

$$
T \tilde{\sigma}_{k}^{2} \geq \sum_{t=1}^{T} \varepsilon_{t}^{2}-6\left(q+\hat{q}_{k}+3\right) \log T
$$

Since $\delta_{T} \geq \underline{C}_{1}(\log T)^{\alpha^{\prime}}$ with $\alpha^{\prime}>1, \Delta_{j^{*}}^{\mathbf{f}} \sqrt{\left\lfloor\delta_{T} / 2\right\rfloor} \geq \underline{C}_{2} \sqrt{\left\lfloor\delta_{T} / 2\right\rfloor} \geq 2 \sqrt{6 \log T}$ for large enough $T$. Consequently, combining the previous two displayed equations lead to

$$
T \hat{\sigma}_{k}^{2} \geq \sum_{t=1}^{T} \varepsilon_{t}^{2}-6\left(q+\hat{q}_{k}+6\right) \log T+\frac{1}{2}\left(\underline{C}_{2} \sqrt{\left\lfloor\delta_{T} / 2\right\rfloor}-2 \sqrt{6 \log T}\right)^{2}
$$

for large enought $T$. This means that for all $k$ with $\hat{q}_{k}<q$,

$$
\begin{aligned}
\operatorname{sSIC}(k)-\operatorname{sSIC}\left(k^{*}\right)= & T\left(\hat{\sigma}_{k}^{2}-\hat{\sigma}_{k^{*}}^{2}\right)+2\left(\hat{q}_{k}-q\right) \log ^{\alpha}(T) \\
\geq \geq & \frac{1}{2}\left(\underline{C}_{2} \sqrt{\left\lfloor\delta_{T} / 2\right\rfloor}-2 \sqrt{6 \log T}\right)^{2} \\
& \quad-\left(4 q \bar{C} \sqrt{6 C^{\prime}}+4 q C^{\prime} \bar{C}^{2}+6 q+6 \hat{q}_{k}+36\right) \log T-2 q \log ^{\alpha}(T) \\
> & 0
\end{aligned}
$$

for sufficiently large $T$, where we again used that fact that $\delta_{T} \geq \underline{C}_{1}(\log T)^{\alpha^{\prime}}$ with $\alpha^{\prime}>\alpha>$ 1, so $\frac{1}{2}\left(\underline{C}_{2} \sqrt{\left\lfloor\delta_{T} / 2\right\rfloor}-2 \sqrt{6 \log T}\right)^{2}$ is at least of order $(\log T)^{\alpha^{\prime}}$. This implies $\mathbb{P}(\hat{q}<q) \rightarrow 0$.

In conclusion, we have established $\mathbb{P}(\hat{q}=q) \rightarrow 1$.

## Part III. Estimation of the change-point locations

In view of the conclusion of Part II, in the rest of the proof we could assume that $E_{T}$ holds and $\hat{q}=q$. Suppose that the model picked via NOT with the sSIC is $\hat{\tau}_{1}, \ldots, \hat{\tau}_{q} \in \mathcal{T}\left(\zeta^{(\hat{k})}\right)$. Furthermore, let

$$
j^{*}=\operatorname{argmax}_{j=1, \ldots, q} \min _{i=1, \ldots, q}\left|\hat{\tau}_{i}-\tau_{j}\right| \quad \text { and } \quad C:=\frac{\min \left(\left\lfloor\delta_{T} / 2\right\rfloor, \min _{i=1, \ldots, q}\left|\hat{\tau}_{i}-\tau_{j^{*}}\right|\right)}{\log T} .
$$

Our aim is to show that $C$ is finite (more precisely, has an upper bound independent of $T)$. Now consider a "near-saturated oracle" candidate model with $2 q+1$ change-points at

$$
\left\{\hat{\tau}_{1}, \ldots, \hat{\tau}_{q}, \tau_{1}, \ldots, \tau_{j^{*}-1}, \tau_{j^{*}+1}, \ldots, \hat{\tau}_{q}, \tau_{j^{*}}-C \log T, \tau_{j^{*}}+C \log T\right\}
$$

with the corresponding estimated variance of the errors denoted as $\dot{\sigma}_{\hat{k}}^{2}$. So here instead of adding all the true change-points to the set of estimated change-points as before (which generates the so-called "saturated oracle"), we add all true change-points apart from $\tau_{j^{*}}$, and replace it by $\tau_{j^{*}} \pm C \log T$.

Note that by construction (i.e. via $\delta_{T}$ in the definition of $C$ ), $f_{t}$ is constant over $\left(\tau_{j^{*}}-C \log T, \tau_{j^{*}}\right]$ and $\left(\tau_{j^{*}}, \tau_{j^{*}}+C \log T\right]$. In addition, $\Delta_{j^{*}}^{\mathrm{f}}=\left|f_{\tau_{j^{*}+1}}-f_{\tau_{j^{*}}}\right| \geq \underline{f}_{T}$. Write

$$
\bar{\varepsilon}_{*}=\frac{1}{2 C \log T} \sum_{t=\tau_{j^{*}}-C \log T+1}^{\tau_{j^{*}}+C \log T} \varepsilon_{t} .
$$

Without loss of generality, assume that $f_{\tau_{j^{*}+1}}>f_{\tau_{j^{*}}}$. Now using the arguments similar to those in Part II, we see that

$$
\begin{aligned}
T \hat{\sigma}_{\hat{k}}^{2} \geq T \dot{\sigma}_{\hat{k}}^{2} \geq & \sum_{t=1}^{\tau_{j^{*}}-C \log T} \varepsilon_{t}^{2}+\sum_{t=\tau_{j^{*}}+C \log T+1}^{T} \varepsilon_{t}^{2}-(2 q) 6 \log T \\
& +\sum_{t=\tau_{j^{*}}-C \log T+1}^{\tau_{j^{*}}}\left(\varepsilon_{t}-\Delta_{j^{*}}^{\mathrm{f}} / 2-\bar{\varepsilon}_{*}\right)^{2}+\sum_{t=\tau_{j^{*}+1}}^{\tau_{j^{*}}+C \log T}\left(\varepsilon_{t}+\Delta_{j^{*}}^{\mathrm{f}} / 2-\bar{\varepsilon}_{*}\right)^{2} \\
= & \sum_{t=1}^{T} \varepsilon_{t}^{2}-12 q \log T+\Delta_{j^{*}}^{\mathrm{f}}\left(\sum_{t=\tau_{j^{*}+1}}^{\tau_{j^{*}}+C \log T} \varepsilon_{t}-\sum_{t=\tau_{j^{*}}-C \log T+1}^{\tau_{j^{*}}} \varepsilon_{t}\right) \\
& \quad+\left(\Delta_{j^{*}}^{\mathrm{f}} / 2\right)^{2}(2 C \log T)-(2 C \log T) \bar{\varepsilon}_{*}^{2} \\
= & \sum_{t=1}^{T} \varepsilon_{t}^{2}-12 q \log T+\Delta_{j^{*}}^{\mathrm{f}} \sqrt{C \log T}\left\{\left\langle\varepsilon, \mathbf{1}_{\left.\left.\tau_{j^{*}+1, \tau_{j^{*}}+C \log T}\right\rangle-\left\langle\varepsilon, \mathbf{1}_{\tau_{j^{*}}-C \log T+1, \tau_{j^{*}}}\right\rangle\right\}} \quad+\left(\Delta_{j^{*}}^{\mathrm{f}} / 2\right)^{2}(2 C \log T)-\left\langle\varepsilon, \mathbf{1}_{\left.\tau_{j^{*}}-C \log T+1, \tau_{j^{*}}+C \log T\right\rangle^{2}}\right.\right.\right. \\
\geq & \sum_{t=1}^{T} \varepsilon_{t}^{2}-\left\{6(2 q+1)+2 \sqrt{6 C} \Delta_{j^{*}}^{\mathrm{f}}\right\} \log T+\left(\Delta_{j^{*}}^{\mathrm{f}} / 2\right)^{2}(2 C \log T)
\end{aligned}
$$

However,

$$
T \hat{\sigma}_{\hat{k}}^{2} \leq T \hat{\sigma}_{k^{*}}^{2} \leq \sum_{t=1}^{T} \varepsilon_{t}^{2}+\left(4 q \bar{C} \sqrt{6 C^{\prime}}+4 q C^{\prime} \bar{C}^{2}\right) \log T
$$

Combining the above two inequalities, and after some algebraic manipulations, we get

$$
2 q \bar{C} \sqrt{6 C^{\prime}}+2 q C^{\prime} \bar{C}^{2} \geq C\left(\Delta_{j^{*}}^{\mathbf{f}} / 2\right)^{2}-3(2 q+1)-\sqrt{6 C} \Delta_{j^{*}}^{\mathrm{f}},
$$

and thus

$$
2 q \bar{C} \sqrt{6 C^{\prime}}+2 q C^{\prime} \bar{C}^{2}+3(2 q+1)+6 \geq\left(\sqrt{C} \Delta_{j^{*}}^{\mathrm{f}} / 2-\sqrt{6}\right)^{2},
$$

which entails

$$
C \leq 4\left[\left\{2 q \bar{C} \sqrt{6 C^{\prime}}+2 q C^{\prime} \bar{C}^{2}+3(2 q+1)+6\right\}^{1 / 2}+\sqrt{6}\right]^{2} / \underline{C}_{2}^{2} .
$$

Finally, we remark that since $\delta_{T}=\min _{j=1, \ldots, q+1}\left(\tau_{j}-\tau_{j-1}\right) \geq \underline{C}_{1}(\log T)^{\alpha^{\prime}}$, for sufficiently large $T$,

$$
C \log T \geq \min \left(\left\lfloor\delta_{T} / 2\right\rfloor, \max _{j=1, \ldots, q} \min _{i=1, \ldots, q}\left|\hat{\gamma}_{i}-\tau_{j}\right|\right)=\max _{j=1, \ldots, q}\left|\hat{\tau}_{j}-\tau_{j}\right| .
$$

Therefore, $\mathbb{P}\left(\max _{j=1, \ldots, q}\left|\hat{\tau}_{j}-\tau_{j}\right| \leq C \log T\right) \rightarrow 1$, as required.

## I.5. Proof of Theorem 4

First, we strengthen Theorem 2 in the scenario where the true signal has finitely many kinks with spacing $\sim T$.

Lemma 9. Under the assumptions of Theorem 4, there exist constants $C^{\prime}$ and $\tilde{C}$ such that by setting $\zeta_{T}=\tilde{C} \sqrt{T}$ and $M \geq 36 \underline{C}_{1}^{-2} \log \left(\underline{C}_{1}^{-1} T\right)$, we have that

$$
\begin{equation*}
\mathbb{P}\left(\hat{q}=q, \max _{j=1, \ldots, q}\left|\hat{\tau}_{j}-\tau_{j}\right| \leq C^{\prime} \sqrt{T \log T}\right) \rightarrow 1 \tag{31}
\end{equation*}
$$

as $T \rightarrow \infty$.
Proof. Let $\underline{C}, C_{1}, C_{2}, C_{3}>0$ be the constants upon applying Theorem 2. For simplicity, here we shall take

$$
\widetilde{C}=C_{2} \underline{C}_{1}^{3 / 2} \underline{C}_{2} / 2 \quad \text { and } \quad C^{\prime}=\frac{32 \sqrt{3}(\sqrt{2}+1)}{\underline{C}_{2}\left\{\sqrt{3} \underline{C}_{1} \widetilde{C} / \bar{C}\right\}^{1 / 3}}
$$

First, we verify that the conditions in Theorem 2 are satisfied. Specifically, we note that under the additional assumptions of Theorem 4, for sufficiently large $T$,
(a) $\delta_{T}^{3 / 2} \underline{f}_{T} \geq \underline{C}_{1}^{3 / 2} \underline{C}_{2} \sqrt{T}>\underline{C} \sqrt{\log T}$,
(b) $\zeta_{T}=\widetilde{C} \sqrt{T} \in\left[C_{1} \sqrt{\log T}, C_{2} \delta_{T}^{3 / 2} \underline{f}_{T}\right)$,
(c) $M \geq 36 \underline{C}_{1}^{-2} \log \left(\underline{C}_{1}^{-1} T\right) \geq 36\left(T / \delta_{T}\right)^{2} \log \left\{\left(T / \delta_{T}\right) T\right\}$.

This means that

$$
\mathbb{P}\left(\hat{q}=q, \max _{j=1, \ldots ., q}\left|\hat{\tau}_{j}-\tau_{j}\right| \leq C_{3} \underline{C}_{2}^{-2 / 3}\left(T^{2} \log T\right)^{1 / 3}\right) \rightarrow 1
$$

Second, to strengthen the convergence rate of $\max _{j=1, \ldots, q}\left|\hat{\tau}_{j}-\tau_{j}\right|$, we make some minor modifications to Step Four in the proof of Theorem 2.

We still let $m^{*}=\operatorname{argmin}_{m \in \mathcal{O}_{(s, e]}}\left(e_{m}-s_{m}\right)$ and $b^{*}=\operatorname{argmax}_{s_{m^{*}}+1<b<e_{m^{*}}} \mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}(\mathbf{Y})$, where $\left(s_{m^{*}}+1, e_{m^{*}}\right)$ must contain exactly one change-point. Again, we consider $\tau_{j} \in$ $\left(s_{m^{*}}+1, e_{m^{*}}\right)$, and let $\eta_{L}=\tau_{j}-s_{m^{*}}-1$ and $\eta_{R}=e_{m^{*}}-\tau_{j}$. Note that

$$
\max _{j=1, \ldots, q} \Delta_{j}^{\mathbf{f}} \leq \frac{4 \max _{i=1, \ldots, T}\left|f_{i}\right|}{\delta_{T}} \leq \frac{4 \bar{C}}{\underline{C}_{1}} \frac{1}{T}
$$

By setting $\eta_{T}=\left\{\sqrt{3} \underline{C}_{1} \widetilde{C} /(8 \bar{C})\right\}^{2 / 3} T-1$ (different from the proof of Theorem 2 ), we observe that $\min \left(\eta_{L}, \eta_{R}\right)>\eta_{T}$ for sufficiently large $T$ (satisfying $8 \log T<\widetilde{C}^{2} T / 4$ ). It is because otherwise $\min \left(\eta_{L}, \eta_{R}\right) \leq \eta_{T}$ and Lemma 5 would imply that

$$
\begin{aligned}
\mathcal{C}_{\left(s_{\left.m^{*}, e_{m^{*}}\right]}^{b^{*}}\right.}(\mathbf{Y}) & \leq \mathcal{C}_{\left(s_{\left.m^{*}, e_{m^{*}}\right]}^{b^{*}}\right.}(\mathbf{f})+\lambda_{T} \leq \mathcal{C}_{\left(s_{\left.m^{*}, e_{m^{*}}\right]}^{\tau_{j}}\right.}(\mathbf{f})+\lambda_{T} \leq \frac{1}{\sqrt{3}}\left(\eta_{T}+1\right)^{3 / 2} \frac{4 \bar{C}}{\underline{C}_{1}} \frac{1}{T}+\lambda_{T} \\
& =\frac{\widetilde{C}}{2} \sqrt{T}+\sqrt{8 \log T}<\widetilde{C} \sqrt{T}=\zeta_{T}
\end{aligned}
$$

which leads to a contradiction.
We are now in the position to prove that $\left|b^{*}-\tau_{j}\right| \leq C^{\prime} \sqrt{T \log T}:=\epsilon_{T}$. Note that in view of Theorem 2, it suffices to only consider
$b \in\left\{s_{m^{*}}+2, \ldots, e_{m^{*}}-1\right\} \cap\left\{\tau_{j}-\left\lceil C_{3}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2 / 3}(\log T)^{1 / 3}\right\rceil, \ldots, \tau_{j}+\left\lceil C_{3}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2 / 3}(\log T)^{1 / 3}\right\rceil\right\}$

Our aim is to show that given $\left|b-\tau_{j}\right|>\epsilon_{T}$ (as well as $\left|b-\tau_{j}\right| \leq C_{3}\left(\Delta_{j}^{\mathbf{f}}\right)^{-2 / 3}(\log T)^{1 / 3}$ in view of Theorem 2),

$$
\begin{equation*}
\left(\mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{\tau_{j}}(\mathbf{Y})\right)^{2}-\left(\mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b}(\mathbf{Y})\right)^{2}>0 \tag{32}
\end{equation*}
$$

Inequality (32) does not hold for $b=b^{*}$ by the definition of $b^{*}$, so proving this claim would demonstrate that $\left|b^{*}-\tau_{j}\right| \leq \epsilon_{T}$.

Using arguments as those in Step Four of the proof of Theorem 1 (or Theorem 2), we can show that $(32)$ is implied by $\kappa>(\sqrt{2}+1)^{2} \lambda_{T}^{2}$, where $\kappa=\left(\mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{\tau_{j}}(\mathbf{f})\right)^{2}-\left(\mathcal{C}_{\left(s_{m^{*}}, e_{m^{*}}\right]}^{b}(\mathbf{f})\right)^{2}$. By Lemma $8, \kappa>(\sqrt{2}+1)^{2} \lambda_{T}^{2}$ is implied by

$$
\begin{equation*}
\frac{\left(\Delta_{j}^{\mathbf{f}}\right)^{2}}{48}\left\{\min \left(\eta_{L}, \eta_{R}\right)-1\right\}\left|b-\tau_{j}\right|^{2}>(\sqrt{2}+1)^{2} \lambda_{T}^{2} \tag{33}
\end{equation*}
$$

In view of the fact that

$$
\min \left(\eta_{L}, \eta_{R}\right)-1>\eta_{T}-2=\left\{\sqrt{3} \underline{C}_{1} \widetilde{C} /(8 \bar{C})\right\}^{2 / 3} T-2>\left\{\sqrt{3} \underline{C}_{1} \widetilde{C} /(8 \bar{C})\right\}^{2 / 3} T / 2
$$

for sufficiently large $T,(33)$ is further implied by

$$
\left|b-\tau_{j}\right|>\frac{8 \sqrt{6}(\sqrt{2}+1) \sqrt{\log T}}{\underline{C}_{2} / T\left\{\sqrt{3} \underline{C}_{1} \widetilde{C} /(8 \bar{C})\right\}^{1 / 3} \sqrt{T / 2}}=\frac{32 \sqrt{3}(\sqrt{2}+1)}{\underline{C}_{2}\left\{\sqrt{3} \underline{C}_{1} \widetilde{C} / \bar{C}\right\}^{1 / 3}} \sqrt{T \log T}=C^{\prime} \sqrt{T \log T}
$$

In conclusion, $\left|b-\tau_{j}\right|>\epsilon_{T}$ implies (32), leading to a contradiction. So it must hold that $\left|b^{*}-\tau_{j}\right| \leq \epsilon_{T}$ for large $T$.

Finally, since $\mathbb{P}(\hat{q}=q) \rightarrow 1$, we have that

$$
\mathbb{P}\left(\hat{q}=q, \max _{j=1, \ldots, q}\left|\hat{\tau}_{j}-\tau_{j}\right| \leq C^{\prime} \sqrt{T \log T}\right) \rightarrow 1
$$

as required.

Now we are in the position to prove Theorem 4.
Proof. The proof proceeds in analogy to the proof of Theorem 3. In the following, we present details of the main steps.

Again, thanks to the standard Gaussianity of the noise, for any candidate $\mathcal{T}\left(\zeta^{(k)}\right)$ on the NOT solution path, the sSIC criterion function in (S2) can be written as

$$
T \hat{\sigma}_{k}^{2}+\left(2 \hat{q}_{k}+2\right) \log ^{\alpha}(T)+\text { constant }
$$

where $\hat{\sigma}_{k}^{2}$ is the estimated variance of the noise (i.e. the residual sum of squares divided by $T$ ) based on $\mathcal{T}\left(\zeta^{(k)}\right)$, and $\hat{q}_{k}$ is the estimated number of kinks.

Part I. About a particular model candidate on the NOT solution path
By Lemma 9, we know that with arbitrarily high probability for sufficiently large $T$, there exists $k^{*}$ such that $\mathcal{T}\left(\zeta^{\left(k^{*}\right)}\right)$ on the NOT solution path is a "good" candidate with $\hat{\tau}_{1}, \ldots, \hat{\tau}_{\hat{q}_{k^{*}}} \in \mathcal{T}\left(\zeta^{\left(k^{*}\right)}\right)$ satisfying $\hat{q}_{k^{*}}=q$ and $\max _{i=1}^{q}\left|\hat{\tau}_{i}-\tau_{i}\right| \leq C^{\prime} \sqrt{T \log T}$ for some $C^{\prime}>0$. In the rest of the proof, for presentational convenience, we assume the existence of such $k^{*}$.

Define the set

$$
E_{T}=\left\{\max _{s, e: 0 \leq s<e \leq T} \max \left(\left|\left\langle\gamma_{(s, e]}, \boldsymbol{\varepsilon}\right\rangle\right|,\left|\left\langle\mathbf{1}_{(s, e]}, \varepsilon\right\rangle\right|\right) \leq \sqrt{6 \log T}\right\} .
$$

Using the Bonferroni bound, we see that $\mathbb{P}\left(E_{T}^{c}\right)=O\left(T^{-1}\right)$. Again, in the following, we could assume that $E_{T}$ holds.

Let $\left\{\hat{f}_{t}\right\}_{t=1}^{T}$ be the fitted values using the candidate on the solution path with $\hat{\tau}_{1}, \ldots, \hat{\tau}_{\hat{\tau}_{k^{*}}} \in$ $\mathcal{T}\left(\zeta^{\left(k^{*}\right)}\right)$, and define $\tilde{f}_{1}=\hat{f}_{1}, \tilde{f}_{t+1}=\tilde{f}_{t}+\left(f_{\tau_{j}+1}-f_{\tau_{j}}\right)$ for $t=\hat{\tau}_{j}, \ldots, \hat{\tau}_{j+1}-1$ for every $j=0,1, \ldots q$. Again, here for notational convenience, we suppressed the dependence of $\left\{\hat{f}_{t}\right\}_{t=1}^{T}$ and $\left\{\tilde{f}_{t}\right\}_{t=1}^{T}$ on $k^{*}$. It is easy to see that $f_{t}-\tilde{f}_{t}$ is piecewise-linear and continuous, with at most $2 q$ kinks and

$$
\max _{t=1, \ldots, T}\left|f_{t}-\tilde{f}_{t}\right| \leq q \max _{j}\left(\Delta_{j}^{\mathbf{f}}\right) C^{\prime} \sqrt{T \log T} \leq \frac{4 \bar{C}}{\underline{C}_{1} T} C^{\prime} q \sqrt{T \log T}=\frac{4 q \bar{C} C^{\prime}}{\underline{C}_{1}} \sqrt{\log T / T} .
$$

Write $\tilde{\mathbf{f}}=\left(\tilde{f}_{1}, \ldots, \tilde{f}_{T}\right)^{\prime}$, then $\|\mathbf{f}-\tilde{\mathbf{f}}\|^{2} \leq\left(4 q \bar{C} C^{\prime} / \underline{C}_{1}\right)^{2} \log T$. Furthermore, it is easy to verify (under $E_{T}$ ) that

$$
\begin{aligned}
T \hat{\sigma}_{k^{*}}^{2} & =\sum_{t=1}^{T}\left(\varepsilon_{t}+f_{t}-\hat{f}_{t}\right)^{2} \leq \sum_{t=1}^{T}\left(\varepsilon_{t}+f_{t}-\tilde{f}_{t}\right)^{2}=\sum_{t=1}^{T} \varepsilon_{t}^{2}+2\langle\varepsilon, \mathbf{f}-\tilde{\mathbf{f}}\rangle+\|\mathbf{f}-\tilde{\mathbf{f}}\|^{2} \\
& \leq \sum_{t=1}^{T} \varepsilon_{t}^{2}+M^{\prime} \log T
\end{aligned}
$$

for some positive constant $M^{\prime}$ that does not depend on $T$. Consequently, as $T \rightarrow \infty$, it follows that $\mathbb{P}\left(\hat{\sigma}_{k^{*}}^{2}<1+\delta / 2\right)=1$ for any $\delta>0$.

## Part II. Estimation of the number of change-points

Our aim in this part is to show that $\mathbb{P}(\hat{q}=q) \rightarrow 1$ as $T \rightarrow \infty$. We accomplish this by showing separately that (i) $\mathbb{P}(\hat{q}<q) \rightarrow 0$ and (ii) $\mathbb{P}(\hat{q}>q) \rightarrow 0$.

First, we note that it follows from Lemma 5.3 and 5.4 of Liu et al. (1997) that there exists $\delta>0$ such that as $T \rightarrow \infty$,

$$
\min _{k: \hat{q}_{k}<q} \mathbb{P}\left(\hat{\sigma}_{k}^{2}>1+\delta\right) \rightarrow 1
$$

This means that for all $k$ with $\hat{q}_{k}<q$,

$$
\operatorname{sSIC}(k)-\operatorname{sSIC}\left(k^{*}\right)=T\left(\hat{\sigma}_{k}^{2}-\hat{\sigma}_{k^{*}}^{2}\right)+2\left(\hat{q}_{k}-q\right) \log ^{\alpha}(T) \geq \delta T / 2-2 q \log ^{\alpha}(T)>0
$$

for large enough $T$, which implies $\mathbb{P}(\hat{q}<q) \rightarrow 0$.
Second, for all $k$ with $\hat{q}_{k}>q$ and $\hat{\tau}_{1}, \ldots, \hat{\tau}_{\hat{q}_{k}} \in \mathcal{T}\left(\zeta^{(k)}\right)$, we consider a "saturated oracle" candidate model with $\hat{q}_{k}+q$ kinks at $\hat{\tau}_{1}, \ldots, \hat{\tau}_{\hat{q}_{k}}, \tau_{1}, \ldots, \tau_{q}$ respectively. We reorder these $\hat{q}_{k}+q$ locations as $0=\grave{\tau}_{0}<\stackrel{\tau}{\tau}_{1} \leq \ldots \leq \dot{\tau}_{\hat{q}_{k}+q}<{\stackrel{\circ}{\hat{q}_{k}+q+1}}=T$, and denote by $\dot{\sigma}_{k}^{2}$ the estimated variance of the errors corresponding to a piecewise-linear model with features at these locations but without the continuity constraint (so effectively the way of estimating this quantity under Scenario (S3)). Let $\boldsymbol{\varepsilon}=\left(\varepsilon_{1}, \ldots, \varepsilon_{T}\right)^{\prime}$,

$$
\boldsymbol{\Gamma}_{(s, e]}:=\left[\mathbf{1}_{(s, e]}, \boldsymbol{\gamma}_{(s, e]}\right] \quad \text { and } \quad \mathbf{H}_{(s, e]}=\boldsymbol{\Gamma}_{(s, e]}\left(\boldsymbol{\Gamma}_{(s, e]}^{\prime} \boldsymbol{\Gamma}_{(s, e]}\right)^{-1} \boldsymbol{\Gamma}_{(s, e]}^{\prime}
$$

for $0 \leq s<e \leq T$, where $\boldsymbol{\Gamma}_{(s, e]}$ is a $T \times 2$ matrix and $\mathbf{H}_{(s, e]}$ is a $T \times T$ matrix. Furthermore, denote by $\mathbf{D}_{(s, e]}$ a $T \times T$ diagonal matrix with 1 in the $(s+1, s+1)$-th to the $(e, e)$-th entries and zero elsewhere. Here both $\mathbf{H}_{(s, e]}$ and $\mathbf{H}_{(s, e]}-\mathbf{D}_{(s, e]}$ are idempotent matrices.

Then the residual sum of squares for fitting a linear line over $\left(\stackrel{\circ}{\tau}_{j}, \stackrel{\circ}{\tau}_{j+1}\right]$ (on which $f_{t}$ is linear as well) is

$$
(\mathbf{f}+\boldsymbol{\varepsilon})^{\prime}\{\mathbf{D}_{(\overbrace{j}, \tilde{\tau}_{j+1}]}-\mathbf{H}_{\left(\tilde{\tau}_{j}, \tilde{\tau}_{j+1}\right]}\}(\mathbf{f}+\boldsymbol{\varepsilon})=\boldsymbol{\varepsilon}^{\prime}\left\{\mathbf{D}_{\left(\tilde{\tau}_{j}, \tilde{\tau}_{j+1}\right]}-\mathbf{H}_{\left(\tilde{\tau}_{j}, \tilde{\tau}_{j+1}\right]}\right\} \varepsilon .
$$

It then follows that

$$
\begin{aligned}
T \hat{\sigma}_{k}^{2} \geq T \stackrel{\circ}{\sigma}_{k}^{2} & =\sum_{j=0}^{\hat{q}_{k}+q} \varepsilon^{\prime}\left\{\mathbf{D}_{\left(\tau_{j}, \tilde{\tau}_{j+1}\right]}-\mathbf{H}_{\left(\stackrel{\tau}{j}_{j}, \tilde{\tau}_{j+1}\right]}\right\} \varepsilon . \\
& =\sum_{t=1}^{T} \varepsilon_{t}^{2}-\sum_{j=0}^{\hat{q}_{k}+q} \varepsilon^{\prime} \mathbf{H}_{\left(\stackrel{\tau}{\tau}_{j}, \tilde{\tau}_{j+1}\right]} \varepsilon
\end{aligned}
$$

Note that $\varepsilon^{\prime} \mathbf{H}_{(s, e]} \varepsilon$ follows a $\chi_{2}^{2}$ distribution. For any $Z \sim \chi_{2}^{2}, \mathbb{P}(Z>z) \leq e^{-z / 2}$. Therefore, by defining the set

$$
G_{T}=\left\{\max _{s, e: 0 \leq s<e \leq T} \varepsilon^{\prime} \mathbf{H}_{(s, e]} \varepsilon \leq 6 \log T\right\}
$$

we have that $\mathbb{P}\left(G_{T}^{c}\right)=O\left(T^{-1}\right)$ using the Bonferroni bound. Now assume that $G_{T}$ holds, it follows that

$$
T \hat{\sigma}_{k}^{2} \geq \sum_{t=1}^{T} \varepsilon_{t}^{2}-6\left(\hat{q}_{k}+q+1\right) \log T
$$

This means that for all $k$ with $\hat{q}_{k}>q$,

$$
\begin{aligned}
\operatorname{sSIC}(k)-\operatorname{sSIC}\left(k^{*}\right) & \geq T\left(\stackrel{\circ}{\sigma}_{k}^{2}-\hat{\sigma}_{k^{*}}^{2}\right)+2\left(\hat{q}_{k}-q\right) \log ^{\alpha}(T) \\
& \geq 2\left(\hat{q}_{k}-q\right) \log ^{\alpha}(T)-\left\{6\left(\hat{q}_{k}+q+1\right)+M^{\prime}\right\} \log T \\
& =2\left(\hat{q}_{k}-q\right)\left\{\log ^{\alpha}(T)-3 \log T\right\}-\left(12 q+6+M^{\prime}\right) \log T \\
& \geq 2 \log ^{\alpha}(T)-\left(12 q+12+M^{\prime}\right) \log T>0
\end{aligned}
$$

for large enough $T$, which in turn implies $\mathbb{P}(\hat{q}>q) \rightarrow 0$.
In conclusion, we have established that $\mathbb{P}(\hat{q}=q) \rightarrow 1$.

Part III. Estimation of the change-point locations
In view of the conclusion of Part II, in the rest of the proof we could assume that $A_{T} \cap$ $B_{T} \cap D_{T} \cap E_{T} \cap G_{T}$ holds and $\hat{q}=q$.

Suppose that the model picked via NOT with the sSIC is $\hat{\tau}_{1}, \ldots, \hat{\tau}_{q} \in \mathcal{T}\left(\zeta^{(\hat{k})}\right)$. Comparing the residual sum of squares of this candidate with $\mathcal{T}\left(\zeta^{\left(k^{*}\right)}\right)$ yields that $\hat{\tau}_{j} \in\left\{\tau_{j}-\right.$ $\left.\left\lfloor\delta_{T} / 6\right\rfloor+1, \ldots, \tau_{j}+\left\lfloor\delta_{T} / 6\right\rfloor-1\right\}$. It is because otherwise one could find an interval of length roughly $\delta_{T} / 3$ (i.e. $\sim T$ ) with a true kink in the middle of but with no kinks in its estimates, leading to $\hat{\sigma}^{2}>1+\delta$ for some $\delta>0$ (see Lemma 5.3 and 5.4 of Liu et al. (1997)), and thus a contradiction (as the sSIC would clearly prefer $\mathcal{T}\left(\zeta^{\left(k^{*}\right)}\right)$ ). Likewise, since $\hat{q}=q$, it is easy to see that $\hat{\tau}_{j}$ is the only estimated kink over $\left(\tau_{j}-\left\lceil\delta_{T} / 3\right\rceil-2, \tau_{j}+\left\lceil\delta_{T} / 3\right\rceil\right\rceil$ for every $j=1, \ldots, q$.

Let

$$
j^{*}=\operatorname{argmax}_{j=1, \ldots, q}\left|\hat{\tau}_{j}-\tau_{j}\right|
$$

Now consider a "near-saturated oracle" candidate model with $2 q+1$ kinks at

$$
\left\{\hat{\tau}_{1}, \ldots, \hat{\tau}_{q}, \tau_{1}, \ldots, \tau_{j^{*}-1}, \tau_{j^{*}+1}, \ldots, \hat{\tau}_{q}, \tau_{j^{*}}-\left\lceil\delta_{T} / 3\right\rceil-2, \tau_{j^{*}}+\left\lceil\delta_{T} / 3\right\rceil+1\right\}
$$

with the corresponding estimated variance of the errors denoted as $\dot{\sigma}_{\hat{k}}^{2}$. So again, instead of adding all the true kinks to the set of estimated kinks as before (which generates the so-called "saturated oracle"), we add all true kinks apart from $\tau_{j^{*}}$, and replace it by $\tau_{j^{*}}-\left(\left\lceil\delta_{T} / 3\right\rceil+2\right)$ and $\tau_{j^{*}}+\left(\left\lceil\delta_{T} / 3\right\rceil+1\right)$.

Note that $\dot{\sigma}_{\hat{k}}^{2}$ is no smaller than the estimated variance of the errors from a model with the features at the same $2 q+1$ locations, but with the continuity constraint only enforced at $\hat{\tau}_{j^{*}}$. More precisely, in the rest of the proof we could effectively follow a model with the signal following Scenario (S2) over $\left\{\tau_{j^{*}}-\left\lceil\delta_{T} / 3\right\rceil-1, \ldots, \tau_{j^{*}}+\left\lceil\delta_{T} / 3\right\rceil+1\right\}$ and Scenario (S3) elsewhere.

In addition, for any $1 \leq s+1<b<e \leq T$,

$$
\begin{array}{r}
\left\|\left.\mathbf{Y}\right|_{(s, e]}-\left\langle\mathbf{Y}, \boldsymbol{\phi}_{(s, e]}^{b}\right\rangle \boldsymbol{\phi}_{(s, e]}^{b}-\left\langle\mathbf{Y}, \boldsymbol{\gamma}_{(s, e]}\right\rangle \boldsymbol{\gamma}_{(s, e]}-\left\langle\mathbf{Y}, \mathbf{1}_{(s, e]}\right) \mathbf{1}_{(s, e]}\right\|^{2} \\
=\left\|\left.\mathbf{Y}\right|_{(s, e]}-\left\langle\mathbf{Y}, \boldsymbol{\gamma}_{(s, e]}\right\rangle \boldsymbol{\gamma}_{(s, e]}-\left\langle\mathbf{Y}, \mathbf{1}_{(s, e]}\right\rangle \mathbf{1}_{(s, e]}\right\|^{2}-\left\langle\mathbf{Y}, \boldsymbol{\phi}_{(s, e]}^{b}\right\rangle^{2} \\
=\left\|\left.\mathbf{Y}\right|_{(s, e]}-\left\langle\mathbf{Y}, \boldsymbol{\gamma}_{(s, e]}\right\rangle \boldsymbol{\gamma}_{(s, e]}-\left\langle\mathbf{Y}, \mathbf{1}_{(s, e]}\right\rangle \mathbf{1}_{(s, e]}\right\|^{2}-\left(\mathcal{C}_{(s, e]}^{b}(\mathbf{Y})\right)^{2}
\end{array}
$$

Applying this result on $s=\tau_{j^{*}}-\left\lceil\delta_{T} / 3\right\rceil-2, e=\tau_{j^{*}}+\left\lceil\delta_{T} / 3\right\rceil+1, b=\tau_{j^{*}}$ or $\hat{\tau}_{j^{*}}$, and using the argument similar to that in Part II, we obtain that

$$
\begin{aligned}
T \hat{\sigma}_{\hat{k}}^{2} \geq T \dot{\sigma}_{\hat{k}}^{2} \geq & \sum_{t=1}^{\tau_{j^{*}}-\left\lceil\delta_{T} / 3\right\rceil-2} \varepsilon_{t}^{2}+\sum_{t=\tau_{j^{*}}+\left\lceil\delta_{T} / 3\right\rceil+2}^{T} \varepsilon_{t}^{2}-(2 q) 6 \log T \\
& +\left(\mathcal{C}_{(s, e]}^{\tau_{j^{*}}}(\mathbf{Y})\right)^{2}-\left(\mathcal{C}_{(s, e]}^{\hat{\mathcal{T}}_{j^{*}}}(\mathbf{Y})\right)^{2}+\left(\sum_{\tau_{j^{*}}-\left\lceil\delta_{T} / 3\right\rceil-1}^{\tau_{j^{*}}+\left\lceil\delta_{T} / 3\right\rceil+1} \varepsilon_{t}^{2}-12 \log T\right)
\end{aligned}
$$

where $\sum_{\tau_{j^{*}}-\left\lceil\delta_{T} / 3\right\rceil-1}^{\tau_{j^{*}}+\left\lceil\delta_{T} / 3\right\rceil+1} \varepsilon_{t}^{2}-12 \log T$ is the lower-bound of the residual sum of squares for fitting a piecewise-linear function over $\left\{\tau_{j^{*}}-\left\lceil\delta_{T} / 3\right\rceil-1, \ldots, \tau_{j^{*}}+\left\lceil\delta_{T} / 3\right\rceil+1\right\}$ with only one feature at $\tau_{j^{*}}$. Consequently, it follows from an argument similar to that in Step Four of the proof of Theorem 1 that

$$
\begin{aligned}
T \hat{\sigma}_{\hat{k}}^{2} \geq & \sum_{t=1}^{T} \varepsilon_{t}^{2}-6(2 q+2) \log T+\left(\mathcal{C}_{(s, e]}^{\tau_{j^{*}}}(\mathbf{f})\right)^{2}-\left(\mathcal{C}_{(s, e]}^{\hat{\tau}_{j^{*}}}(\mathbf{f})\right)^{2} \\
& \quad-2 \sqrt{8 \log T} \sqrt{\left(\mathcal{C}_{(s, e]}^{\tau_{j^{*}}}(\mathbf{f})\right)^{2}-\left(\mathcal{C}_{(s, e]}^{\hat{\tau}_{j^{*}}}(\mathbf{f})\right)^{2}}-8 \log T \\
= & \sum_{t=1}^{T} \varepsilon_{t}^{2}-6(2 q+2) \log T+\left(\sqrt{\left(\mathcal{C}_{(s, e]}^{\tau_{j^{*}}}(\mathbf{f})\right)^{2}-\left(\mathcal{C}_{(s, e]}^{\hat{\tau}_{j^{*}}}(\mathbf{f})\right)^{2}}-\sqrt{8 \log T}\right)^{2}-16 \log T
\end{aligned}
$$

Using the fact that $\left|\hat{\tau}_{j^{*}}-\tau_{j^{*}}\right|<\delta_{T} / 6 \leq\left(\left\lceil\delta_{T} / 3\right\rceil+1\right) / 2$ and Lemma 8 , we have that in the
case where $\frac{C_{1} C_{2}^{2}}{144 T}\left|\hat{\tau}_{j^{*}}-\tau_{j^{*}}\right|^{2} \geq 8 \log T$,

$$
\begin{aligned}
T \hat{\sigma}_{\hat{k}}^{2} & \geq \sum_{t=1}^{T} \varepsilon_{t}^{2}-(12 q+28) \log T+\left(\frac{\underline{C}_{2}}{\sqrt{48} T}\left(\underline{C}_{1} T / 3+1-1\right)^{1 / 2}\left|\hat{\tau}_{j^{*}}-\tau_{j^{*}}\right|-\sqrt{8 \log T}\right)^{2} \\
& =\sum_{t=1}^{T} \varepsilon_{t}^{2}-(12 q+28) \log T+\left(\sqrt{\frac{C_{1} \underline{C}_{2}^{2}}{144 T}}\left|\hat{\tau}_{j^{*}}-\tau_{j^{*}}\right|-\sqrt{8 \log T}\right)^{2}
\end{aligned}
$$

However,

$$
T \hat{\sigma}_{\stackrel{k}{k}}^{2} \leq T \hat{\sigma}_{k^{*}}^{2} \leq \sum_{t=1}^{T} \varepsilon_{t}^{2}+M^{\prime} \log T
$$

Combining the above two inequalities, and after some algebraic manipulations, we get

$$
\left|\hat{\tau}_{j^{*}}-\tau_{j^{*}}\right| \leq \frac{12}{\sqrt{\underline{C}_{1} \underline{C}_{2}}}\left(\sqrt{M^{\prime}+12 q+28}+\sqrt{8}\right) \sqrt{T \log T}=: C \sqrt{T \log T}
$$

On the other hand, in the case where $\frac{C_{1} C_{2}^{2}}{144 T}\left|\hat{\tau}_{j^{*}}-\tau_{j^{*}}\right|^{2}<8 \log T$, we directly have that

$$
\left|\hat{\tau}_{j^{*}}-\tau_{j^{*}}\right|<\frac{24 \sqrt{2}}{\sqrt{\underline{C}_{1} \underline{C}_{2}}} \sqrt{T \log T}<C \sqrt{T \log T}
$$

Therefore, $\mathbb{P}\left(\max _{j=1, \ldots, q}\left|\hat{\tau}_{j}-\tau_{j}\right| \leq C \sqrt{T \log T}\right) \rightarrow 1$, as required.

### 1.6. Proof of Corollary 1

Proof. We set $P:=\sum_{k=-\infty}^{\infty}\left|\rho_{k}\right|$, where $\rho_{k}$ is the auto-correlation function of $\left\{\varepsilon_{t}\right\}$. Now we modify our proof of Theorem 1 as follows:

## Step One and Two

Let $\lambda_{T}=\sqrt{8 P \log T}$ and define the set $A_{T}$ as before. Denote the autocorrelation matrix of $\left\{\varepsilon_{t}\right\}$ by $\mathbf{P}_{T}=\left[\rho_{i-j}\right]_{i, j=1, \ldots, T}$ (which is also the autocovariance matrix, since $\varepsilon_{t}$ has unit-variance). Then since $\mathbf{P}_{T}$ is symmetric, we have that

$$
\left\|\mathbf{P}_{T}\right\|_{\infty}=\left\|\mathbf{P}_{T}\right\|_{1}=\max _{j} \sum_{i}\left|P_{i j}\right| \leq P
$$

where $\|\cdot\|_{\infty}$ and $\|\cdot\|_{1}$ are the operator norms of a matrix. Consequently, by Hölder's inequality, $\left\|\mathbf{P}_{T}\right\|_{2} \leq \sqrt{\left\|\mathbf{P}_{T}\right\|_{1}\left\|\mathbf{P}_{T}\right\|_{\infty}} \leq P$, i.e., the largest eigenvalue of $\mathbf{P}_{T}$ is bounded above by $P$, which is irrelevant of $T$.

For any $s, b, e$ such that $0 \leq s<b<e \leq T$, since $\left\langle\boldsymbol{\psi}_{(s, e]}^{b}, \boldsymbol{\varepsilon}\right\rangle$ has a normal distribution, with zero-mean and

$$
\operatorname{Var}\left(\left\langle\boldsymbol{\psi}_{(s, e]}^{b}, \boldsymbol{\varepsilon}\right\rangle\right)=\left(\boldsymbol{\psi}_{(s, e]}^{b}\right)^{T} \mathbf{P}_{T} \boldsymbol{\psi}_{(s, e]}^{b} \leq P\left\|\boldsymbol{\psi}_{(s, e]}^{b}\right\|_{2}^{2} \leq P,
$$

we have that

$$
\mathbb{P}\left(\left|\mathcal{C}_{(s, e]}^{b}(\varepsilon)\right| \geq \lambda_{T}\right)=\mathbb{P}\left(\left|\mathcal{C}_{(s, e]}^{b}(\varepsilon)\right| / \sqrt{P} \geq \sqrt{8 \log T}\right) \leq \frac{2 e^{-8 \log T / 2}}{\sqrt{8 \log T \sqrt{2 \pi}}}
$$

It follows from the Bonferroni bound that $\mathbb{P}\left(A_{T}^{c}\right) \leq 12 \sqrt{\pi} T^{-1}$.
Using the same argument as above, we can show that for any $0 \leq s<b<e \leq T$, $\frac{\left\langle\boldsymbol{\psi}_{(s, e]}^{b}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e)}^{b}\right\rangle-\boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\right\rangle, \boldsymbol{\varepsilon}\right\rangle}{\left\|\boldsymbol{\psi}_{(s, e]}^{b}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{b}\right\rangle-\boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\left\langle\mathbf{f}, \boldsymbol{\psi}_{(s, e]}^{\tau_{j}}\right\rangle\right\|_{2}}$ is normal distributed, with zero-mean and variance bounded above by $P$. Thus, $\mathbb{P}\left(B_{T}^{c}\right) \leq 12 \sqrt{\pi} T^{-1}$.

Step Three, Four and Five
The rest of the proof goes through by simply changing the constants as
$\underline{C}=\sqrt{6}\left(2 \sqrt{C_{3}}+\sqrt{32 P}\right)+1, \quad C_{1}=2 \sqrt{C_{3}}+\sqrt{8 P}, \quad C_{2}=\frac{1}{\sqrt{6}}-\frac{\sqrt{8 P}}{\underline{C}}, \quad C_{3}=(32 \sqrt{2}+48) P$
and setting

$$
\eta_{T}=\left(C_{1}-\sqrt{8 P}\right)^{2}
$$

Finally, we remark that the proof of Corollary 2 is similar to that of Corollary 1, so is omitted for brevity.

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