Learning from Censored and Dependent Data: 
The case of Linear Dynamics

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Abstract

Observations from dynamical systems often exhibit irregularities, such as censoring, where values are recorded only if they fall within a certain range. Censoring is ubiquitous in practice, due to saturating sensors, limit-of-detection effects, image frame effects, and combined with temporal dependencies within the data, makes the task of system identification particularly challenging.

In light of recent developments on learning linear dynamical systems (LDSs), and on censored statistics with independent data, we revisit the decades-old problem of learning an LDS, from censored observations (Lee and Maddala (1985); Zeger and Brookmeyer (1986)). Here, the learner observes the state $x_t \in \mathbb{R}^d$ if and only if $x_t$ belongs to some set $S_t \subseteq \mathbb{R}^d$. We develop the first computationally and statistically efficient algorithm for learning the system, assuming only oracle access to the sets $S_t$. Our algorithm, Stochastic Online Newton with Switching Gradients, is a novel second-order method that builds on the Online Newton Step (ONS) of Hazan et al. (2007). Our Switching-Gradient scheme does not always use (stochastic) gradients of the function we want to optimize, which we call censor-aware function. Instead, in each iteration, it performs a simple test to decide whether to use the censor-aware, or another censor-oblivious function, for getting a stochastic gradient.

In our analysis, we consider a “generic” Online Newton method, which uses arbitrary vectors instead of gradients, and we prove an error-bound for it. This can be used to appropriately design these vectors, leading to our Switching-Gradient scheme. This framework significantly deviates from the recent long line of works on censored statistics (e.g. Daskalakis et al. (2018); Kontonis et al. (2019); Daskalakis et al. (2019)), which apply Stochastic Gradient Descent (SGD), and their analysis reduces to establishing conditions for off-the-shelf SGD-bounds. Our approach enables to relax these conditions, and gives rise to phenomena that might appear counterintuitive, given the previous works. Specifically, our method makes progress even when the current “survival probability” is exponentially small. We believe that our analysis framework will have applications in more settings where the data are subject to censoring.

Keywords: Truncated/Censored Data, Linear Dynamical Systems, Time-Series, Online Learning

1. Introduction

System identification is the problem of learning the evolution equations of a dynamical system from data. Mathematically, we have a sequence of system states $(x_t)_t$, and observations $(y_t)_t$ evolving as

$$x_{t+1} = f(x_t, u_t, w_t), \quad y_t = g(x_t, v_t),$$

where $u_t$'s are inputs to the system, $w_t$ is process noise, and $v_t$ is sensor noise. At each step $t$, the learner observes the input $u_t$, and the resulting output $y_{t+1}$. The goal is to learn the functions $f$...
and $g$, from an observed trajectory.\footnote{Some works consider access to multiple trajectories.} In this paper, we consider system identification with \textit{censored} observations, where the learner observes $y_t$ if and only if $y_t$ belongs to some set $S_t$, which we call the \textit{observable} set.

Dynamical data with missing observations are ubiquitous in practice (Honaker and King, 2010). This is often due to censoring, which frequently manifests in the fields of signal processing and control theory (Yang and Li, 2009), time series analysis (Lee and Maddala, 1985), business (Hausman and Wise, 1977), economics (Johannsen and Mertens, 2018), in medicine, and in physical sciences. For example, consider learning the dynamics of some target dynamical system, using observations from cameras. Here, censoring naturally arises due to occlusions blocking visibility, or from the system exiting the camera frame. Despite the numerous applications, the problem is not statistically and computationally understood even for linear, fully-observable dynamics:

$$x_{t+1} = A_\ast x_t + B_\ast u_t + w_t, \quad y_t = x_t,$$

where $x_t \in \mathbb{R}^d$, $u_t \in \mathbb{R}^m$, and $w_t \sim i.i.d. N(0, I)$. Actually, even for one-dimensional linear systems ($d = 1$), with no inputs ($u_t = 0$), and observable sets $S_t$ being half-lines, no known efficient algorithm learns (the scalar) $A_\ast$, not even asymptotically. Specifically, the existing methods fall in two categories: 1) iterative methods trying to maximize the non-concave log-likelihood, as in Lee and Maddala (1985) and Zeger and Brookmeyer (1986), and 2) EM-based imputation methods (e.g., Park et al. (2007)). Both of them are not guaranteed to recover the underlying system. On top of that, for large dimension $d$, the likelihood-based approach is inefficient to implement because it requires computing high-dimensional integrals over $S_t$'s.

In this work, we study the multidimensional linear case with no inputs: $x_{t+1} = A_\ast x_t + w_t$. We allow \textit{arbitrary} observable sets $S_t$, with only requirement being that we have oracle access to each $S_t$, namely, given a point $x$, the oracle efficiently computes $\mathbb{1}\{x \in S_t\}$. For this model, we obtain the first computationally and statistically efficient algorithm for learning $A_\ast$, under the following assumptions (stated informally here, and in detail in Section 3):

\textbf{Assumption 1:} The system is \textit{stable}, i.e., the spectral radius $\rho(A_\ast)$ is less than one. Stability is a classical assumption in linear dynamical systems (LDSs). If $\rho(A_\ast) > 1$, then the state explodes exponentially, with high probability.

\textbf{Assumption 2:} The number of times we observe the state is at least a constant fraction (say 1\%) of the trajectory-length.

\textbf{Assumption 3:} For most of the timesteps $t$, given that we observed $x_t$, the probability of observing $x_{t+1}$ is at least a constant (say 1\%).

To motivate the last two assumptions, we note that for the simpler problems of censored Gaussian estimation and censored linear regression (with independent data), the only known computationally and statistically efficient algorithms assume that the probability of observing a sample is at least a constant (see Daskalakis et al. (2018), Daskalakis et al. (2019)). Assumptions 2 and 3 are the adaption of this, to fit our dynamical setting. Under these assumptions, our estimation error bound matches (up to logarithmic factors) the best known bound for learning uncensored LDSs (Simchowitz et al. (2018)), with respect to the dimension $d$, the trajectory-length $T$, and the spectrum...
of $A_\ast$. It is also the first, to the best of our knowledge, estimator that can accommodate arbitrary observable sets, since previous works considered intervals, half-lines, and products of these (e.g., Zeger and Brookmeyer (1986), Yang and Li (2009)). Key for tackling this decades-old problem are recent advances in 1) (uncensored) linear system identification theory, and 2) censored/truncated statistics (CTS) for independent data.

Learning Linear Dynamics. The difficulty in learning an LDS, compared to standard linear regression, is that the observations are dependent. A classical approach that avoids this issue is based on the system’s mixing time, which is roughly $\tau_{\text{mix}} = \frac{1}{1 - \rho(A_\ast)}$ steps. Thus, the learner can use one $x_t$ every $\tau_{\text{mix}}$ steps, and reduce the analysis to standard linear regression (e.g., Yu (1994)). However, the resulting bounds get worse with larger $\tau_{\text{mix}}$, and as pointed in Simchowitz et al. (2018), this behavior is qualitatively incorrect. Intuitively, larger $A_\ast$ gives larger states, which implies larger signal-to-noise ratio, i.e., easier estimation. The authors provided bounds that express this intuition, and that was the first sharp analysis for stable systems.\footnote{The authors also considered $\rho(A_\ast) = 1$, known as the marginally stable regime.} This progress initiated a large line of work on learning LDSs (see Section 1.1).

Censored/Truncated Statistics (CTS). Consider the problem of learning a Gaussian distribution, having access only to samples from some set $S$. These are called truncated samples. Censoring is when we also know the number of unobserved samples, as in our case, where this number can be inferred from the lengths of time-intervals during which we do not observe anything. Truncation and censoring go back to at least Galton (1898) and Pearson (1902), and there has been a large volume of research devoted to them (see Cohen (1991)). Nevertheless, the first provably computationally and statistically efficient algorithm, for learning a truncated Gaussian, was only recently discovered by Daskalakis et al. (2018). The authors developed a general algorithmic framework, based on Stochastic Gradient Descent (SGD), which bypasses the computation of high-dimensional integrals over $S$. The result and the generality of the approach created a lot of excitement, and a large number of subsequent works applied the SGD framework to other problems in CTS (for independent data), e.g., linear regression by Daskalakis et al. (2019).

Our Contributions

We build on the above advances, by introducing new algorithmic and technical ideas, which we now overview.

Our algorithm. Our first observation is that non-convexity of the negative log-likelihood can be bypassed by focusing on “paired observations”, i.e., $(x_t, x_{t+1})$ such that $x_t \in S_t$ and $x_{t+1} \in S_{t+1}$. By ignoring the other terms in the objective, we get a convex function. The second observation is that for LDSs, if an SGD-based algorithm relies on the Markovian property $\mathbb{E}[x_{t+1}|x_t, x_{t-1}, \ldots, x_1] = A_\ast x_t$ to produce unbiased gradient estimates, then the algorithm should process the data in temporal order. The reason is that if we “see” $x_t$ and $x_{t+s}$ (for some $s \geq 2$), then the expectation of $x_{t+1}$ is not $A_\ast x_t$, i.e., $\mathbb{E}[x_{t+1}|x_t, x_{t+s}] \neq A_\ast x_t$. In other words, we need an online algorithm. Unfortunately, for censored linear regression, the SGD framework of Daskalakis et al. (2019) processes the data in random order, and also does multiple passes over them. This is not a technicality; making that algorithm online will lead to slow statistical rates, as we explain later. For this reason, we design a new stochastic second-order method, building on the Online Newton Step method (ONS) of Hazan.
et al. (2007). The crucial difference with ONS, and all recent works on CTS, is that we do not always use (stochastic) gradients of the function we want to optimize, which we call censor-aware function. Instead, in each iteration, we perform a simple test, based on which we decide whether to use the censor-aware function, or another censor-oblivious function, for getting a stochastic gradient. We call our method Stochastic Online Newton with Switching Gradients (SON-SG, given in Algorithm 1), and we show that it can be combined with a least-squares-based warmup procedure, to learn censored LDSs (Algorithm 4). We also show that SON-SG can be applied to an even broader setting, i.e., general linear-response time-series (Section 4).

Algorithm-design framework. We came up with this switching-gradient scheme, by considering a “generic” ONS method, where instead of gradients we have an arbitrary vector sequence \( g_1, g_2, \ldots \). We proved a general estimation error bound for this method, which serves as a guideline for designing the \( g_i \)’s. A similar framework for designing \( g_i \)’s has been proposed in the context of non-convex optimization (for first-order methods), by Arora et al. (2015). This approach gives a lot of freedom for algorithm design, compared to previous works on CTS which apply off-the-shelf SGD-bounds, and we believe it will have more applications in censored and truncated statistics.

Technical contributions and insights. In all recent CTS papers (see Section 1.1 for an extensive list), a crucial step for proving parameter recovery by SGD is establishing anti-concentration of truncated Gaussians \( \mathcal{N}(\mu, \Sigma, S) \), where \( \mathcal{N}(\mu, \Sigma) \) is the Gaussian \( \mathcal{N}(\mu, \Sigma) \), given that the sample is in \( S \). Daskalakis et al. (2018) reduced this task to showing that the “survival probability” is large, i.e., \( \mathbb{P}_{x \sim \mathcal{N}(\mu, \Sigma)}[x \in S] \geq \Omega(1) \). Hence, this and all follow-up works have focused on lower bounding survival probabilities. In our case, this methodology does not apply, due to temporal dependencies. However, our approach (generic ONS bound) enables to significantly relax the high survival probability condition, and gives rise to phenomena that might appear counterintuitive, given the intuition built in the recent literature. Specifically, in high dimensions \( (d \to \infty) \) our analysis deals with cases where the survival probability is exponentially small \( (e^{-\Omega(d)}) \), while at the same time anti-concentration tends to infinity. Our other technical contribution is a lower-bound on the covariance matrix of the observed states. Simchowitz et al. (2018) proved such a bound when all states are observed. Here, due to censoring, we observe only a subset of the whole trajectory. For independent data, Daskalakis et al. (2018) address this issue using union-bound over all possible subsets. Unfortunately, for LDSs union-bound gives vacuous guarantees, because for a fixed subset, the concentration degrades with larger mixing times. We resolve this difficulty, by generalizing the “small-ball” technique of Simchowitz et al. (2018), and we lower-bound the covariance matrices of all subsets of size \( \Omega(T) \) simultaneously, where \( T \) is the trajectory-length.

1.1. Further Related Work

Censored and Truncated Statistics. As we mentioned, there has been a long line of recent works on several settings within CTS: Gaussian parameter estimation (Daskalakis et al., 2018; Kontonis et al., 2019), linear, logistic and probit regression (Daskalakis et al., 2019; Ilyas et al., 2020), compressed sensing (Daskalakis et al., 2020), sparse graphical models (Bhattacharyya et al., 2020), estimation of boolean product distributions (Fotakis et al., 2020), mixtures of Gaussians (Nagarajan

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3. This is implicit in their analysis. They use a result of Diakonikolas et al. (2019), which applies union bound over subsets.
and Panageas, 2020). All these works consider independent data and apply the SGD framework introduced in Daskalakis et al. (2018).

**Linear System Identification.** Even though linear system identification is a decades-old field (Ljung, 1999), a sharp non-asymptotic theory was only recently developed (Simchowitz et al., 2018; Sarkar and Rakhlin, 2019; Tsiamis and Pappas, 2019; Oymak and Ozay, 2019; Simchowitz et al., 2020; Ghai et al. (2020); Ljung, 1999), a sharp non-asymptotic theory was only recently developed (Simchowitz et al., 2018; Sarkar and Rakhlin, 2019; Tsiamis and Pappas, 2019; Oymak and Ozay, 2019; Simchowitz et al., 2020).

**Online Convex Optimization.** To design our algorithm, we build on ideas from online convex optimization (OCO). In the recent years, OCO has been extensively used for learning and controlling LDSs (e.g., Agarwal et al. (2019); Hazan et al. (2020); Simchowitz et al. (2020); Ghai et al. (2020); Simchowitz et al. 2020). For a general overview of OCO see Hazan (2019).

2. Notation

For every vector \( x \), we use \( \| x \| \) to denote the \( \ell_2 \) norm \( \| x \|_2 \). Also, we use \( \langle A, B \rangle = \text{tr}(A^\top B) \) to denote the matrix inner product. For a matrix \( A \) and a \( \Sigma \succ 0 \), we define the norm \( \| A \|_{\Sigma} = \sqrt{\langle A^\top A, \Sigma \rangle} \). The covariance matrix between two random vectors \( x, y \) is \( \text{Cov}(x, y) \). For a sequence \( (x_t)_{t=1}^\infty \), we denote by \( x_{<\tau} \), \( x_{\tau} \) and \( x_{-\tau} \) the subsequences \( (x_t)_{t<\tau}, (x_t)_{t<\tau}, (x_t)_{t\neq \tau} \) respectively.

For a Gaussian distribution \( \mathcal{N}(\mu, \Sigma) \) in \( \mathbb{R}^d \), and a measurable \( \mathcal{S} \subseteq \mathbb{R}^d \), we define the survival probability \( \mathcal{N}(\mu, \Sigma; \mathcal{S}) = \Pr_{x \sim \mathcal{N}(\mu, \Sigma)}[x \in \mathcal{S}] \). We also define the truncated Gaussian \( \mathcal{N}(\mu, \Sigma, \mathcal{S}) \) to be \( \mathcal{N}(\mu, \Sigma) \) conditioned on taking values in \( \mathcal{S} \). Finally, whenever we say that a set \( \mathcal{S} \) is “revealed” to the learner, we mean that she has access to a membership oracle \( M_\mathcal{S} \), i.e., an efficient procedure that computes the \( 1 \{ x \in \mathcal{S} \} \), for any point \( x \).

3. Censored Linear Dynamics: Model, and Main Theorem

We study the system \( x_{t+1} = A_s x_t + w_t \), where \( x_t \in \mathbb{R}^d \), \( A_s \in \mathbb{R}^{d \times d} \) and \( w_t \overset{i.i.d.}{\sim} \mathcal{N}(0, I) \). Starting from \( x_0 = 0 \), consider the trajectory \( x_1, x_2, \ldots, x_{T+1} \). The learner has access to censored observations, i.e., there is a process of observable sets \( (\mathcal{S}_t)_t \) and she observes \( x_t \) if and only if \( x_t \in \mathcal{S}_t \). Also, at time \( t \), the set \( \mathcal{S}_t \) is revealed to her. Now, \( \mathcal{S}_t \)’s may depend on the state-trajectory, but we assume that given \( x_t \), the state \( x_{t+1} \) and the set \( \mathcal{S}_{t+1} \) are statistically independent. To see why this is a natural assumption, consider the camera-based example (Section 1), and think of \( x_t \in \mathbb{R}^3 \) as the position of some object, and \( \mathcal{S}_t \) as the visible part of the space at time \( t \). Having observed \( x_t \), the camera could adapt its frame (affecting \( \mathcal{S}_{t+1} \)), to improve the chances for observing \( x_{t+1} \), but without knowing the next “excitation” \( w_t \). We now formally state our assumptions, sketched in Section 1.

**Assumption 1** \( A_s \) is diagonalizable and stable, i.e., \( A_s = U D U^{-1} \), where \( D \) is diagonal and \( \rho(A_s) = \max_i |D_{ii}| < 1 \).

Let \( \mathcal{O} \) be the set of observation times: \( \mathcal{O} = \{ t \in [T] : x_t \in \mathcal{S}_t \} \). We assume that we observe “enough” data:

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4. Exception is Nagarajan and Panageas (2020) who consider the EM-algorithm.
5. Usually, \( \| A \|_\Sigma \) is used to denote \( \sqrt{\langle AA^\top, \Sigma \rangle} \). However, this definition is more appropriate for our setting.
6. We assume \( x_0 = 0 \) to simplify the exposition. Our proofs generalize for any \( x_0 \), by paying a \( \log \| x_0 \| \) factor in the bound.
7. Note that \( \mathcal{S} \) and \( D \) can have complex entries.
Assumption 2  For a known constant $\beta \in (0, 1)$, with probability $1 - o(1)$, we have $|O| \geq \beta T$, where $o(1)$ denotes a $\delta_T \to 0$, as $T \to \infty$.

Let $B(a)$ be the set of timesteps $t$, at which 1) we observe $x_t$, and 2) given $x_t$ and $S_{t+1}$, the probability of observing $x_{t+1}$ is less than $a$, i.e., $B(a) = \{ t \in O : N(A_s x_t, I; S_{t+1}) < a \}$.

Assumption 3  For a known constant $\alpha \in (0, 1)$, and some bound $L > 0$, with probability $1 - o(1)$, we have $|B(\alpha)| \leq L$. Also, $\mathbb{E}[|B(\alpha)|] \leq L$.

Our bounds will depend on $L$, and will match the uncensored case if $L \leq \tilde{O}(d)$. As we mentioned in the introduction, Assumptions 2 and 3 are the adaptation of the $\Omega(1)$-survival-probability assumption in Daskalakis et al. (2018). We further motivate Assumption 3 with a natural one-dimensional example.

Example 1  Let $x_{t+1} = a_s x_t + w_t$, where $a_s \in [0, 1)$. The observable set is a static half-line: $S_t = S = \{ x \in \mathbb{R} : x \geq \lambda \}$, $\lambda > 0$. We claim that here, Assumption 2 implies Assumption 3 with $\alpha = \Omega(1)$ and $L = 0$. This is clear for $a_s \approx 1$, since if $x_t \geq \lambda$, then $a_s x_t + w_t \geq x_t \geq \lambda$ with probability almost $1/2$. For general $a_s \in [0, 1)$, the implication is less obvious (see Appendix B).

As in Simchowitz et al. (2018), our bounds depend on the controllability Gramian $\Gamma_T$, defined as $\Gamma_T := \sum_{s=0}^{T-1} A_s^\top A_s$. This matrix quantifies how much the noise process excites the system. In the theorem that follows, we use $\hat{\Theta}(1)$ to denote polylogarithmic factors in $T$ and in $\text{cond}(U)$, where $\text{cond}(U)$ is the condition number of the eigenvector-matrix $U$.

Theorem 1  Under Assumptions 1, 2, 3, there exist $C_{\alpha, \beta}, C_{\alpha, \beta}^\prime = \hat{\Theta}(1) \cdot \text{poly}\left(\frac{1}{\alpha \beta}\right)$, such that if

$$T \geq C_{\alpha, \beta}^\prime \cdot \left( d^2 + \frac{d}{1 - \rho(A_s)} + dL \right),$$

then with probability at least 99%, Algorithm 4 runs in polynomial time, and outputs an $\hat{A}$ such that

$$\| \hat{A} - A_s \|_{\Gamma_T} \leq C_{\alpha, \beta} \sqrt{\frac{d^2 + dL}{T}}. \tag{1}$$

Remark 2  For uncensored LDSs, and stable-diagonalizable $A_s$, the best known (spectral-norm) bound is $\| (\hat{A} - A_s) \|_{\Gamma_T}^{1/2} \leq \tilde{O} \left( \sqrt{d/T} \right)$. So, the Frobenius version is $\| \hat{A} - A_s \|_{\Gamma_T} \leq \tilde{O} \left( \sqrt{d^2/T} \right)$, which matches (1), if $\alpha, \beta = \Omega(1)$, and $L = \tilde{O}(d)$. We leave the spectral-norm bound for censored LDSs for future work.

Remark 3  It is possible to drop the assumption that $A_s$ is diagonalizable by paying an exponential dependence in the size of the largest Jordan block. This dependence (for a bound on $\| \cdot \|_{\Gamma_T}$) appears even in the uncensored case (see Ghai et al. (2020) for a discussion on this).

As we mentioned, our algorithm is SON-SG, preceded by a least-squares warmup procedure. Let $P$ be the set of of pairs $(x_t, x_{t+1})$, such that both $t, t+1 \in O$. We will ignore all “isolated” observations, i.e., the ones not participating in any pair of $P$. As we will see, by ignoring them this decision is justified by the following proposition:

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8. $\tilde{O}(\cdot)$ hides logarithmic factors.

9. Note that for $a_s \approx 1$, we did not need Assumption 2, to show Assumption 3. However, we need it for general $a_s$. 

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Proposition 4  Let $M := |\mathcal{P}|$. With probability $1 - o(1)$, we have $M \geq \alpha \beta T/2$.

We prove Proposition 4 in Appendix C, using Assumptions 2, 3, and that $T >> L$. Let $\mathcal{P}_0$ be the first (in time-order) $\lfloor M/2 \rfloor$ pairs of $\mathcal{P}$, and $\mathcal{P}_1 = \mathcal{P} \setminus \mathcal{P}_0$. Our algorithm first uses $\mathcal{P}_0$ to create a (rough) confidence ellipsoid $K$ that includes $A_*$, with high probability. Then, it uses $K$ as constraint-set for SON-SG, which will operate on $\mathcal{P}_1$. In the next section, we present and analyze SON-SG in a more general setting, which reveals the key structure that our method exploits. Then, in Section 6, we give the full-algorithm and conclude the proof of Theorem 1.

4. Truncated Time Series with Linear Responses

For timesteps $1 \leq t \leq T$, consider a covariate-response process $(x_t, y_t)_t$, where $x_t \in \mathbb{R}^d$, $y_t \in \mathbb{R}^n$. The learner only observes a subset of the data, based on a Bernoulli process $(o_t)_t$, i.e., if $o_t = 1$, then she observes $(x_t, y_t)$, otherwise the pair is hidden. In untruncated linear-response time series, we have $y_t = A_* x_t + w_t$, where $w_t \sim \mathcal{N}(0, I)$, and $A_* \in \mathbb{R}^{n \times d}$. Here, when the learner gets to see a datapoint, the noise will be biased due to truncation. Formally, consider a process of observable sets $(S_t)_t$, where $S_t \subseteq \mathbb{R}^n$. Let $\mathcal{F}_t$ be the $\sigma$-algebra generated by $x_{\leq t}$, $y_{< t}$, $o_{\leq t}$ and $S_{\leq t}$ (note that $y_t$ is not in $\mathcal{F}_t$). We assume that given $\mathcal{F}_t$ and $o_t = 1$, the set $S_t$ is revealed to the learner, and $y_t \sim \mathcal{N}(A_* x_t, I, S_t)$. This model has three notable special cases:

1. **Untruncated, linear-response time series.** This model was studied by Simchowitz et al. (2018), and corresponds to $o_t = 1$ and $S_t = \mathbb{R}^n$, for all $t$.

2. **Truncated linear regression.** First considered by Tobin (1958), this model was revisited by Daskalakis et al. (2019), and corresponds to $o_t = 1$ for all $t$, and it requires independent data, i.e., given $x_t$, the response $y_t$ is independent of $x_{-t}$.\(^{10}\)

3. **LDS with censored observations.** Here, $y_{t} = x_{t+1}$, $o_t = 1 \{ x_t \in S_t \land x_{t+1} \in S_{t+1} \}$, and $S_t = S_{t+1}$. Note that we pretend we do not observe the “isolated” observations, which is aligned with what our algorithm does.

**Initial confidence ellipsoid.** SON-SG receives as input a rough initial estimate $A_0 \in \mathbb{R}^{n \times d}$, and a $\Sigma_0 > 0$ that represents a confidence ellipsoid $K = \{ A \in \mathbb{R}^{n \times d} : \| A - A_0 \|_{\Sigma_0} \leq 1 \}$. For LDSs, we will later show how to use $\mathcal{P}_0$ to construct a $K$ with the following property:

**Definition 5** A confidence ellipsoid $K = \{ A \in \mathbb{R}^{n \times d} : \| A - A_0 \|_{\Sigma_0} \leq 1 \}$ is $(R, \omega)$-accurate, for some $R, \omega > 0$, if (a) $A_* \in K$, (b) $\Sigma_0 \succeq \omega \cdot I$, and (c) for all $t$, $\| \Sigma_0^{-1/2} x_t \| \leq R^2$-subgaussian.\(^{11}\)

**Theorem 6** Fix $R, R_w, R_x, L, \alpha, \delta > 0$, and let $B(\alpha) := \{ t \in [T] : o_t = 1, N(A_* x_t, I, S_t) < \alpha \}$. Suppose that (a) we are given an $(R, \omega)$-accurate confidence ellipsoid $K$, (b) $\mathbb{E}[\| B(\alpha) \|_2] \leq L$, (c) for all $t$, the noise $w_t = y_t - A_* x_t$ has norm $\| w_t \|$ that is $R_w^2$-subgaussian, and (d) $\mathbb{E}[\| x_t \|^2] \leq R_x^2$. Let $t_1 < t_2 < \cdots < t_N$ be the observation times ($o_{t_i} = 1$). If the total number of steps $T \geq \text{poly}(1/\alpha)$, then SON-SG (Algorithm 1) outputs an $\hat{A}$ such that

\[
\mathbb{E} \left[ \left\| \hat{A} - A_* \right\|_{\Sigma}^2 \right] \leq \frac{D + LD'}{N},
\]

\(^{10}\) To be precise, Daskalakis et al. (2019) consider one-dimensional responses and fixed truncation set. However, the extension to multidimensional $y_t$’s and time-varying truncation sets is relatively straightforward.

\(^{11}\) A random variable $X$ is called $\sigma^2$-subgaussian if $\mathbb{P}[|X| \geq \delta \cdot \sigma] \leq \exp(-\delta^2/2C)$, where $C = O(1)$. 

where $\Sigma = \frac{1}{N} \sum_{i=1}^{N} x_i x_i^T$, and $D = \tilde{O}(1) \cdot dD'$, $D' = \tilde{O}(1) \cdot \text{poly}(1/\alpha) \cdot (d + R^2 + R^2_w)$. Here, $\tilde{O}(1)$ denotes a polylogarithmic factor in the parameters defined in this section.

**Remark 7** Note that in the LDS case, $R_w = O(\sqrt{d})$ by standard concentration of Gaussian norm. However, for slowly mixing systems ($\rho(A_x) \to 1$), $R_x$ can grow polynomially with $T$. Our bounds only degrade in the logarithm of $R_x$, which is absorbed in $\tilde{O}(\cdot)$, and the same happens with $1/\omega$.

Before presenting SON-SG and its analysis, we give some background on existing techniques for CTS.

### 4.1. Existing Techniques and their Limitations

Even though Daskalakis et al. (2019) consider truncated linear regression with no temporal dependencies, we will use the same likelihood-based objective. Specifically, let

$$
\ell_S(\mu; y) := -\frac{1}{2} \| y - \mu \|^2 - \log \left( \int_S \exp \left( -\frac{1}{2} \| z - \mu \|^2 \right) dz \right),
$$

and observe that given $F_t$ and that $o_t = 1$, we have that for a candidate matrix $A$, the log-likelihood for $y_t$ is $\ell_S(Ax_t; y_t)$. Also, let $f_t(A)$ be the (negative) population log-likelihood:

$$
f_t(A) := -\mathbb{E}_y \left[ \ell_S(Ax_t; y) \mid F_t, o_t = 1 \right],
$$

and observe that $f_t(A) = -\mathbb{E}_{y \sim \mathcal{N}(Ax_t, I, S_t)} \left[ \ell_S(Ax_t; y) \right]$. Given the observed data, our goal will be to minimize $f(A) := \frac{1}{N} \sum_{i=1}^{N} f_t(A)$. To see why $f(A)$ is a “good” objective, we take the first and second derivatives of $f_t(A)$:

$$
\nabla f_t(A) = \mathbb{E}_{z \sim \mathcal{N}(Ax_t, I, S_t)} [z] x_t^T - \mathbb{E}_{y \sim \mathcal{N}(A_x x_t, I, S_t)} [y] x_t^T
$$

$$
\nabla^2 f_t(A) = \text{Cov}_{z \sim \mathcal{N}(Ax_t, I, S_t)} [z, z] \otimes \left( x_t x_t^T \right), \tag{3}
$$

where $\nabla^2 f_t(A)$ is used to denote $\nabla^2 f_t (\text{vec} (A^T))$, and vec($\cdot$) is the standard vectorization. Now, note that $\nabla f_t(A_x) = 0$, and so $\nabla f(A_x) = 0$. Also, since the Kronecker product of positive semidefinite matrices (PSD) is PSD, $f_t(A)$ is convex, and so $f(A)$ is also convex. Now, if $f$ was strongly-convex ($\nabla^2 f(A) \succ 0$), then $A_x$ would be the unique optimal solution, justifying the use of the objective. Even though strong-convexity does not hold, Daskalakis et al. (2019) show how to address this for independent data, by restricting $A$ in some set that contains $A_x$, and $\nabla^2 f(A) \succeq \Omega(1) \cdot I$, inside the set. Let’s assume (for now) that here, $\nabla^2 f(A) \succeq \Omega(1) \cdot I$ holds. Now, note that a priori, it is not clear how to optimize $f(A)$, since we do not have a closed-form expression. An important conceptual contribution of Daskalakis et al. (2019) is the observation that by sampling $z_t \sim \mathcal{N}(A x_t, I, S_t)$, and computing $v_t = (y_t - z_t) x_t^T$, we have $\mathbb{E} \left[ v_t \mid F_t, o_t = 1 \right] = \nabla f_t(A)$, i.e,

---

12. $\otimes$ denotes Kronecker product.
13. As we mentioned in the introduction, for censored LDSs, the negative log-likelihood is non-convex. However, here we have convexity, because $f(A)$ corresponds to a part of the overall log-likelihood.
14. This is done via rejection sampling, using the membership oracle for $S_t$. 
we get an unbiased gradient estimate. Based on this observation, they employ a variant of SGD to minimize \( f(A) \), but their algorithm is tailored to independent data. The reason is that it processes the data in random order, and also does multiple passes over them. Thus, by the time it computes \( v_t \), it is very likely that before that, it had processed \((x'_t, y'_t)\), for \( t' > t \). Because of this, if the data are temporally dependent, \( v_t \) can be a biased estimate of \( \nabla f_t(A) \).

4.2. Our Approach

To avoid the above issue, we need to process the data in temporal order. This is exactly the case for Online Convex Optimization (OCO). In OCO though, the goal is not to recover some parameter, but to minimize regret. However, regret bounds can often be transformed to statistical recovery rates via “online-to-batch” conversions (e.g., Cesa-Bianchi et al. (2004)). In our setting, this conversion can be done, but is trickier than usual, and we will deal with it later. Now, since we are aiming for a fast \( \mathcal{O}(1/N) \)-rate, the first natural attempt is online SGD, which has only logarithmic (in \( N \)) regret, provided all \( f_t \)'s are strongly-convex (Hazan (2019))\(^{15}\). Unfortunately, this is not true for any \( f_t \), due to the rank-one component \( x_t x_t^\top \) in the Hessian \( \nabla^2 f_t(A) \) (3). However, there is still structure we can exploit. Notice that every row of \( \nabla f_t(A) \) has the same direction as \( x_t \), which corresponds exactly to that “problematic” rank-one component in the Hessian. This structure is reminiscent of the exp-concavity property (Hazan et al. (2007)), which essentially is strong-convexity, in the direction of the gradient:

**Definition 8** A function \( f \) is called \( \lambda \)-exp-concave, if for all \( x \), \( \nabla^2 f(x) \succcurlyeq \lambda \cdot \nabla f(x) \nabla f(x)^\top \).\(^{16}\)

For exp-concave functions, the Online Newton Step (ONS) algorithm, introduced in Hazan et al. (2007), has regret that depends logarithmically in \( N \). Unfortunately, this result does not apply here:

**Obstacles for ONS.** First, \( f_t \)'s are not necessarily exp-concave (unless \( \lambda \) is exponentially small, which is not useful). This is because of the covariance term in 3. To the best of our knowledge, the idea used in Daskalakis et al. (2019) to restrict \( A \) is some appropriate set, does not resolve this issue, due to temporal dependencies. The second obstacle, is that the regret bound of ONS (Hazan et al. (2007)) will have linear dependence in \( R_\varepsilon \), which as we said can grow as \( \text{poly}(T) \).

4.3. Stochastic Online Newton with Switching Gradients

We now describe SON-SG. First, we use as projection-set the ellipsoid \( \mathcal{K} \). Second, we use pre-conditioning as in ONS, but while in ONS the preconditioner has outer-products of the gradients, here we use outer-products of the covariates. This is done to simplify the analysis.

**Switching Gradients.** The crucial difference with ONS is the choice of the \( g_t \)'s, by the “Switch-Grad” function. To ease notation, suppose we are at time \( t \), and \( t = t_i \) for some \( i \). We define \( A(t) = A_i \), and \( g(t) = g_i \). In ONS, \( g(t) \) would simply be \( \nabla f_t(A(t)) \). Of course, we do not have access to this gradient, but as we said we can get a stochastic gradient by sampling \( z_t \sim \mathcal{N}(A(t)x_t, I, S_t) \), and setting \( g(t) = (z_t - y_t)x_t^\top \). Sampling from \( \mathcal{N}(A(t)x_t, I, S_t) \) can be done via rejection sampling, using the membership oracle. However, to be efficient, the mass \( \gamma_t := \mathcal{N}(A(t)x_t, I; S_t) \) should be sufficiently large. Unfortunately, here \( \gamma_t \) can be exponentially small.

\(^{15}\) In online-to-batch conversions, a regret-bound \( R_N > 0 \) often translates to \( R_N/N \) statistical recovery rate.

\(^{16}\) These functions are called exp-concave, because this property is equivalent to \( e^{-\lambda f(x)} \) being concave.
Algorithm 1 Stochastic Online Newton with Switching Gradients

**Input:** $A_0 \in \mathbb{R}^{n \times d}$, PSD matrix $\Sigma_0 \in \mathbb{R}^{d \times d}$, data $(x_{t_i}, y_{t_i})_{i=1}^N$.

$\eta = (2/\alpha)^c_0$ $\triangleright c_0 \geq 0$ is a large constant.

$K = \{ A \in \mathbb{R}^{n \times d} : \| A - A_0 \|_{\Sigma_0} \leq 1 \}$

$A_1 = A_0$

For $i = 1$ to $N$

$g_i = \text{SwitchGrad}(A_i x_{t_i}, x_{t_i}, y_{t_i}, S_{t_i})$

$\Sigma_i = \Sigma_{i-1} + x_{t_i}x_{t_i}^\top$

$\tilde{A}_{i+1} = A_i - \eta \cdot g_i \Sigma_i^{-1}$

$A_{i+1} = \arg\min_{A \in K} \| A - \tilde{A}_{i+1} \|_{\Sigma_i}$

End

Return $\hat{A} = A_{N+1}$

Algorithm 2 SwitchGrad

**Input:** $\mu, x, y, S$

$z = \mu$

If $\text{Test}(\mu, S)$ then

Sample $z' \sim N(\mu, I, S)$ via rejection sampling using the membership oracle $M_S$.

$z = z'$

End

Return $g = (z - y)x^T$

Algorithm 3 Test

**Input:** $\mu, S$.

$\gamma = (\alpha/2)^c_\gamma$, $k = \frac{4}{\gamma} \log T$. $\triangleright c_\gamma \geq 0$ is a large constant.

Sample $\xi_1, \ldots, \xi_k \sim N(\mu, I)$

$p = \frac{1}{k} \sum_{j=1}^k 1\{\xi_j \in S\}$

Return $(p \geq 2\gamma)$

However, the case of small $\gamma_t$ is easily recognizable, by estimating $\gamma_t$ via sampling from the normal $N(A(t)x_t, I)$, and counting how many times we hit $S_t$, using the membership oracle. This is done by the “Test” function. If the Test returns “True”, then with high probability (w.h.p), $\gamma_t \geq \alpha O(1)$, and so we can efficiently sample a stochastic gradient, and assign it to $g(t)$. If it returns “False”, then w.h.p, $\gamma_t \leq \alpha O(1)$. They key idea here is that, as we will show, $\gamma_t$ being small is actually an “easy” case, and simply choosing $g(t) = (A(t)x_t - y_t)x_t^\top$ suffices to make progress towards $A_\ast$. Observe that here $g(t)$ is a stochastic gradient of $\tilde{f}(A) = -E_{y \sim N(Ax_t, I, S_t)} \ell_{\mathbb{R}^n}(Ax_t; y)$. In other words, the Test is a “switch” between $\nabla f_t$ and $\nabla \tilde{f}_t$. 


5. Proof of Theorem 6

In this section, we give an overview of the proof of Theorem 6, with an emphasis on the novel technical components. Our goal is to convey the key ideas, and so at some steps we are slightly informal. We provide the formal and detailed proof in Appendix D.

The Generic Bound. We first prove a bound on \( \| \hat{A} - A_s \|_\Sigma^2 \) (remember that \( \Sigma = \Sigma_N \) in Algorithm 1), which holds for any sequence \( g_1, g_2, \ldots, g_N \). This bound will serve as a guideline for choosing the \( g_i \)’s.

Lemma 9 Independently of how \( g_i \)’s are chosen,

\[
\| \hat{A} - A_s \|_\Sigma^2 \leq 1 - \sum_{i=1}^N \left( 2\eta (g_i, A_i - A_s) - \| (A_i - A_s)x_{t_i} \|_2 \right) + \eta^2 \sum_{i=1}^N \text{tr} \left( g_i \Sigma_i^{-1} g_i^\top \right). \tag{4}
\]

The proof of the lemma is along the lines of the analysis of ONS in Hazan et al. (2007), and we provide it in Appendix D.1. Let \( E_1 \) be the first sum in 4, and \( E_2 \) the second. The rest of the proof is about showing that for our choice of \( g_i \)’s, \( E_1 \) is (almost) non-negative, and \( E_2 \) is not too large (both in expectation). In the main text we present the crux of the proof, which is the bound on \( E_1 \). Furthermore, in the main text we assume that \( L = 0 \), since the extension for general \( L \) is straightforward.

Bound on \( E_1 \).

Fix a time \( t_i \), condition on \( F_{t_i} \), and suppose that \( t = t_i \), for some \( i \). To ease notation, we define \( A(t) = A_{t_i} \), \( g(t) = g_i \), \( \mu_t^* = A_{t_i} x_{t_i} \), \( \mu_t = A(t)x_{t_i} \), and \( V_t = 2\eta (g(t), A(t) - A_s) - \| (A(t) - A_s)x_{t_i} \|_2 \). To bound \( \mathbb{E}[E_1] \), we show (roughly) that \( \mathbb{E}[V_t | F_{t_i}, t = t_i] \geq 0 \). Observe that this immediately implies \( \mathbb{E}[E_1] \geq 0 \). We consider two cases, based on what the Test function returns.

Large Survival Probability. Suppose at time \( t \), Test returns “True”. We call this event \( T_t \), and here we assume that Test returning “True” implies \( N(\mu_t^*, I; S_t) \geq \gamma \) (this is correct w.h.p.).\(^{17} \)

Furthermore, \( z_t \sim N(\mu, I, S_t) \), and \( g(t) = (z_t - y_t)x_t^\top \). Also, since \( y_t \sim N(\mu^*_t, I, S_t) \), we have

\[
\mathbb{E}[V_t | F_{t_i}, t = t_i, T_t] = 2\eta (\nu_t - \nu_t^*, \mu_t - \mu_t^*) - \| \mu_t - \mu_t^* \|^2, \tag{5}
\]

where \( \nu_t = \mathbb{E}_{z \sim N(\mu_t, I, S_t)}[z] \) and \( \nu_t^* = \mathbb{E}_{y \sim N(\mu_t^*, I, S_t)}[y] \). Now, note that since \( t = t_i \) and \( L = 0 \), we have \( N(\mu_t^*, I; S_t) \geq \alpha \geq \gamma \). The key component of our proof is the following lemma, which we prove here in detail.

Lemma 10 Let \( \mathcal{L}_S(\mu; \mu^*) = -\mathbb{E}_{y \sim N(\mu^*, I, S)}[\ell_S(\mu; y)] \), where \( \mu, \mu^* \in \mathbb{R}^n \) and \( S \subseteq \mathbb{R}^n \). Suppose that \( N(\mu, I; S) \), \( N(\mu^*, I; S) \geq \gamma \). Then, there exists an absolute constant \( c > 0 \) such that

\[
\langle \nabla_\mu \mathcal{L}_S(\mu; \mu^*), \mu - \mu^* \rangle \geq (\gamma/2)^c \cdot \| \mu - \mu^* \|^2. \tag{6}
\]

\(^{17} \) The parameter \( \gamma \) is defined in Algorithm 3. Also, in Appendix D, where we give the details for handling the low-probability events, we show that \( \mathbb{E}[E_1] \geq \Delta \), for some \( \Delta > 0 \), but small.
Lemma 10 implies that, given large enough $c_\eta$,\(^\text{18}\) the expectation in (5) is non-negative. Indeed, note that $\nabla_{\mu} L_S(\mu; \mu^*) = \nu - \nu^*$, where $\nu = \mathbb{E}_{z \sim N(\mu, I, S)}[z]$ and $\nu^* = \mathbb{E}_{y \sim N(\mu^*, I, S)}[z]$, so

$$
\mathbb{E}[V_t \mid F_t, t = i, T_i] \geq (2\eta(\gamma/2)^c - 1) \cdot \|\mu_t - \mu_t^*\|^2 = (2\eta(\gamma/2)^c - 1) \cdot \|\mu_t - \mu_t^*\|^2
$$

for any constant $c_\eta \geq 2c \cdot c_\gamma$. We now prove Lemma 10.

**Proof** Let $s(\gamma) := \sqrt{2 \log(1/\gamma)} + 1$. We will need two technical claims, which hold for any $\mu, \mu^* \in \mathbb{R}^n$ and $S \subseteq \mathbb{R}^n$.

**Claim 1** Let $\nu = \mathbb{E}_{z \sim N(\mu, I, S)}[z]$. If $N(\mu, I, S) \geq \gamma$, then $\|\nu - \mu\| \leq s(\gamma)$.

**Claim 2** Suppose $N(\mu^*, I, S) \geq \gamma$, and for some $\tilde{\mu}$ we have $\|\tilde{\mu} - \mu^*\| \leq c \cdot s(\gamma)$. Then,

$$
\text{Cov}_{z \sim N(\tilde{\mu}, I, S)}[z, z] \succeq (\gamma/2)^{\text{poly}(c)} \cdot I.
$$

Both Claims are mainly from Daskalakis et al. (2018). Claim 1 is Lemma 6 in that paper. Claim 2 is not explicitly stated, but can be derived by slightly adapting their proof (see Appendix H.1). We consider two cases:

**Case 1:** $\|\mu - \mu^*\| \geq 4s(\gamma)$. Then,

$$
\langle \nu - \nu^*, \mu - \mu^* \rangle = \|\mu - \mu^*\|^2 + \langle \nu - \mu, \mu - \mu^* \rangle + \langle \mu^* - \nu^*, \mu - \mu^* \rangle
$$

$$
\geq \|\mu - \mu^*\|^2 - \|\nu - \mu\| \cdot \|\mu - \mu^*\| - \|\nu^* - \mu^*\| \cdot \|\mu - \mu^*\|
$$

Since $N(\mu, I, S), N(\mu^*, I, S) \geq \gamma$, Claim 1 implies $\|\nu - \mu\|, \|\nu^* - \mu^*\| \leq s(\gamma)$. Being in Case 1,

$$
\langle \nu - \nu^*, \mu - \mu^* \rangle \geq \|\mu - \mu^*\| \cdot (\|\mu - \mu^*\| - 2s(\gamma)) \geq \|\mu - \mu^*\|^2/2.
$$

**Case 2:** $\|\mu - \mu^*\| < 4s(\gamma)$. To ease notation, fix $\mu^*, S$, and let $\mathcal{L}(\mu) = L_S(\mu; \mu^*)$. From fundamental theorem of calculus,

$$
\nabla \mathcal{L}(\mu) - \nabla \mathcal{L}(\mu^*) = \int_0^1 \nabla^2 \mathcal{L}(\mu(\theta)) \, d\theta \cdot (\mu - \mu^*),
$$

where $\mu(\theta) := \mu^* + \theta(\mu - \mu^*)$. Since $\nabla \mathcal{L}(\mu^*) = 0$,

$$
\langle \nabla \mathcal{L}(\mu), \mu - \mu^* \rangle = \int_0^1 \langle \nabla^2 \mathcal{L}(\mu(\theta)) (\mu - \mu^*), \mu - \mu^* \rangle \, d\theta.
$$

For $\theta \in (0, 1)$, $\|\mu(\theta) - \mu^*\| = \theta \|\mu - \mu^*\| \leq O(s(\gamma))$. Using $N(\mu^*, I, S) \geq \gamma$ and Claim 2,

$$
\nabla^2 \mathcal{L}(\mu(\theta)) = \text{Cov}_{z \sim N(\mu(\theta), I, S)}[z, z] \succeq (\gamma/2)^{O(1)} \cdot I.
$$

Using (7), we finish the proof. \(\blacksquare\)

\(^{18}\) Remember that $\eta = (2/\alpha)^c$. 

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Remark. We want to highlight an important qualitative difference between Lemma 10 and all recent works in CTS. Suppose $\gamma = \Omega(1)$ and let $v = \frac{\mu - \mu^*}{\|\mu - \mu^*\|}$. From the argument in Case 2, 6 can be rephrased as

\[
\mathbb{E}_{\tilde{\mu} \sim [\mu^*, \mu]} \left[ \langle \text{Cov}_{z \sim \mathcal{N}(\tilde{\mu}, I, S)}[z, z] \cdot v, v \rangle \right] \geq \Omega(1),
\]

where $\tilde{\mu}$ is distributed uniformly on the segment $[\mu, \mu^*]$. In other words, the average variance in the $v$-direction is $\Omega(1)$. This is an anti-concentration bound. As we mentioned, proving anti-concentration bounds for truncated normals $\mathcal{N}(\mu, I, S)$ is a core component in all recent works in CTS, where this task is reduced to lower bounding $\mathcal{N}(\mu^*, I; S)$. However, inequality (8) cannot be proven with this methodology. We illustrate this with an insightful example (Figure 1). In the example, we consider $\mu, \mu^* \in \mathbb{R}^d$ and $\|\mu - \mu^*\| = \sqrt{d}$, because the guaranteed bound in the LDS-case for $\|\mu_t - \mu_t^*\|$ will be $\Theta(\sqrt{d})$. Let $\mu^* = 0$, $\mu = \sqrt{d} \cdot e_1$ ($e_1$ is the standard-basis vector), $S = \{x \in \mathbb{R}^d : x_1 \in (-\infty, 0] \cup [\sqrt{d}, +\infty)\}$. Observe that the conditions of Lemma 10 are satisfied with $\gamma = 1/2$. However, for most $\tilde{\mu} \in [\mu^*, \mu]$, the mass $\mathcal{N}(\tilde{\mu}, I; S)$ is exponentially small (in $d$). Consider now $\tilde{\mu}$ exactly in the middle of $\mu^*, \mu$. Even though $\mathcal{N}(\tilde{\mu}, I; S) = \exp(-\Omega(d))$, the variance in the $e_1$ direction is $\Omega(d)$, due to symmetry. It can actually be shown that precisely this $\tilde{\mu}$ and small perturbations of it make the average variance $\Omega(1)$ in (8). Note that it is necessary for both corners ($\mu$ and $\mu^*$) to have high survival probability, e.g., in Figure 2 (for the same distance-scales) it can be shown that the variance is $O(1/d)$.

Small Survival Probability. Suppose that at time $t$, Test returns “False” ($\neg T_t$). Suppose also that this implies $\mathcal{N}(\mu^*_t, I; S_t) \leq 4\gamma$ (again this is w.h.p). We will show that $\mathbb{E}[V_t | F_t, t = t_i, \neg T_t] \geq 0$. Observe that given $\neg T_t$, we have $g(t) = (\mu_t - y_t)x_t^\top$, and so

\[
\mathbb{E}[V_t | F_t, t = t_i, \neg T_t] = 2\eta\langle \nu_t^* - \mu_t^*, \mu_t - \mu_t^* \rangle - \|\mu_t - \mu_t^*\|^2.
\]

Again, since $t = t_i$ and $L = 0$, we have $\mathcal{N}(\mu^*_t, I; S_t) \geq \alpha$, so by Claim 1, $\|\nu_t^* - \mu_t^*\| \leq s(\alpha)$. We need one more technical claim.

Claim 3. If $\mathcal{N}(\mu^*_t, I; S) \geq \alpha$, and $\|\mu_t - \mu^*_t\| \leq c \cdot s(\alpha)$, then $\mathcal{N}(\mu, I; S) \geq (\alpha/2)^{\text{poly}(c)}$.

Again, Claim 3 is proven in Appendix H.2, by slightly adapting the analysis of Daskalakis et al. (2018). Now, we claim that $\|\mu_t - \mu_t^*\| > 2s(\alpha)$, provided that $c_\gamma = O(1)$ is sufficiently large. Indeed, suppose $\|\mu_t - \mu_t^*\| \leq 2s(\alpha)$. Combining with $\mathcal{N}(\mu^*_t, I; S_t) \geq \alpha$ and Claim 3, we get $4(\alpha/2)^{c_\gamma} = 4\gamma \geq \mathcal{N}(\mu_t, I; S_t) \geq (\alpha/2)^{O(1)}$, which is a contradiction for large enough $c_\gamma = O(1)$. Using that $\|\nu_t^* - \mu_t^*\| \leq s(\alpha)$, we get

\[
\langle \mu_t - \nu_t^*, \mu_t - \mu_t^* \rangle = \|\mu_t - \mu_t^*\|^2 + \langle \nu_t^*, \mu_t - \mu_t^* \rangle \geq \|\mu_t - \mu_t^*\|^2 - \|\nu_t^* - \mu_t^*\| \cdot \|\mu_t - \mu_t^*\|
\]

\[
\geq \|\mu_t - \mu_t^*\| \cdot (\|\mu_t - \mu_t^*\| - s(\alpha)) \geq \|\mu_t - \mu_t^*\|^2/2.
\]
Since \( \eta \geq 1 \), Equation 23 implies \( \mathbb{E}[V_t \mid \mathcal{F}_t, t = t_i, \neg \mathcal{T}_t] \geq 0 \). Thus, \( \mathbb{E}[V_t \mid \mathcal{F}_t, t = t_i] \geq 0 \), and \( \mathbb{E}[E_1] \geq 0 \).

6. The Initial Ellipsoid and Proof of Theorem 1

Here, we provide the overall algorithm (Algorithm 4), and an outline of the proof of Theorem 1. Let \( \mathcal{I}_0 = \{ t \in [T] : (x_t, x_{t+1}) \in \mathcal{P}_0 \} \), and \( \mathcal{I}_1 = \{ t \in [T] : (x_t, x_{t+1}) \in \mathcal{P}_1 \} \).

Algorithm 4 SON-SG for censored linear dynamics

\[
A_0 = \arg \min_{A \in \mathbb{R}^{d \times d}} \sum_{t \in \mathcal{I}_0} \|x_{t+1} - Ax_t\|^2 \\
\Sigma_0 = \frac{1}{|\mathcal{I}_0|} \sum_{t \in \mathcal{I}_0} x_t x_t^T, \text{ where } s = c_s \left( \sqrt{\log \left( \frac{1}{\alpha} \right)} + 1 \right). \quad \triangleright c_s \geq 0 \text{ is a large constant.}
\]

Get \( \hat{A} \) by running SON-SG (Algorithm 1) with \( A_0, \Sigma_0 \), and dataset \( \mathcal{P}_1 \).

return \( \hat{A} \)

The least-squares method when all \( x_t \)'s are observed was analyzed in Simchowitz et al. (2018). Here, we observe only a subset the \( x_t \)'s. In Appendix G, we generalize the “small-ball” technique of Simchowitz et al. (2018), to prove the following theorem.

**Theorem 11** There exist \( c_1, c_2, c_3 = \text{poly} \left( \frac{1}{\alpha \beta} \right) \), such that if \( T \geq \tilde{\Theta}(c_1) \cdot \left( d^2 + \frac{d}{1 - \rho(A_*)^T} + dL \right) \), then with probability \( 1 - o(1) \), the ellipsoid \( \tilde{K} = \{ A \in \mathbb{R}^{d \times d} : \|A - A_0\|_{\Sigma_0} \leq 1 \} \) is \((R, \omega)\)-accurate with \( R = c_2 \sqrt{d}, \omega = 1/c_2 \), and also \( \frac{1}{|\mathcal{I}_1|} \sum_{t \in \mathcal{I}_1} x_t x_t^T \geq \frac{1}{c_3} \Gamma_T \).

We now prove Theorem 1 using Theorems 6 and 11.

**Proof** We use the time-series model with \( y_t = x_{t+1}, S_t = S_{t+1}, o_t = \mathbb{I}(t \in \mathcal{I}_1) \). Thus, \( B(\alpha) \subseteq B(\alpha) \), so \( \mathbb{E}[|B(\alpha)|] \leq L \). Let \( \mathcal{E}_0 \) be the event that all guarantees in Theorem 11 hold, so \( \mathbb{P}[\neg \mathcal{E}_0] \leq o(1) \). As we mentioned, \( \|w_t\| \) is \( O(d) \)-subgaussian. Also, in Appendix H.3, we show that \( \mathbb{E}[\|x_t\|^2] \) is polynomial in all parameters. Let \( \Sigma = \frac{1}{|\mathcal{I}_1|} \sum_{t \in \mathcal{I}_1} x_t x_t^T \). By conditioning on \( \mathcal{E}_0 \) and applying Theorem 6, using Markov’s inequality in (2), we have that with probability at least \( 1 - \delta \),

\[
\|\hat{A} - A_*\|^2_{\Sigma} \leq \frac{\tilde{O}(C)}{\delta N} \cdot \left( d^2 + R^2 + R_w^2 \right) + L(d + R^2 + R_w^2) \leq \tilde{O}(C) \cdot \frac{d^2 + dL}{\delta N}.
\]

where \( C = \text{poly} \left( \frac{1}{\alpha \beta} \right) \). Now, on \( \mathcal{E}_0, \Sigma \gtrsim \Gamma_T / \text{poly} \left( \frac{1}{\alpha \beta} \right) \). From Proposition 4, \( \mathbb{P} \left[ N < \frac{\alpha \beta T}{4} \right] \leq o(1) \). Thus, by choosing small enough \( \delta = \Omega(1) \), we are done.

**References**


19. Observe that our bounds depend on \( 1/\delta \). It is possible to achieve \( \log(1/\delta) \) dependence, but with more technical proofs. To keep the presentation simpler, we chose not to optimize this dependence.


Benjamin Kramer Johannsen and Elmar Mertens. A time series model of interest rates with the effective lower bound. 2018.


