A Deterministic Algorithm for the Capacity of Finite-State Channels *[†]

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Abstract

We propose two modified versions of the classical gradient ascent method to compute the capacity of finite-state channels with Markovian inputs. For the case that the channel mutual information is strongly concave in a parameter taking values in a compact convex subset of some Euclidean space, our first algorithm proves to achieve polynomial accuracy in polynomial time and, moreover, for some special families of finite-state channels our algorithm can achieve exponential accuracy in polynomial time under some technical conditions. For the case that the channel mutual information may not be strongly concave, our second algorithm proves to be at least locally convergent.

Index Terms: channel capacity, finite-state channels, gradient ascent, hidden Markov process.

1 Introduction

As opposed to a discrete memoryless channel, which can be characterized by the conditional distribution of the output given the input, in a finite-state channel this conditional distribution depends on an underlying state variable which evolves with time. Encompassing

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discrete memoryless channels as special cases, finite-state channels have long been used in a wide range of communication scenarios where the current behavior of the channel may be affected by its past. Among many others, conventional examples of such channels include inter-symbol interference channels [8], partial response channels [22, 23] and Gilbert-Elliott channels [21].

While it is well-known that the Blahut-Arimoto algorithm [2, 4] can be used to efficiently compute the capacity of a discrete memoryless channel, the computation of the capacity of a general finite-state channel has long been a notoriously difficult problem, which has been open for decades. The difficulty of this problem may be justified by the widely held (yet not proven) belief that the capacity of a finite-state channel may not be achieved by any finite-order Markovian input, and an increase of the memory of the input may lead to an increase of the channel capacity.

We are mainly concerned with finite-state channels with Markov processes of a fixed order as their inputs. Possibly an unavoidable compromise we have to make in exchange for progress in computing the capacity, the extra fixed-order assumption imposed on the input processes is also necessary for the situation where the channel input has to satisfy certain constraints, notably finite-type constraints [19] that are commonly used in magnetic and optical recording. On the other hand, the focus on Markovian inputs can also be justified by the known fact that the Shannon capacity of an indecomposable finite-state channel [9] can be approximated by the Markov capacity with increasing orders (see Theorem 2.1 of [18]). Recently, there has been some progress in computing the capacity of finite-state channels with such input constraints. Below we only list the most relevant work in the literature, and we refer the reader to [12] for a comprehensive list of references. In [15], the Blahut-Arimoto algorithm was reformulated into a stochastic expectation-maximization procedure and a similar algorithm for computing a lower bound on the capacity of finite-state channels was proposed, which led to a generalized Blahut-Arimoto algorithm [24] that proves to compute the capacity under some concavity assumptions. More recently, inspired by ideas in stochastic approximation, a randomized algorithm was proposed [12] to compute the capacity under weaker concavity assumptions, which can be verified to hold true for several families of practical channels [14, 16]. Both of the above-mentioned algorithms, however, are of a randomized nature (a feasible implementation of the generalized Blahut-Arimoto algorithm will necessitate a randomization procedure). By comparison, among many other advantages, our algorithms, which are deterministic in nature, can be used to derive accurate estimates on the channel capacity, as evidenced by the tight bounds in Section 3.2.

In this paper, we first deal with the case that the mutual information of the finite-state channel is strongly concave in a parameter taking values in a compact convex subset of some Euclidean space, for which we propose our first algorithm that proves to converge to the channel capacity exponentially fast. This algorithm largely follows the spirit of the classical gradient ascent method. However, unlike the classical case, the lack of an explicit expression for our target function and the boundedness of the variable domain (without an explicit description of the boundary) pose additional challenges. To overcome the first issue, a convergent sequence of approximating functions (to the original target function) is used instead in our treatment; meanwhile, an additional check condition is also added to ensure that the iterates stay inside the given variable domain. A careful convergence analysis has been carried out to deal with the difficulties caused by such modifications. This algorithm is efficient in the sense that, for a general finite-state channel (satisfying the above-mentioned concavity condition and some additional technical conditions), it achieves polynomial accuracy in polynomial time (see Theorem 3.12), and for some special families of finite-state channels it achieves exponential accuracy in polynomial time (see Section 3.2).

It is well known that the mutual information rate of a finite-state channel may not be concave under the natural parametrization in several examples; see, e.g., [14, 16]. Another modification of the classical gradient ascent method is proposed to handle this challenging scenario. Similar to our first algorithm, our second one replaces the original target function with a sequence of approximating functions, which unfortunately renders conventional methods such as the Frank-Wolfe method (see, e.g., [3]) or methods using the Lojasiewicz inequality (see, e.g., [1]) inapplicable. To address this issue, among other subtle modifications, we impose an extra check in the algorithm to slow down the pace "a bit" to avoid an immature convergence to a non-stationary point but "not too much" to ensure the local convergence.

As variants of the classical gradient ascent method, our algorithms can be applied to any sequence of convergent functions, so they can be of particular interest in information theory since many information-theoretic quantities are defined as the limit of their finite-block versions. On the other hand though, we would like to add that our algorithms are actually stated in a much more general setting and may have potential applications in optimization scenarios where the target functions are difficult to compute but amenable to approximations.

The remainder of this paper is organized as follows. In Section 2, we describe our channel model in greater detail. Then, we present our first algorithm (Algorithm 3.3) in Section 3 and analyze its convergence behavior in Section 3.1 under some strong concavity assumptions. Applications of this algorithm for computing the capacity of finite-state channels under concavity assumptions will be discussed in Section 3.2. In particular, in this section, we show that the estimation of the channel capacity can be improved by increasing the Markov order of the input process in some examples. In Section 4, our second algorithm (Algorithm 4.2) is presented, which proves to be at least locally convergent. Finally, in Section 4.2, our second algorithm is applied to a Gilbert-Elliott channel where the concavity of the channel mutual information rate in the natural parametrization is not known, and yet fast convergence behavior is observed.

In the remainder of this paper, the base of the logarithm is assumed to be e.

2 Channel Model and Problem Formulation

In this section, we introduce the channel model considered in this paper, which is essentially the same as that in [12, 24].

As mentioned before, we are concerned with a discrete-time finite-state channel with a Markovian channel input. Let $X = \{X_n : n = 1, 2, ...\}$ denote the channel input process, which is often assumed to be a first-order stationary Markov chain ¹ over a finite alphabet \mathcal{X} , and let $Y = \{Y_n : n = 1, 2, ...\}$ and $S = \{S_n : n = 0, 1, ...\}$ denote the channel output and state processes over finite alphabets \mathcal{Y} and \mathcal{S} , respectively.

¹The assumption that X is a first-order Markov chain is for notational convenience only: through a usual "reblocking" technique, the higher-order Markov case can be boiled down to the first-order case.

Let Π be the set of all the stochastic matrices of dimension $|\mathcal{X}| \times |\mathcal{X}|$. For any finite set $F \subseteq \mathcal{X}^2$ and any $\delta > 0$, define

$$\Pi_{F,\delta} \triangleq \{A \in \Pi : A_{ij} = 0, \text{ for } (i,j) \in F \text{ and } A_{ij} \geq \delta \text{ otherwise} \}.$$

It can be easily verified that if one of the matrices from $\Pi_{F,\delta}$ is primitive, then all matrices from $\Pi_{F,\delta}$ will be primitive, in which case, as elaborated on in [12], F gives rise to a so-called mixing finite-type constraint. Such a constraint has been widely used in data storage and magnetic recoding [20], the best known example being the so-called (d, k)-run length limited (RLL) constraint over the alphabet $\{0, 1\}$, which forbids any sequence with fewer than d or more than k consecutive zeros in between two successive 1's.

The following conditions will be imposed on the finite-state channel described above:

- (2.a) There exist $F \subseteq \mathcal{X}^2$ and $\delta > 0$ such that the transition probability matrix of X belongs to $\Pi_{F,\delta}$, each element of which is a primitive matrix.
- (2.b) (X, S) is a first-order stationary Markov chain whose transition probabilities satisfy

$$p(x_n, s_n | x_{n-1}, s_{n-1}) = p(x_n | x_{n-1}) p(s_n | x_n, s_{n-1}), \quad n = 1, 2, \dots,$$

where $p(s_n | x_n, s_{n-1}) > 0$ for any s_{n-1}, s_n, x_n .

(2.c) The channel is stationary and characterized by

$$p(y_n|y_1^{n-1}, x_1^n, s_1^{n-1}) = p(y_n|x_n, s_{n-1}) > 0, \quad n = 1, 2, \dots,$$

that is, conditioned on the pair (x_n, s_{n-1}) , the output Y_n is statistically independent of all inputs, outputs and states prior to X_n, Y_n and S_{n-1} , respectively.

As elaborated on in Remark 4.1 of [12], a finite-state channel specified as above is indecomposable. Therefore, assuming that the input X (or, more precisely, the transition probability matrix of X) is analytically parameterized by a finite-dimensional parameter θ in the interior of a compact convex subset Θ of some Euclidean space, the parametrization being continuous at the boundary (such a parameterization exists thanks to the stationarity of X, and we will simply say that X is analytically parametrized by θ , for convenience), we can express the capacity of the above channel as

$$C = \max_{\theta \in \Theta} I(X(\theta); Y(\theta)) = \max_{\theta \in \Theta} \lim_{k \to \infty} I_k(X(\theta); Y(\theta)),$$
(1)

where

$$I_k(X(\theta); Y(\theta)) \triangleq \frac{H(X_1^k(\theta)) + H(Y_1^k(\theta)) - H(X_1^k(\theta), Y_1^k(\theta))}{k}.$$
(2)

Moreover, it has also been shown in [12] that $I_k(X(\theta); Y(\theta))$ (resp., its derivatives) converges to $I(X(\theta); Y(\theta))$ (resp., the corresponding derivatives) exponentially fast in k under Assumptions (2.a), (2.b) and (2.c). Hence, although the value of the target function $I(X(\theta); Y(\theta))$ cannot be exactly computed, it can be approximated by the function $I_k(X(\theta); Y(\theta))$, which has an explicit expression, within an error exponentially decreasing in k. Instead of merely solving (1), we will deal with the following slightly more general problem

$$\max f(\theta) \triangleq \lim_{k \to \infty} f_k(\theta)$$

subject to $\theta \in \Theta$, (3)

under the following assumptions:

- A.1. Θ is a compact convex subset of \mathbb{R}^d for some $d \in \mathbb{N}$ with nonempty interior Θ° and boundary $\partial \Theta$;
- A.2. $f(\theta)$ and all $f_k(\theta)$, $k \ge 0$, are continuous on Θ and twice continuously differentiable in Θ° ;
- A.3. there exist $M_0 > 0$, N > 0 and $0 < \rho < 1$ such that for all $k \ge 1$, $\theta \in \Theta^{\circ}$ and $\ell = 0, 1, 2$, it holds true that $||f_0^{(\ell)}(\theta)||_2 \le M_0$ and

$$||f_k^{(\ell)}(\theta) - f_{k-1}^{(\ell)}(\theta)||_2 \le N\rho^k, \qquad ||f_k^{(\ell)}(\theta) - f^{(\ell)}(\theta)||_2 \le N\rho^k, \tag{4}$$

where the superscript $^{(\ell)}$ denotes the ℓ -th order derivative and $||\cdot||_2$ denotes the Frobenius norm of a vector/matrix.

Obviously, if we set $f_k(\theta) = I_k(X(\theta); Y(\theta))$ and assume that $X(\theta)$ analytically parameterized by some $\theta \in \Theta$, then (3) boils down to (1).

When the target function $f(\theta)$ has an explicit expression and Θ is specified by finitely many inequalities with twice differentiable terms, the optimization problem (3) can be effectively dealt with via, for example, the classical gradient ascent method [5] or the Frank-Wolfe method [3] or their numerous variants. However, feasible implementations and executions of these algorithms usually hinge on explicit descriptions of Θ and ∇f , both of which can be rather intricate in our setting.

Before moving to the next two sections to present our algorithms, we make some observations about the sequence $\{f_k(\theta)\}_{k=0}^{\infty}$. It immediately follows from the uniform boundedness of $||f_0^{(\ell)}(\theta)||_2$ and the inequality (4) that there exists M > 0 such that for all $k \ge 0$, $\ell = 0, 1, 2$ and $\theta \in \Theta^\circ$,

$$||f_k^{(\ell)}(\theta)||_2 \le M. \tag{5}$$

In particular, for any $\theta \in \Theta^{\circ}$, when $\ell = 2$, $f_k^{(\ell)}(\theta) = \nabla^2 f_k(\theta)$ is a symmetric matrix whose spectral norm is given by

$$|||\nabla^2 f_k(\theta)|||_2 \triangleq \sup_{\mathbf{x} \neq 0} \frac{||\nabla^2 f_k(\theta) \cdot \mathbf{x}||_2}{||\mathbf{x}||_2} = |\lambda_1(\theta)|_2$$

where λ_1 denotes the largest (in modulus) eigenvalue of $\nabla^2 f_k(\theta)$. Hence, the inequality (5) and the easily verifiable fact that $|||\nabla^2 f_k(\theta)||_2 \leq ||\nabla^2 f_k(\theta)||_2$ imply

$$-M\mathbb{I}_d \preceq \nabla^2 f_k(\theta) \preceq M\mathbb{I}_d \tag{6}$$

for any k and any $\theta \in \Theta^{\circ}$, where \mathbb{I}_d denotes the $d \times d$ identity matrix, and for two matrices A, B of the same dimension, by $A \preceq B$, we mean that B - A is a positive semidefinite matrix. The existence of the constant M in (6) will be crucial for implementing our algorithms.

3 The First Algorithm: with Concavity

Throughout this section, we assume that $f(\theta)$ is strongly concave, i.e., there exists m > 0 such that for all $\theta \in \Theta^{\circ}$,

$$\nabla^2 f(\theta) \preceq -m \mathbb{I}_d,\tag{7}$$

and moreover

f achieves its unique maximum in Θ° . (8)

We will present our first algorithm to solve the optimization problem (3). As mentioned before, the algorithm is in fact a modified version of the classical gradient ascent algorithm, whereas its convergence analysis is more intricate than the classical one. To overcome the issue that the target function $f(\theta)$ may not have an explicit expression we capitalize on the fact that it can be well approximated by $\{f_k(\theta)\}_{k=0}^{\infty}$, which will be used instead to compute the estimates in each iteration.

Before presenting our algorithm, we need the following lemma, which, as evidenced later, is important in initializing and analyzing our first algorithm.

Lemma 3.1. There exists a non-negative integer k_0 such that

(a)
$$\frac{(N+M)M\rho^{k_0+1}+2N\rho^{k_0+1}}{1-\rho} \leq \frac{\delta}{8} \text{ and } N\rho^{k_0} \leq \frac{\delta}{8}, \text{ where } \delta \triangleq \max_{\theta \in \Theta} f(\theta) - \max_{\theta \in \partial \Theta} f(\theta) > 0.$$

(b) For any $k \ge k_0$, $f_k(\theta)$ is strongly concave and has a unique maximum in Θ° ; and moreover, we have

$$\sup_{k \ge k_0} ||\theta_k^* - \theta^*||_2 + \frac{d^{1/2} \rho^{k_0}}{1 - \rho} < dist(\theta^*, \partial \Theta),$$
(9)

where θ^* denotes the unique maximum point of f and θ^*_k denotes the unique maximum point of f_k .

(c) There exists $y_0 \in \mathbb{R}$ such that for all $k \geq k_0$,

$$\emptyset \subsetneq B_k \subseteq C_k \subseteq \Theta^\circ \quad and \quad dist(C_k, \partial \Theta) > 0,$$

where

$$B_k \triangleq \{x \in \Theta : f_k(x) \ge y_0\}$$
 and $C_k \triangleq \left\{x \in \Theta : f_k(x) \ge y_0 - \frac{\delta}{8}\right\}.$

Proof. Since (a) trivially holds for sufficiently large k_0 , we will omit its proof and proceed to prove (b). Towards this end, note that according to (4) and (7), it holds true that for sufficiently large k, each f_k is strongly concave. Noting that $f(\theta^*) - \max_{\theta \in \partial \Theta} f(\theta) = \delta$, we deduce from (a) and (4) that for k large enough,

$$\max_{\theta \in \Theta} f_k(\theta) - \max_{\theta \in \partial \Theta} f_k(\theta) \ge f_k(\theta^*) - \max_{\theta \in \partial \Theta} f(\theta) - \frac{\delta}{8} \ge f(\theta^*) - \max_{\theta \in \partial \Theta} f(\theta) - \frac{\delta}{4} = \frac{3\delta}{4} > 0.$$
(10)

Hence, for k sufficiently large, f_k achieves its unique maximum at $\theta_k^* \in \Theta^\circ$.

We now prove that $\theta_k^* \to \theta^*$ as $k \to \infty$. To see this, observe that (4) implies the uniform convergence of f_k to f, i.e., for any $\varepsilon > 0$, there exists K > 0 such that for any k > K and any $\theta \in \Theta$, $f(\theta) - \varepsilon \leq f_k(\theta) \leq f(\theta) + \varepsilon$. In particular, for k > K, we have

$$f(\theta^*) - \varepsilon \le f_k(\theta^*) \le f_k(\theta^*_k) \le f(\theta^*_k) + \varepsilon \le f(\theta^*) + \varepsilon$$

which further implies that $f_k(\theta_k^*) \to f(\theta^*)$ as $k \to \infty$. It then follows from the triangle inequality that

$$f(\theta_k^*) \to f(\theta^*), \quad \text{as } k \to \infty.$$
 (11)

Now, by the Taylor series expansion, there exists some $\tilde{\theta} \in \Theta^{\circ}$ such that

$$f(\theta_k^*) - f(\theta^*) = \nabla f(\theta^*)^T (\theta_k^* - \theta^*) + (\theta_k^* - \theta^*)^T \nabla^2 f(\tilde{\theta}) (\theta_k^* - \theta^*).$$
(12)

Since $\nabla f(\theta^*) = 0$ and $\nabla^2 f(\tilde{\theta}) \preceq -m\mathbb{I}_d$ according to (7), it follows from (11) and (12) that $\theta_k^* \to \theta^*$ as $k \to \infty$, as desired.

It then immediately follows that $||\theta_k^* - \theta^*||_2 + d^{1/2}\rho^k/(1-\rho) \to 0$ as $k \to \infty$. Observing that $dist(\theta^*, \partial\Theta) > 0$ (since $\theta^* \in \Theta^\circ$), we infer that (9) holds for sufficiently large k. Hence, (b) will be satisfied as long as k_0 is sufficiently large.

We now show that (c) also holds for sufficiently large k_0 . From the definition of δ , there exists y_0 such that $\max_{\theta \in \partial \Theta} f(\theta) + \frac{\delta}{4} < y_0 < \max_{\theta \in \Theta} f(\theta) - \frac{\delta}{4}$. From (4), using the same logic as that used to derive (10), we infer that for sufficiently large k,

$$\max_{\theta \in \partial \Theta} f_k(\theta) < y_0 - \frac{\delta}{8} < y_0 < \max_{\theta \in \Theta} f_k(\theta).$$
(13)

According to (b) and the fact that $\theta_k^* \in \Theta^\circ$, which follows from (13), we deduce that $\emptyset \subseteq B_k \subseteq C_k \subseteq \Theta^\circ$ and $dist(C_k, \partial \Theta) > 0$ with

$$C_k \triangleq \left\{ x : f_k(x) \ge y_0 - \frac{\delta}{8} \right\}$$
 and $B_k \triangleq \{ x : f_k(x) \ge y_0 \}.$

Therefore, (c) is valid as long as k_0 is sufficiently large. Finally, choosing a larger k_0 if necessary, we conclude that there exists k_0 such that (a), (b) and (c) are all satisfied.

Remark 3.2. We remark that, for any $k \ge k_0$, each B_k specified as above has a non-empty interior, which is due to the rightmost strict inequality in (13) and the continuity of f_k .

We are now ready to present our first algorithm, which modifies the classical gradient ascent method in the following manner: Instead of using ∇f to find a feasible direction, we use ∇f_k as the ascent direction in the k-th iteration and then pose additional check conditions for a careful choice of the step size. Note that such modifications make the convergence analysis more difficult compared to the classical case, as elaborated on in the next subsection. Algorithm 3.3. (The first modified gradient ascent algorithm)

Step 0. Choose k_0 such that Lemma 3.1 (a)-(c) hold. Set k = 0, $g_0 = f_{k_0}$ and choose $\alpha \in (0, 0.5), \beta \in (0, 1)$ and $\theta_0 \in \Theta^\circ$ such that $\theta_0 \in B_{k_0}$ and $\nabla g_0(\theta_0) \neq 0$.

Step 1. Increase k by 1, and set t = 1, $g_k = f_{k_0+k}$. **Step** 2. If $\nabla g_{k-1}(\theta_{k-1}) = 0$, set

$$\tau = \theta_{k-1} + t\nabla g_{k-1}(\theta_{k-1} + \rho^{k+k_0}\mathbf{1}),$$

where 1 denotes the all-one vector in \mathbb{R}^d ; otherwise, set

$$\tau = \theta_{k-1} + t \nabla g_{k-1}(\theta_{k-1}).$$

If $\tau \notin \Theta$ or

$$g_k(\tau) < g_k(\theta_{k-1}) + \alpha t ||\nabla g_{k-1}(\theta_{k-1})||_2^2 - (N+M)Mt\rho^{k+k_0},$$

set $t = \beta t$ and go to Step 2, otherwise set $\theta_k = \tau$ and go to Step 1.

Remark 3.4. It is obvious from the definition of g_k that as k tends to infinity, g_k (resp., its first and second order derivatives) converges to f (resp., its first and second order derivatives) exponentially fast with the same constants N and ρ as in (4).

Remark 3.5. According to Lemma 3.1, the choice of k_0 depends on practical evaluations of constants N, ρ, M , and is different from case to case. Moreover, the existence of θ_0 can also be justified by Lemma 3.1 (c).

Remark 3.6. We point out that in Step 2 of Algorithm 3.3, for any $k \ge 1$, when $\nabla g_{k-1}(\theta_{k-1}) = 0$, the point $\theta_{k-1} + \rho^{k+k_0} \mathbf{1}$ will always lie in Θ° . To see this, note that if θ_{k-1} is the maximum point of $g_{k-1} = f_{k+k_0-1}$, then $\theta_{k-1} = \theta^*_{k+k_0-1}$. However, by Lemma 3.1 (b), $\theta^*_{k+k_0-1}$ satisfies (9), which immediately implies that $\theta_{k-1} + \rho^{k+k_0} \mathbf{1} \in \Theta^{\circ}$ when $\nabla g_{k-1}(\theta_{k-1}) = 0$, for any $k \ge 1$.

Remark 3.7. For technical reasons that will be made clear in the next section, α is chosen within (0, 0.5) to ensure the convergence of the algorithm. In Step 2 of Algorithm 3.3, the case that $\nabla g_{k-1}(\theta_{k-1}) = 0$ is singled out for special treatment to prevent the algorithm from getting trapped at the maximum point of f_{k-1} for a fixed k, which may be still far away from the maximum point of f.

3.1 Convergence Analysis

As mentioned earlier, compared to the classical gradient ascent method, Algorithm 3.3 poses additional challenges for convergence analysis. The main difficulties come from the two check conditions in Step 2: the "perturbed" Armijo condition (see, e.g., Chapter 2 of [3] for more details)

$$g_k(\tau) \ge g_k(\theta_{k-1}) + \alpha t ||\nabla g_{k-1}(\theta_{k-1})||_2^2 - (N+M)Mt\rho^{k+k_0}$$

may break the monotonicity of the sequence $\{g_k(\theta_k)\}_{k=0}^{\infty}$ which would have been used to simplify the convergence analysis in the classical case; and the extra check condition $\tau \in \Theta$

 $(\tau \text{ depends on } k)$ forces us to seek uniform control (over all k) of the time used to ensure the validity of this condition in each iteration. In the remainder of this section, we deal with these problems and examine the convergence behavior of Algorithm 3.3. In a nutshell, we will prove that our algorithm converges exponentially fast under some strong concavity assumptions.

Note that the variable k as in Algorithm 3.3 actually records the number of times that Step 1 has been executed at the present moment. To facilitate the analysis of our algorithm, we will put it into an equivalent form, where an additional variable n is used to record the number of times that Step 2 has been executed.

Below is Algorithm 3.3 rewritten with the additional variable n.

Algorithm 3.8. (An equivalent form of Algorithm 3.3)

Step 0. Choose k_0 such that Lemma 3.1 (a)-(c) hold. Set $n = 0, k = 0, \hat{g}_0 = g_0 = f_{k_0}$, and choose $\alpha \in (0, 0.5), \beta \in (0, 1)$ and $\hat{\theta}_0 \in \Theta^\circ$ such that $\hat{\theta}_0 \in B_{k_0}$ and $\nabla \hat{g}_0(\hat{\theta}_0) \neq 0$.

Step 1. Increase k by 1, and set t = 1, $g_k = f_{k_0+k}$. Step 2. Increase n by 1. If $\nabla \hat{g}_{n-1}(\hat{\theta}_{n-1}) = 0$, set

$$\tau = \hat{\theta}_{n-1} + t \nabla \hat{g}_{n-1} (\hat{\theta}_{n-1} + \rho^{k+k_0} \mathbf{1}); \tag{14}$$

otherwise, set

$$\tau = \hat{\theta}_{n-1} + t\nabla \hat{g}_{n-1}(\hat{\theta}_{n-1}).$$
(15)

If $\tau \notin \Theta^{\circ}$ or

$$g_k(\tau) < g_k(\hat{\theta}_{n-1}) + \alpha t ||\nabla \hat{g}_{n-1}(\hat{\theta}_{n-1})||_2^2 - (N+M)Mt\rho^{k+k_0},$$
(16)

then set $\hat{\theta}_n = \hat{\theta}_{n-1}, \hat{g}_n = \hat{g}_{n-1}, t = \beta t$ and go to Step 2; otherwise, set $\hat{\theta}_n = \tau, \hat{g}_n = g_k$ and go to Step 1.

Remark 3.9. Let $n_0 = 0$, and for any $k \ge 1$, recursively define

$$n_k \triangleq \inf\{n > n_{k-1} : \hat{\theta}_n \neq \hat{\theta}_{n-1}\}.$$

Then, one verifies that for any $k \ge 0$, it holds true that $\hat{\theta}_{n_k} = \theta_k$, $\hat{g}_{n_k} = g_k = f_{k+k_0}$ and moreover, $\hat{\theta}_l = \hat{\theta}_{l+1}$, $\hat{g}_l = \hat{g}_{l+1}$ for any l with $n_{k-1} \le l \le n_k - 1$, which justify the equivalence between Algorithm 3.3 and Algorithm 3.8.

The following theorem establishes the exponential convergence of Algorithm 3.8 with respect to n.

Theorem 3.10. Suppose, as in (7) and (8), that the strongly concave function f achieves its unique maximum in Θ° . Then there exist $\hat{M} > 0$ and $0 < \hat{\xi} < 1$ such that for all $n \ge 0$,

$$|\hat{g}_n(\hat{\theta}_n) - f(\theta^*)| \le \hat{M}\hat{\xi}^n,\tag{17}$$

where $\hat{g}_n(\hat{\theta}_n)$ is obtained by executing Algorithm 3.8.

Proof. For simplicity, we only deal with the case $\nabla \hat{g}_{n-1}(\hat{\theta}_{n-1}) \neq 0$ in Step 2 of Algorithm 3.8 (and therefore (15) is actually executed), since the opposite case follows from a similar argument by replacing $\hat{\theta}_{n-1}$ with $\hat{\theta}_{n-1} + \rho^{k+k_0} \mathbf{1}$.

Let $T_1(k)$ denote the smallest non-negative integer p such that

$$\hat{\theta}_{n_{k-1}} + \beta^p \nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}}) \in \Theta^\circ, \tag{18}$$

T(k) denote the smallest non-negative integer q such that $q \ge T_1(k)$ and

$$g_k(\hat{\theta}_{n_{k-1}} + \beta^q \nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})) \ge g_k(\hat{\theta}_{n_{k-1}}) + \alpha \beta^q ||\nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})||_2^2 - (N+M)M\beta^q \rho^{k+k_0}.$$

Note that the well-definedness of $T_1(k)$ and T(k) follows from the observation that if (18) holds for some non-negative integer p, then it also holds for any integer p' > p. Adopting these definitions, we can immediately verify that

$$T(k) = n_k - n_{k-1},$$

which corresponds to the number of times Step 2 (of Algorithm 3.3) has been executed to obtain $\hat{\theta}_{n_k}$ from $\hat{\theta}_{n_{k-1}}$.

The remainder of the proof consists of the following three steps.

Step 1: Uniform boundedness of T(k). In this step, we show that there exists $A \ge 0$ such that, for all $k, T(k) \le A$.

Since Θ° is open and $\hat{\theta}_0 \in \Theta^{\circ}$, we have $T_1(k) < \infty$ for any $k \ge 0$. Note that we haven't show that $T_1(k)$ is uniformly bounded as this stage.

For any $q \geq T_1(k)$, letting

$$\tau = \hat{\theta}_{n_{k-1}} + \beta^q \nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})$$

we have $\tau \in \Theta^{\circ}$ and so both $f_k(\tau)$ and $f(\tau)$ are well-defined. Recalling from (6) that

$$\nabla^2 g_k(\theta) = \nabla^2 f_{k+k_0}(\theta) \succeq -M \mathbb{I}_d$$

for any $k \ge 0$ and any $\theta \in \Theta^{\circ}$, we derive from the Taylor series expansion that

$$g_{k}(\tau) = g_{k}(\hat{\theta}_{n_{k-1}}) + \beta^{q} \nabla g_{k}(\hat{\theta}_{n_{k-1}})^{T} \nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}}) + \frac{\beta^{2q}}{2} \nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})^{T} \nabla^{2} g_{k}(\tilde{\theta}_{k}) \nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}}) \\ \geq g_{k}(\hat{\theta}_{n_{k-1}}) + \beta^{q} \nabla g_{k}(\hat{\theta}_{n_{k-1}})^{T} \nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}}) - \frac{M\beta^{2q}}{2} ||\nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})||_{2}^{2},$$
(19)

where $\tilde{\theta}_k \in \Theta^{\circ}$. According to (4), we have

$$\begin{split} \nabla g_k(\hat{\theta}_{n_{k-1}})^T \nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}}) \\ &= \nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})^T \nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}}) + (\nabla g_k(\hat{\theta}_{n_{k-1}})^T \nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}}) - \nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})^T \nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})) \\ &\geq ||\nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})||_2^2 - N\rho^{k+k_0} ||\nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})||_2. \end{split}$$

This, together with (19), implies

$$g_{k}(\tau) \geq g_{k}(\hat{\theta}_{n_{k-1}}) + \beta^{q} ||\nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})||_{2}^{2} - \frac{M\beta^{2q}}{2} ||\nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})||_{2}^{2} - N\beta^{q} \rho^{k+k_{0}} ||\nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})||_{2}^{2} \\ \geq g_{k}(\hat{\theta}_{n_{k-1}}) + \beta^{q} ||\nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})||_{2}^{2} - \frac{M\beta^{2q}}{2} ||\nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})||_{2}^{2} - NM\beta^{q} \rho^{k+k_{0}},$$

where the last inequality follows from (5). Note that for any non-negative integer $q \ge -\log M / \log \beta$, we have

$$\beta^q - \frac{M\beta^{2q}}{2} \ge \frac{1}{2}\beta^q > \alpha\beta^q,$$

which immediately implies that (16) fails; in other words, for any non-negative integer $q \ge T_1(k)$, we have

$$g_k(\hat{\theta}_{n_{k-1}} + \beta^q \nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})) \ge g_k(\hat{\theta}_{n_{k-1}}) + \alpha \beta^q ||\nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})||_2^2 - (N+M)M\beta^q \rho^{k+k_0}$$

as long as $q \ge -\log M / \log \beta$. It then follows that for any integer $k \ge 1$, T(k) can be bounded as

$$T(k) \le \begin{cases} A_2 & \text{if } T_1(k) \le A_2 \\ T_1(k) & \text{if } T_1(k) > A_2, \end{cases}$$
(20)

where $A_2 \triangleq \max\{0, -\log M/\log \beta + 1\}$ is a constant independent of k. Now, to prove the uniform boundedness of T(k), what remains is to show that there exists $A_1 \ge 0$ such that for all $k, T_1(k) \le A_1$.

From the definition of T(k), we have

$$g_k(\hat{\theta}_{n_k}) \ge g_k(\hat{\theta}_{n_{k-1}}) + \alpha \beta^{T(k)} ||\nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}})||_2^2 - (N+M)M\beta^{T(k)}\rho^{k+k_0}.$$
 (21)

Note that (4) and (21) together imply that

$$g_k(\hat{\theta}_{n_k}) \ge g_{k-1}(\hat{\theta}_{n_{k-1}}) - (N+M)M\rho^{k+k_0} - N\rho^{k+k_0}$$

By summation, we obtain that

$$g_k(\hat{\theta}_{n_k}) \ge g_0(\hat{\theta}_0) - \sum_{i=0}^{k-1} \left[(N+M)M\rho^{i+k_0+1} + N\rho^{i+k_0+1} \right]$$
$$\ge g_0(\hat{\theta}_0) - \left[\frac{(N+M)M\rho^{k_0+1}}{1-\rho} + \frac{N\rho^{k_0+1}}{1-\rho} \right].$$

It then follows from (4) that for all $k \ge 0$ we have

$$g_{0}(\hat{\theta}_{n_{k}}) \geq g_{0}(\hat{\theta}_{0}) - \left[\frac{(N+M)M\rho^{k_{0}+1}}{1-\rho} + \frac{N\rho^{k_{0}+1}}{1-\rho}\right] - \sum_{i=1}^{k} N\rho^{i+k_{0}}$$
$$\geq g_{0}(\hat{\theta}_{0}) - \left[\frac{(N+M)M\rho^{k_{0}+1}}{1-\rho} + \frac{2N\rho^{k_{0}+1}}{1-\rho}\right]$$
$$\geq g_{0}(\hat{\theta}_{0}) - \frac{\delta}{8}, \tag{22}$$

where the last inequality follows from Lemma 3.1 (a). Now, letting y_0, B_{k_0} and C_{k_0} be defined as in Lemma 3.1, we infer from (22) and Lemma 3.1 (c) that $\{\hat{\theta}_{n_k}\}_{k=0}^{\infty} \subseteq C_{k_0} \subseteq \Theta^{\circ}$. Hence, for any non-negative integer $p \ge \log(dist(C_{k_0}, \partial \Theta)/M)/\log \beta$, we have $\hat{\theta}_{n_{k-1}} + \beta^p \nabla \hat{g}_{n_{k-1}}(\hat{\theta}_{n_{k-1}}) \in \Theta^\circ$ and it then follows that $T_1(k) \le A_1$, where A_1 is defined as

$$A_1 \triangleq \max\left\{0, \frac{\log(dist(C_{k_0}, \partial\Theta)/M)}{\log\beta} + 1\right\}.$$
(23)

Finally, it follows from (20) and (23) that

$$T(k) \le A \triangleq \max\{A_1, A_2\},\tag{24}$$

as desired.

Step 2: Exponential convergence of $\{f(\hat{\theta}_{n_k})\}$. From (21), using (4), (5) and the definition of $\{\hat{g}_{n_k}\}_{k=0}^{\infty}$, we deduce that

$$f(\hat{\theta}_{n_k}) \ge f(\hat{\theta}_{n_{k-1}}) + \alpha \beta^{T(k)} ||\nabla f(\hat{\theta}_{n_{k-1}})||_2^2 - [(N+M)M\beta^{T(k)} + 2N + 2NM\rho]\rho^{k+k_0}.$$

According to (7), we have

$$f(\theta^*) \le f(\hat{\theta}_{n_{k-1}}) + \nabla f(\hat{\theta}_{n_{k-1}})^T (\theta^* - \hat{\theta}_{n_{k-1}}) - \frac{m}{2} ||\theta^* - \hat{\theta}_{n_{k-1}}||_2^2$$

which, coupled with some straightforward estimates, yields

$$2m(f(\theta^*) - f(\hat{\theta}_{n_{k-1}})) \le ||\nabla f(\hat{\theta}_{n_{k-1}})||_2^2$$

It then follows that

$$\begin{aligned} f(\theta^{*}) - f(\hat{\theta}_{n_{k}}) \\ &\leq f(\theta^{*}) - f(\hat{\theta}_{n_{k-1}}) - \alpha \beta^{T(k)} ||\nabla f(\hat{\theta}_{n_{k-1}})||_{2}^{2} + [(N+M)M\beta^{T(k)} + 2N + 2NM\rho]\rho^{k+k_{0}} \\ &\leq (1 - 2m\alpha\beta^{T(k)})(f(\theta^{*}) - f(\hat{\theta}_{n_{k-1}})) + [(N+M)M + 2N + 2NM\rho]\rho^{k+k_{0}} \\ &\stackrel{(d)}{\leq} \left(1 - \min\left\{2m\alpha\beta^{A_{1}}, 2m\alpha\beta^{A_{2}}\right\}\right)(f(\theta^{*}) - f(\hat{\theta}_{n_{k-1}})) + \left(\frac{NM + M^{2} + 2N}{\rho} + 2NM\right)\rho^{k+k_{0}+1} \\ &= \eta(f(\theta^{*}) - f(\hat{\theta}_{n_{k-1}})) + \gamma_{k}, \end{aligned}$$
(25)

where

$$\eta = 1 - \min\left\{2m\alpha, \frac{dist(C_{k_0}, \partial\Theta)}{M} 2m\alpha\beta, \frac{2m\alpha\beta}{M}\right\}, \quad \gamma_k = \left(\frac{NM + M^2 + 2N}{\rho} + 2NM\right)\rho^{k+k_0+1}$$

and (d) follows from (24). Recursively applying inequality (25) and noting that $0 < \eta < 1$, we infer that there exist $0 < \xi < 1$ and M' > 0 such that

$$f(\theta^*) - f(\hat{\theta}_{n_k}) \le M' \xi^k.$$
(26)

Step 3: Exponential convergence of $\{\hat{g}_n(\hat{\theta}_n)\}$. In this step, we establish (17) and thereby finish the proof.

First, note that for any positive integer $n \ge 0$, there exists an integer $k' \ge 0$ such that

$$n_{k'} \le n \le n_{k'+1}, \quad n \le (k'+1)A, \quad \hat{\theta}_n = \hat{\theta}_{n_{k'}}, \quad \hat{g}_n(\hat{\theta}_n) = \hat{g}_{n_{k'}}(\hat{\theta}_{n_{k'}}),$$

where A is defined in (24). These four inequalities, together with (4) and (26), imply the existence of $\hat{M} > 0$ and $0 < \hat{\xi} < 1$ such that for any $n \ge 0$,

$$\begin{aligned} |\hat{g}_n(\hat{\theta}_n) - f(\theta^*)| &\leq |\hat{g}_{n_{k'}}(\hat{\theta}_{n_{k'}}) - f(\hat{\theta}_{n_{k'}})| + |f(\hat{\theta}_{n_{k'}}) - f(\theta^*)| \\ &\leq N\rho^{k'+k_0} + M'\xi^{k'} \\ &\leq N\rho^{k_0}\rho^{\lfloor n/A \rfloor - 1} + M'\xi^{\lfloor n/A \rfloor - 1} \\ &\leq \hat{M}\hat{\xi}^n, \end{aligned}$$

which completes the proof of the theorem.

Theorem 3.10, together with the uniform boundedness of T(k) established in its proof, immediately implies that Algorithm 3.3 converges exponentially in k. More precisely, we have the following theorem.

Theorem 3.11. For a strongly concave function f whose unique maximum is achieved in Θ° , given in terms of the approximating sequence of functions $\{f_k\}_{k=0}^{\infty}$ as in (3), satisfying assumptions A.1, A.2 and A.3 in Section 2, there exist $\tilde{M} > 0$ and $0 < \tilde{\xi} < 1$ depending on m, M, N and ρ such that for all k,

$$|g_k(\theta_k) - f(\theta^*)| \le \tilde{M}\xi^k, \tag{27}$$

where $g_k(\theta_k)$ is defined as in Algorithm 3.3.

3.2 Applications of Algorithm 3.3

In this section, we discuss some applications of Algorithm 3.3 in information theory.

Consider a finite-state channel satisfying (2.a)-(2.c) and assume that all the matrices in $\Pi_{F,\delta}$ are analytically parameterized by $\theta \in \Theta^{\circ}$, where Θ is a compact convex subset of \mathbb{R}^d , $d \in \mathbb{N}$. Setting

$$f(\theta) = I(X(\theta); Y(\theta))$$

and

$$f_k(\theta) = H(X_2(\theta)|X_1(\theta)) + H(Y_{k+1}(\theta)|Y_1^k(\theta)) - H(X_{k+1}(\theta), Y_{k+1}(\theta)|X_1^k(\theta), Y_1^k(\theta)),$$

we derive from [11] that (4) holds. So, when $f(\theta)$ is strongly concave with respect to θ (this may hold true for some special channels, see, for example, [14] and [16]) as in (7), our algorithm applied to $\{f_k(\theta)\}_{k=0}^{\infty}$ converges exponentially fast in the number of steps to the maximum value of $f(\theta)$. This, and the easily verifiable fact that the computational complexity of $f_k(\theta)$ is at most exponential in k, leads to the conclusion that Algorithm 3.3, when applied to $\{f_k(\theta)\}_{k=0}^{\infty}$ as above, achieves exponential accuracy in exponential time. We now trade exponential time for polynomial time at the expense of accuracy. For any fixed $r \in \mathbb{R}_+$ and any $k \geq \lceil r \log 2 \rceil$, choose the largest $l \in \mathbb{N}$ such that $k = \lceil r \log l \rceil$. Substituting this into (27), we have

$$|g_{\lceil r \log l\rceil}(\theta_{\lceil r \log l\rceil}) - f(\theta^*)| \le \tilde{M}l^{r \log \xi}.$$

In other words, as summarized in the following theorem, we have shown that Algorithm 3.3, when used to compute the channel capacity as above, achieves polynomial accuracy in polynomial time.

Theorem 3.12. For a general finite-state channel satisfying (2.a)-(2.c) and parameterized as above, if $I(X(\theta); Y(\theta))$ is strongly concave with respect to $\theta \in \Theta$ and achieves its unique maximum in Θ° , then there exists an algorithm computing its fixed order Markov capacity that achieves polynomial accuracy in polynomial time.

In the following, we show that for certain special families of finite-state channels, we get a stronger convergence result than that in Theorem 3.12. In particular, for the following two examples, Algorithm 3.3 can be used to compute the channel capacity, achieving exponential accuracy in polynomial time.

3.2.1 A noisy channel with one state

In this section, we consider the Markov capacity of a binary erasure channel (BEC) under the $(1, \infty)$ -RLL constraint. This channel can be mathematically characterized by the inputoutput equation

$$Y_n = X_n \cdot E_n,\tag{28}$$

where $\{X_n\}_{n=1}^{\infty}$ is the input stationary Markov chain taking values in $\{1, 2\}$ such that $\{22\}$ is a *forbidden set* (see, e.g., [19]), and $\{E_n\}_{n=1}^{\infty}$ is an i.i.d. process taking values in $\{0, 1\}$ with

$$P(E_n = 0) = \varepsilon, \quad P(E_n = 1) = 1 - \varepsilon$$

for $0 < \varepsilon < 1$. Here we note that the BEC given above can be viewed as a degenerate finitestate channel with only one state. In the following, we will compare the channel capacity when $\{X_n\}_{n=1}^{\infty}$ is a first-order stationary Markov chain with the capacity when $\{X_n\}_{n=1}^{\infty}$ is a second-order stationary Markov chain. In particular, Algorithm 3.3 will be used to evaluate the first-order Markov capacity, which, compared to a lower bound for the second-order Markov capacity, will lead to the conclusion that higher order memory in the channel input may increase the Markov capacity.

For the first case, suppose that $\{X_n\}_{n=1}^{\infty}$ is a first-order stationary Markov chain with the transition probability matrix (indexed by 1, 2)

$$\Pi = \left[\begin{array}{cc} 1 - \theta & \theta \\ 1 & 0 \end{array} \right]$$

for $0 < \theta < 1$. It has been established in [17] that the mutual information rate $I(X(\theta); Y(\theta))$ of the BEC channel (28) can be computed as

$$I(X(\theta); Y(\theta)) = (1 - \varepsilon)^2 \sum_{l=0}^{\infty} H(X_{l+2}(\theta) | X_1(\theta)) \varepsilon^l,$$

which is strictly concave with respect to θ . Now, setting $f(\theta) = I(X(\theta); Y(\theta))$, one verifies, through straightforward computation, that

$$f(\theta) = \lim_{k \to \infty} f_k(\theta),$$

where

$$\begin{split} f_0(\theta) &= f_1(\theta) \triangleq (1-\varepsilon)^2 \frac{-\theta \log \theta - (1-\theta) \log(1-\theta)}{1+\theta}, \\ f_k(\theta) \triangleq (1-\varepsilon)^2 \frac{-\theta \log \theta - (1-\theta) \log(1-\theta)}{1+\theta} \\ &+ (1-\varepsilon)^2 \sum_{l=2}^k \left\{ \frac{1}{1+\theta} H\left(\frac{1-(-\theta)^{l+1}}{1+\theta}\right) \right\} \varepsilon^{l-1} + (1-\varepsilon)^2 \sum_{l=2}^k \left\{ \frac{\theta}{1+\theta} H\left(\frac{1-(-\theta)^l}{1+\theta}\right) \right\} \varepsilon^{l-1} \end{split}$$

for $k \ge 2$ and $H(p) \triangleq -p \log p - (1-p) \log(1-p)$ is the binary entropy function. In what follows, assuming $\varepsilon = 0.1$, we will show that Algorithm 3.3 can be applied to compute the first-order Markov capacity of the channel (28), i.e., the maximum of $f(\theta)$ over all $\theta \in [0, 1]$.

First of all, we claim that $f(\theta)$ achieves its unique maximum within the interval [0.25, 0.55] and therefore in the interior of $\Theta \triangleq [0.2, 0.6]$. To see this, noting that $f_k(\theta) \leq f(\theta)$ for any θ and through evaluating the elementary function $f_{100}(\theta)$, we have

$$0.442239 < \max_{\theta \in [0.25, 0.55]} f_{100}(\theta) < 0.442240$$

and therefore

$$\max_{\theta \in [0.25, 0.55]} f(\theta) \ge 0.442239,\tag{29}$$

where (29) follows from the fact that $f_k(\theta)$ is monotonically increasing in k. On the other hand, using the stationarity of $\{Y_n\}_{n=1}^{\infty}$ and the fact that conditioning reduces entropy, we have

$$f(\theta) = I(X(\theta); Y(\theta)) = H(Y) - H(\varepsilon) \le H(Y_3(\theta)|Y_1(\theta), Y_2(\theta)) - H(\varepsilon),$$

where H(Y) is the entropy rate of $\{Y_n\}_{n=1}^{\infty}$. Then, by straightforward computation, we deduce that

$$\max_{\theta \in [0,0.25] \cup [0.55,1]} f(\theta) \le \max_{\theta \in [0,0.25] \cup [0.55,1]} H(Y_3(\theta)|Y_1(\theta), Y_2(\theta)) - H(\varepsilon) < 0.414483,$$

which, together with (29), yields

$$\max_{\theta \in [0,0.25] \cup [0.55,1]} f(\theta) < \max_{\theta \in [0.25,0.55]} f(\theta),$$

as desired.

Next, we will verify that (4), (5) and (7) are satisfied for all $\theta \in [0.2, 0.6]$. Note that for $k \geq 2$ we have

$$f_k(\theta) - f_{k-1}(\theta) = (1-\varepsilon)^2 \left[\frac{1}{1+\theta} H\left(\frac{1-(-\theta)^{k+1}}{1+\theta}\right) + \frac{\theta}{1+\theta} H\left(\frac{1-(-\theta)^k}{1+\theta}\right) \right] \varepsilon^{k-1}.$$

This implies that for any $k \ge 5$ and any $\theta \in [0.2, 0.6]$,

$$|f_k(\theta) - f_{k-1}(\theta)| \le (1-\varepsilon)^2 \varepsilon^{k-1} = 8.1 \times 0.1^k.$$

This, together with the easily verifiable fact that $0.378 \leq f_5(\theta) \leq 0.443$ for $\theta \in [0.2, 0.6]$, further implies that

$$|f_k(\theta) - f(\theta)| \le 0.9 \times 0.1^k$$
 and $0.37 \le f_k(\theta) \le 0.45$

for all $k \geq 5$ and $\theta \in [0.2, 0.6]$.

Going through similar arguments, we obtain that, for any $k \ge 13$ and any $\theta \in [0.2, 0.6]$,

$$|f'_k(\theta) - f'_{k-1}(\theta)| \le 72.9 \times 0.1^k, \quad |f'_k(\theta) - f'(\theta)| \le 8.1 \times 0.1^k,$$

and

$$-0.44 \le f'_k(\theta) \le 0.76,$$

and, for any $k \ge 18$ and any $\theta \in [0.2, 0.6]$,

$$|f_k''(\theta) - f_{k-1}''(\theta)| \le 370.575 \times 0.1^k, \quad |f_k''(\theta) - f''(\theta)| \le 41.175 \times 0.1^k,$$

and

$$-5.81 \le f_k''(\theta) \le -1.88.$$

To sum up, we have shown that (4) is satisfied with N = 371 and $\rho = 0.1$, (5) is satisfied with M = 5.81 and (7) is satisfied with m = 1.88. Under these choices of the constants, direct calculation shows that $k_0 = 18$ is sufficient for Lemma 3.1. As a result, Algorithm 3.3 is applicable to the channel (28). Observing that, by its definition, the computational complexity of $f_k(\theta)$ is polynomial in k, we conclude that Algorithm 3.3 achieves exponential accuracy in polynomial time.

Now, applying Algorithm 3.3 to the sequence $\{f_k(\theta) : k \ge 18\}$ over $\Theta = [0.2, 0.6]$ with $\alpha = 0.4, \beta = 0.9$ and the initial point $\theta_0 = 0.5$, we obtain that

$$\theta_{110} \approx 0.395485, \quad f_{110}(\theta_{110}) \approx 0.442239.$$

Furthermore, under the settings given above, ξ and η can be chosen such that $\xi = \eta < 0.767$. It now follows from (4), (26) and $\hat{\theta}_{n_k} = \theta_k$ (see Remark 3.9) that

$$|f_{110}(\theta_{110}) - f(\theta^*)| \le |f_{110}(\theta_{110}) - f(\theta_{110})| + |f(\theta_{110}) - f(\theta^*)| \le 2.621 \times 10^{-7},$$

which further implies that when the input is a first-order Markov chain, the capacity of the BEC channel (28) can be bounded as

$$0.4422382 \le f(\theta^*) \le 0.4422398. \tag{30}$$

We now consider the case when the input is a second-order stationary Markov chain, whose transition probability matrix (indexed by 11, 12 and 21 only since 22 is prohibited by the $(1, \infty)$ -RLL constraint) is given by

$$\left[\begin{array}{ccc} p & 1-p & 0 \\ 0 & 0 & 1 \\ q & 1-q & 0 \end{array} \right],$$

where 0 < p, q < 1. For this case, from the Birch lower bound (see, e.g., Lemma 4.5.1 of [7]), we have

$$H(Y_6|Y_5, Y_4, Y_3, X_2, X_1) - H(\varepsilon) \le H(Y) - H(\varepsilon) = I(X; Y).$$

It can then be verified by direct computation that, when $p \approx 0.597275$ and $q \approx 0.614746$,

$$H(Y_6|Y_5, Y_4, Y_3, X_2, X_1) - H(\varepsilon) \approx 0.442329,$$

which is a lower bound on the second-order Markov capacity yet strictly larger than the upper bound on the first-order Markov capacity given in (30). Hence we can draw the conclusion that for the BEC channel with Markovian inputs under the $(1, \infty)$ -RLL constraint, an increase of the Markov order of the input process from 1 to 2 does increase the channel capacity.

3.2.2 A noiseless channel with two states

In this section, we consider a noiseless finite-state channel with two channel states, for which we show that Algorithm 3.3 can be applied to show that higher order memory can yield larger Markov capacity.

More precisely, the channel input $\{X_n\}_{n=1}^{\infty}$ is a first-order stationary Markov chain taking values from the alphabet $\mathcal{A} = \{0, 1\}$ and, except at time 0, the channel state $\{S_n\}_{n=1}^{\infty}$ is determined by the channel input, that is, $S_n = X_n$, $n = 1, 2, \ldots$ The channel is characterized by the following input-output equation:

$$Y_n = \phi(S_{n-1}, X_n), \quad n = 1, 2, \dots,$$
(31)

where ϕ is a deterministic function with $\phi(0,0) = 1, \phi(0,1) = 0, \phi(1,0) = 0$ and $\phi(1,1) = 0$. Note that ϕ naturally induces a sliding block code that maps the full \mathcal{A} -shift \mathcal{S} to the shift of finite type $\mathcal{S}_{\mathcal{F}}$, where the forbidden set \mathcal{F} is {101}. It can be readily verified that the Shannon capacity of (31) is equal to its stationary capacity [10], which can be computed as the largest eigenvalue of the adjacency matrix of the 3rd higher block shift of $\mathcal{S}_{\mathcal{F}}$ and is approximately equal to 0.562399 (see Chapter 4 and 13 of [19] for more details). In what follows, we will focus on the Markov capacity of (31); more specifically, we will compute the Markov capacity when the input $\{X_n\}_{n=1}^{\infty}$ is an i.i.d. process and a first-order stationary Markov chain, which will be compared with the Shannon capacity.

It can be easily verified that the mutual information rate of (31) can be computed as

$$I(X;Y) = \lim_{k \to \infty} H(Y_{k+1}|Y_1^k) - \frac{1}{k}H(Y_1^k|X_1^k) = \lim_{k \to \infty} H(Y_{k+1}|Y_1^k) = H(Y).$$

When $\{X_n\}_{n=1}^{\infty}$ is a stationary Markov chain, the output $\{Y_n\}_{n=1}^{\infty}$ is a hidden Markov chain with an unambiguous symbol whose entropy rate can be computed by the following formula [13]:

$$H(Y) = \sum_{n=1}^{\infty} P(Y_1^n = (1, \underbrace{0, \dots, 0}_{n-1})) H(Y_{n+1} | Y_1^n = (1, \underbrace{0, \dots, 0}_{n-1})).$$
(32)

This formula will play a key role in our analysis detailed below.

We first consider the degenerate case that $\{X_n\}_{n=1}^{\infty}$ is an i.i.d. process. Letting θ denote $P(X_1 = 0)$, we note that the Markov chain $\{(X_{n-1}, X_n)\}_{n=2}^{\infty}$ has the following transition probability matrix (indexed by 00, 01, 10, 11)

$$\left[\begin{array}{cccc} \theta & 1-\theta & 0 & 0 \\ 0 & 0 & \theta & 1-\theta \\ \theta & 1-\theta & 0 & 0 \\ 0 & 0 & \theta & 1-\theta \end{array} \right],$$

whose left eigenvector corresponding to the largest eigenvalue is

$$(\pi_1(\theta), \pi_2(\theta), \pi_3(\theta), \pi_4(\theta)) = (\theta^2, \theta(1-\theta), \theta(1-\theta), (1-\theta)^2).$$

Using (32), we have

$$H(Y) = -\sum_{l=0}^{\infty} \pi_1(\theta) \mathbf{r}(B_{\theta})^l \mathbf{1} \log \frac{\mathbf{r}(B_{\theta})^l \mathbf{1}}{\mathbf{r}(B_{\theta})^{l-1} \mathbf{1}} - \sum_{l=0}^{\infty} \pi_1(\theta) \mathbf{r}(B_{\theta})^{l-1} \mathbf{c} \log \frac{\mathbf{r}(B_{\theta})^{l-1} \mathbf{c}}{\mathbf{r}(B_{\theta})^{l-1} \mathbf{1}},$$

where $\mathbf{r} = (1 - \theta, 0, 0), \ \mathbf{c} = (0, \theta, 0)^T, \ \mathbf{1} = (1, 1, 1)^T,$

$$B_{\theta} = \left[\begin{array}{ccc} 0 & \theta & 1-\theta \\ 1-\theta & 0 & 0 \\ 0 & \theta & 1-\theta \end{array} \right],$$

and both $\mathbf{r}(B_{\theta})^{-1}\mathbf{1}, \mathbf{r}(B_{\theta})^{-1}\mathbf{c}$ should be interpreted as 1.

Setting $f(\theta) \triangleq H(Y)$, we note that

$$f(\theta) = \lim_{k \to \infty} f_k(\theta),$$

where

$$f_k(\theta) \triangleq -\sum_{l=0}^k \pi_1(\theta) \mathbf{r}(B_\theta)^l \mathbf{1} \log \frac{\mathbf{r}(B_\theta)^l \mathbf{1}}{\mathbf{r}(B_\theta)^{l-1} \mathbf{1}} - \sum_{l=0}^k \pi_1(\theta) \mathbf{r}(B_\theta)^{l-1} \mathbf{c} \log \frac{\mathbf{r}(B_\theta)^{l-1} \mathbf{c}}{\mathbf{r}(B_\theta)^{l-1} \mathbf{1}}, \quad k \ge 0.$$

Similarly as in the previous example, we can show that

$$\max_{\theta \in [0,0.41] \cup [0.89,1]} f(\theta) < \max_{\theta \in [0.41,0.89]} f(\theta),$$

which means that $f(\theta)$ will achieve its maximum within the interior of [0.4, 0.9]. Moreover, through tedious but similar evaluations as in the previous example, we can choose (below, rather than a constant, N is a polynomial in k, but the proof of Theorem 3.10 carries over almost verbatim)

$$k_0 = 120, \quad N = (374.945k^2 + 6207.73k + 46587.2), \quad \rho = 0.875, \quad m = 1.2, \quad M = 10.37.$$

Though the function $f(\theta)$ is not concave near $\theta = 0$, tedious yet straightforward computation indicates that $f''(\theta) \leq f''_{120}(\theta) + N\rho^{120} < 0$ for any $\theta \in [0.4, 0.9]$, which immediately implies

that $f(\theta)$ is strongly concave within the interior of the interval [0.4, 0.9]. Then, similarly as in Section 3.2.1, one verifies that, when applied to the channel in (31), Algorithm 3.3 achieves exponential accuracy in polynomial time.

Letting $\alpha = 0.4, \beta = 0.9$, we apply our algorithm to the sequence $\{f_k(\theta) : k \ge 120\}$ with $\Theta \triangleq [0.4, 0.9], \theta_0 = 0.5, \eta = \xi = 0.901061$, and we obtain that

$$\theta_{450} \approx 0.6257911, \quad f_{450}(\theta_{450}) \approx 0.4292892.$$

Now from (4), (26) and the fact that $\hat{\theta}_{n_k} = \theta_k$, we conclude

$$|f_{450}(\theta_{450}) - f(\theta^*)| \le |f_{450}(\theta_{450}) - f(\theta_{450})| + |f(\theta_{450}) - f(\theta^*)| \le 0.0001745$$

which further implies

$$0.4291146 \le f(\theta^*) \le 0.4294638 \tag{33}$$

for the i.i.d. case.

Now, we consider the case that $\{X_n\}_{n=1}^{\infty}$ is a genuine first-order stationary Markov process, and assume the Markov chain $\{(X_{n-1}, X_n)\}_{n=2}^{\infty}$ has the following transition probability matrix (indexed by 00, 01, 10, 11)

$$\begin{pmatrix} p & 1-p & 0 & 0 \\ 0 & 0 & q & 1-q \\ p & 1-p & 0 & 0 \\ 0 & 0 & q & 1-q \end{pmatrix},$$

where 0 < p, q < 1. Again, straightforward computation shows that for $p \approx 0.674521, q \approx 0.595176$, $H(Y_4|Y_3, X_2, X_1)$ is approximately 0.513259, which gives a lower bound on H(Y). Comparing this lower bound with the upper bound in (33), we conclude that the capacity is increased when increasing the Markov order of the input from 0 to 1. Finally, we also point out that direct evaluation of a trivial upper bound (for the first order Markov capacity of (31)) gives

$$\max_{p,q} H(Y_6|Y_5, Y_4, Y_3, Y_2, Y_1) \approx 0.548481 \quad \text{for} \quad p \approx 0.629902, q \approx 0.734121.$$

Comparing this upper bound with 0.562399, the Shannon capacity given at the beginning of this section, we also conclude that the Shannon capacity of (31) cannot be achieved by any first-order Markov input.

4 The Second Algorithm: without Concavity

In this section, we consider the optimization problem (3) for the case when f may not be concave.

For a non-convex optimization problem with a continuously differentiable target function f and a bounded domain, conventionally there are two major methods for finding its solution: the Frank-Wolfe method [3] and the method through the Lojasiewicz inequality (see, e.g., [1]). However, both of these methods in general tend to fail in our setting: for the Frank-Wolfe method, the computation for finding the feasible ascent direction and the verification

of the relevant gradient condition (which is necessary for the convergence of this method) both depend on the existence of an exact formula for ∇f and a tractable description of Θ , which is however not available in our case; on the other hand, due to the fact that our target function is the limit of a sequence of approximating functions, the method through the Lojasiewicz inequality necessitates a "uniform" version of the Lojasiewicz inequality over all sequences of approximating functions, which does not seem to hold true in our setting.

Motivated by Algorithm 3.3, we propose in the following our second algorithm to efficiently solve the optimization problem (3) whose target function may not be concave. Except for using the sequence $\{\nabla f_k\}_{k=0}^{\infty}$ as the ascent direction in each iteration, an additional check condition is proposed for the choice of the step size. This check condition is chosen carefully to ensure an appropriate pace for the decay of ∇f_k , which turns out to be crucial for the convergence of this algorithm.

Similarly as in Section 3, we need the following lemma before presenting our second algorithm.

Lemma 4.1. Assume the function f has s stationary points $\{\theta_i^*\}_{i=1}^s$ which are all contained in Θ° , and that f achieves its maximum in Θ° . If, for each k, f_k also has finitely many stationary points which are all contained in Θ° , then there exists a non-negative integer k_0 such that

(a)
$$\rho^{1/3} + \rho^{2k_0/3} < 1$$
 and $\frac{2N\rho^{k_0}}{1-\rho} \leq \frac{\delta}{8}$, where $\delta \triangleq \max_{\substack{\theta_i^*: 1 \leq i \leq s}} f(\theta_i^*) - \max_{\theta \in \partial \Theta} f(\theta) > 0$;

(b) There exists $y_0 \in \mathbb{R}$ such that for any fixed b with 0 < b < 1, we have

$$\emptyset \subsetneq B_{k_0} \subseteq C_{k_0} \subseteq \Theta^\circ, \quad A_{k_0} \cap B_{k_0} \neq \emptyset \quad and \quad dist(C_{k_0}, \partial \Theta) > 0,$$

where

$$A_{k_0} \triangleq \left\{ x \in \Theta^\circ : ||\nabla f_{k_0}(x)||_2 \ge \frac{2N\rho^{k_0/3}}{1-b} \right\},$$

$$B_{k_0} \triangleq \left\{ x \in \Theta : f_{k_0}(x) \ge y_0 \right\},$$

$$C_{k_0} \triangleq \left\{ x \in \Theta : f_{k_0}(x) \ge y_0 - \frac{\delta}{8} \right\}.$$

Note that A_{k_0} depends on b, whereas B_{k_0} and C_{k_0} do not.

Proof. By replacing what was assumed to be the unique maximum of f with $\max_{\theta_i^*:1 \le i \le s} f(\theta_i^*)$, a similar argument as in the proof of Lemma 3.1(*a*) yields that there exists $y_0 < y^* - \frac{\delta}{4}$ such that for all sufficiently large $k, \emptyset \subsetneq B_k \subseteq C_k \subseteq \Theta^\circ$ and $dist(C_k, \Theta^c) > 0$, where

$$y^* = \max_{\theta_i^*: 1 \le i \le s} f(\theta_i^*), \quad B_k \triangleq \{x \in \Theta : f_k(x) \ge y_0\} \quad \text{and} \quad C_k \triangleq \left\{x \in \Theta : f_k(x) \ge y_0 - \frac{\delta}{8}\right\}.$$

Now, for any k and any fixed 0 < b < 1, let

$$A_k \triangleq \left\{ x \in \Theta^\circ : ||\nabla f_k(x)||_2 \ge \frac{2N\rho^{k/3}}{1-b} \right\}.$$

We claim that for all large enough $k, A_k \cap B_k \neq \emptyset$. To see this, define

$$D_k \triangleq \left\{ x \in \Theta^\circ : ||\nabla f(x)||_2 \ge \frac{2N\rho^{k/3}}{1-b} + N\rho^k \right\} \quad \text{and} \quad B' \triangleq \left\{ x \in \Theta : f(x) \ge y_0 + \frac{\delta}{8} \right\}.$$

It then follows from (4), the continuity of f and the fact $y_0 + \delta/8 < y^*$ that $D_k \subseteq A_k, B' \subseteq B_k$ for all large enough k and B' has a non-empty interior. Observing that D_k^c converges to the finite set consisting of all stationary points of f, we deduce that $D_k \cap B' \neq \emptyset$ and therefore $A_k \cap B_k \neq \emptyset$ for sufficiently large k and therefore establish the claim. Finally, it immediately follows from this claim and the observation that (a) trivially holds for k_0 sufficiently large that there exists k_0 such that (a) and (b) are both satisfied.

Recalling that f and each f_k are assumed to have finitely many stationary points in Θ° , we now present our second algorithm.

Algorithm 4.2. (The second modified gradient ascent algorithm)

Step 0. Choose k_0 , y_0 and 0 < b < 1 such that the conditions in Lemma 4.1 are satisfied. Set k = 0, $g_0 = f_{k_0}$ and choose $\alpha \in (0, 0.5)$, $\beta \in (0, 1)$, $\theta_0 \in A_{k_0} \cap B_{k_0}$ where A_{k_0} and B_{k_0} are defined as in Lemma 4.1.

Step 1. Increase k by 1. Set t = 1 and $g_k = f_{k+k_0}$. Step 2. Set

$$\tau = \theta_{k-1} + t\nabla g_{k-1}(\theta_{k-1}).$$

If $\tau \notin \Theta^{\circ}$ or

$$||\nabla g_k(\tau)||_2 < \frac{2N\rho^{k/3}}{1-b}$$

or

$$g_k(\tau) < g_k(\theta_{k-1}) + \alpha t ||\nabla g_{k-1}(\theta_{k-1})||_2^2$$

set $t = \beta t$ and go to Step 2, otherwise set $\theta_k = \tau$ and go to Step 1.

Remark 4.3. The constants in Step 0 are chosen to ensure the convergence of the algorithm. And the existence of θ_0 follows from Lemma 4.1 (b).

Remark 4.4. In Step 2, for any feasible k, one of the necessary conditions for updating the value of θ_k is

$$\|\nabla g_k(\tau)\|_2 \ge \frac{2N\rho^{k/3}}{1-b}.$$

This is a key condition imposed to make sure that $\|\nabla g_k(\tau)\|$ is not too small and thereby the algorithm will not prematurely converge to a non-stationary point.

4.1 Convergence Analysis

To conduct the convergence analysis of Algorithm 4.2, we need to reformulate the algorithm via possible relabelling of the functions $\{g_k\}_{k=0}^{\infty}$ and iterates $\{\theta_k\}_{k=0}^{\infty}$ similarly as in Section 3.1. For ease of presentation only, we assume in the remainder of this section that such

a relabelling is not needed and thereby k actually records the number of times that Step 2 has been executed.

The following theorem asserts the convergence of Algorithm 4.2 under some regularity conditions.

Theorem 4.5. Under the same assumptions as in Lemma 4.1,

$$\lim_{k \to \infty} g_k(\theta_k) \text{ exists and } \|\nabla g_k(\theta_k)\|_2 \to 0,$$

where $g_k(\theta_k)$ is defined in Algorithm 4.2.

Proof. Similarly as in Section 3.1, define

$$T_{1}(k) \triangleq \inf\{p \in \mathbb{Z} : \theta_{k-1} + \beta^{p} \nabla g_{k-1}(\theta_{k-1}) \in \Theta^{\circ}\},$$

$$\hat{T}(k) \triangleq \inf\left\{q \in \mathbb{Z} : q \ge T_{1}(k), \ ||\nabla g_{k}(\theta_{k-1} + \beta^{q} \nabla g_{k-1}(\theta_{k-1}))||_{2} \ge \frac{2N\rho^{(k+k_{0})/3}}{1-b}\right\},$$

$$T(k) \triangleq \inf\{r \in \mathbb{Z} : r \ge \hat{T}(k), \ g_{k}(\theta_{k-1} + \beta^{r} \nabla g_{k-1}(\theta_{k-1})) \ge g_{k}(\theta_{k-1}) + \alpha\beta^{r} ||\nabla g_{k-1}(\theta_{k-1})||_{2}^{2}\},$$

and

$$T_2(k) \triangleq \hat{T}(k) - T_1(k), \quad T_3(k) \triangleq T(k) - \hat{T}(k).$$

In other words, for each k, $T_1(k)$ can be regarded as the number of times that Step 2 of Algorithm 4.2 has been executed before the condition $\tau \in \Theta^{\circ}$ is met; $T_2(k)$ can be regarded as the number of additional times that Step 2 of Algorithm 4.2 has been executed before the condition

$$||\nabla g_k(\theta_{k-1} + \beta^q \nabla g_{k-1}(\theta_{k-1}))||_2 \ge \frac{2N\rho^{(k+k_0)/3}}{1-b}$$

is also met while $T_3(k)$ can be regarded as the number of additional times that Step 2 of Algorithm 4.2 has been executed before the Armijo condition

$$g_k(\theta_{k-1} + \beta^r \nabla g_{k-1}(\theta_{k-1})) \ge g_k(\theta_{k-1}) + \alpha \beta^r ||\nabla g_{k-1}(\theta_{k-1})||_2^2$$

is also met. The well-definedness of $\hat{T}(k)$ is based on the fact that if $\theta_{k-1} + \beta^p \nabla g_{k-1}(\theta_{k-1}) \in \Theta^\circ$ for some non-negative integer p, then the same inequality also holds for any integer p' > p; and the well-definedness of T(k) will be postponed to Step 2 of this proof detailed below.

The remainder of the proof consists of 5 steps, with the first three devoted to establishing the uniform boundedness of $T_1(k)$, $T_2(k)$ and $T_3(k)$ and thus that of T(k).

Step 1: Uniform boundedness of $T_2(k)$. As in the proof of Theorem 3.10, it can be readily verified that $T_1(k) < \infty$ for all $k \ge 0$. Hence, when considering $T_2(k)$, we assume that $\tau = \theta_{k-1} + \beta^q \nabla g_{k-1}(\theta_{k-1})$ is already in Θ° .

In order to prove the uniform boundedness of $T_2(k)$, we proceed by way of induction. First of all, by the definition of g_0 and the choice of θ_0 , we have $||\nabla g_0(\theta_0)||_2 \ge 2N\rho^{k_0/3}/(1-b)$. Now, assuming that for some $k = 1, 2, \ldots$,

$$||\nabla g_{k-1}(\theta_{k-1})||_2 \ge \frac{2N\rho^{(k_0+k-1)/3}}{1-b},\tag{34}$$

we will derive a sufficient condition on β^q such that $||\nabla g_k(\tau)||_2 \ge 2N\rho^{(k_0+k)/3}/(1-b)$, where we recall that τ is defined as

$$\tau = \theta_{k-1} + \beta^q \nabla g_{k-1}(\theta_{k-1}). \tag{35}$$

To this end, we first note that by the Taylor series expansion, there exist ξ and $\hat{\xi}$ in Θ° such that

$$g_k(\tau) - g_k(\theta_{k-1}) = \nabla g_k(\tau)^T (\tau - \theta_{k-1}) - (\theta_{k-1} - \tau)^T \frac{\nabla^2 g_k(\xi)}{2} (\theta_{k-1} - \tau)$$

and

$$g_k(\tau) - g_k(\theta_{k-1}) = \nabla g_k(\theta_{k-1})^T (\tau - \theta_{k-1}) + (\tau - \theta_{k-1})^T \frac{\nabla^2 g_k(\hat{\xi})}{2} (\tau - \theta_{k-1}),$$

which immediately imply that

$$\nabla g_{k}(\tau)^{T}(\tau - \theta_{k-1}) - (\theta_{k-1} - \tau)^{T} \frac{\nabla^{2} g_{k}(\xi)}{2} (\theta_{k-1} - \tau)$$
$$= \nabla g_{k}(\theta_{k-1})^{T} (\tau - \theta_{k-1}) + (\tau - \theta_{k-1})^{T} \frac{\nabla^{2} g_{k}(\hat{\xi})}{2} (\tau - \theta_{k-1}).$$
(36)

Noting that $||\nabla^2 g_k(\xi)||_2 \leq M$ for all $\xi \in \Theta^\circ$ and

$$||\nabla g_k(\theta) - \nabla g_{k-1}(\theta)||_2 = ||\nabla f_{k+k_0}(\theta) - \nabla f_{k+k_0-1}(\theta)||_2 \le N\rho^{k+k_0}$$
(37)

for all $\theta \in \Theta^{\circ}$, we deduce from (36) that

$$\begin{aligned} ||\nabla g_k(\tau)||_2 \cdot ||\tau - \theta_{k-1}||_2 &\geq \nabla g_k(\theta_{k-1})^T (\tau - \theta_{k-1}) - M ||\tau - \theta_{k-1}||_2^2 \\ &\geq \nabla g_{k-1}(\theta_{k-1})^T (\tau - \theta_{k-1}) - N\rho^{k+k_0} ||\tau - \theta_{k-1}||_2 - M ||\tau - \theta_{k-1}||_2^2. \end{aligned}$$
(38)

Clearly, it follows from (35) that the vectors $\nabla g_{k-1}(\theta_{k-1})$ and $\tau - \theta_{k-1}$ have the same direction, which means that (38) can be rewritten as

$$||\nabla g_k(\tau)||_2 \cdot ||\tau - \theta_{k-1}||_2 \ge ||\nabla g_{k-1}(\theta_{k-1})||_2 \cdot ||\tau - \theta_{k-1}||_2 - N\rho^{k+k_0}||\tau - \theta_{k-1}||_2 - M||\tau - \theta_{k-1}||_2^2.$$

Simplifying this inequality, we have

$$||\nabla g_{k}(\tau)||_{2} \geq (1 - M\beta^{q})||\nabla g_{k-1}(\theta_{k-1})||_{2} - N\rho^{k+k_{0}}$$

$$\geq (1 - M\beta^{q})\frac{2N\rho^{(k_{0}+k-1)/3}}{1 - b} - N\rho^{k+k_{0}}.$$
(39)

Now, using the fact $1 - \rho^{1/3} - \rho^{2k_0/3} > 0$ (see Lemma 4.1 (a)), (34) and (39), we conclude that the condition

$$\beta^q \le \frac{1 - \rho^{1/3} - \rho^{2k_0/3}}{M} \tag{40}$$

is sufficient for $||\nabla g_k(\tau)||_2 \ge 2N\rho^{(k+k_0)/3}/(1-b)$. In other words, the induction argument successfully proceeds as long as (40) holds, and therefore $T_2(k)$ can be uniformly bounded as below:

$$T_2(k) \le \max\left\{0, \frac{\log\left(\left(1-\rho^{1/3}-\rho^{2k_0/3}\right)/M\right)}{\log\beta}+1\right\}.$$
 (41)

Step 2: Uniform boundedness of $T_3(k)$. First note that from (40), if the inequality

$$||\nabla g_k(\theta_{k-1} + \beta^q \nabla g_{k-1}(\theta_{k-1}))||_2 \ge \frac{2N\rho^{(k+k_0)/3}}{1-b}$$

holds for some non-negative integer q, then it remains true for all integers q' > q. This observation justifies the well-definedness of $T_3(k)$. Moreover, due to the boundedness of $T_2(k)$ for each k (in fact, it is uniformly bounded), we can assume without loss of generality that $||\nabla g_k(\tau)||_2 \geq 2N\rho^{(k+k_0)/3}/(1-b)$ is already satisfied, where $\tau = \theta_{k-1} + \beta^r \nabla g_{k-1}(\theta_{k-1})$, before we proceed to establish the uniform boundedness of $T_3(k)$.

By the Taylor series expansion formula and (37), we have

$$g_{k}(\tau) \geq g_{k}(\theta_{k-1}) + \beta^{r} \nabla g_{k}(\theta_{k-1})^{T} \nabla g_{k-1}(\theta_{k-1}) - \frac{M\beta^{2r}}{2} ||\nabla g_{k-1}(\theta_{k-1})||_{2}^{2}$$

$$\geq g_{k}(\theta_{k-1}) + \beta^{r} ||\nabla g_{k-1}(\theta_{k-1})||_{2}^{2} - \frac{M\beta^{2r}}{2} ||\nabla g_{k-1}(\theta_{k-1})||_{2}^{2} - N\rho^{k+k_{0}}\beta^{r} ||\nabla g_{k-1}(\theta_{k-1})||_{2},$$

where $\tau = \theta_{k-1} + \beta^r \nabla g_{k-1}(\theta_{k-1})$. It then follows that the condition

$$\beta^{r} \leq \frac{1}{M} - \frac{2N\rho^{k+k_{0}-1}}{M||\nabla g_{k-1}(\theta_{k-1})||_{2}}$$

is sufficient to ensure that

$$g_k(\tau) \ge g_k(\theta_{k-1}) + \alpha \beta^r ||\nabla g_{k-1}(\theta_{k-1})||_2^2.$$
(42)

Recalling that

$$||\nabla g_{k-1}(\theta_{k-1})||_2 \ge \frac{2N\rho^{(k+k_0-1)/3}}{1-b} \ge \frac{2N\rho^{k+k_0-1}}{1-b}$$

with 0 < b < 1, we deduce that the condition $\beta^r \leq b/M$ is sufficient for (42). In other words, we have

$$T_3(k) \le \max\left\{0, \frac{\log b - \log M}{\log \beta} + 1\right\}.$$
(43)

Step 3: Uniform boundedness of $T_1(k)$ and T(k). In this step, we will show that $T_1(k)$ is uniformly bounded over all k. This, together with the established fact that $T_2(k)$ and $T_3(k)$ are both uniformly bounded, immediately implies the uniform boundedness of T(k) over all k.

From Algorithm 4.2, $g_k(\theta_k) \ge g_k(\theta_{k-1}) + \alpha t ||\nabla g_{k-1}(\theta_{k-1})||_2^2$ for all $k \ge 0$, where $\theta_k = \theta_{k-1} + \beta^{T(k)} \nabla g_{k-1}(\theta_{k-1})$. Using (4), we have

$$g_0(\theta_k) \ge g_k(\theta_k) - \frac{N\rho^{k_0+1}}{1-\rho} \ge g_k(\theta_{k-1}) - \frac{N\rho^{k_0+1}}{1-\rho} \ge g_{k-1}(\theta_{k-1}) - N\rho^{k+k_0} - \frac{N\rho^{k_0+1}}{1-\rho},$$

from which we arrive at

$$g_0(\theta_k) \ge g_0(\theta_0) - \sum_{k=1}^{\infty} N\rho^{k+k_0} - \frac{N\rho^{k_0+1}}{1-\rho} \ge g_0(\theta_0) - \frac{2N\rho^{k_0}}{1-\rho},$$
(44)

for all $k \ge 0$. Recalling from Lemma 4.1 and Step 0 of Algorithm 4.2 that

$$\theta_0 \in B_{k_0} = \{ x \in \Theta : f_{k_0}(x) \ge y_0 \} = \{ x \in \Theta : g_0(x) \ge y_0 \},\$$

we deduce from (44) and Lemma 4.1 that for all $k \ge 0$,

$$\theta_k \in \left\{ x : g_0(x) \ge y_0 - \frac{2N\rho^{k_0}}{1-\rho} \right\} \subseteq C_{k_0} \subseteq \Theta^\circ,$$

where C_{k_0} is defined in Lemma 4.1 (b) and $dist(C_{k_0}, \partial \Theta) > 0$. Hence, for any non-negative integer p such that $p \geq \log(dist(C_{k_0}, \partial \Theta)/M)/\log\beta$, we have $\theta_{k-1} + \beta^p \nabla g_k(\theta_{k-1}) \in \Theta^\circ$, establishing the following uniform bound

$$T_1(k) \le \max\left\{0, \frac{\log(dist(C_{k_0}, \Theta^c)/M)}{\log \beta} + 1\right\}.$$
(45)

Finally, it is clear from (41), (43), (45) and the definition of T(k) that there exists a non-negative integer B such that, for all k,

$$T(k) \le B. \tag{46}$$

Step 4: Convergence of $g_k(\theta_k)$.

It follows from (4), (46) and the fact $||\nabla g_{k-1}(\theta_{k-1})||_2 \ge 2N\rho^{(k+k_0-1)/3}/(1-b)$ that

$$g_{k}(\theta_{k}) \geq g_{k}(\theta_{k-1}) + \alpha \beta^{B+1} ||\nabla g_{k-1}(\theta_{k-1})||_{2}^{2}$$

$$\geq g_{k-1}(\theta_{k-1}) + \alpha \beta^{B+1} ||\nabla g_{k-1}(\theta_{k-1})||_{2}^{2} - N \rho^{k+k_{0}}$$

$$\geq g_{k-1}(\theta_{k-1}) + \frac{4\alpha \beta^{B+1} N^{2} \rho^{2(k+k_{0}-1)/3}}{(1-b)^{2}} - N \rho^{k+k_{0}}.$$

Observing that if k is large enough,

$$\frac{4\alpha\beta^{B+1}N^2\rho^{2(k+k_0-1)/3}}{(1-b)^2} \ge N\rho^{k+k_0},$$

we deduce that $g_k(\theta_k) \ge g_{k-1}(\theta_{k-1})$ for sufficiently large k. Noting that (4) and the definition of g_k imply that there exists C > 0 such that $g_k(\theta_k) \le C$ for all k, we conclude that $\lim_{k\to\infty} g_k(\theta_k)$ exists.

Step 5:
$$\|\nabla g_k(\theta_k)\|_2 \to 0$$
.
Since $g_k(\theta_k) \ge g_{k-1}(\theta_{k-1}) + \alpha \beta^{B+1} ||\nabla g_{k-1}(\theta_{k-1})||_2^2 - N \rho^{k+k_0}$, we have
 $\sum_{k=1}^{n-1} \alpha \beta^{B+1} ||\nabla g_{k-1}(\theta_{k-1})||_2^2 \le g_n(\theta_n) - g_0(\theta_0) + \sum_{k=1}^{n-1} N \rho^{k+k_0}$,

which, together with the uniform boundedness of $\{g_k(\theta_k)\}_{k=0}^{\infty}$, yields

$$\sum_{k=1}^{\infty} \alpha \beta^{B+1} ||\nabla g_{k-1}(\theta_{k-1})||_2^2 < \infty.$$

Hence, $\lim_{n\to\infty} ||\nabla g_{k-1}(\theta_{k-1})||_2 = 0$. The proof of the theorem is thus complete.

4.2 A noisy channel with two states: Gilbert-Elliott Channel

In this section, we consider a Gilbert-Elliott channel with a first-order Markovian input under the $(1, \infty)$ -RLL constraint. To be more specific, let \oplus denote binary addition and $\{S_n\}_{n=0}^{\infty}$ be the state process which is a binary stationary Markov chain with the transition probability matrix

$$\left[\begin{array}{rrr} 0.7 & 0.3 \\ 0.3 & 0.7 \end{array}\right].$$

We focus on the Gilbert-Elliott channel characterized by the input-output equation

$$Y_n = X_n \oplus E_n,\tag{47}$$

where $\{X_n\}_{n=1}^{\infty}$ is a binary first-order stationary Markov chain independent of $\{S_n\}_{n=1}^{\infty}$ with the transition probability matrix

$$\begin{bmatrix} 1-\theta & \theta \\ 1 & 0 \end{bmatrix},$$

and $\{E_n\}_{n=1}^{\infty}$ is the noise process given by

$$E_n = \begin{cases} 0, & \text{with probability 0.99,} \\ 1, & \text{with probability 0.01,} \end{cases}$$

when $S_{n-1} = 0$ and

$$E_n = \begin{cases} 0, & \text{with probability } 0.9, \\ 1, & \text{with probability } 0.1, \end{cases}$$

when $S_{n-1} = 1$. In other words, at time *n*, if the channel state takes the value 0, the channel is a binary symmetric channel (BSC) with crossover probability 0.01, and if the channel state takes the value 1, it is a BSC with crossover probability 0.1. It is worth noting that E_n and E_{n-1} are not statistically independent for this channel.

It can be readily checked that the aforementioned channel is a finite-state channel characterized by

 $p(y_n, s_n | x_n, s_{n-1}) = p(y_n | x_n, s_{n-1}) p(s_n | s_{n-1})$

and the mutual information rate can be computed as

$$I(X(\theta); Y(\theta)) = \lim_{k \to \infty} H(Y_k(\theta) | Y_1^{k-1}(\theta)) - H(E_k(\theta) | E_1^{k-1}(\theta)).$$

The concavity of $I(X(\theta); Y(\theta))$ with respect to θ is not known to the best of our knowledge, yet Algorithm 4.2 can be applied to effectively maximize it. More specifically, setting

$$f_k(\theta) = H(Y_k(\theta)|Y_1^{k-1}(\theta)) - H(E_k(\theta)|E_1^{k-1}(\theta)),$$

we have applied Algorithm 4.2 with the initial point $\theta_0 = 0.2$ and we have obtained the following simulation results, from which one can observe fast convergence of the algorithm:



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