# Analysis of the Convergence Speed of the Arimoto-Blahut Algorithm by the Second-Order Recurrence Formula 

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#### Abstract

In this paper, we investigate the convergence speed of the Arimoto-Blahut algorithm. For many channel matrices, the convergence speed is exponential, but for some channel matrices it is slower than exponential. By analyzing the Taylor expansion of the defining function of the Arimoto-Blahut algorithm, we will make the conditions clear for the exponential or slower convergence. The analysis of the slow convergence in this paper is new. Based on this analysis, we will compare the convergence speeds of the Arimoto-Blahut algorithm numerically with the values obtained in our theorems for several channel matrices. The purpose of this paper is to obtain a complete understanding of the convergence speed of the Arimoto-Blahut algorithm.


Index Terms-channel capacity, discrete memoryless channel, Arimoto-Blahut algorithm, convergence speed, Hessian matrix, second-order recurrence formula.

## I. Introduction

We consider a discrete memoryless channel with the input and output alphabets of finite sizes. The channel capacity $C$ of the channel is defined as the maximum value of the mutual information between the input and output, where the maximum is taken over the entire input probability distributions. By Shannon's channel coding theorem, $C$ is equal to the maximum information rate that can be transmitted with arbitrarily small error probability.

Methods for calculating the value of the channel capacity $C$ fall into two categories. One is the direct method of solving the mutual information maximization problem [14], [15]. This is a convex optimization problem regarding the input distribution and is generally considered to be solved by the method of Lagrange multipliers. However, the constraint that the probability must be between 0 and 1 should be carefully considered, hence the channels where the channel capacity can be calculated explicitly are limited.

The other is a sequential computation algorithm represented by the Arimoto-Blahut algorithm [2], [3], [4]. In the ArimotoBlahut algorithm, a series of input probability distributions

[^0]are computed by a recurrence formula, and they converge to the input probability distribution that achieves the channel capacity $C$. The Arimoto-Blahut algorithm can be applied to any channel, and its capacity is computed.

Based on the original Arimoto-Blahut algorithm, much work has been done. In [12], [17], [20], applicable channels were extended to include more than just discrete memoryless channels. In [11], [21], algorithms to accelerate the ArimotoBlahut algorithm were proposed. Furthermore, in [8], [11], [13], [14], the Arimoto-Blahut algorithm is characterized by divergence geometry.
In this paper, we study the convergence speed of the Arimoto-Blahut algorithm. Although the convergence speed is exponential for many channels, we found that for some channels, the convergence speed is slower than exponential. Exponential convergence has been studied in [2], [3], [21], but no study of convergence speed other than exponential has been conducted.

## A. Purpose of Study

The purpose of our study is to classify discrete memoryless channels and to determine when the Arimoto-Blahut algorithm converges quickly and when it converges slowly.

First, we show that only two kinds of convergence speeds occur. One is the exponential convergence, which is a fast convergence, and the other is the convergence of order $O(1 / N)$, which is a slow convergence. Then, we evaluate the convergence speed theoretically in each case.

Exponential convergence is evaluated by the eigenvalues of the Jacobian matrix of the defining function $F$ of the ArimotoBlahut algorithm. Its analysis is not difficult. On the other hand, a convergence speed of order $O(1 / N)$ is determined not only by the Jacobian matrix of $F$ but also requires the analysis of the Hessian matrix, i.e., the second derivative of $F$. The convergence speed of the order $O(1 / N)$ is analyzed for the first time in this paper. The Hessian matrix of $F$ yields a second-order nonlinear recurrence formula, and then we will investigate the convergence speed of the recurrence formula. In general, a nonlinear recurrence formula can rarely be solved in an explicit form, and the recurrence formula in this paper also cannot be solved explicitly. Therefore, a method to investigate the speed of convergence must be devised. The key lemma is Lemma 6 on the convergence speed of a recurrence formula defined by a quadratic polynomial in one variable. In this
paper, we also investigate the speed at which the mutual information converges to the channel capacity $C$.

Based on these analyses, we numerically calculate the convergence speeds of the Arimoto-Blahut algorithm for several channels and compare them with the theoretical values obtained in this paper.

## II. Arimoto-Blahut algorithm

## A. Channel matrix and channel capacity

Consider a discrete memoryless channel $X \rightarrow Y$ with the input source $X$ and the output source $Y$. Let $\mathcal{X}=$ $\left\{x_{1}, \ldots, x_{m}\right\}$ be the input alphabet and $\mathcal{Y}=\left\{y_{1}, \ldots, y_{n}\right\}$ be the output alphabet. The conditional probability that the output symbol $y_{j}$ is received when the input symbol $x_{i}$ is transmitted is denoted by $P_{j}^{i}=P\left(Y=y_{j} \mid X=x_{i}\right), i=$ $1, \ldots, m, j=1, \ldots, n$, and the row vector $P^{i}$ is defined by $P^{i}=\left(P_{1}^{i}, \ldots, P_{n}^{i}\right), i=1, \ldots, m$. The channel matrix $\Phi$ is defined by

$$
\Phi=\left(\begin{array}{c}
P^{1}  \tag{1}\\
\vdots \\
P^{m}
\end{array}\right)=\left(\begin{array}{ccc}
P_{1}^{1} & \ldots & P_{n}^{1} \\
\vdots & & \vdots \\
P_{1}^{m} & \ldots & P_{n}^{m}
\end{array}\right)
$$

We assume that for any $j(j=1, \ldots, n)$, there exists at least one $i(i=1, \ldots, m)$ with $P_{j}^{i}>0$. This assumption means that there are no useless output symbols.

The set of input probability distributions on the input alphabet $\mathcal{X}$ is denoted by $\Delta(\mathcal{X}) \equiv\left\{\boldsymbol{\lambda}=\left(\lambda_{1}, \ldots, \lambda_{m}\right) \mid \lambda_{i} \geq\right.$ $\left.0, i=1, \ldots, m, \sum_{i=1}^{m} \lambda_{i}=1\right\}$. The interior of $\Delta(\mathcal{X})$ is denoted by $\Delta(\mathcal{X})^{\circ} \equiv\left\{\boldsymbol{\lambda}=\left(\lambda_{1}, \ldots, \lambda_{m}\right) \in \Delta(\mathcal{X}) \mid \lambda_{i}>\right.$ $0, i=1, \ldots, m\}$. Similarly, the set of output probability distributions on the output alphabet $\mathcal{Y}$ is denoted by $\Delta(\mathcal{Y}) \equiv$ $\left\{Q=\left(Q_{1}, \ldots, Q_{n}\right) \mid Q_{j} \geq 0, j=1, \ldots, n, \sum_{j=1}^{n} Q_{j}=1\right\}$, and its interior $\Delta(\mathcal{Y})^{\circ}$ is similarly defined as $\Delta(\mathcal{X})^{\circ}$.

Let $Q=\lambda \Phi$ be the output distribution for an input distribution $\boldsymbol{\lambda} \in \Delta(\mathcal{X})$ and write its components as $Q_{j}=$ $\sum_{i=1}^{m} \lambda_{i} P_{j}^{i}, j=1, \ldots, n$. Then, the mutual information is defined by $I(\boldsymbol{\lambda}, \Phi)=\sum_{i=1}^{m} \sum_{j=1}^{n} \lambda_{i} P_{j}^{i} \log \left(P_{j}^{i} / Q_{j}\right)$, where $\log$ is the natural logarithm. The channel capacity $C$ is defined by

$$
\begin{equation*}
C=\max _{\boldsymbol{\lambda} \in \Delta(\mathcal{X})} I(\boldsymbol{\lambda}, \Phi) \tag{2}
\end{equation*}
$$

The unit of $C$ is nat/symbol.
The Kullback-Leibler divergence $D\left(Q \| Q^{\prime}\right)$ for two output distributions $Q=\left(Q_{1}, \ldots, Q_{n}\right), Q^{\prime}=\left(Q_{1}^{\prime}, \ldots, Q_{n}^{\prime}\right) \in \Delta(\mathcal{Y})$ is defined [7] by

$$
\begin{equation*}
D\left(Q \| Q^{\prime}\right)=\sum_{j=1}^{n} Q_{j} \log \frac{Q_{j}}{Q_{j}^{\prime}} \tag{3}
\end{equation*}
$$

An important proposition for investigating the convergence speed of the Arimoto-Blahut algorithm is the Kuhn-Tucker condition on the input distribution $\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}$ that achieves the maximum of (2).

Theorem: (Kuhn-Tucker condition [6]) In the maximization problem (2), a necessary and sufficient condition for the input
distribution $\boldsymbol{\lambda}^{*}=\left(\lambda_{1}^{*}, \ldots, \lambda_{m}^{*}\right) \in \Delta(\mathcal{X})$ to achieve the maximum is that there is a certain constant $\tilde{C}$ with

$$
D\left(P^{i} \| \boldsymbol{\lambda}^{*} \Phi\right) \begin{cases}=\tilde{C}, & \text { for } i \text { with } \lambda_{i}^{*}>0  \tag{4}\\ \leq \tilde{C}, & \text { for } i \text { with } \lambda_{i}^{*}=0\end{cases}
$$

In (4), $\tilde{C}$ is equal to the channel capacity $C$.
Since the Kuhn-Tucker condition is a necessary and sufficient condition, all the information about the capacityachieving input distribution $\boldsymbol{\lambda}^{*}$ can be derived from this condition.

## B. Definition of the Arimoto-Blahut algorithm

A sequence $\left\{\boldsymbol{\lambda}^{N}=\left(\lambda_{1}^{N}, \ldots, \lambda_{m}^{N}\right)\right\}_{N=0,1, \ldots} \subset \Delta(\mathcal{X})$ of input distributions is defined by the Arimoto-Blahut algorithm as follows [2], [4]. First, let $\boldsymbol{\lambda}^{0}=\left(\lambda_{1}^{0}, \ldots, \lambda_{m}^{0}\right)$ be an initial distribution taken in $\Delta(\mathcal{X})^{\circ}$, i.e., $\lambda_{i}^{0}>0, i=1, \ldots, m$. Then, the Arimoto-Blahut algorithm is given by the recurrence formula

$$
\begin{align*}
\lambda_{i}^{N+1} & =\frac{\lambda_{i}^{N} \exp D\left(P^{i} \| \boldsymbol{\lambda}^{N} \Phi\right)}{\sum_{k=1}^{m} \lambda_{k}^{N} \exp D\left(P^{k} \| \boldsymbol{\lambda}^{N} \Phi\right)}  \tag{5}\\
i & =1, \ldots, m, N=0,1, \ldots
\end{align*}
$$

On the convergence of this Arimoto-Blahut algorithm, the following results were obtained by Arimoto in [2], [3].

In [2], [3], defining

$$
\begin{align*}
& C(N+1, N) \equiv-\sum_{i=1}^{m} \lambda_{i}^{N+1} \log \lambda_{i}^{N+1} \\
&+\sum_{i=1}^{m} \sum_{j=1}^{n} \lambda_{i}^{N+1} P_{j}^{i} \log \frac{\lambda_{i}^{N} P_{j}^{i}}{\sum_{k=1}^{m} \lambda_{k}^{N} P_{j}^{k}}  \tag{6}\\
&=-D\left(\boldsymbol{\lambda}^{N+1} \| \boldsymbol{\lambda}^{N}\right)+\sum_{i=1}^{m} \lambda_{i}^{N+1} D\left(P^{i} \| Q^{N}\right) \tag{7}
\end{align*}
$$

Arimoto obtained the following results.
Theorem A1: [2] If the initial input distribution $\boldsymbol{\lambda}^{0}$ is in $\Delta(\mathcal{X})^{\circ}$, then

$$
\begin{equation*}
\lim _{N \rightarrow \infty} C(N+1, N)=C \tag{8}
\end{equation*}
$$

Theorem A2: [2] If $\boldsymbol{\lambda}^{0}$ is the uniform distribution, then

$$
\begin{equation*}
0 \leq C-C(N+1, N) \leq \frac{\log m-h\left(\boldsymbol{\lambda}^{*}\right)}{N} \tag{9}
\end{equation*}
$$

where $\boldsymbol{\lambda}^{*}$ is the capacity-achieving input distribution, and $h\left(\boldsymbol{\lambda}^{*}\right)$ is the entropy of $\boldsymbol{\lambda}^{*}$.

Theorem A3: [3] Assume that $\boldsymbol{\lambda}^{*}$ is unique and belongs to $\Delta(\mathcal{X})^{\circ}$. Then,
i) $\boldsymbol{\lambda}^{N} \rightarrow \boldsymbol{\lambda}^{*}, N \rightarrow \infty$, and
ii) for sufficiently small arbitrary $\epsilon>0$, there exists $N_{0}=$ $N_{0}(\epsilon)$ such that

$$
\begin{equation*}
0 \leq C-C(N+1, N) \leq \epsilon \cdot(\theta)^{N-N_{0}}, N \geq N_{0} \tag{10}
\end{equation*}
$$

where $\theta$ is a constant with $0 \leq \theta<1$ and is unrelated to $\epsilon$ and $N_{0}$. Further, in the expression $\epsilon \cdot(\theta)^{N-N_{0}}$, denotes multiplication, and $(\theta)^{N-N_{0}}$ means $\theta$ to the power of $N-N_{0}$.
We can see from Theorem A3, under the assumptions of Theorem A3, that the first term of (7) converges to 0 , and the second term converges to $C$.

## C. Function from $\Delta(\mathcal{X})$ to $\Delta(\mathcal{X})$

Let $F_{i}(\boldsymbol{\lambda})$ be the defining function of the Arimoto-Blahut algorithm (5), i.e.,

$$
\begin{equation*}
F_{i}(\boldsymbol{\lambda})=\frac{\lambda_{i} \exp D\left(P^{i} \| \boldsymbol{\lambda} \Phi\right)}{\sum_{k=1}^{m} \lambda_{k} \exp D\left(P^{k} \| \boldsymbol{\lambda} \Phi\right)}, i=1, \ldots, m \tag{11}
\end{equation*}
$$

Define $F(\boldsymbol{\lambda}) \equiv\left(F_{1}(\boldsymbol{\lambda}), \ldots, F_{m}(\boldsymbol{\lambda})\right)$. Then, $F(\boldsymbol{\lambda})$ is a differentiable function from $\Delta(\mathcal{X})$ to $\Delta(\mathcal{X})$, and (5) is represented by $\boldsymbol{\lambda}^{N+1}=F\left(\boldsymbol{\lambda}^{N}\right)$.

In this paper, for the analysis of the convergence speed, we assume

$$
\begin{equation*}
\operatorname{rank} \Phi=m \tag{12}
\end{equation*}
$$

Concerning this assumption, we see that in [2], [3], for the analysis of the convergence speed, the uniqueness of the capacity-achieving $\boldsymbol{\lambda}^{*}$ is assumed, which is a necessary condition for (12). In fact, we have

Lemma 1: The capacity-achieving input distribution $\boldsymbol{\lambda}^{*}$ is unique under the assumption (12).

Proof: By Csiszár[7], p.137, Eq. (37), for arbitrary $Q \in$ $\Delta(\mathcal{Y})$,

$$
\begin{equation*}
\sum_{i=1}^{m} \lambda_{i} D\left(P^{i} \| Q\right)=I(\boldsymbol{\lambda}, \Phi)+D(\boldsymbol{\lambda} \Phi \| Q) \tag{13}
\end{equation*}
$$

By the assumption (12), we see that there exists $Q^{0} \in \Delta(\mathcal{Y})$ [15] with

$$
\begin{equation*}
D\left(P^{1} \| Q^{0}\right)=\ldots=D\left(P^{m} \| Q^{0}\right) \equiv C^{0} \tag{14}
\end{equation*}
$$

Substituting $Q=Q^{0}$ into (13), we have $C^{0}=I(\boldsymbol{\lambda}, \Phi)+$ $D\left(\boldsymbol{\lambda} \Phi \| Q^{0}\right)$. Because $C^{0}$ is a constant,

$$
\begin{equation*}
\max _{\boldsymbol{\lambda} \in \Delta(\mathcal{X})} I(\boldsymbol{\lambda}, \Phi) \Longleftrightarrow \min _{\boldsymbol{\lambda} \in \Delta(\mathcal{X})} D\left(\boldsymbol{\lambda} \Phi \| Q^{0}\right) \tag{15}
\end{equation*}
$$

Define $W \equiv\{\boldsymbol{\lambda} \Phi \mid \boldsymbol{\lambda} \in \Delta(\mathcal{X})\}$. Then, $W$ is a closed convex set, and thus by Cover [6], p.297, Theorem 12.6.1, $Q=Q^{*}$, which achieves $\min _{Q \in W} D\left(Q \| Q^{0}\right)$, exists and is unique. By the assumption (12), the mapping $\Delta \ni \boldsymbol{\lambda} \mapsto \boldsymbol{\lambda} \Phi \in W$ is one to one; therefore, $\boldsymbol{\lambda}^{*}$ with $Q^{*}=\boldsymbol{\lambda}^{*} \Phi$ is unique.

Remark 1: Because of the equivalence (15), the ArimotoBlahut algorithm can be obtained by Csiszár [8], Chapter 4, "Minimizing information distance from a single measure", Theorem 5.

Lemma 2: The capacity-achieving input distribution $\boldsymbol{\lambda}^{*}$ is the fixed point of the function $F(\boldsymbol{\lambda})$. That is, $F\left(\boldsymbol{\lambda}^{*}\right)=\boldsymbol{\lambda}^{*}$.

Proof: In the Kuhn-Tucker condition (4), let us define $m_{1}$ as the number of indices $i$ with $\lambda_{i}^{*}>0$, i.e.,

$$
\lambda_{i}^{*}\left\{\begin{align*}
>0, & i=1, \ldots, m_{1},  \tag{16}\\
=0, & i=m_{1}+1, \ldots, m,
\end{align*}\right.
$$

by reordering the input symbols (if necessary), then

$$
D\left(P^{i} \| \boldsymbol{\lambda}^{*} \Phi\right) \begin{cases}=C, & i=1, \ldots, m_{1}  \tag{17}\\ \leq C, & i=m_{1}+1, \ldots, m\end{cases}
$$

We have

$$
\begin{equation*}
\sum_{k=1}^{m} \lambda_{k}^{*} \exp D\left(P^{k} \| \boldsymbol{\lambda}^{*} \Phi\right)=\sum_{k=1}^{m_{1}} \lambda_{k}^{*} e^{C}=e^{C} \tag{18}
\end{equation*}
$$

hence, by (11), (16), and (18),

$$
\begin{align*}
F_{i}\left(\boldsymbol{\lambda}^{*}\right) & = \begin{cases}e^{-C} \lambda_{i}^{*} e^{C}, & i=1, \ldots, m_{1} \\
0, & i=m_{1}+1, \ldots, m\end{cases}  \tag{19}\\
& =\lambda_{i}^{*}, i=1, \ldots, m \tag{20}
\end{align*}
$$

which shows $F\left(\boldsymbol{\lambda}^{*}\right)=\boldsymbol{\lambda}^{*}$.
The sequence $\left\{\boldsymbol{\lambda}^{N}\right\}_{N=0,1, \ldots}$ of the Arimoto-Blahut algorithm converges to the fixed point $\boldsymbol{\lambda}^{*}$, i.e., $\boldsymbol{\lambda}^{N} \rightarrow \boldsymbol{\lambda}^{*}, N \rightarrow$ $\infty$. We will investigate the convergence speed by using the Taylor expansion of $F(\boldsymbol{\lambda})$ about $\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}$.

## D. Convergence speed of $\boldsymbol{\lambda}^{N} \rightarrow \boldsymbol{\lambda}^{*}$

Now, we define two kinds of convergence speeds for investigating $\boldsymbol{\lambda}^{N} \rightarrow \boldsymbol{\lambda}^{*}$.
(i) Exponential convergence
$\boldsymbol{\lambda}^{N} \rightarrow \boldsymbol{\lambda}^{*}$ is the exponential convergence if

$$
\begin{align*}
& \left\|\boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}\right\|<K \cdot(\theta)^{N}, N=0,1, \ldots  \tag{21}\\
& \quad K>0,0 \leq \theta<1
\end{align*}
$$

where $\|\boldsymbol{\lambda}\|$ denotes the Euclidean norm, i.e., $\|\boldsymbol{\lambda}\|=$ $\left(\lambda_{1}^{2}+\ldots+\lambda_{m}^{2}\right)^{1 / 2}$
(ii) Convergence of order $O(1 / N)$
$\boldsymbol{\lambda}^{N} \rightarrow \boldsymbol{\lambda}^{*}$ is the convergence of order $O(1 / N)$ if

$$
\begin{equation*}
\lim _{N \rightarrow \infty} N\left(\lambda_{i}^{N}-\lambda_{i}^{*}\right)=K_{i} \neq 0, i=1, \ldots, m \tag{22}
\end{equation*}
$$

## E. Type of index

Now, we classify the indices $i(i=1, \ldots, m)$ in the KuhnTucker condition (4) in more detail into the following 3 types.

$$
D\left(P^{i} \| \boldsymbol{\lambda}^{*} \Phi\right) \begin{cases}=C, & \text { for } i \text { with } \lambda_{i}^{*}>0 \text { [type-I] }  \tag{23}\\ =C, & \text { for } i \text { with } \lambda_{i}^{*}=0 \text { [type-II] } \\ <C, & \text { for } i \text { with } \lambda_{i}^{*}=0 \text { [type-III]. }\end{cases}
$$

Let us define the sets of indices as follows.

$$
\begin{align*}
& \text { all the indices: } \mathcal{I} \equiv\{1, \ldots, m\}  \tag{24}\\
& \text { type-I indices: } \mathcal{I}_{\mathrm{I}} \equiv\left\{1, \ldots, m_{1}\right\}  \tag{25}\\
& \text { type-II indices: } \mathcal{I}_{\mathrm{II}} \equiv\left\{m_{1}+1, \ldots, m_{1}+m_{2}\right\}  \tag{26}\\
& \text { type-III indices: } \mathcal{I}_{\mathrm{III}} \equiv\left\{m_{1}+m_{2}+1, \ldots, m\right\} \tag{27}
\end{align*}
$$

We have $|\mathcal{I}|=m,\left|\mathcal{I}_{\mathrm{I}}\right|=m_{1},\left|\mathcal{I}_{\mathrm{II}}\right|=m_{2},\left|\mathcal{I}_{\mathrm{III}}\right|=m-m_{1}-$ $m_{2} \equiv m_{3}, \mathcal{I}=\mathcal{I}_{\mathrm{I}} \cup \mathcal{I}_{\mathrm{II}} \cup \mathcal{I}_{\mathrm{III}}$ and $m=m_{1}+m_{2}+m_{3} . \mathcal{I}_{\mathrm{I}}$ is not empty, and $\left|\mathcal{I}_{\mathrm{I}}\right|=m_{1} \geq 2$ for any channel matrix, but $\mathcal{I}_{\text {II }}$ and $\mathcal{I}_{\text {III }}$ may be empty for some channel matrices.


Fig. 1. Positional relation of row vectors $P^{1}, P^{2}, P^{3}$ of $\Phi^{(1)}$ and $Q^{*}$ in Example 1.

## F. Examples of convergence speed

Let us consider the dependence of the convergence speed of the Arimoto-Blahut algorithm on the channel matrices.

For many channel matrices $\Phi$, the convergence is exponential, but for some special $\Phi$, the convergence is very slow. Let us consider the following three examples with input alphabet size $m=3$ and output alphabet size $n=3$, taking types-I, -II and -III into account.

Example 1: (only type-I) If only type-I indices exist, then $\lambda_{i}^{*}>0, i=1,2,3$; hence, $Q^{*} \equiv \lambda^{*} \Phi$ is in the interior of $\triangle P^{1} P^{2} P^{3}$. See Fig. 1. As a concrete channel matrix of this example, let us consider

$$
\Phi^{(1)} \equiv\left(\begin{array}{ccc}
0.800 & 0.100 & 0.100  \tag{28}\\
0.100 & 0.800 & 0.100 \\
0.250 & 0.250 & 0.500
\end{array}\right)
$$

For this $\Phi^{(1)}$, we have $C=0.323$ [nat/symbol], $\lambda^{*}=$ $(0.431,0.431,0.138)$ and $Q^{*}=(0.422,0.422,0.156)$. The vertices of the large triangle in Fig. 1 are $\boldsymbol{e}_{1}=(1,0,0), \boldsymbol{e}_{2}=$ $(0,1,0), \boldsymbol{e}_{3}=(0,0,1)$. We have $D\left(P^{i} \| Q^{*}\right)=C, i=$ $1,2,3$, then considering the analogy to Euclidean geometry, $\triangle P^{1} P^{2} P^{3}$ can be regarded as an "acute triangle."

Example 2: (types-I and -II) If there are type-I and type-II indices, we can assume $\lambda_{1}^{*}>0, \lambda_{2}^{*}>0, \lambda_{3}^{*}=0$ without loss of generality; hence, $Q^{*}$ is on the side $P^{1} P^{2}$ and $D\left(P^{i} \| Q^{*}\right)=$ $C, i=1,2,3$. See Fig. 2. As a concrete channel matrix of this example, let us consider

$$
\Phi^{(2)} \equiv\left(\begin{array}{ccc}
0.800 & 0.100 & 0.100  \tag{29}\\
0.100 & 0.800 & 0.100 \\
0.300 & 0.300 & 0.400
\end{array}\right)
$$

For this $\Phi^{(2)}$, we have $C=0.310$ [nat/symbol], $\boldsymbol{\lambda}^{*}=$ $(0.500,0.500,0.000)$ and $Q^{*}=(0.450,0.450,0.100)$. Considering the analogy to Euclidean geometry, $\triangle P^{1} P^{2} P^{3}$ can be regarded as a "right triangle."

Example 3: (types-I and -III) If there are type-I and typeIII indices, we can assume $\lambda_{1}^{*}>0, \lambda_{2}^{*}>0, \lambda_{3}^{*}=0$ without


Fig. 2. Positional relation of row vectors $P^{1}, P^{2}, P^{3}$ of $\Phi^{(2)}$ and $Q^{*}$ in Example 2.


Fig. 3. Positional relation of row vectors $P^{1}, P^{2}, P^{3}$ of $\Phi^{(3)}$ and $Q^{*}$ in Example 3.
loss of generality; hence, $Q^{*}$ is on the side $P^{1} P^{2}$ and $C=$ $D\left(P^{1} \| Q^{*}\right)=D\left(P^{2} \| Q^{*}\right)>D\left(P^{3} \| Q^{*}\right)$. See Fig. 3. As a concrete channel matrix of this example, let us consider

$$
\Phi^{(3)} \equiv\left(\begin{array}{ccc}
0.800 & 0.100 & 0.100  \tag{30}\\
0.100 & 0.800 & 0.100 \\
0.350 & 0.350 & 0.300
\end{array}\right)
$$

For this $\Phi^{(3)}$, we have $C=0.310\left[n a t /\right.$ symbol], $\boldsymbol{\lambda}^{*}=$ $(0.500,0.500,0.000)$ and $Q^{*}=(0.450,0.450,0.100)$. Considering the analogy to Euclidean geometry, $\triangle P^{1} P^{2} P^{3}$ can be regarded as an "obtuse triangle."

For the above $\Phi^{(1)}, \Phi^{(2)}, \Phi^{(3)}$, we show in Fig. 4 the speed of convergence of $\left|\lambda_{1}^{N}-\lambda_{1}^{*}\right| \rightarrow 0$. By Fig. 4, we see that in Examples 1 and 3, the convergence is exponential, while in Example 2, the convergence is slower than exponential.
From the above three examples, it is inferred that the Arimoto-Blahut algorithm converges exponentially when type-


Fig. 4. Comparison of the convergence speed in Examples 1,2,3.

II indices do not exist and slowly when type-II indices exist. We will analyze this phenomenon in this paper.

Remark 2: Related to the above investigation of triangles, we have given a geometric investigation of the channel capacity in [15] based on divergence geometry.

## G. Analysis of convergence speed

We will show in this paper for which channel the convergence speed of the Arimoto-Blahut algorithm is exponential and for which channel it is of order $O(1 / N)$.

In previous studies [2], [21], the exponential convergence has been investigated. In [2], [21], the authors proved that the convergence speed is exponential if the capacity-achieving input distribution $\boldsymbol{\lambda}^{*}=\left(\lambda_{1}^{*}, \ldots, \lambda_{m}^{*}\right)$ satisfies $\lambda_{i}^{*}>0, i=$ $1, \ldots, m$. In this case, all the indices are of type-I and there are no type-II indices. When there are no type-II indices, the convergence speed is determined by the Jacobian matrix of the defining function $F(\boldsymbol{\lambda})$ of the Arimoto-Blahut algorithm.

In this paper, we will investigate the case in which type-II indices exist. If there exist type-II indices, the convergence speed of the Arimoto-Blahut algorithm is determined not only by the Jacobian matrix, but also by the Hessian matrix. We will analyze the nonlinear recurrence formula defined by a quadratic polynomial, which is obtained from the Hessian matrix. Then, we will show, under some assumptions, that the convergence speed of the recurrence formula is of order $O(1 / N)$.

## III. TAYLOR EXPANSION OF $F(\boldsymbol{\lambda})$ about $\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}$

We will examine the convergence speed of the ArimotoBlahut algorithm by the Taylor expansion of $F(\boldsymbol{\lambda})$ about the
fixed point $\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}$. The Taylor expansion of the function $F(\boldsymbol{\lambda})=\left(F_{1}(\boldsymbol{\lambda}), \ldots, F_{m}(\boldsymbol{\lambda})\right)$ about $\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}$ is

$$
\begin{align*}
& F(\boldsymbol{\lambda})=F\left(\boldsymbol{\lambda}^{*}\right)+\left(\boldsymbol{\lambda}-\boldsymbol{\lambda}^{*}\right) J\left(\boldsymbol{\lambda}^{*}\right) \\
& \quad+\frac{1}{2!}\left(\boldsymbol{\lambda}-\boldsymbol{\lambda}^{*}\right) H\left(\boldsymbol{\lambda}^{*}\right)^{t}\left(\boldsymbol{\lambda}-\boldsymbol{\lambda}^{*}\right)+o\left(\left\|\boldsymbol{\lambda}-\boldsymbol{\lambda}^{*}\right\|^{2}\right) \tag{31}
\end{align*}
$$

where ${ }^{t} \boldsymbol{\lambda}$ denotes the transpose of $\boldsymbol{\lambda}$, and $o\left(\|\boldsymbol{\lambda}\|^{2}\right)$ means $\lim _{\|\boldsymbol{\lambda}\| \rightarrow 0} o\left(\|\boldsymbol{\lambda}\|^{2}\right) /\|\boldsymbol{\lambda}\|^{2}=0$.

In (31), $J\left(\boldsymbol{\lambda}^{*}\right)$ is the Jacobian matrix of $F(\boldsymbol{\lambda})$ at $\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}$, i.e.,

$$
\begin{equation*}
J\left(\boldsymbol{\lambda}^{*}\right)=\left(\left.\frac{\partial F_{i}}{\partial \lambda_{i^{\prime}}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}}\right)_{i^{\prime}, i=1, \ldots, m} \tag{32}
\end{equation*}
$$

In this paper, we assume that the input probability distribution $\boldsymbol{\lambda}$ is a row vector; thus, the Jacobian matrix $J\left(\boldsymbol{\lambda}^{*}\right)$ is

$$
J\left(\boldsymbol{\lambda}^{*}\right)=\begin{gather*}
\uparrow\left(\begin{array}{ccc}
\left.\frac{\partial F_{1}}{\partial \lambda_{1}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}} & \cdots & \left.\frac{\partial F_{m}}{\partial \lambda_{1}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}} \\
\vdots & & \vdots \\
\left.\frac{\partial F_{1}}{}\right|_{i^{\prime}} & & \left.\frac{\partial F_{m}}{}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}} \\
\cdots & \left.\frac{\partial \lambda_{m}}{}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}}
\end{array}\right) \in \mathbb{R}^{m \times m}, .
\end{gather*}
$$

i.e., $\partial F_{i} /\left.\partial \lambda_{i^{\prime}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}}$ is the $\left(i^{\prime}, i\right)$ component.

Lemma 3: Every row sum of $J\left(\boldsymbol{\lambda}^{*}\right)$ is equal to 0 .
Proof: By (11), we have $\sum_{i=1}^{m} F_{i}(\boldsymbol{\lambda})=1$; thus, by (32), the lemma holds.

In $(31), H\left(\boldsymbol{\lambda}^{*}\right) \equiv\left(H_{1}\left(\boldsymbol{\lambda}^{*}\right), \ldots, H_{m}\left(\boldsymbol{\lambda}^{*}\right)\right)$, where $H_{i}\left(\boldsymbol{\lambda}^{*}\right)$ is the Hessian matrix of $F_{i}(\boldsymbol{\lambda})$ at $\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}$, i.e.,

$$
\begin{equation*}
H_{i}\left(\boldsymbol{\lambda}^{*}\right)=\left(\left.\frac{\partial^{2} F_{i}}{\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}}\right)_{i^{\prime}, i^{\prime \prime}=1, \ldots, m} \tag{34}
\end{equation*}
$$

and $\left(\boldsymbol{\lambda}-\boldsymbol{\lambda}^{*}\right) H\left(\boldsymbol{\lambda}^{*}\right)^{t}\left(\boldsymbol{\lambda}-\boldsymbol{\lambda}^{*}\right)$ is an abbreviated expression of the $m$ dimensional row vector
$\left(\left(\boldsymbol{\lambda}-\boldsymbol{\lambda}^{*}\right) H_{1}\left(\boldsymbol{\lambda}^{*}\right)^{t}\left(\boldsymbol{\lambda}-\boldsymbol{\lambda}^{*}\right), \ldots,\left(\boldsymbol{\lambda}-\boldsymbol{\lambda}^{*}\right) H_{m}\left(\boldsymbol{\lambda}^{*}\right)^{t}\left(\boldsymbol{\lambda}-\boldsymbol{\lambda}^{*}\right)\right)$.
Remark 3: $\lambda_{1}, \ldots, \lambda_{m}$ satisfy the constraint $\sum_{i=1}^{m} \lambda_{i}=1$, but in (31), (32), (34) we consider $\lambda_{1}, \ldots, \lambda_{m}$ as independent (or constraint free) variables to have the Taylor series approximation (31). This approximation is justified as follows. By the Kuhn-Tucker condition (4), $D\left(P^{i} \| Q^{*}\right) \leq C<$ $\infty, i=1, \ldots, m$; hence, by the assumption put below (1), we have $Q_{j}^{*}>0, j=1, \ldots, n$. See [2]. For $\epsilon>0$, define $\mathcal{Q}_{\epsilon}^{*} \equiv\left\{Q=\left(Q_{1}, \ldots, Q_{n}\right) \in \mathbb{R}^{n} \mid\left\|Q-Q^{*}\right\|<\epsilon\right\}$, i.e., $\mathcal{Q}_{\epsilon}^{*}$ is an open ball in $\mathbb{R}^{n}$ centered at $Q^{*}$ with radius $\epsilon$. Note that $Q \in \mathcal{Q}_{\epsilon}^{*}$ is free from the constraint $\sum_{j=1}^{n} Q_{j}=1$. Taking $\epsilon>0$ sufficiently small, we can have $Q_{j}>0, j=1, \ldots, n$, for any $Q \in \mathcal{Q}_{\epsilon}^{*}$. The function $F(\boldsymbol{\lambda})$ is defined for $\boldsymbol{\lambda}$ with $Q_{j}=(\boldsymbol{\lambda} \Phi)_{j}>0, j=1, \ldots, n$, even if some $\lambda_{i}<0$. Therefore, the domain of definition of $F(\boldsymbol{\lambda})$ can be extended to $\Phi^{-1}\left(\mathcal{Q}_{\epsilon}^{*}\right) \subset \mathbb{R}^{m}$, where $\Phi^{-1}\left(\mathcal{Q}_{\epsilon}^{*}\right)$ is the inverse image of $\mathcal{Q}_{\epsilon}^{*}$ by the mapping $\mathbb{R}^{m} \ni \boldsymbol{\lambda} \mapsto \boldsymbol{\lambda} \Phi \in \mathbb{R}^{n} . \Phi^{-1}\left(\mathcal{Q}_{\epsilon}^{*}\right)$ is an open neighborhood of $\boldsymbol{\lambda}^{*}$ in $\mathbb{R}^{m}$. Then, $F(\boldsymbol{\lambda})$ is a function of $\boldsymbol{\lambda}=\left(\lambda_{1}, \ldots, \lambda_{m}\right) \in \Phi^{-1}\left(\mathcal{Q}_{\epsilon}^{*}\right)$ as independent variables (free from the constraint $\sum_{i=1}^{m} \lambda_{i}=1$ ). We can consider (31) to be the Taylor expansion by independent variables $\lambda_{1}, \ldots, \lambda_{m}$;
then, substituting $\boldsymbol{\lambda} \in \Delta(\mathcal{X}) \cap \Phi^{-1}\left(\mathcal{Q}_{\epsilon}^{*}\right)$ into (31), we obtain the approximation for $F(\boldsymbol{\lambda})$ about $\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}$.

Now, substituting $\boldsymbol{\lambda}=\boldsymbol{\lambda}^{N}$ into (31) and using $F\left(\boldsymbol{\lambda}^{*}\right)=\boldsymbol{\lambda}^{*}$ and $F\left(\boldsymbol{\lambda}^{N}\right)=\boldsymbol{\lambda}^{N+1}$, we have

$$
\begin{align*}
& \boldsymbol{\lambda}^{N+1}=\boldsymbol{\lambda}^{*}+\left(\boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}\right) J\left(\boldsymbol{\lambda}^{*}\right) \\
& \quad+\frac{1}{2!}\left(\boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}\right) H\left(\boldsymbol{\lambda}^{*}\right)^{t}\left(\boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}\right)+o\left(\left\|\boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}\right\|^{2}\right) \tag{35}
\end{align*}
$$

Then, by substituting $\boldsymbol{\mu}^{N} \equiv \boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}$, (35) becomes

$$
\begin{equation*}
\boldsymbol{\mu}^{N+1}=\boldsymbol{\mu}^{N} J\left(\boldsymbol{\lambda}^{*}\right)+\frac{1}{2!} \boldsymbol{\mu}^{N} H\left(\boldsymbol{\lambda}^{*}\right)^{t} \boldsymbol{\mu}^{N}+o\left(\left\|\boldsymbol{\mu}^{N}\right\|^{2}\right) . \tag{36}
\end{equation*}
$$

Then, we will investigate the convergence $\boldsymbol{\mu}^{N} \rightarrow \mathbf{0}, N \rightarrow \infty$, based on the Taylor expansion (36). Let $\mu_{i}^{N} \equiv \lambda_{i}^{N}-\lambda_{i}^{*}, i=$ $1, \ldots, m$, denote the components of $\boldsymbol{\mu}^{N}=\boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}$; then, we can write $\boldsymbol{\mu}^{N}$ by components as $\boldsymbol{\mu}^{N}=\left(\mu_{1}^{N}, \ldots, \mu_{m}^{N}\right)$. We have $\sum_{i=1}^{m} \mu_{i}^{N}=0, N=0,1, \ldots$, because $\sum_{i=1}^{m} \lambda_{i}^{N}=$ $\sum_{i=1}^{m} \lambda_{i}^{*}=1$.

## A. The Jacobian matrix $J\left(\boldsymbol{\lambda}^{*}\right)$

Let us consider the Jacobian matrix $J\left(\boldsymbol{\lambda}^{*}\right)$. We are assuming $\operatorname{rank} \Phi=m$ in (12), hence $m \leq n$. We will calculate the components (32) of $J\left(\boldsymbol{\lambda}^{*}\right)$.

Defining $D_{i} \equiv D\left(P^{i} \| \boldsymbol{\lambda} \Phi\right)$ and $F_{i} \equiv F_{i}(\boldsymbol{\lambda}), i=1, \ldots, m$, we can write (11) as

$$
\begin{equation*}
F_{i}=\frac{\lambda_{i} e^{D_{i}}}{\sum_{k=1}^{m} \lambda_{k} e^{D_{k}}}, i=1, \ldots, m \tag{37}
\end{equation*}
$$

From (37), it follows that

$$
\begin{equation*}
F_{i} \sum_{k=1}^{m} \lambda_{k} e^{D_{k}}=\lambda_{i} e^{D_{i}} \tag{38}
\end{equation*}
$$

Then, differentiating both sides of (38) with respect to $\lambda_{i^{\prime}}$, we have

$$
\begin{equation*}
\frac{\partial F_{i}}{\partial \lambda_{i^{\prime}}} \sum_{k=1}^{m} \lambda_{k} e^{D_{k}}+F_{i} \frac{\partial}{\partial \lambda_{i^{\prime}}} \sum_{k=1}^{m} \lambda_{k} e^{D_{k}}=\delta_{i^{\prime} i} e^{D_{i}}+\lambda_{i} e^{D_{i}} \frac{\partial D_{i}}{\partial \lambda_{i^{\prime}}}, \tag{39}
\end{equation*}
$$

where $\delta_{i^{\prime} i}$ is the Kronecker delta.
Before substituting $\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}=\left(\lambda_{1}^{*}, \ldots, \lambda_{m}^{*}\right)$ into both sides of (39), we define the following symbols. Remember that the integer $m_{1}$ was defined in (16). See also (25).

Let us define

$$
\begin{align*}
Q^{*} & \equiv Q\left(\boldsymbol{\lambda}^{*}\right)=\boldsymbol{\lambda}^{*} \Phi  \tag{40}\\
Q_{j}^{*} & \equiv Q\left(\boldsymbol{\lambda}^{*}\right)_{j}=\sum_{i=1}^{m} \lambda_{i}^{*} P_{j}^{i}=\sum_{i=1}^{m_{1}} \lambda_{i}^{*} P_{j}^{i}, j=1, \ldots, n  \tag{41}\\
D_{i}^{*} & \equiv D\left(P^{i} \| Q^{*}\right), i=1, \ldots, m  \tag{42}\\
D_{i^{\prime}, i}^{*} & \left.\equiv \frac{\partial D_{i}}{\partial \lambda_{i^{\prime}}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}}, i^{\prime}, i=1, \ldots, m  \tag{43}\\
F_{i}^{*} & \equiv F_{i}\left(\boldsymbol{\lambda}^{*}\right), i=1, \ldots, m \tag{44}
\end{align*}
$$

Lemma 4: We have

$$
\begin{align*}
& \left.\sum_{k=1}^{m} \lambda_{k} e^{D_{k}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}}=e^{C}  \tag{45}\\
& \frac{\partial D_{i}}{\partial \lambda_{i^{\prime}}}=-\sum_{j=1}^{n} \frac{P_{j}^{i^{\prime}} P_{j}^{i}}{Q_{j}}, i^{\prime}, i=1, \ldots, m  \tag{46}\\
& \left.\frac{\partial}{\partial \lambda_{i^{\prime}}} \sum_{k=1}^{m} \lambda_{k} e^{D_{k}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}}=e^{D_{i^{\prime}}^{*}}-e^{C}, i^{\prime}=1, \ldots, m  \tag{47}\\
& F_{i}^{*}=\lambda_{i}^{*}, i=1, \ldots, m \tag{48}
\end{align*}
$$

Proof: Eq. (45) was proved in (18). Eq. (46) is proved by a simple calculation. Eq. (47) is proved as follows.

$$
\begin{align*}
& \left.\frac{\partial}{\partial \lambda_{i^{\prime}}} \sum_{k=1}^{m} \lambda_{k} e^{D_{k}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}}=\left.\sum_{k=1}^{m}\left(\delta_{i^{\prime} k} e^{D_{k}}+\lambda_{k} e^{D_{k}} \frac{\partial D_{k}}{\partial \lambda_{i^{\prime}}}\right)\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}} \\
& =e^{D_{i^{\prime}}^{*}}+\sum_{k=1}^{m_{1}} \lambda_{k}^{*} e^{C}\left(-\sum_{j=1}^{n} \frac{P_{j}^{k} P_{j}^{i^{\prime}}}{Q_{j}^{*}}\right) \\
& =e^{D_{i^{\prime}}^{*}}-e^{C} \sum_{j=1}^{n} P_{j}^{i^{\prime}} \frac{1}{Q_{j}^{*}} \sum_{k=1}^{m_{1}} \lambda_{k}^{*} P_{j}^{k} \\
& =e^{D_{i^{\prime}}^{*}}-e^{C} \tag{49}
\end{align*}
$$

Note that $Q_{j}^{*}>0, j=1, \ldots, n$, from Remark 3. Eq. (48) is the result of Lemma 2.

Substituting the results of Lemma 4 into (39), we have

$$
\begin{equation*}
\left.\frac{\partial F_{i}}{\partial \lambda_{i^{\prime}}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}} e^{C}+\lambda_{i}^{*}\left(e^{D_{i^{\prime}}^{*}}-e^{C}\right)=\delta_{i^{\prime} i} e^{D_{i}^{*}}+\lambda_{i}^{*} e^{D_{i}^{*}} D_{i^{\prime}, i}^{*} \tag{50}
\end{equation*}
$$

Consequently, we have
Theorem 1: The components of the Jacobian matrix $J\left(\boldsymbol{\lambda}^{*}\right)$ are given as follows.

$$
\begin{align*}
& \left.\frac{\partial F_{i}}{\partial \lambda_{i^{\prime}}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}} \\
& =e^{D_{i}^{*}-C}\left(\delta_{i^{\prime} i}+\lambda_{i}^{*} D_{i^{\prime}, i}^{*}\right)+\lambda_{i}^{*}\left(1-e^{D_{i^{\prime}}^{*}-C}\right), i^{\prime}, i \in \mathcal{I},  \tag{51}\\
& =\left\{\begin{array}{l}
\delta_{i^{\prime} i}+\lambda_{i}^{*}\left(D_{i^{\prime}, i}^{*}+1-e^{D_{i^{\prime}}^{*}-C}\right), i^{\prime} \in \mathcal{I}, i \in \mathcal{I}_{\mathrm{I}}, \\
\delta_{i^{\prime} i}, i^{\prime} \in \mathcal{I}, i \in \mathcal{I}_{\mathrm{II}}, \\
e^{D_{i}^{*}-C} \delta_{i^{\prime} i}, i^{\prime} \in \mathcal{I}, i \in \mathcal{I}_{\mathrm{III}},
\end{array}\right. \tag{52}
\end{align*}
$$

where the sets of indices $\mathcal{I}, \mathcal{I}_{\mathrm{I}}, \mathcal{I}_{\mathrm{II}}, \mathcal{I}_{\text {III }}$ were defined in (24)(27). Note that $D_{i}^{*}=C$ for $i \in \mathcal{I}_{\mathrm{I}} \cup \mathcal{I}_{\mathrm{II}}$ and $\lambda_{i}^{*}=0$ for $i \in \mathcal{I}_{\text {II }} \cup \mathcal{I}_{\text {III }}$.

## B. Eigenvalues of the Jacobian matrix $J\left(\boldsymbol{\lambda}^{*}\right)$

From (52), we see that the Jacobian matrix $J\left(\boldsymbol{\lambda}^{*}\right)$ is of the form

$$
\begin{align*}
J\left(\boldsymbol{\lambda}^{*}\right) & \equiv\left(\begin{array}{ccc}
J^{\mathrm{I}} & O & O \\
* & J^{\mathrm{II}} & O \\
* & O & J^{\mathrm{III}}
\end{array}\right)  \tag{53}\\
J^{\mathrm{I}} & \equiv\left(\partial F_{i} /\left.\partial \lambda_{i^{\prime}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}}\right)_{i, i^{\prime} \in \mathcal{I}_{\mathrm{I}}} \in \mathbb{R}^{m_{1} \times m_{1}}  \tag{54}\\
J^{\mathrm{II}} & \equiv\left(\partial F_{i} /\left.\partial \lambda_{i^{\prime}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}}\right)_{i, i^{\prime} \in \mathcal{I}_{\mathrm{II}}} \\
& =I(\text { the identity matrix }) \in \mathbb{R}^{m_{2} \times m_{2}}  \tag{55}\\
J^{\mathrm{III}} & \equiv\left(\partial F_{i} /\left.\partial \lambda_{i^{\prime}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}}\right)_{i, i^{\prime} \in \mathcal{I}_{\mathrm{III}}} \\
& =\operatorname{diag}\left(e^{D_{i}^{*}-C}, i \in \mathcal{I}_{\mathrm{III}}\right) \in \mathbb{R}^{m_{3} \times m_{3}} \tag{56}
\end{align*}
$$

$O$ denotes the all-zero matrix of appropriate size, * denotes an appropriate matrix.

In (56), $\operatorname{diag}\left(e^{D_{i}^{*}-C}, i \in \mathcal{I}_{\text {III }}\right)$ denotes the diagonal matrix with diagonal components $e^{D_{i}^{*}-C}, i \in \mathcal{I}_{\text {IIII }}$. Here, $e^{D_{i}^{*}-C}<$ $1, i \in \mathcal{I}_{\text {III }}$ holds from type-III in (23).

Let $\left\{\theta_{1}, \ldots, \theta_{m}\right\} \equiv\left\{\theta_{i} \mid i \in \mathcal{I}\right\}$ be the set of eigenvalues of $J\left(\boldsymbol{\lambda}^{*}\right)$. By (53), the eigenvalues of $J\left(\boldsymbol{\lambda}^{*}\right)$ are the eigenvalues of $J^{\mathrm{I}}, J^{\mathrm{II}}$, and $J^{\mathrm{III}}$, hence we can put
$\left\{\theta_{i} \mid i \in \mathcal{I}_{\mathrm{I}}\right\}:$ the set of eigenvalues of $J^{\mathrm{I}}$,
$\left\{\theta_{i} \mid i \in \mathcal{I}_{\text {II }}\right\}$ : the set of eigenvalues of $J^{\mathrm{II}}$,
$\left\{\theta_{i} \mid i \in \mathcal{I}_{\text {III }}\right\}$ : the set of eigenvalues of $J^{\text {IIII }}$.
We will evaluate the eigenvalues of $J^{\mathrm{I}}, J^{\mathrm{II}}$, and $J^{\mathrm{III}}$ as follows.

## C. Eigenvalues of $J^{\mathrm{I}}$

First, we consider the eigenvalues of $J^{\mathrm{I}}$. Let $J_{i^{\prime} i}^{\mathrm{I}}$ be the $\left(i^{\prime}, i\right)$ component of $J^{\mathrm{I}}$, then by (52),

$$
\begin{equation*}
J_{i^{\prime} i}^{\mathrm{I}}=\delta_{i^{\prime} i}+\lambda_{i}^{*} D_{i^{\prime}, i}^{*}, \quad i^{\prime}, i \in \mathcal{I}_{\mathrm{I}} . \tag{57}
\end{equation*}
$$

Let $I \in \mathbb{R}^{m_{1} \times m_{1}}$ denote the identity matrix and define $B \equiv$ $I-J^{\mathrm{I}}$. Let $B_{i^{\prime} i}$ be the $\left(i^{\prime}, i\right)$ component of $B$. Then, from (57),

$$
\begin{align*}
B_{i^{\prime} i} & =-\lambda_{i}^{*} D_{i^{\prime}, i}^{*}  \tag{58}\\
& =\lambda_{i}^{*} \sum_{j=1}^{n} \frac{P_{j}^{i^{\prime}} P_{j}^{i}}{Q_{j}^{*}}, i^{\prime}, i \in \mathcal{I}_{\mathrm{I}} . \tag{59}
\end{align*}
$$

Let $\left\{\beta_{i} \mid i \in \mathcal{I}_{\mathrm{I}}\right\}$ be the set of eigenvalues of $B$, then we have $\theta_{i}=1-\beta_{i}, i \in \mathcal{I}_{\mathrm{I}}$. To calculate the eigenvalues of $B$, we will define the following matrices. Similar calculations are performed in [21].

Let us define

$$
\begin{align*}
\Phi_{1} & \equiv\left(\begin{array}{c}
P^{1} \\
\vdots \\
P^{m_{1}}
\end{array}\right) \in \mathbb{R}^{m_{1} \times n}  \tag{60}\\
\Gamma & \equiv\left(-D_{i^{\prime}, i}^{*}\right)=\left(\sum_{j=1}^{n} \frac{P_{j}^{i^{\prime}} P_{j}^{i}}{Q_{j}^{*}}\right)_{i^{\prime}, i \in \mathcal{I}_{\mathrm{I}}} \in \mathbb{R}^{m_{1} \times m_{1}}  \tag{61}\\
\Lambda & \equiv \operatorname{diag}\left(\lambda_{1}^{*}, \ldots, \lambda_{m_{1}}^{*}\right) \in \mathbb{R}^{m_{1} \times m_{1}} \tag{62}
\end{align*}
$$

Furthermore, define

$$
\begin{align*}
\sqrt{\Lambda} & \equiv \operatorname{diag}\left(\sqrt{\lambda_{1}^{*}}, \ldots, \sqrt{\lambda_{m_{1}}^{*}}\right) \in \mathbb{R}^{m_{1} \times m_{1}}  \tag{63}\\
\Omega & \equiv \operatorname{diag}\left(\left(Q_{1}^{*}\right)^{-1}, \ldots,\left(Q_{n}^{*}\right)^{-1}\right) \in \mathbb{R}^{n \times n}  \tag{64}\\
\sqrt{\Omega} & \equiv \operatorname{diag}\left(\left(Q_{1}^{*}\right)^{-1 / 2}, \ldots,\left(Q_{n}^{*}\right)^{-1 / 2}\right) \in \mathbb{R}^{n \times n} \tag{65}
\end{align*}
$$

Then, we have, by calculation,

$$
\begin{align*}
\sqrt{\Lambda} B \sqrt{\Lambda}^{-1} & =\sqrt{\Lambda} \Gamma \sqrt{\Lambda}  \tag{66}\\
& =\sqrt{\Lambda} \Phi_{1} \Omega^{t} \Phi_{1}{ }^{t} \sqrt{\Lambda}  \tag{67}\\
& =\sqrt{\Lambda} \Phi_{1} \sqrt{\Omega}{ }^{t} \sqrt{\Omega}{ }^{t} \Phi_{1}{ }^{t} \sqrt{\Lambda}  \tag{68}\\
& =\sqrt{\Lambda} \Phi_{1} \sqrt{\Omega}{ }^{t}\left(\sqrt{\Lambda} \Phi_{1} \sqrt{\Omega}\right) \tag{69}
\end{align*}
$$

From (16), $\sqrt{\Lambda}$ is a regular matrix and from the assumption (12), $\operatorname{rank} \Phi_{1}=m_{1}$. Therefore, by $m_{1} \leq m \leq n$, we have $\operatorname{rank} \sqrt{\Lambda} \Phi_{1} \sqrt{\Omega}=m_{1}$, and thus from (69), $\sqrt{\Lambda} B \sqrt{\Lambda}^{-1}$ is symmetric and positive definite. In particular, all the eigenvalues $\beta_{1}, \ldots, \beta_{m_{1}}$ of $B$ are positive. Without loss of generality, let $\beta_{1} \geq \ldots \geq \beta_{m_{1}}>0$. By (59), every component of $B$ is nonnegative and by Lemma 3, every row sum of $B$ is equal to 1 . Hence, by the Perron-Frobenius theorem [10]

$$
\begin{equation*}
1=\beta_{1} \geq \beta_{2} \geq \ldots \geq \beta_{m_{1}}>0 \tag{70}
\end{equation*}
$$

Because $\theta_{i}=1-\beta_{i}, i \in \mathcal{I}_{\mathrm{I}}$, we have

$$
\begin{equation*}
0=\theta_{1} \leq \theta_{2} \leq \ldots \leq \theta_{m_{1}}<1 \tag{71}
\end{equation*}
$$

therefore,
Theorem 2: The eigenvalues of $J^{\mathrm{I}}$ satisfy

$$
\begin{equation*}
0 \leq \theta_{i}<1, i \in \mathcal{I}_{\mathrm{I}} \tag{72}
\end{equation*}
$$

## D. Eigenvalues of $J^{\mathrm{II}}$

Second, we consider the eigenvalues of $J^{\mathrm{II}}$. From (53), (55), we have

Theorem 3: The eigenvalues of $J^{\mathrm{II}}$ satisfy

$$
\begin{equation*}
\theta_{i}=1, i \in \mathcal{I}_{\mathrm{II}} \tag{73}
\end{equation*}
$$

## E. Eigenvalues of $J^{\mathrm{III}}$

Third, we consider the eigenvalues of $J^{\text {IIII }}$. From (53), (56), we have

Theorem 4: The eigenvalues of $J^{\mathrm{III}}$ are $\theta_{i}=e^{D_{i}^{*}-C}, D_{i}^{*}<$ $C, i \in \mathcal{I}_{\mathrm{III}}$, hence

$$
\begin{equation*}
0<\theta_{i}<1, i \in \mathcal{I}_{\mathrm{III}} \tag{74}
\end{equation*}
$$

Remark 4: From the above consideration, we know that all the eigenvalues of the Jacobian matrix $J\left(\boldsymbol{\lambda}^{*}\right)$ are real.

Lemma 5: If the eigenvalues of $J^{\mathrm{I}}$ and $J^{\mathrm{III}}$ are distinct, i.e.,

$$
\begin{equation*}
\theta_{i} \neq \theta_{i^{\prime}}, i \in \mathcal{I}_{\mathrm{I}}, i^{\prime} \in \mathcal{I}_{\mathrm{III}} \tag{75}
\end{equation*}
$$

then $J\left(\boldsymbol{\lambda}^{*}\right)$ is diagonalizable. In particular, if $\mathcal{I}_{\text {III }}=\emptyset$, then (75) holds, and thus $J\left(\boldsymbol{\lambda}^{*}\right)$ is diagonalizable.

Proof: See Appendix A.

## IV. On the exponential convergence

We obtained in Theorems 2, 3, and 4 the evaluation for the eigenvalues of $J\left(\boldsymbol{\lambda}^{*}\right)$. Let $\theta_{\text {max }} \equiv \max _{i \in \mathcal{I}} \theta_{i}$ be the maximum eigenvalue of $J\left(\boldsymbol{\lambda}^{*}\right)$, then by Theorems 2,3 , and 4 , we have $0 \leq \theta_{\text {max }}<1$ if $\mathcal{I}_{\text {II }}$ is empty, and $\theta_{\max }=1$ if $\mathcal{I}_{\text {II }}$ is not empty. First, we show that the convergence is exponential if $\mathcal{I}_{\text {II }}$ is empty.

Theorem 5: Assume $\mathcal{I}_{\mathrm{II}}=\emptyset$, then for any $\theta$ with $\theta_{\max }<$ $\theta<1$, there exist $\delta>0$ and $K>0$, such that for arbitrary initial distribution $\boldsymbol{\lambda}^{0}$ with $\left\|\boldsymbol{\lambda}^{0}-\boldsymbol{\lambda}^{*}\right\|<\delta$, we have

$$
\begin{equation*}
\left\|\boldsymbol{\mu}^{N}\right\|=\left\|\boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}\right\|<K \cdot(\theta)^{N}, N=0,1, \ldots \tag{76}
\end{equation*}
$$

i.e., the convergence is exponential.

Proof: See Appendix B.

## V. On the convergence of order $O(1 / N)$

We will consider the second-order recurrence formula obtained by truncating the Taylor expansion of $F(\boldsymbol{\lambda})$ up to the second-order term and analyze the convergence of order $O(1 / N)$ of the sequence defined by the second-order recurrence formula.

## A. The Hessian matrix $H_{i}\left(\boldsymbol{\lambda}^{*}\right)$

If $0 \leq \theta_{\max }<1$, then the convergence speed of $\boldsymbol{\lambda}^{N} \rightarrow$ $\boldsymbol{\lambda}^{*}$ is determined by the Jacobian matrix $J\left(\boldsymbol{\lambda}^{*}\right)$ because of Theorem 5. However, if $\theta_{\max }=1$, the convergence speed is not determined only by $J\left(\boldsymbol{\lambda}^{*}\right)$; hence, we must investigate the Hessian matrix. In [2], [21], the Jacobian matrix is considered, but the Hessian matrix is neglected in the literature.

Now, we will calculate the Hessian matrix

$$
\begin{equation*}
H_{i}\left(\boldsymbol{\lambda}^{*}\right)=\left(\left.\frac{\partial^{2} F_{i}}{\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}}\right)_{i^{\prime}, i^{\prime \prime} \in \mathcal{I}}, i \in \mathcal{I} \tag{77}
\end{equation*}
$$

of $F_{i}(\boldsymbol{\lambda})$ at $\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}$. We have
Theorem 6: The components of the Hessian matrix $H_{i}\left(\boldsymbol{\lambda}^{*}\right), i \in \mathcal{I}$, are given as follows.

$$
\begin{aligned}
& \left.\frac{\partial^{2} F_{i}}{\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}} \\
& =e^{D_{i}^{*}-C}\left\{\delta_{i i^{\prime}} D_{i, i^{\prime \prime}}^{*}+\delta_{i i^{\prime \prime}} D_{i, i^{\prime}}^{*}+\lambda_{i}^{*}\left(D_{i, i^{\prime}}^{*} D_{i, i^{\prime \prime}}^{*}+D_{i, i^{\prime}, i^{\prime \prime}}^{*}\right)\right. \\
& \quad+\left(\delta_{i i^{\prime}}+\lambda_{i}^{*} D_{i, i^{\prime}}^{*}\right)\left(1-e^{D_{i^{\prime \prime}}^{*}-C}\right) \\
& \left.\quad+\left(\delta_{i i^{\prime \prime}}+\lambda_{i}^{*} D_{i, i^{\prime \prime}}^{*}\right)\left(1-e^{D_{i^{\prime}}^{*}-C}\right)\right\} \\
& \quad+2 \lambda_{i}^{*}\left(1-e^{D_{i^{\prime}}^{*}-C}\right)\left(1-e^{D_{i^{\prime \prime}}^{*}-C}\right) \\
& \quad-\lambda_{i}^{*}\left(e^{D_{i^{\prime}}^{*}-C} D_{i^{\prime}, i^{\prime \prime}}^{*}+e^{D_{i^{\prime \prime}}^{*}-C} D_{i^{\prime}, i^{\prime \prime}}^{*}+E_{i^{\prime}, i^{\prime \prime}}-D_{i^{\prime}, i^{\prime \prime}}^{*}\right) \\
& i, i^{\prime}, i^{\prime \prime} \in \mathcal{I}
\end{aligned}
$$

where $D_{i, i^{\prime}, i^{\prime \prime}}^{*} \equiv \partial^{2} D_{i} /\left.\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}}$ and $E_{i^{\prime}, i^{\prime \prime}} \equiv$ $\sum_{k=1}^{m_{1}} \lambda_{k}^{*} D_{k, i^{\prime}}^{*} D_{k, i^{\prime \prime}}^{*}$.

In particular, if $i \in \mathcal{I}_{\text {II }}$, then $\lambda_{i}^{*}=0$ by (23), and thus

$$
\begin{align*}
& \left.\frac{\partial^{2} F_{i}}{\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}} \\
& =\delta_{i i^{\prime \prime}}\left(1-e^{D_{i^{\prime}}^{*}-C}+D_{i, i^{\prime}}^{*}\right)+\delta_{i i^{\prime}}\left(1-e^{D_{i^{\prime \prime}}^{*}-C}+D_{i, i^{\prime \prime}}^{*}\right) \tag{78}
\end{align*}
$$

Further, if $i \in \mathcal{I}_{\text {II }}$ and $D_{i^{\prime}}^{*}=C$ holds for all $i^{\prime} \in \mathcal{I}$, then by (78), we have

$$
\begin{equation*}
\left.\frac{\partial^{2} F_{i}}{\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}}=\delta_{i i^{\prime}} D_{i, i^{\prime \prime}}^{*}+\delta_{i i^{\prime \prime}} D_{i, i^{\prime}}^{*}, i^{\prime}, i^{\prime \prime} \in \mathcal{I} \tag{79}
\end{equation*}
$$

Proof: See Appendix C.

## B. Analysis of the convergence of order $O(1 / N)$

We consider a recurrence formula obtained by truncating the Taylor expansion (36) up to the second-order term and write the variables as $\overline{\boldsymbol{\mu}}^{N}=\left(\bar{\mu}_{1}^{N}, \ldots, \bar{\mu}_{m}^{N}\right)$. That is, we have

$$
\begin{equation*}
\overline{\boldsymbol{\mu}}^{N+1}=\overline{\boldsymbol{\mu}}^{N} J\left(\boldsymbol{\lambda}^{*}\right)+\frac{1}{2!} \overline{\boldsymbol{\mu}}^{N} H\left(\boldsymbol{\lambda}^{*}\right)^{t} \overline{\boldsymbol{\mu}}^{N} \tag{80}
\end{equation*}
$$

The recurrence formula (80) is called the second-order recurrence formula of the Taylor expansion (36). We investigate the convergence speed of $\overline{\boldsymbol{\mu}}^{N} \rightarrow \mathbf{0}$. The convergence speed of $\overline{\boldsymbol{\mu}}^{N} \rightarrow \mathbf{0}$ seems to be the same as that of the original $\boldsymbol{\mu}^{N} \rightarrow \mathbf{0}$, but the proof has not been obtained. Numerical comparison will be done in Chapter VI. In this chapter, we will prove that, if $\mathcal{I}_{\text {II }} \neq \emptyset$, there exists an initial vector $\overline{\boldsymbol{\mu}}^{0}$ such that $\overline{\boldsymbol{\mu}}^{N} \rightarrow \mathbf{0}$ is the convergence of order $O(1 / N)$. Furthermore, we will consider the condition that $\overline{\boldsymbol{\mu}}^{N} \rightarrow \mathbf{0}$ is the convergence of order $O(1 / N)$ for arbitrary initial vector $\overline{\boldsymbol{\mu}}^{0}$.

The convergence of order $O(1 / N)$ will be proved by the following three steps.

Step 1: Represent $\bar{\mu}_{i}^{N}$ with types-I and -III indices by $\bar{\mu}_{i}^{N}$ with type-II indices.
Step 2: Obtain the recurrence formula satisfied by $\bar{\mu}_{i}^{N}$ with type-II indices.
Step 3: Prove that the convergence of $\bar{\mu}_{i}^{N}$ with type-II indices is of order $O(1 / N)$ for some initial vector $\overline{\boldsymbol{\mu}}^{0}$.

## C. Step 1

Here, we consider the types-I and -III indices together. Then, substitute $\mathcal{I}_{\text {I }} \cup \mathcal{I}_{\text {III }}=\left\{1, \ldots, m^{\prime}\right\}$ and $\mathcal{I}_{\text {II }}=\left\{m^{\prime}+\right.$ $1, \ldots, m\}$. We have $m_{2}=m-m^{\prime}$ and $\left|\mathcal{I}_{\text {II }}\right|=m_{2}$. The purpose of step 1 is to represent $\overline{\boldsymbol{\mu}}_{\mathrm{I}, \mathrm{IIII}}^{N} \equiv\left(\bar{\mu}_{1}^{N}, \ldots, \bar{\mu}_{m^{\prime}}^{N}\right)$ by $\overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N} \equiv\left(\bar{\mu}_{m^{\prime}+1}^{N}, \ldots, \bar{\mu}_{m}^{N}\right)$.

In the Jacobian matrix, by changing the order of $J^{\mathrm{II}}$ and $J^{\text {III }}$, we have

$$
J\left(\boldsymbol{\lambda}^{*}\right)=\left(\begin{array}{ccc}
J^{\mathrm{I}} & O & O  \tag{81}\\
* & J^{\mathrm{III}} & O \\
* & O & J^{\mathrm{II}}
\end{array}\right)
$$

Then, by defining

$$
J^{\prime} \equiv\left(\begin{array}{cc}
J^{\mathrm{I}} & O  \tag{82}\\
* & J^{\mathrm{III}}
\end{array}\right)
$$

we have

$$
J\left(\boldsymbol{\lambda}^{*}\right)=\left(\begin{array}{cc}
J^{\prime} & O  \tag{83}\\
* & J^{\mathrm{II}}
\end{array}\right)
$$

The eigenvalues of $J\left(\boldsymbol{\lambda}^{*}\right)$ are $\theta_{i}, i=1, \ldots, m$, and then the eigenvalues of $J^{\prime}$ are $\theta_{i}, i=1, \ldots, m^{\prime}$ with $0 \leq \theta_{i}<1$, and those of $J^{\mathrm{II}}$ are $\theta_{i}, i=m^{\prime}+1, \ldots, m$ with $\theta_{i}=1$.

Now, let $\boldsymbol{a}_{i}$ be a right eigenvector of $J\left(\boldsymbol{\lambda}^{*}\right)$ for $\theta_{i}$ and define

$$
\begin{equation*}
A \equiv\left(\boldsymbol{a}_{1}, \ldots, \boldsymbol{a}_{m}\right) \in \mathbb{R}^{m \times m} \tag{84}
\end{equation*}
$$

Under the assumption (75), by choosing the eigenvectors $\boldsymbol{a}_{1}, \ldots, \boldsymbol{a}_{m}$ appropriately, we can make $A$ a regular matrix. In fact, because $J\left(\boldsymbol{\lambda}^{*}\right)$ is diagonalizable by Lemma 5 , the direct sum of all the eigenspaces spans the whole $\mathbb{R}^{m}$ (See [19], p.161, Example 4).

For $i=m^{\prime}+1, \ldots, m$, define

$$
\begin{gather*}
\stackrel{i \text { th }}{\vee}  \tag{85}\\
\boldsymbol{e}_{i}=(0, \ldots, 0, \stackrel{1}{1}, 0, \ldots, 0) \in \mathbb{R}^{m} \\
i=m^{\prime}+1, \ldots, m
\end{gather*}
$$

Then, because $\theta_{i}=1$, we can take

$$
\begin{equation*}
\boldsymbol{a}_{i}={ }^{t} \boldsymbol{e}_{i}, i=m^{\prime}+1, \ldots, m \tag{86}
\end{equation*}
$$

Therefore, we have

$$
\begin{align*}
A & =\left(\boldsymbol{a}_{1}, \ldots, \boldsymbol{a}_{m^{\prime}},{ }^{t} \boldsymbol{e}_{m^{\prime}+1}, \ldots,{ }^{t} \boldsymbol{e}_{m}\right)  \tag{87}\\
& =\left(\begin{array}{ccc:ccc}
a_{11} & \ldots & a_{m^{\prime} 1} & \\
\vdots & & \vdots & & O & \\
a_{1 m^{\prime}} & \ldots & a_{m^{\prime} m^{\prime}} & & & \\
\hdashline a_{1, m^{\prime}+1} & \ldots & a_{m^{\prime}, m^{\prime}+1} & 1 & \ldots & 0 \\
\vdots & & \vdots & \vdots & \ddots & \vdots \\
a_{1 m} & \ldots & a_{m^{\prime} m} & 0 & \ldots & 1
\end{array}\right)  \tag{88}\\
& \equiv\left(\begin{array}{cc}
A_{1} & O \\
A_{2} & I
\end{array}\right), \tag{89}
\end{align*}
$$

where

$$
\begin{align*}
A_{1} & \equiv\left(\begin{array}{ccc}
a_{11} & \ldots & a_{m^{\prime} 1} \\
\vdots & & \vdots \\
a_{1 m^{\prime}} & \ldots & a_{m^{\prime} m^{\prime}}
\end{array}\right),  \tag{90}\\
A_{2} & \equiv\left(\begin{array}{ccc}
a_{1, m^{\prime}+1} & \ldots & a_{m^{\prime}, m^{\prime}+1} \\
\vdots & & \vdots \\
a_{1 m} & \ldots & a_{m^{\prime} m}
\end{array}\right) . \tag{91}
\end{align*}
$$

Because $A$ is regular, $A_{1}$ is also regular by (89). $J\left(\boldsymbol{\lambda}^{*}\right)$ is diagonalized by $A$, i.e., $A^{-1} J\left(\boldsymbol{\lambda}^{*}\right) A=\Theta$, where

$$
\begin{align*}
& \Theta=\left(\begin{array}{cc}
\Theta_{1} & O \\
O & I
\end{array}\right), \Theta_{1}=\left(\begin{array}{ccc}
\theta_{1} & & O \\
& \ddots & \\
O & & \theta_{m^{\prime}}
\end{array}\right)  \tag{92}\\
& 0 \leq \theta_{i}<1, i=1, \ldots, m^{\prime} .
\end{align*}
$$

Using only the first-order term of the Taylor expansion (36), we have

$$
\begin{align*}
\boldsymbol{\mu}^{N+1} A & =\boldsymbol{\mu}^{N} J\left(\boldsymbol{\lambda}^{*}\right) A  \tag{93}\\
& =\boldsymbol{\mu}^{N} A \Theta \tag{94}
\end{align*}
$$

thus, by $\boldsymbol{\mu}^{N}=\left(\boldsymbol{\mu}_{\mathrm{I}, \mathrm{III}}^{N}, \boldsymbol{\mu}_{\mathrm{II}}^{N}\right)$, (94), and (89),

$$
\begin{equation*}
\boldsymbol{\mu}_{\mathrm{I}, \mathrm{III}}^{N+1} A_{1}+\boldsymbol{\mu}_{\mathrm{II}}^{N+1} A_{2}=\left(\boldsymbol{\mu}_{\mathrm{I}, \mathrm{III}}^{N} A_{1}+\boldsymbol{\mu}_{\mathrm{II}}^{N} A_{2}\right) \Theta_{1} \tag{95}
\end{equation*}
$$

Hence, if the second- and higher-order terms are negligible, we have

$$
\begin{equation*}
\boldsymbol{\mu}_{\mathrm{I}, \mathrm{III}}^{N} A_{1}+\boldsymbol{\mu}_{\mathrm{II}}^{N} A_{2} \rightarrow \mathbf{0} \text { (exponentially), } N \rightarrow \infty \tag{96}
\end{equation*}
$$

To show that the second- and higher-order terms are negligible, let us consider $\boldsymbol{\mu}^{N+1} \boldsymbol{a}_{i}$ as a function of $\boldsymbol{\mu}^{N} \boldsymbol{a}_{1}, \ldots, \boldsymbol{\mu}^{N} \boldsymbol{a}_{m^{\prime}}$. If the Taylor expansion (36) satisfies the condition that

$$
\begin{equation*}
\boldsymbol{\mu}^{N+1} \boldsymbol{a}_{i} \text { is divisible by } \boldsymbol{\mu}^{N} \boldsymbol{a}_{i} \tag{97}
\end{equation*}
$$

i.e., if $\boldsymbol{\mu}^{N} \boldsymbol{a}_{i}=0$ implies that $\boldsymbol{\mu}^{N+1} \boldsymbol{a}_{i}=0$, then, we have

$$
\begin{align*}
\boldsymbol{\mu}^{N+1} \boldsymbol{a}_{i} & =\theta_{i} \boldsymbol{\mu}^{N} \boldsymbol{a}_{i}\left(1+o\left(\left|\boldsymbol{\mu}^{N} \boldsymbol{a}_{i}\right|\right)\right), N \rightarrow \infty  \tag{98}\\
i & =1, \ldots, m^{\prime}
\end{align*}
$$

and hence (96) holds. However, in general, (97) is difficult to prove. We will show later in Examples 6 and 7 that (97) holds.

In what follows, we assume (96) and regard it as $\boldsymbol{\mu}_{\mathrm{I}, \mathrm{III}}^{N} A_{1}+$ $\boldsymbol{\mu}_{\mathrm{II}}^{N} A_{2}=\mathbf{0}$. Then, we replace $\boldsymbol{\mu}^{N}$ by $\overline{\boldsymbol{\mu}}^{N}$ to have $\overline{\boldsymbol{\mu}}_{\mathrm{I}, \mathrm{III}}^{N} A_{1}+$ $\overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N} A_{2}=\mathbf{0}$ and hence

$$
\begin{equation*}
\overline{\boldsymbol{\mu}}_{\mathrm{I}, \mathrm{III}}^{N}=-\overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N} A_{2} A_{1}^{-1} \tag{99}
\end{equation*}
$$

The validity of (99) will be checked by numerical examples.
Step 1 is achieved by (99).

## D. Step 2

The purpose of step 2 is to obtain a recurrence formula satisfied by $\bar{\mu}_{i}^{N}, i=m^{\prime}+1, \ldots, m$.

The $i$-th component of (80) for $i=m^{\prime}+1, \ldots, m$ is

$$
\begin{gather*}
\bar{\mu}_{i}^{N+1}=\bar{\mu}_{i}^{N}+\frac{1}{2!} \bar{\mu}^{N} H_{i}\left(\boldsymbol{\lambda}^{*}\right)^{t} \overline{\boldsymbol{\mu}}^{N}  \tag{100}\\
i=m^{\prime}+1, \ldots, m
\end{gather*}
$$

We will represent the second term on the right hand side of (100) by $\overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N}$.

Let $H_{i, i^{\prime} i^{\prime \prime}}$ be the $\left(i^{\prime}, i^{\prime \prime}\right)$ component of the Hessian matrix $H_{i}\left(\boldsymbol{\lambda}^{*}\right)$. Then, by (78) of Theorem 6, we have

$$
\begin{align*}
& H_{i, i^{\prime} i^{\prime \prime}}=\left.\frac{\partial^{2} F_{i}}{\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}} \\
& =\delta_{i i^{\prime \prime}}\left(1-e^{D_{i^{\prime}}^{*}-C}+D_{i, i^{\prime}}^{*}\right)+\delta_{i i^{\prime}}\left(1-e^{D_{i^{\prime \prime}}^{*}-C}+D_{i, i^{\prime \prime}}^{*}\right) \tag{101}
\end{align*}
$$

$$
i=m^{\prime}+1, \ldots, m, i^{\prime}, i^{\prime \prime}=1, \ldots, m
$$

Here, for the simplicity of symbols, define

$$
\begin{align*}
& S_{i i^{\prime}} \equiv 1-e^{D_{i^{\prime}}^{*}-C}+D_{i, i^{\prime}}^{*}  \tag{102}\\
& \quad i=m^{\prime}+1, \ldots, m, i^{\prime}=1, \ldots, m
\end{align*}
$$

Then, we can write (101) as

$$
\begin{equation*}
H_{i, i^{\prime} i^{\prime \prime}}=\delta_{i i^{\prime \prime}} S_{i i^{\prime}}+\delta_{i i^{\prime}} S_{i i^{\prime \prime}} \tag{103}
\end{equation*}
$$

Further, writing $A_{1}^{-1} \equiv\left(\zeta_{i^{\prime} i^{\prime \prime}}\right), i^{\prime}, i^{\prime \prime}=1, \ldots, m^{\prime}, T_{i i^{\prime}} \equiv$ $-\sum_{k, k^{\prime}=1}^{m^{\prime}} a_{i^{\prime} k} \zeta_{k k^{\prime}} S_{i k^{\prime}}, i, i^{\prime}=m^{\prime}+1, \ldots, m$, and $r_{i i^{\prime}} \equiv$ $T_{i i^{\prime}}+D_{i, i^{\prime}}^{*}, i, i^{\prime}=m^{\prime}+1, \ldots, m$, we have the following theorem.

Theorem 7: $\left\{\bar{\mu}_{i}^{N}\right\}, i=m^{\prime}+1, \ldots, m$, satisfies the recurrence formula

$$
\begin{align*}
\bar{\mu}_{i}^{N+1} & =\bar{\mu}_{i}^{N}+\bar{\mu}_{i}^{N} \sum_{i^{\prime}=m^{\prime}+1}^{m} r_{i i^{\prime}} \bar{\mu}_{i^{\prime}}^{N}  \tag{104}\\
i & =m^{\prime}+1, \ldots, m .
\end{align*}
$$

## Proof: See Appendix D.

Step 2 is achieved by (104).

## E. Step 3

The purpose of step 3 is to prove that $\bar{\mu}_{i}^{N} \rightarrow 0, i=$ $1, \ldots, m$, is the convergence of order $O(1 / N)$. We will define the canonical form of the recurrence formula (104). Writing $R \equiv\left(r_{i i^{\prime}}\right), i, i^{\prime}=m^{\prime}+1, \ldots, m$, and $\mathbf{1} \equiv(1, \ldots, 1) \in \mathbb{R}^{m_{2}}$, we consider the equation $R^{t} \boldsymbol{\sigma}=-{ }^{t} \mathbf{1}$ for the variables $\boldsymbol{\sigma}=\left(\sigma_{m^{\prime}+1}, \ldots, \sigma_{m}\right)$. Assuming that $R$ is regular, we have

$$
\begin{equation*}
{ }^{t} \boldsymbol{\sigma}=-R^{-1 t} \mathbf{1} \tag{105}
\end{equation*}
$$

and then we assume $\sigma>0$, i.e., $\sigma_{i}>0, i=m^{\prime}+1, \ldots, m$. Further, by substituting

$$
\begin{align*}
\nu_{i}^{N} & \equiv \bar{\mu}_{i}^{N} / \sigma_{i}, i=m^{\prime}+1, \ldots, m  \tag{106}\\
p_{i i^{\prime}} & \equiv-r_{i i^{\prime}} \sigma_{i^{\prime}}, i, i^{\prime}=m^{\prime}+1, \ldots, m \tag{107}
\end{align*}
$$

the recurrence formula (104) becomes

$$
\begin{align*}
\nu_{i}^{N+1} & =\nu_{i}^{N}-\nu_{i}^{N} \sum_{i^{\prime}=1}^{m} p_{i i^{\prime}} \nu_{i^{\prime}}^{N}  \tag{108}\\
i & =m^{\prime}+1, \ldots, m
\end{align*}
$$

where $\boldsymbol{p}_{i} \equiv\left(p_{i, m^{\prime}+1}, \ldots, p_{i, m}\right)$ is a probability vector.
The recurrence formula (108) is called the canonical form of (104).

For the analysis of (108), we prepare the following lemma.
Lemma 6: Let us define a positive sequence $\left\{\nu^{N}\right\}_{N=0,1 \ldots}$ by the recurrence formula

$$
\begin{align*}
& \nu^{N+1}=\nu^{N}-\left(\nu^{N}\right)^{2}, N=0,1, \ldots  \tag{109}\\
& 0<\nu^{0} \leq 1 / 2 \tag{110}
\end{align*}
$$

Then, we have

$$
\begin{equation*}
\lim _{N \rightarrow \infty} N \nu^{N}=1 \tag{111}
\end{equation*}
$$

Proof: Since the function $g(\nu) \equiv \nu-\nu^{2}$ satisfies $0<g(\nu)<\nu$ for $0<\nu \leq 1 / 2$, we see $0<\nu^{N+1}<$ $\nu^{N}, N=0,1, \ldots$ by mathematical induction. Thus, $\nu^{\infty} \equiv$ $\lim _{N \rightarrow \infty} \nu^{N} \geq 0$ exists, and by (109), $\nu^{\infty}=\nu^{\infty}-\left(\nu^{\infty}\right)^{2}$ holds. Then, we have $\nu^{\infty}=0$.

Next, by (109), we have

$$
\begin{align*}
\frac{1}{N}\left(\frac{1}{\nu^{N}}-\frac{1}{\nu^{0}}\right) & =\frac{1}{N} \sum_{l=0}^{N-1}\left(\frac{1}{\nu^{l+1}}-\frac{1}{\nu^{l}}\right)  \tag{112}\\
& =\frac{1}{N} \sum_{l=0}^{N-1} \frac{1}{1-\nu^{l}} \tag{113}
\end{align*}
$$

Applying the proposition that "the arithmetic mean of the first $n$ terms of a convergent sequence converges to the same limit
as the original sequence" ([1], p.37, Exercise 2) to (113), we have

$$
\begin{align*}
\lim _{N \rightarrow \infty} \frac{1}{N \nu^{N}} & =\lim _{N \rightarrow \infty} \frac{1}{1-\nu^{N}}  \tag{114}\\
& =1 \tag{115}
\end{align*}
$$

which proves the lemma.
Lemma 7: For the sequence $\left\{\nu_{i}^{N}\right\}, i=m^{\prime}+1, \ldots, m$, defined by (108) with the initial values $\nu_{i}^{0}=1 / 2$, we have

$$
\begin{equation*}
\lim _{N \rightarrow \infty} N \nu_{i}^{N}=1, i=m^{\prime}+1, \ldots, m \tag{116}
\end{equation*}
$$

Proof: By mathematical induction, we see that $\nu_{m^{\prime}+1}^{N}=$ $\ldots=\nu_{m}^{N}$ holds for $N=0,1, \ldots$, and thus (108) becomes $\nu_{i}^{N+1}=\nu_{i}^{N}-\left(\nu_{i}^{N}\right)^{2}, i=m^{\prime}+1, \ldots, m$. Therefore, (116) holds by Lemma 6.

Theorem 8: For (104) with the initial values $\bar{\mu}_{i}^{0}=\sigma_{i} / 2, i=$ $m^{\prime}+1, \ldots, m$, we have

$$
\begin{equation*}
\lim _{N \rightarrow \infty} N \bar{\mu}_{i}^{N}=\sigma_{i}, i=m^{\prime}+1, \ldots, m \tag{117}
\end{equation*}
$$

Further, under the assumption (99), we have

$$
\begin{equation*}
\lim _{N \rightarrow \infty} N \bar{\mu}_{i}^{N}=-\left(\boldsymbol{\sigma} A_{2} A_{1}^{-1}\right)_{i}, i=1, \ldots, m^{\prime} \tag{118}
\end{equation*}
$$

where $\sigma, A_{1}, A_{2}$ were defined by (105), (91).
Proof: From $\nu_{i}^{N}=\bar{\mu}_{i}^{N} / \sigma_{i}, \quad i=m^{\prime}+1, \ldots, m$, and Lemma 7, we obtain (117). Further, by (99), we obtain (118).

Step 3 is achieved by Theorem 8. Summarizing the above, we have the following theorem.

Theorem 9: If type-II indices exist, then there exists an initial vector $\overline{\boldsymbol{\mu}}^{0}$ such that $\overline{\boldsymbol{\mu}}^{N} \rightarrow \mathbf{0}$ is the convergence of order $O(1 / N)$.

Remark 5: In Theorem 9, we must choose a specific initial vector $\overline{\boldsymbol{\mu}}^{0}$, but we want to prove it for an arbitrary initial vector. If the existence of $\lim _{N \rightarrow \infty} N \mu_{i}^{N}, i=1, \ldots, m$, is proved for arbitrary $\overline{\boldsymbol{\mu}}^{0}$, then we have (117) and (118). However, the analysis for the convergence speed of the recurrence formula (104) for an arbitrary initial vector is very difficult, so it has not been achieved. We will give a proof under the assumption that a conjecture holds.

## F. On the initial vector

In this section, we will investigate the convergence of (108) for an arbitrary initial vector.

Now, for the sake of simplicity, we will change the indices and symbols of the canonical form (108). Noting that $m-m^{\prime}=m_{2}$, we change the indices from $m^{\prime}+1, \ldots, m$ to $1, \ldots, m_{2}$, and define

$$
\begin{align*}
\xi_{i}^{N} & \equiv \nu_{m^{\prime}+i}, i=1, \ldots, m_{2}  \tag{119}\\
q_{i i^{\prime}} & \equiv p_{m^{\prime}+i, m^{\prime}+i^{\prime}}, i, i^{\prime}=1, \ldots, m_{2} \tag{120}
\end{align*}
$$

which are just shifting the indices. By the above change, the canonical form (108) becomes

$$
\begin{equation*}
\xi_{i}^{N+1}=\xi_{i}^{N}-\xi_{i}^{N} \sum_{i^{\prime}=1}^{m_{2}} q_{i i^{\prime}} \xi_{i^{\prime}}^{N}, i=1, \ldots, m_{2} \tag{121}
\end{equation*}
$$

where $\boldsymbol{q}_{i} \equiv\left(q_{i 1}, \ldots, q_{i, m_{2}}\right) \in \mathbb{R}^{m_{2}}$ is a probability vector.

## G. Diagonally dominant condition

For the probability vectors $\boldsymbol{q}_{i}, i=1, \cdots, m_{2}$, we assume

$$
\begin{equation*}
q_{i i}>\sum_{i^{\prime}=1, i^{\prime} \neq i}^{m_{2}} q_{i i^{\prime}}, i=1, \ldots, m_{2} \tag{122}
\end{equation*}
$$

This assumption means that the matrix made by arranging the row vectors $\boldsymbol{q}_{1}, \ldots, \boldsymbol{q}_{m_{2}}$ vertically is diagonally dominant. We call (122) the diagonally dominant condition. By numerical calculation, we confirmed that (122) holds for all our examples.

## H. Conjecture

Our goal is to prove $\lim _{N \rightarrow \infty} N \xi_{i}^{N}=1, i=1, \ldots, m_{2}$, for any initial vector $\boldsymbol{\xi}^{0}$. We found that we can prove this if the following conjecture holds; however, the proof of this conjecture has not yet been obtained.

Conjecture: By reordering the indices if necessary, there exists an $N_{0}$ such that the following inequalities hold.

$$
\begin{equation*}
\xi_{1}^{N} \geq \xi_{2}^{N} \geq \ldots \geq \xi_{m_{2}}^{N}, N \geq N_{0} \tag{123}
\end{equation*}
$$

Many numerical examples we observed seem to support this conjecture.

We have
Theorem 10: We assume that the sequence defined by the recurrence formula (121) satisfies (122) and (123). Then, for arbitrary initial values with $0<\xi_{i}^{0} \leq 1 / 2, i=1, \ldots, m_{2}$, we have

$$
\begin{equation*}
\lim _{N \rightarrow \infty} N \xi_{i}^{N}=1, i=1, \ldots, m_{2} \tag{124}
\end{equation*}
$$

## Proof: See Appendix E.

## I. Special cases where the conjecture holds

Now, we consider some special cases where the conjecture (123) holds.

If $m_{2}=1$, then the variable is only $\xi_{1}^{N}$, hence we do not need to consider the inequality condition.

If $m_{2}=2$, the canonical form (121) becomes

$$
\begin{align*}
\xi_{1}^{N+1} & =\xi_{1}^{N}-q_{11}\left(\xi_{1}^{N}\right)^{2}-q_{12} \xi_{1}^{N} \xi_{2}^{N}  \tag{125}\\
\xi_{2}^{N+1} & =\xi_{2}^{N}-q_{21} \xi_{2}^{N} \xi_{1}^{N}-q_{22}\left(\xi_{2}^{N}\right)^{2}  \tag{126}\\
q_{11} & >0, q_{12}>0, q_{11}+q_{12}=1  \tag{127}\\
q_{21} & >0, q_{22}>0, q_{21}+q_{22}=1 . \tag{128}
\end{align*}
$$

By calculation, we have

$$
\begin{equation*}
\xi_{1}^{N+1}-\xi_{2}^{N+1}=\left(\xi_{1}^{N}-\xi_{2}^{N}\right)\left(1-q_{11} \xi_{1}^{N}-q_{22} \xi_{2}^{N}\right) \tag{129}
\end{equation*}
$$

By Lemma 12 in Appendix E, we have $1-q_{11} \xi_{1}^{N}-q_{22} \xi_{2}^{N}>0$ for any $N=0,1, \ldots$. Hence, if $\xi_{1}^{0} \geq \xi_{2}^{0}$, then $\xi_{1}^{N} \geq \xi_{2}^{N}, N=$ $0,1, \ldots$. Thus, the conjecture (123) holds with $N_{0}=0$. Note that, when $m_{2}=2$, the diagonally dominant condition (122) is not necessary, because it is not used in the above proof.

In the case of $m_{2} \geq 3$, this problem is very difficult. We have a sufficient condition for (123) in the following lemma.

Lemma 8: For $i^{\prime}$, consider $q_{i i^{\prime}}, i=1, \ldots, m_{2}$. Assume that $q_{i i^{\prime}}$ are equal for all $i$ except $i^{\prime}$, i.e., $q_{1 i^{\prime}}=\ldots=q_{i^{\prime}-1, i^{\prime}}=$ $q_{i^{\prime}+1, i^{\prime}}=\ldots=q_{m_{2}, i^{\prime}}$. Then, the conjecture (123) holds.

Proof: By calculation, we have

$$
\begin{align*}
& \xi_{1}^{N+1}-\xi_{2}^{N+1} \\
& =\left(\xi_{1}^{N}-\xi_{2}^{N}\right)\left(1-q_{11} \xi_{1}^{N}-q_{22} \xi_{2}^{N}-\sum_{i^{\prime}=3}^{m_{2}} q_{1 i^{\prime}} \xi_{i^{\prime}}^{N}\right) \tag{130}
\end{align*}
$$

The second factor on the right hand side of (130) is positive for all $N=0,1, \ldots$ by Lemma 12. Hence, if $\xi_{1}^{0} \geq \xi_{2}^{0}$, then $\xi_{1}^{N} \geq \xi_{2}^{N}$ holds for $N=0,1, \ldots$. Similarly, by considering any pair $\xi_{i}^{N}$ and $\xi_{i^{\prime}}^{N}$, we see that (123) holds.
Summarizing the above, we have
Theorem 11: Under the assumptions of Lemma 8 and the diagonally dominant condition (122), the convergence of the recurrence formula (121) is of order $O(1 / N)$ for arbitrary initial vector $\boldsymbol{\xi}^{0}=\left(\xi_{1}^{0}, \ldots, \xi_{m_{2}}^{0}\right)$ with $0<\xi_{i}^{0} \leq 1 / 2, i=$ $1, \ldots, m_{2}$, and

$$
\begin{equation*}
\lim _{N \rightarrow \infty} N \xi_{i}^{N}=1, i=1, \ldots, m_{2} \tag{131}
\end{equation*}
$$

## VI. Numerical Evaluation

Based on the analysis in the previous sections, we will evaluate numerically the convergence speed of the ArimotoBlahut algorithm for several channel matrices.
In Examples 4 and 5 below, we will investigate the exponential convergence, where all the indices are of typeI , in other words, the capacity-achieving $\lambda^{*}$ is in $\Delta(\mathcal{X})^{\circ}$ (the interior of $\Delta(\mathcal{X})$ ). In Example 5, we will discuss how the convergence speed depends on the choice of the initial distribution $\boldsymbol{\lambda}^{0}$. Next, in Examples 6 and 7, we will consider the convergence of order $O(1 / N)$. In these examples, there exist type-II indices. It will be confirmed that the convergence speed is accurately approximated by the values obtained in Theorem 8. In Example 8, we will investigate the exponential convergence, where all the indices are of type-I or -III, i.e., there exist no type-II indices.
Here, in the case of exponential convergence, we will evaluate the values of the function

$$
\begin{equation*}
L(N) \equiv-\frac{1}{N} \log \left\|\boldsymbol{\mu}^{N}\right\| \tag{132}
\end{equation*}
$$

Based on the results of Theorem 5, i.e., $\left\|\boldsymbol{\mu}^{N}\right\|=\left\|\boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}\right\|<$ $K \cdot(\theta)^{N}, \theta \doteqdot \theta_{\max }$, we will compare $L(N)$ for large $N$ with $-\log \theta_{\max }$ (or other values).

On the other hand, in the case of the convergence of order $O(1 / N)$, we will evaluate

$$
\begin{equation*}
N \boldsymbol{\mu}^{N}=\left(N \mu_{1}^{N}, \ldots, N \mu_{m}^{N}\right) \tag{133}
\end{equation*}
$$

We will compare $N \boldsymbol{\mu}^{N}$ for large $N$ with the values obtained in Theorem 8.

## A. Exponential convergence where all the indices are of type-I

Example 4: Consider the channel matrix $\Phi^{(1)}$ of (28), i.e.,

$$
\Phi^{(1)}=\left(\begin{array}{lll}
0.800 & 0.100 & 0.100  \tag{134}\\
0.100 & 0.800 & 0.100 \\
0.250 & 0.250 & 0.500
\end{array}\right)
$$



Fig. 5. Convergence of $L(N)$ in Example 4 with initial distribution $\boldsymbol{\lambda}^{0}=$ $(1 / 3,1 / 3,1 / 3)$.
and an initial distribution $\boldsymbol{\lambda}^{0}=(1 / 3,1 / 3,1 / 3)$. We have

$$
\begin{align*}
C & =0.323[\mathrm{nat} / \mathrm{symbol}]  \tag{135}\\
\boldsymbol{\lambda}^{*} & =(0.431,0.431,0.138),  \tag{136}\\
Q^{*} & =(0.422,0.422,0.156),  \tag{137}\\
J\left(\boldsymbol{\lambda}^{*}\right) & =\left(\begin{array}{ccc}
0.308 & -0.191 & -0.117 \\
-0.191 & 0.308 & -0.117 \\
-0.369 & -0.369 & 0.738
\end{array}\right) . \tag{138}
\end{align*}
$$

By (136), we see that $\lambda_{1}^{*}>0, \lambda_{2}^{*}>0, \lambda_{3}^{*}>0$, thus all the indices are of type-I.

The eigenvalues of $J\left(\boldsymbol{\lambda}^{*}\right)$ are $\left(\theta_{1}, \theta_{2}, \theta_{3}\right)=(0.000,0.500$, $0.855)$. Then, $\theta_{\max }=\theta_{3}=0.855$. We have, for $N=500$,

$$
\begin{equation*}
L(500)=0.161 \doteqdot-\log \theta_{\max }=0.157 \tag{139}
\end{equation*}
$$

We can see from Fig. 5 that $L(N)$ for large $N$ is accurately approximated by the value $-\log \theta_{\max }$.

Example 5: Let us consider another channel matrix. Define

$$
\Phi^{(4)} \equiv\left(\begin{array}{ccc}
0.793 & 0.196 & 0.011  \tag{140}\\
0.196 & 0.793 & 0.011 \\
0.250 & 0.250 & 0.500
\end{array}\right)
$$

We have

$$
\begin{align*}
C & =0.294[\mathrm{nat} / \mathrm{symbol}],  \tag{141}\\
\boldsymbol{\lambda}^{*} & =(0.352,0.352,0.296),  \tag{142}\\
Q^{*} & =(0.422,0.422,0.156),  \tag{143}\\
J\left(\boldsymbol{\lambda}^{*}\right) & =\left(\begin{array}{ccc}
0.443 & -0.260 & -0.183 \\
-0.260 & 0.443 & -0.183 \\
-0.218 & -0.218 & 0.436
\end{array}\right) . \tag{144}
\end{align*}
$$

By (142), we see that all the indices are of type-I.
The eigenvalues of $J\left(\boldsymbol{\lambda}^{*}\right)$ are $\left(\theta_{1}, \theta_{2}, \theta_{3}\right)=(0.000,0.618$, 0.702 ). Then, $\theta_{\max }=\theta_{3}=0.702$. Write the second largest eigenvalue as $\theta_{\mathrm{sec}}$, then $\theta_{\mathrm{sec}}=\theta_{2}=0.618$.

We show in Fig. 6 the graph of $L(N)$ with initial distribution $\boldsymbol{\lambda}_{1}^{0} \equiv(1 / 3,1 / 3,1 / 3)$ by the solid line, and the graph with initial distribution $\boldsymbol{\lambda}_{2}^{0} \equiv(1 / 2,1 / 3,1 / 6)$ by the dotted line.


Fig. 6. Convergence of $L(N)$ in Example 5 with initial distribution $\boldsymbol{\lambda}_{1}^{0}=$ $(1 / 3,1 / 3,1 / 3)$ and $\boldsymbol{\lambda}_{2}^{0}=(1 / 2,1 / 3,1 / 6)$.

The larger $L(N)$ is, the faster the convergence, hence the convergence with $\boldsymbol{\lambda}_{1}^{0}$ is faster than with $\boldsymbol{\lambda}_{2}^{0}$. The convergence speed varies depending on the choice of the initial distribution. What kind of initial distribution yields faster convergence? We will investigate this question below.

Let us define

$$
\begin{align*}
& \boldsymbol{\mu}_{1}^{0} \equiv \boldsymbol{\lambda}_{1}^{0}-\boldsymbol{\lambda}^{*}=(-0.019,-0.019,0.038)  \tag{145}\\
& \boldsymbol{\mu}_{2}^{0} \equiv \boldsymbol{\lambda}_{2}^{0}-\boldsymbol{\lambda}^{*}=(0.148,-0.019,-0.129) \tag{146}
\end{align*}
$$

We will execute the following calculation by assuming $\boldsymbol{\mu}^{N+1}=\boldsymbol{\mu}^{N} J\left(\boldsymbol{\lambda}^{*}\right), N=0,1, \ldots$ holds exactly.

Here, we will investigate for general $m, n$. We assume (75). Let $\boldsymbol{b}_{\text {max }}$ be the left eigenvector of $J\left(\boldsymbol{\lambda}^{*}\right)$ for $\theta_{\max }$, and let $\boldsymbol{b}_{\text {max }}^{\perp}$ be the orthogonal complement of $\boldsymbol{b}_{\max }$, i.e., $\boldsymbol{b}_{\max }^{\perp} \equiv$ $\left\{\boldsymbol{\mu} \in \mathbb{R}^{m} \mid \boldsymbol{\mu}^{t} \boldsymbol{b}_{\max }=0\right\}$.

Lemma 9: If

$$
\begin{equation*}
\boldsymbol{\mu}^{N} \in \boldsymbol{b}_{\max }^{\perp}, N=0,1, \ldots, \tag{147}
\end{equation*}
$$

then for any $\theta$ with $\theta_{\text {sec }}<\theta<1$, we have $\left\|\boldsymbol{\mu}^{N}\right\|<K$. $(\theta)^{N}, K>0, N=0,1, \ldots$.

Proof: See Appendix F.
Because $\theta_{\text {sec }}<\theta_{\max }$, if (147) holds, then the convergence speed is faster than $\theta_{\max }$ by Lemma 9. The next lemma gives a necessary and sufficient condition that guarantees (147).

Lemma 10: $\boldsymbol{\mu} J\left(\boldsymbol{\lambda}^{*}\right) \in \boldsymbol{b}_{\max }^{\perp}$ holds for any $\boldsymbol{\mu} \in \boldsymbol{b}_{\max }^{\perp}$ if and only if ${ }^{t} \boldsymbol{b}_{\text {max }}$ is a right eigenvector for $\theta_{\text {max }}$.

Proof: See Appendix G.
If ${ }^{t} \boldsymbol{b}_{\text {max }}$ is a right eigenvector, then by Lemma 10, any $\boldsymbol{\mu}^{0} \in$ $b_{\text {max }}^{\perp}$ yields (147), hence the convergence becomes faster.

Now, we will evaluate the convergence speed for the initial vectors (145) and (146). For $J\left(\boldsymbol{\lambda}^{*}\right)$ of (144), $\theta_{\max }=0.702$ and $\theta_{\text {sec }}=0.618$. The left eigenvector for $\theta_{\max }$ is $\boldsymbol{b}_{\max }=$ $(-0.500,0.500,0.000)$. We can confirm that ${ }^{t} \boldsymbol{b}_{\text {max }}$ is a right eigenvector for $\theta_{\max }$ and $\boldsymbol{\mu}_{1}^{0}{ }^{t} \boldsymbol{b}_{\text {max }}=0$, thus by Lemmas 9
and 10 , we have $\lim _{N \rightarrow \infty} L(N) \doteqdot-\log \theta_{\text {sec }}$. Then, by the solid line in Fig. 6, we have, for $N=500$,

$$
\begin{equation*}
L(500)=0.489 \doteqdot-\log \theta_{\sec }=0.481 \tag{148}
\end{equation*}
$$

On the other hand, we have $\boldsymbol{\mu}_{2}^{0}{ }^{t} \boldsymbol{b}_{\text {max }} \neq 0$, thus by Lemma 10 , we have $\lim _{N \rightarrow \infty} L(N) \doteqdot-\log \theta_{\max }$. Then, by the dotted line, we have, for $N=500$,

$$
\begin{equation*}
L(500)=0.360 \doteqdot-\log \theta_{\max }=0.353 \tag{149}
\end{equation*}
$$

Checking Example 4 in this way, we can see that $b_{\max }=$ $(-0.431,-0.431,0.862)$ is a left eigenvector for $\theta_{\max }=$ 0.855 , but ${ }^{t} \boldsymbol{b}_{\text {max }}$ is not a right eigenvector. Thus, by Lemma 10, we have $\lim _{N \rightarrow \infty} L(N) \doteqdot-\log \theta_{\max }$ and (139).

## B. Convergence of order $O(1 / N)$

In the case of the convergence of order $O(1 / N)$, there exist type-II indices.

Example 6: Consider the channel matrix $\Phi^{(2)}$ of (29), i.e.,

$$
\Phi^{(2)}=\left(\begin{array}{lll}
0.800 & 0.100 & 0.100  \tag{150}\\
0.100 & 0.800 & 0.100 \\
0.300 & 0.300 & 0.400
\end{array}\right)
$$

and an initial distribution $\boldsymbol{\lambda}^{0}=(1 / 3,1 / 3,1 / 3)$. We have

$$
\begin{align*}
C & =0.310[\mathrm{nat} / \mathrm{symbol}]  \tag{151}\\
\lambda^{*} & =(0.500,0.500,0.000),  \tag{152}\\
Q^{*} & =(0.450,0.450,0.100),  \tag{153}\\
D_{1,1}^{*} & =-1.544, D_{1,2}^{*}=-0.456, D_{1,3}^{*}=-1,  \tag{154}\\
D_{2,2}^{*} & =-1.544, D_{2,3}^{*}=-1, D_{3,3}^{*}=-2  \tag{155}\\
J\left(\boldsymbol{\lambda}^{*}\right) & =\left(\begin{array}{ccc}
1+\lambda_{1}^{*} D_{1,1}^{*} & \lambda_{2}^{*} D_{1,2}^{*} & 0 \\
\lambda_{1}^{*} D_{2,1}^{*} & 1+\lambda_{2}^{*} D_{2,2}^{*} & 0 \\
\lambda_{1}^{*} D_{3,1}^{*} & \lambda_{2}^{*} D_{3,2}^{*} & 1
\end{array}\right)  \tag{156}\\
& =\left(\begin{array}{ccc}
0.228 & -0.228 & 0.000 \\
-0.228 & 0.228 & 0.000 \\
-0.500 & -0.500 & 1.000
\end{array}\right) \tag{157}
\end{align*}
$$

By (152) and (153), we see that $\lambda_{3}^{*}=0$ and $C=$ $D\left(P^{1} \| Q^{*}\right)=D\left(P^{2} \| Q^{*}\right)=D\left(P^{3} \| Q^{*}\right)$, hence, $i=3$ is a type-II index.

The eigenvalues of $J\left(\boldsymbol{\lambda}^{*}\right)$ are $\left(\theta_{1}, \theta_{2}, \theta_{3}\right)=(0.000,0.456$, 1.000).

$$
\begin{align*}
& A=\left(\begin{array}{ccc}
1 & 1 & 0 \\
1 & -1 & 0 \\
1 & 0 & 1
\end{array}\right)  \tag{158}\\
& A_{1}=\left(\begin{array}{cc}
1 & 1 \\
1 & -1
\end{array}\right), A_{2}=\left(\begin{array}{ll}
1 & 0
\end{array}\right) \tag{159}
\end{align*}
$$

By (159), the eigenvectors $\boldsymbol{a}_{1}$ for $\theta_{1}=0$ and $\boldsymbol{a}_{2}$ for $\theta_{2}=$ 0.456 are

$$
\boldsymbol{a}_{1}=\left(\begin{array}{l}
1  \tag{160}\\
1 \\
1
\end{array}\right), \boldsymbol{a}_{2}=\left(\begin{array}{c}
1 \\
-1 \\
0
\end{array}\right)
$$

We will prove (97). First, for $\boldsymbol{a}_{1}, \boldsymbol{\mu}^{N+1} \boldsymbol{a}_{1}=\boldsymbol{\mu}^{N} \boldsymbol{a}_{1}=\mathbf{0}$, thus (97) is trivial. Next, for $\boldsymbol{a}_{2}$, we will prove that $\boldsymbol{\mu}^{N+1} \boldsymbol{a}_{2}=$
$\mu_{1}^{N+1}-\mu_{2}^{N+1}$ is divisible by $\boldsymbol{\mu}^{N} \boldsymbol{a}_{2}=\mu_{1}^{N}-\mu_{2}^{N}$. In fact, by $\lambda_{1}^{*}=\lambda_{2}^{*}$ we have

$$
\begin{align*}
\mu_{1}^{N+1}-\mu_{2}^{N+1} & =\lambda_{1}^{N+1}-\lambda_{2}^{N+1}  \tag{161}\\
= & F_{1}\left(\boldsymbol{\lambda}^{N}\right)-F_{2}\left(\boldsymbol{\lambda}^{N}\right)  \tag{162}\\
= & \frac{\lambda_{1}^{N} \exp D\left(P^{1} \| \boldsymbol{\lambda}^{N} \Phi\right)-\lambda_{2}^{N} \exp D\left(P^{2} \| \boldsymbol{\lambda}^{N} \Phi\right)}{\sum_{k=1}^{3} \lambda_{k}^{N} \exp D\left(P^{k} \| \boldsymbol{\lambda}^{N} \Phi\right)} \tag{163}
\end{align*}
$$

If we substitute $\lambda_{1}^{N}=\lambda_{2}^{N}$ in (163), we have by calculation, $\mu_{1}^{N+1}-\mu_{2}^{N+1}=0$. Because $\lambda_{1}^{N}=\lambda_{2}^{N}$ is equivalent to $\mu_{1}^{N}=$ $\mu_{2}^{N}$, we see that $\mu_{1}^{N}-\mu_{2}^{N}=0$ implies $\mu_{1}^{N+1}-\mu_{2}^{N+1}=0$, thus (97) holds, and then we can consider (99).

For $\overline{\boldsymbol{\mu}}_{\mathrm{I}, \mathrm{III}}^{N}=\left(\bar{\mu}_{1}^{N}, \bar{\mu}_{2}^{N}\right), \overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N}=\left(\bar{\mu}_{3}^{N}\right)$, we have $\overline{\boldsymbol{\mu}}_{\mathrm{I}, \mathrm{III}}^{N}=$ - $\bar{\mu}_{\text {II }}^{N} A_{2} A_{1}^{-1}$ by (99), hence

$$
\begin{equation*}
\bar{\mu}_{1}^{N}=\bar{\mu}_{2}^{N}=-(1 / 2) \bar{\mu}_{3}^{N} . \tag{164}
\end{equation*}
$$

Further, the Hessian matrix $H_{3}\left(\boldsymbol{\lambda}^{*}\right)$ is

$$
\begin{align*}
H_{3}\left(\boldsymbol{\lambda}^{*}\right) & =\left(\begin{array}{ccc}
0 & 0 & D_{3,1}^{*} \\
0 & 0 & D_{3,2}^{*} \\
D_{3,1}^{*} & D_{3,2}^{*} & 2 D_{3,3}^{*}
\end{array}\right)  \tag{165}\\
& =\left(\begin{array}{ccc}
0.000 & 0.000 & -1.000 \\
0.000 & 0.000 & -1.000 \\
-1.000 & -1.000 & -4.000
\end{array}\right), \tag{166}
\end{align*}
$$

then, we have by (164),

$$
\begin{equation*}
\frac{1}{2} \overline{\boldsymbol{\mu}}^{N} H_{3}^{t} \overline{\boldsymbol{\mu}}^{N}=-\left(\bar{\mu}_{3}^{N}\right)^{2} \tag{167}
\end{equation*}
$$

and the second-order recurrence formula

$$
\begin{equation*}
\bar{\mu}_{3}^{N+1}=\bar{\mu}_{3}^{N}-\left(\bar{\mu}_{3}^{N}\right)^{2} \tag{168}
\end{equation*}
$$

By Lemma 6 and (164), we have $\lim _{N \rightarrow \infty} N \bar{\mu}_{3}^{N}=$ 1, $\lim _{N \rightarrow \infty} N \bar{\mu}_{1}^{N}=\lim _{N \rightarrow \infty} N \bar{\mu}_{2}^{N}=-1 / 2$.

By the numerical simulation, $N \boldsymbol{\mu}^{N}$ for $N=500$ is

$$
\begin{align*}
& N \boldsymbol{\mu}^{N}=(-0.510,-0.510,1.019)  \tag{169}\\
& \doteqdot \lim _{N \rightarrow \infty} N \overline{\boldsymbol{\mu}}^{N}=(-1 / 2,-1 / 2,1) \tag{170}
\end{align*}
$$

See Fig. 7. We can confirm that $N \boldsymbol{\mu}^{N}$ for large $N$ is near the values obtained in Theorem 8.

Example 7: Consider a channel matrix

$$
\Phi^{(5)} \equiv\left(\begin{array}{ccccc}
0.6 & 0.1 & 0.1 & 0.1 & 0.1  \tag{171}\\
0.1 & 0.6 & 0.1 & 0.1 & 0.1 \\
s & s & t & 0.1 & 0.1 \\
s & s & 0.1 & t & 0.1 \\
s & s & 0.1 & 0.1 & t
\end{array}\right)
$$

where $s \equiv 0.238, t \equiv 0.324,(2 s+t+0.2=1)$,


Fig. 7. Convergence of $N \mu_{i}^{N}$ in Example 6.
and an initial distribution $\boldsymbol{\lambda}^{0}=(1 / 5,1 / 5,1 / 5,1 / 5,1 / 5)$. For this $\Phi^{(5)}$, we have

$$
\begin{align*}
& C=0.198[\mathrm{nat} / \mathrm{symbol}]  \tag{173}\\
& \boldsymbol{\lambda}^{*}=(0.5,0.5,0,0,0)  \tag{174}\\
& Q^{*}=(0.35,0.35,0.1,0.1,0.1),  \tag{175}\\
& D_{1,1}^{*}=-19 / 14, D_{1,2}^{*}=-9 / 14, D_{1,3}^{*}=-1 \\
& D_{1,4}^{*}=-1, D_{1,5}^{*}=-1  \tag{176}\\
& D_{2,2}^{*}=-19 / 14, D_{2,3}^{*}=-1, D_{2,4}^{*}=-1 \\
& D_{2,5}^{*}=-1  \tag{177}\\
& D_{3,3}^{*}=-1.576 \equiv-\alpha, D_{3,4}^{*}=-1.072 \equiv-\beta \\
& D_{3,5}^{*}=-\beta  \tag{178}\\
& D_{4,4}^{*}=-\alpha, D_{4,5}^{*}=-\beta, D_{5,5}^{*}=-\alpha  \tag{179}\\
& J\left(\boldsymbol{\lambda}^{*}\right)=\left(\begin{array}{ccccc}
9 / 28 & -9 / 28 & 0 & 0 & 0 \\
-9 / 28 & 9 / 28 & 0 & 0 & 0 \\
-1 / 2 & -1 / 2 & 1 & 0 & 0 \\
-1 / 2 & -1 / 2 & 0 & 1 & 0 \\
-1 / 2 & -1 / 2 & 0 & 0 & 1
\end{array}\right) \tag{180}
\end{align*}
$$

By (174) and (175), we see that $\lambda_{3}^{*}=\lambda_{4}^{*}=\lambda_{5}^{*}=0$ and $C=D\left(P^{1} \| Q^{*}\right)=\ldots=D\left(P^{5} \| Q^{*}\right)$, thus $i=3,4,5$ are type-II indices.

The eigenvalues of $J\left(\boldsymbol{\lambda}^{*}\right)$ are $\left(\theta_{1}, \theta_{2}, \theta_{3}, \theta_{4}, \theta_{5}\right)=$ $(0,9 / 14,1,1,1)$ and

$$
\begin{align*}
& A=\left(\begin{array}{ccccc}
1 & 1 & 0 & 0 & 0 \\
1 & -1 & 0 & 0 & 0 \\
1 & 0 & 1 & 0 & 0 \\
1 & 0 & 0 & 1 & 0 \\
1 & 0 & 0 & 0 & 1
\end{array}\right),  \tag{181}\\
& A_{1}=\left(\begin{array}{cc}
1 & 1 \\
1 & -1
\end{array}\right), A_{2}=\left(\begin{array}{ll}
1 & 0 \\
1 & 0 \\
1 & 0
\end{array}\right) . \tag{182}
\end{align*}
$$

We can prove (97) in a similar way as in Example 6.

For $\overline{\boldsymbol{\mu}}_{\mathrm{I}, \mathrm{III}}^{N}=\left(\bar{\mu}_{1}^{N}, \bar{\mu}_{2}^{N}\right), \overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N}=\left(\bar{\mu}_{3}^{N}, \bar{\mu}_{4}^{N}, \bar{\mu}_{5}^{N}\right)$, we have $\overline{\boldsymbol{\mu}}_{\mathrm{I}, \mathrm{III}}^{N}=-\overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N} A_{2} A_{1}^{-1}$ by (99), hence

$$
\begin{equation*}
\bar{\mu}_{1}^{N}=\bar{\mu}_{2}^{N}=-\left(\bar{\mu}_{3}^{N}+\bar{\mu}_{4}^{N}+\bar{\mu}_{5}^{N}\right) / 2 . \tag{183}
\end{equation*}
$$

Further, the Hessian matrix $H_{3}\left(\boldsymbol{\lambda}^{*}\right)$ is

$$
H_{3}\left(\boldsymbol{\lambda}^{*}\right)=\left(\begin{array}{ccccc}
0 & 0 & -1 & 0 & 0  \tag{184}\\
0 & 0 & -1 & 0 & 0 \\
-1 & -1 & -2 \alpha & -\beta & -\beta \\
0 & 0 & -\beta & 0 & 0 \\
0 & 0 & -\beta & 0 & 0
\end{array}\right)
$$

where $\alpha$ and $\beta$ were defined in (178). Then, we have by (183)

$$
\begin{align*}
\frac{1}{2} \overline{\boldsymbol{\mu}}^{N} H_{3}\left(\boldsymbol{\lambda}^{*}\right)^{t} \overline{\boldsymbol{\mu}}^{N}= & -(\alpha-1)\left(\bar{\mu}_{3}^{N}\right)^{2}-(\beta-1) \bar{\mu}_{3}^{N} \bar{\mu}_{4}^{N} \\
& -(\beta-1) \bar{\mu}_{3}^{N} \bar{\mu}_{5}^{N} \tag{185}
\end{align*}
$$

Similarly,

$$
\begin{align*}
\frac{1}{2} \overline{\boldsymbol{\mu}}^{N} H_{4}\left(\boldsymbol{\lambda}^{*}\right)^{t} \overline{\boldsymbol{\mu}}^{N}= & -(\beta-1) \bar{\mu}_{3}^{N} \bar{\mu}_{4}^{N}-(\alpha-1)\left(\bar{\mu}_{4}^{N}\right)^{2} \\
& -(\beta-1) \bar{\mu}_{4}^{N} \bar{\mu}_{5}^{N},  \tag{186}\\
\frac{1}{2} \overline{\boldsymbol{\mu}}^{N} H_{5}\left(\boldsymbol{\lambda}^{*}\right)^{t} \overline{\boldsymbol{\mu}}^{N}= & -(\beta-1) \bar{\mu}_{3}^{N} \bar{\mu}_{5}^{N}-(\beta-1) \bar{\mu}_{4}^{N} \bar{\mu}_{5}^{N} \\
& -(\alpha-1)\left(\bar{\mu}_{5}^{N}\right)^{2} . \tag{187}
\end{align*}
$$

Therefore, by substituting $\alpha^{\prime} \equiv \alpha-1, \beta^{\prime} \equiv \beta-1$, we have

$$
\begin{align*}
& \bar{\mu}_{3}^{N+1}=\bar{\mu}_{3}^{N}-\alpha^{\prime}\left(\bar{\mu}_{3}^{N}\right)^{2}-\beta^{\prime} \bar{\mu}_{3}^{N} \bar{\mu}_{4}^{N}-\beta^{\prime} \bar{\mu}_{3}^{N} \bar{\mu}_{5}^{N},  \tag{188}\\
& \bar{\mu}_{4}^{N+1}=\bar{\mu}_{4}^{N}-\beta^{\prime} \bar{\mu}_{3}^{N} \bar{\mu}_{4}^{N}-\alpha^{\prime}\left(\bar{\mu}_{4}^{N}\right)^{2}-\beta^{\prime} \bar{\mu}_{4}^{N} \bar{\mu}_{5}^{N},  \tag{189}\\
& \bar{\mu}_{5}^{N+1}=\bar{\mu}_{5}^{N}-\beta^{\prime} \bar{\mu}_{3}^{N} \bar{\mu}_{5}^{N}-\beta^{\prime} \bar{\mu}_{4}^{N} \bar{\mu}_{5}^{N}-\alpha^{\prime}\left(\bar{\mu}_{5}^{N}\right)^{2} . \tag{190}
\end{align*}
$$

From (105), we have $\boldsymbol{\sigma}=\left(\sigma_{3}, \sigma_{4}, \sigma_{5}\right)=(1.389,1.389$, 1.389), so the canonical form for (188), (189), (190) is

$$
\begin{align*}
& \nu_{3}^{N+1}=\nu_{3}^{N}-0.8\left(\nu_{3}^{N}\right)^{2}-0.1 \nu_{3}^{N} \mu_{4}^{N}-0.1 \nu_{3}^{N} \nu_{5}^{N}  \tag{191}\\
& \nu_{4}^{N+1}=\nu_{4}^{N}-0.1 \nu_{3}^{N} \nu_{4}^{N}-0.8\left(\nu_{4}^{N}\right)^{2}-0.1 \nu_{4}^{N} \nu_{5}^{N}  \tag{192}\\
& \nu_{5}^{N+1}=\nu_{5}^{N}-0.1 \nu_{3}^{N} \nu_{5}^{N}-0.1 \nu_{4}^{N} \nu_{5}^{N}-0.8\left(\nu_{5}^{N}\right)^{2} \tag{193}
\end{align*}
$$

Eqs. (191), (192), (193) satisfy the assumptions of Lemma 8 and the diagonally dominant condition (122). Then, by Theorem 11, they converge for arbitrary initial values and

$$
\begin{equation*}
\lim _{N \rightarrow \infty} N \nu_{i}^{N}=1, i=3,4,5 . \tag{194}
\end{equation*}
$$

Therefore, by (106)

$$
\begin{equation*}
\lim _{N \rightarrow \infty} N \bar{\mu}_{i}^{N}=\sigma_{i}=1.389, i=3,4,5 \tag{195}
\end{equation*}
$$

and by (183),

$$
\begin{equation*}
\lim _{N \rightarrow \infty} N \bar{\mu}_{i}^{N}=-3 \sigma_{3} / 2=-2.083, i=1,2 \tag{196}
\end{equation*}
$$

We will show in Fig. 8 the comparison of the numerical results and the values of (195), (196).


Fig. 8. Convergence of $N \mu_{i}^{N}$ in Example 7.

## C. Exponential convergence where indices of types-I and -III

 existExample 8: Consider the channel matrix $\Phi^{(3)}$ of (30)

$$
\Phi^{(3)}=\left(\begin{array}{lll}
0.800 & 0.100 & 0.100  \tag{197}\\
0.100 & 0.800 & 0.100 \\
0.350 & 0.350 & 0.300
\end{array}\right)
$$

and an initial distribution $\boldsymbol{\lambda}^{0}=(1 / 3,1 / 3,1 / 3)$. We have

$$
\begin{align*}
C & =0.310[\text { nat/symbol }],  \tag{198}\\
\boldsymbol{\lambda}^{*} & =(0.500,0.500,0.000),  \tag{199}\\
Q^{*} & =(0.450,0.450,0.100),  \tag{200}\\
J\left(\boldsymbol{\lambda}^{*}\right) & =\left(\begin{array}{ccc}
0.228 & -0.228 & 0.000 \\
-0.228 & 0.228 & 0.000 \\
-0.428 & -0.428 & 0.856
\end{array}\right) . \tag{201}
\end{align*}
$$

By (199) and (200), we see that $\lambda_{1}^{*}>0, \lambda_{2}^{*}>0, \lambda_{3}^{*}=0$ and $C=D\left(P^{1} \| Q^{*}\right)=D\left(P^{2} \| Q^{*}\right)>D\left(P^{3} \| Q^{*}\right)$, thus $i=1,2$ are type-I indices, and $i=3$ is a type-III index.

The eigenvalues of $J\left(\boldsymbol{\lambda}^{*}\right)$ are $\left(\theta_{1}, \theta_{2}, \theta_{3}\right)=(0.000,0.456$, $0.856)$. Then, $\theta_{\max }=\theta_{3}=0.856$. We have for $N=500$

$$
\begin{equation*}
L(500)=0.159 \doteqdot-\log \theta_{\max }=0.155 \tag{202}
\end{equation*}
$$

See Fig. 9.
Extending this result, we have the following lemma.
Lemma 11: Assume that type-II indices do not exist and $\theta_{\max }=\max _{i \in \mathcal{I}_{\text {III }}} \theta_{i}$, i.e., the maximum eigenvalue of $J\left(\boldsymbol{\lambda}^{*}\right)$ is achieved in $J^{I I I}$. Then, the convergence speed does not depend on the choice of initial distribution. In other words, $\lim _{N \rightarrow \infty} L(N)=-\log \theta_{\max }$ holds for arbitrary initial distribution, hence the convergence speed cannot be increased any more.

Proof: Let $\theta_{\max }=\theta_{i^{*}}, i^{*} \in \mathcal{I}_{\text {III }}$. We have $\theta_{i^{*}}>0$ by Theorem 4. $J^{\text {III }}$ is diagonal by (56), thus we can take ${ }^{t} \boldsymbol{e}_{i^{*}}={ }^{t}\left(0, \ldots, 0, \stackrel{i^{*}}{\vee}\right.$ th $\left., 0, \ldots, 0\right)$ as a right eigenvector for


Fig. 9. Convergence of $L(N)$ in Example 8 with initial distribution $\boldsymbol{\lambda}^{0}=$ ( $1 / 3,1 / 3,1 / 3$ ).
$\theta_{i^{*}}$. However, $\boldsymbol{e}_{i^{*}}$ is not a left eigenvector for $\theta_{i^{*}}$. In fact, since every row sum of $J\left(\boldsymbol{\lambda}^{*}\right)$ is 0 by Lemma 3, putting $\mathbf{1} \equiv$ $(1, \ldots, 1) \in \mathbb{R}^{m}$, we have $J\left(\boldsymbol{\lambda}^{*}\right)^{t} \mathbf{1}=\mathbf{0}$. If $\boldsymbol{e}_{i^{*}}$ were a left eigenvector for $\theta_{i^{*}}$, then $0=\boldsymbol{e}_{i^{*}} J\left(\boldsymbol{\lambda}^{*}\right)^{t} \mathbf{1}=\theta_{i^{*}} \boldsymbol{e}_{i^{*}}{ }^{t} \mathbf{1}=\theta_{i^{*}}>$ 0 , which is a contradiction. Therefore, by Lemmas 9 and 10 , the convergence speed does not depend on the choice of initial distribution and $\lim _{N \rightarrow \infty} L(N)=-\log \theta_{\text {max }}$.

## VII. Convergence speed of $I\left(\boldsymbol{\lambda}^{N}, \Phi\right) \rightarrow C$

Based on the results obtained so far, we will consider the convergence speed that the mutual information $I\left(\boldsymbol{\lambda}^{N}, \Phi\right)$ tends to $C$ as $N \rightarrow \infty$. We will show that if $\mathcal{I}_{\text {III }}=\emptyset$ and $\boldsymbol{\lambda}^{N} \rightarrow \boldsymbol{\lambda}^{*}$ is the convergence of order $O(1 / N)$, then $I\left(\boldsymbol{\lambda}^{N}, \Phi\right) \rightarrow C$ is of order $O\left(1 / N^{2}\right)$. Including this fact, we have the following theorem.

Theorem 12: Let $\mathcal{I}_{\mathrm{III}}=\emptyset$.
If $\left\|\boldsymbol{\mu}^{N}\right\|<K_{1} \cdot(\theta)^{N}, 0 \leq \theta_{\max }<\theta<1, K_{1}>0, N=$ $0,1, \ldots$, then

$$
\begin{gather*}
0<C-I\left(\boldsymbol{\lambda}^{N}, \Phi\right)<K_{2} \cdot(\theta)^{2 N}  \tag{203}\\
K_{2}>0, N=0,1, \ldots
\end{gather*}
$$

Meanwhile, if $\lim _{N \rightarrow \infty} N \mu_{i}^{N}=\sigma_{i} \neq 0, i=1, \ldots, m$, then

$$
\begin{equation*}
\lim _{N \rightarrow \infty} N^{2}\left(C-I\left(\boldsymbol{\lambda}^{N}, \Phi\right)\right)=\frac{1}{2} \sum_{j=1}^{n} \frac{1}{Q_{j}^{*}}\left(\sum_{i=1}^{m} \sigma_{i} P_{j}^{i}\right)^{2} \tag{204}
\end{equation*}
$$

Next, let $\mathcal{I}_{\text {III }} \neq \emptyset$.

$$
\text { If }\left\|\boldsymbol{\mu}^{N}\right\|<K_{1} \cdot(\theta)^{N}, 0 \leq \theta_{\max }<\theta<1, K_{1}>0, N=
$$ $0,1, \ldots$, then

$$
\begin{gather*}
0<C-I\left(\boldsymbol{\lambda}^{N}, \Phi\right)<K_{2} \cdot(\theta)^{N}  \tag{205}\\
K_{2}>0, N=0,1, \ldots
\end{gather*}
$$

Meanwhile, if $\lim _{N \rightarrow \infty} N \mu_{i}^{N}=\sigma_{i} \neq 0, i=1, \ldots, m$, then

$$
\begin{gather*}
0<C-I\left(\boldsymbol{\lambda}^{N}, \Phi\right)<K / N  \tag{206}\\
K>0, N=0,1, \ldots
\end{gather*}
$$

## Proof: See Appendix H.

We will show in the following tables the evaluation of the convergence speed of $I\left(\boldsymbol{\lambda}^{N}, \Phi^{(k)}\right) \rightarrow C$ for $k=1,2,3$, where $\Phi^{(1)}, \Phi^{(2)}$, and $\Phi^{(3)}$ were defined in Examples 1, 2, and 3, and examined in Examples 4, 6, and 8, respectively.

TABLE I
CONVERGENCE SPEED OF $I\left(\boldsymbol{\lambda}^{N}, \Phi^{(1)}\right) \rightarrow C$.

| $-\left.(1 / N) \log \left(C-I\left(\boldsymbol{\lambda}^{N}, \Phi^{(1)}\right)\right)\right\|_{N=500}$ | 0.324 |
| :--- | :--- |
| $-2 \log \theta_{\max }$ | 0.313 |

TABLE II
Convergence speed of $I\left(\boldsymbol{\lambda}^{N}, \Phi^{(2)}\right) \rightarrow C$.

| $\left.N^{2}\left(C-I\left(\boldsymbol{\lambda}^{N}, \Phi^{(2)}\right)\right)\right\|_{N=500}$ | 0.516 |
| :--- | :---: |
| Eq. (204) | 0.500 |

TABLE III
Convergence speed of $I\left(\boldsymbol{\lambda}^{N}, \Phi^{(3)}\right) \rightarrow C$.

| $-\left.(1 / N) \log \left(C-I\left(\boldsymbol{\lambda}^{N}, \Phi^{(3)}\right)\right)\right\|_{N=500}$ | 0.163 |
| :--- | :--- |
| $-\log \theta_{\max }$ | 0.155 |

We can see from these tables that the convergence speed of $I\left(\boldsymbol{\lambda}^{N}, \Phi\right) \rightarrow C$ is accurately approximated by the results of Theorem 12.

## VIII. Conclusion

In this paper, we investigated the convergence speed of the Arimoto-Blahut algorithm. We showed that the capacityachieving input distribution $\boldsymbol{\lambda}^{*}$ is the fixed point of $F(\boldsymbol{\lambda})$ and analyzed the convergence speed by the Taylor expansion of $F(\boldsymbol{\lambda})$ about $\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}$. We concretely calculated the Jacobian matrix $J$ of the first-order term of the Taylor expansion and the Hessian matrix $H$ of the second-order term. The analysis of the convergence speed by the Hessian matrix $H$ was done for the first time in this paper.

We showed that if type-II indices do not exist, then the convergence of $\boldsymbol{\lambda}^{N} \rightarrow \boldsymbol{\lambda}^{*}$ is exponential, and if type-II indices exist, then the convergence of the second-order recurrence formula obtained by truncating the Taylor expansion is of order $O(1 / N)$ for some initial vector. Further, we considered the condition for the convergence of order $O(1 / N)$ for an arbitrary initial vector. Next, we considered the convergence speed of $I\left(\boldsymbol{\lambda}^{N}, \Phi\right) \rightarrow C$ and showed that the type-III indices concern the convergence speed. In particular, if there exist
no type-III indices and $\boldsymbol{\lambda}^{N} \rightarrow \boldsymbol{\lambda}^{*}$ is of order $O(1 / N)$, then $I\left(\boldsymbol{\lambda}^{N}, \Phi\right) \rightarrow C$ is $O\left(1 / N^{2}\right)$.

Based on these analyses, the convergence speeds for several channel matrices were numerically evaluated. As a result, it was confirmed that the convergence speed of the ArimotoBlahut algorithm is accurately approximated by the values obtained by our theorems.
In this paper and in [2], [21], the capacity-achieving $\boldsymbol{\lambda}^{*}$ is used to analyze the convergence speed of the sequence $\left\{\boldsymbol{\lambda}^{N}\right\}$. Using the limit vector is a common method for the analysis of the convergence speed. For example, see [9], Proposition 4.4. However, $\boldsymbol{\lambda}^{*}$ is not known before it is calculated by the Arimoto-Blahut algorithm, etc. Determining whether the convergence speed is fast or slow by a simple calculation from the channel matrix $\Phi$ without using $\boldsymbol{\lambda}^{*}$ would be effective, so we would like to consider finding this condition as a future work.

## REFERENCES

[1] L.V. Ahlfors, Complex Analysis (third edition), McGraw-Hill, 1979.
[2] S. Arimoto, "An algorithm for computing the capacity of arbitrary discrete memoryless channels", IEEE Trans. Inf. Theory, vol. 18, pp. 14-20, Jan. 1972.
[3] S. Arimoto, Information Theory, Kyoritsu Shuppan, Tokyo, 1976. (in Japanese)
[4] R. E. Blahut, "Computation of channel capacity and rate-distortion functions", IEEE Trans. Inf. Theory, vol. 18, no. 4, pp. 460-473, Jul. 1972.
[5] T. J. I'A. Bromwich, An Introduction to the Theory of Infinite Series, third edition, Chelsea Publishing Company, New York, 1991.
[6] T. Cover and J. Thomas, Elements of Information Theory, Wiley, June 2006.
[7] I. Csiszár and J. Körner, Information Theory: Coding Theorems for Discrete Memoryless Systems, Academic Press, Orlando, 1982.
[8] I. Csiszár and G. Tusnády, "Information geometry and alternating minimization procedures", Statistics and Decisions, Supplement Issue No.1, 205-237, 1984.
[9] R. L. Devaney, An Introduction to Chaotic Dynamical Systems, AddisonWesley, 1987,
[10] R. A. Horn and C. R. Johnson, Matrix Analysis, Cambridge University Press, Second Ed., 1985.
[11] G. Matz and P. Duhamel, "Information geometric formulation and interpretation of accelerated Blahut-Arimoto-Type algorithms", in proceedings of ITW2004, 2004.
[12] I. Naiss and H. Permuter, "Extension of the Blahut-Arimoto algorithm for maximizing directed information", IEEE Trans. Inf. Theory, vol. 59, no. 1, pp. 204-222, Jan. 2013.
[13] Z. Naja, F. Alberge, and P. Duhamel, "Geometrical interpretation and improvements of the Blahut-Arimoto's algorithm", in proceedings of IEEE International Conference on Acoustics, Speech and Signal Processing (ICASSP), 2009.
[14] K. Nakagawa and F. Kanaya, "A new geometric capacity characterization of a discrete memoryless channel", IEEE Trans. Inf. Theory, vol. 34, no. 2, pp. 318-321, 1988.
[15] K. Nakagawa, K. Watabe and T. Sabu, "On the search algorithm for the output distribution that achieves the channel capacity", IEEE Trans. Inf. Theory, vol. 63, no. 2, pp. 1043-1062, Feb. 2017.
[16] K. Nakagawa, Y. Takei and K. Watabe, "Analysis for the slow convergence in Arimoto algorithm", in proceedings of ISITA2018, pp. 764-768, Oct. 2018.
[17] M. Rezaeian and A. Grant, "A generalization of Arimoto-Blahut algorithm", in proceedings of IEEE ISIT 2004.
[18] C. Robinson, Dynamical Systems: Stability, Symbolic Dynamics, and Chaos, CRC Press, 1998.
[19] I. Satake, Linear Algebra, Marcel Dekker, 1975.
[20] P. O. Vontobel, A. Kav̌cić, D. M. Arnold, H.-A. Loeliger, "A generalization of the Blahut-Arimoto algorithm to finite-state channels", IEEE Trans. Inf. Theory, vol. 54, no. 5, pp. 1887-1918, 2008.
[21] Y. Yu, "Squeezing the Arimoto-Blahut algorithm for faster convergence", IEEE Trans. Inf. Theory, vol. 56, no. 7, pp. 3149-3157, Jul. 2010.

## Appendix A

## Proof of Lemma 5

Since $\sqrt{\Lambda} B \sqrt{\Lambda}^{-1}$ is symmetric by (69), it is diagonalizable, hence $J^{\mathrm{I}}$ is diagonalizable because $J^{\mathrm{I}}=I-B$. Therefore, there exists a regular matrix $\Pi \in \mathbb{R}^{m_{1} \times m_{1}}$ with

$$
\begin{equation*}
\Pi^{-1} J^{\mathrm{I}} \Pi=\Theta_{1} \tag{207}
\end{equation*}
$$

where $\Theta_{1}$ is a diagonal matrix whose diagonal components are the eigenvalues of $J^{\mathrm{I}}$, i.e., $\Theta_{1}=\operatorname{diag}\left(\theta_{i}, i \in \mathcal{I}_{\mathrm{I}}\right) \in \mathbb{R}^{m_{1} \times m_{1}}$. We can write it by components as $\Theta_{1}=\left(\theta_{i} \delta_{i^{\prime} i}\right), i^{\prime}, i \in \mathcal{I}_{\mathrm{I}}$.

Next, let $\Theta_{2}$ be a diagonal matrix whose diagonal components are the eigenvalues of $J^{\mathrm{II}}$ and $J^{\mathrm{III}}$, i.e., $\Theta_{2}=$ $\operatorname{diag}\left(\theta_{i}, i \in \mathcal{I}_{\text {II }} \cup \mathcal{I}_{\text {III }}\right) \in \mathbb{R}^{\left(m_{2}+m_{3}\right) \times\left(m_{2}+m_{3}\right)}$. We have $\Theta_{2}=\left(\theta_{i} \delta_{i^{\prime} i}\right), i^{\prime}, i \in \mathcal{I}_{\text {II }} \cup \mathcal{I}_{\text {III }}$. Then, by (53), we have

$$
J\left(\boldsymbol{\lambda}^{*}\right)=\left(\begin{array}{cc}
J^{\mathrm{I}} & O  \tag{208}\\
U & \Theta_{2}
\end{array}\right)
$$

where $U \in \mathbb{R}^{\left(m_{2}+m_{3}\right) \times m_{1}}$ is an appropriate matrix.
Now, we will prove that there exists a unique matrix $V \in$ $\mathbb{R}^{\left(m_{2}+m_{3}\right) \times m_{1}}$ that satisfies

$$
\begin{equation*}
V \Theta_{1}-\Theta_{2} V=U \Pi \tag{209}
\end{equation*}
$$

Define the components of $V$ by $V \equiv\left(v_{i^{\prime} i}\right), i^{\prime} \in \mathcal{I}_{\mathrm{II}} \cup \mathcal{I}_{\mathrm{III}}, i \in$ $\mathcal{I}_{\text {I }}$. Then,

$$
\begin{align*}
& \left(V \Theta_{1}\right)_{i^{\prime} i}=\sum_{k \in \mathcal{I}_{\mathrm{I}}} v_{i^{\prime} k} \theta_{i} \delta_{k i}=v_{i^{\prime} i} \theta_{i},  \tag{210}\\
& i^{\prime} \in \mathcal{I}_{\mathrm{II}} \cup \mathcal{I}_{\mathrm{III}}, i \in \mathcal{I}_{\mathrm{I}}, \\
& \left(\Theta_{2} V\right)_{i^{\prime} i}=\sum_{k \in \mathcal{I}_{\mathrm{II}} \cup \mathcal{I}_{\mathrm{III}}} \theta_{i^{\prime}} \delta_{k i^{\prime}} v_{k i}=v_{i^{\prime} i} \theta_{i^{\prime}},  \tag{211}\\
& i^{\prime} \in \mathcal{I}_{\mathrm{II}} \cup \mathcal{I}_{\mathrm{III}}, i \in \mathcal{I}_{\mathrm{I}} .
\end{align*}
$$

Further, defining the components of $U \Pi$ by $U \Pi=\left(u_{i^{\prime} i}^{\Pi}\right), i^{\prime} \in$ $\mathcal{I}_{\mathrm{II}} \cup \mathcal{I}_{\mathrm{III}}, i \in \mathcal{I}_{\mathrm{I}}$, both sides of (209) are represented by components as

$$
\begin{equation*}
\left(\theta_{i}-\theta_{i^{\prime}}\right) v_{i^{\prime} i}=u_{i^{\prime} i}^{\Pi}, i^{\prime} \in \mathcal{I}_{\mathrm{II}} \cup \mathcal{I}_{\mathrm{III}}, i \in \mathcal{I}_{\mathrm{I}} \tag{212}
\end{equation*}
$$

By Theorem 2, the eigenvalues of $J^{\mathrm{I}}$ are less than 1 , hence different from the eigenvalues 1 of $J^{\mathrm{II}}$. Further, $\theta_{i} \neq \theta_{i^{\prime}}, i \in$ $\mathcal{I}_{\mathrm{I}}, i^{\prime} \in \mathcal{I}_{\text {III }}$ by the assumption (75), so we have

$$
\begin{equation*}
v_{i^{\prime} i}=\frac{u_{i^{\prime} i}^{\Pi}}{\theta_{i}-\theta_{i^{\prime}}}, i^{\prime} \in \mathcal{I}_{\mathrm{II}} \cup \mathcal{I}_{\mathrm{III}}, i \in \mathcal{I}_{\mathrm{I}} \tag{213}
\end{equation*}
$$

which shows the existence and uniqueness of $V \in$ $\mathbb{R}^{\left(m_{2}+m_{3}\right) \times m_{1}}$ that satisfies (209).

Now, define

$$
\tilde{\Pi} \equiv\left(\begin{array}{cc}
\Pi & O  \tag{214}\\
V & I
\end{array}\right) \in \mathbb{R}^{m \times m}
$$

Then, by noting (207), (208), (209), we have

$$
\tilde{\Pi}^{-1} J\left(\boldsymbol{\lambda}^{*}\right) \tilde{\Pi}=\left(\begin{array}{cc}
\Theta_{1} & O  \tag{215}\\
O & \Theta_{2}
\end{array}\right)
$$

which proves the lemma.

## Appendix B <br> Proof of Theorem 5

Consider the line segment with the starting point $\boldsymbol{\lambda}^{*}$ and the end point $\boldsymbol{\lambda}^{N}$, i.e.,

$$
\begin{equation*}
\boldsymbol{\lambda}(t) \equiv(1-t) \boldsymbol{\lambda}^{*}+t \boldsymbol{\lambda}^{N}, 0 \leq t \leq 1 \tag{216}
\end{equation*}
$$

The components of (216) are $\lambda_{i}(t)=(1-t) \lambda_{i}^{*}+t \lambda_{i}^{N}, i=$ $1, \ldots, m$. Let us define

$$
\begin{equation*}
f(t) \equiv F(\boldsymbol{\lambda}(t)) \in \Delta(\mathcal{X}) \tag{217}
\end{equation*}
$$

and write its components as $f(t)=\left(f_{1}(t), \ldots, f_{m}(t)\right)$. We have

$$
\begin{align*}
\frac{d f_{i}(t)}{d t} & =\left.\sum_{i^{\prime}=1}^{m} \frac{d \lambda_{i^{\prime}}(t)}{d t} \frac{\partial F_{i}}{\partial \lambda_{i^{\prime}}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}(t)}  \tag{218}\\
& =\left.\sum_{i^{\prime}=1}^{m}\left(\lambda_{i^{\prime}}^{N}-\lambda_{i^{\prime}}^{*}\right) \frac{\partial F_{i}}{\partial \lambda_{i^{\prime}}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}(t)}  \tag{219}\\
& =\left(\left(\boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}\right) J(\boldsymbol{\lambda}(t))\right)_{i}, i=1, \ldots, m \tag{220}
\end{align*}
$$

thus

$$
\begin{equation*}
\frac{d f(t)}{d t}=\left(\boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}\right) J(\boldsymbol{\lambda}(t)) \tag{221}
\end{equation*}
$$

Now, by the relation between the matrix norm and the maximum eigenvalue [10], p.347, Lemma 5.6.10, for $\epsilon \equiv \theta-$ $\theta_{\max }>0$, there exists a vector norm $\|\cdot\|^{\prime}$ in $\mathbb{R}^{m}$ whose associated matrix norm $\|\cdot\|^{\prime}$ satisfies

$$
\begin{equation*}
\theta_{\max } \leq\left\|J\left(\boldsymbol{\lambda}^{*}\right)\right\|^{\prime}<\theta_{\max }+\epsilon \tag{222}
\end{equation*}
$$

(Note that ' does not denote the derivative.) By the continuity of norm, for any $\epsilon_{1}$ with $0<\epsilon_{1}<\theta_{\max }+\epsilon-$ $\left\|J\left(\boldsymbol{\lambda}^{*}\right)\right\|^{\prime}$, there exists $\delta^{\prime}>0$ such that if $\left\|\boldsymbol{\lambda}-\boldsymbol{\lambda}^{*}\right\|^{\prime}<$ $\delta^{\prime}$, then $\left|\|J(\boldsymbol{\lambda})\|^{\prime}-\left\|J\left(\boldsymbol{\lambda}^{*}\right)\right\|^{\prime}\right|<\epsilon_{1}$, especially, $\|J(\boldsymbol{\lambda})\|^{\prime}<$ $\left\|J\left(\boldsymbol{\lambda}^{*}\right)\right\|^{\prime}+\epsilon_{1}$. Thus,

$$
\begin{align*}
\|J(\boldsymbol{\lambda})\|^{\prime} & <\left\|J\left(\boldsymbol{\lambda}^{*}\right)\right\|^{\prime}+\theta_{\max }+\epsilon-\left\|J\left(\boldsymbol{\lambda}^{*}\right)\right\|^{\prime}  \tag{223}\\
& =\theta<1 \tag{224}
\end{align*}
$$

By the mean value theorem, there exists $t^{N} \in[0,1]$, which satisfies

$$
\begin{align*}
\left\|\boldsymbol{\lambda}^{N+1}-\boldsymbol{\lambda}^{*}\right\|^{\prime} & =\left\|F\left(\boldsymbol{\lambda}^{N}\right)-F\left(\boldsymbol{\lambda}^{*}\right)\right\|^{\prime}  \tag{225}\\
& =\|f(1)-f(0)\|^{\prime}  \tag{226}\\
& \leq\left\|\left.\frac{d f(t)}{d t}\right|_{t=t^{N}}\right\|^{\prime}(1-0)  \tag{227}\\
& =\left\|\left(\boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}\right) J\left(\boldsymbol{\lambda}\left(t^{N}\right)\right)\right\|^{\prime}  \tag{228}\\
& \leq\left\|\boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}\right\|^{\prime}\left\|J\left(\boldsymbol{\lambda}\left(t^{N}\right)\right)\right\|^{\prime} \tag{229}
\end{align*}
$$

Here, if $\left\|\boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}\right\|^{\prime}<\delta^{\prime}$, then we have $\left\|J\left(\boldsymbol{\lambda}^{N}\right)\right\|^{\prime}<$ $\theta<1$ by (224), so $\left\|\boldsymbol{\lambda}^{N+1}-\boldsymbol{\lambda}^{*}\right\|^{\prime}<\delta^{\prime}$ by (229). Thus, by mathematical induction, if the initial distribution $\boldsymbol{\lambda}^{0}$ satisfies $\left\|\boldsymbol{\lambda}^{0}-\boldsymbol{\lambda}^{*}\right\|^{\prime}<\delta^{\prime}$, then $\left\|\boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}\right\|^{\prime}<\delta^{\prime}$ for all $N$, and so $\left\|J\left(\boldsymbol{\lambda}^{N}\right)\right\|^{\prime}<\theta<1$ by (224).

Therefore, by (224), (229), $\left\|\boldsymbol{\lambda}^{N+1}-\boldsymbol{\lambda}^{*}\right\|^{\prime}<\theta\left\|\boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}\right\|^{\prime}<$ $\ldots<\theta^{N+1}\left\|\boldsymbol{\lambda}^{0}-\boldsymbol{\lambda}^{*}\right\|^{\prime}$, so we have

$$
\begin{equation*}
\left\|\boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}\right\|^{\prime}<(\theta)^{N}\left\|\boldsymbol{\lambda}^{0}-\boldsymbol{\lambda}^{*}\right\|^{\prime}, \quad N=0,1, \ldots \tag{230}
\end{equation*}
$$

By the equivalence of norms in the finite dimensional vector space [10], [18], we can change the norm from $\|\cdot\|^{\prime}$ to the Euclidean norm $\|\cdot\|$ to have

$$
\begin{equation*}
\left\|\boldsymbol{\lambda}^{N}-\boldsymbol{\lambda}^{*}\right\| \leq K \cdot(\theta)^{N}, K>0, N=0,1, \ldots \tag{231}
\end{equation*}
$$

## Appendix C

## Proof of Theorem 6 (Calculation of Hessian

 MATRIX $\left.H_{i}\left(\boldsymbol{\lambda}^{*}\right)\right)$We will calculate the Hessian matrix $H_{i}\left(\boldsymbol{\lambda}^{*}\right)$ of $F_{i}(\boldsymbol{\lambda})$ at $\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}$, i.e., $H_{i}\left(\boldsymbol{\lambda}^{*}\right)=\left(\partial^{2} F_{i} /\left.\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}}\right)$.

Differentiating both sides of (39) with respect to $\lambda_{i^{\prime \prime}}$, we have

$$
\begin{align*}
& +\frac{\partial F_{i}}{\partial \lambda_{i^{\prime \prime}}} \frac{\partial}{\partial \lambda_{i^{\prime}}} \sum_{k=1}^{m} \lambda_{k} e^{D_{k}}+F_{i} \frac{\partial^{2}}{\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}} \sum_{k=1}^{m} \lambda_{k} e^{D_{k}} \\
& =\delta_{i i^{\prime}} e^{D_{i}} \frac{\partial D_{i}}{\partial \lambda_{i^{\prime \prime}}}+\delta_{i i^{\prime \prime}} e^{D_{i}} \frac{\partial D_{i}}{\partial \lambda_{i^{\prime}}}+\lambda_{i} e^{D_{i}} \frac{\partial D_{i}}{\partial \lambda_{i^{\prime \prime}}} \frac{\partial D_{i}}{\partial \lambda_{i^{\prime}}} \\
& +\lambda_{i} e^{D_{i}} \frac{\partial^{2} D_{i}}{\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}} . \tag{232}
\end{align*}
$$

On the left hand side of (232), part $\star 1$ is evaluated at $\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}$ by (45). Part $\star 2$ is the component of the Jacobian matrix which is evaluated by (51). Part $\star 3$ is evaluated by (47). Part $\star 4$ is evaluated as follows. We have

$$
\begin{align*}
& \frac{\partial^{2}}{\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}} \sum_{k=1}^{m} \lambda_{k} e^{D_{k}}=\frac{\partial}{\partial \lambda_{i^{\prime \prime}}}\left(e^{D_{i^{\prime}}}+\sum_{k=1}^{m} \lambda_{k} e^{D_{k}} \frac{\partial D_{k}}{\partial \lambda_{i^{\prime}}}\right) \\
& =e^{D_{i^{\prime}}} \frac{\partial D_{i^{\prime}}}{\partial \lambda_{i^{\prime \prime}}}+\sum_{k=1}^{m}\left(\delta_{k i^{\prime \prime}} e^{D_{k}} \frac{\partial D_{k}}{\partial \lambda_{i^{\prime}}}+\lambda_{k} e^{D_{k}} \frac{\partial D_{k}}{\partial \lambda_{i^{\prime \prime}}} \frac{\partial D_{k}}{\partial \lambda_{i^{\prime}}}\right. \\
& \left.\quad+\lambda_{k} e^{D_{k}} \frac{\partial^{2} D_{k}}{\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}}\right) \\
& =e^{D_{i^{\prime}}} \frac{\partial D_{i}^{\prime}}{\partial \lambda_{i^{\prime \prime}}^{\prime}}+e^{D_{i^{\prime \prime}}} \frac{\partial D_{i^{\prime \prime}}}{\partial \lambda_{i^{\prime}}}+\sum_{k=1}^{m} \lambda_{k} e^{D_{k}} \frac{\partial D_{k}}{\partial \lambda_{i^{\prime}}} \frac{\partial D_{k}}{\partial \lambda_{i^{\prime \prime}}} \\
& \quad+\sum_{k=1}^{m} \lambda_{k} e^{D_{k}} \frac{\partial^{2} D_{k}}{\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}} \tag{233}
\end{align*}
$$

and part $\star 5$ becomes, at $\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}$,

$$
\begin{align*}
\left.\sum_{k=1}^{m} \lambda_{k} e^{D_{k}} \frac{\partial^{2} D_{k}}{\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}}\right|_{\lambda=\lambda^{*}} & =e^{C} \sum_{k=1}^{m_{1}} \lambda_{k}^{*} \sum_{j=1}^{n} \frac{P_{j}^{k} P_{j}^{i^{\prime}} P_{j}^{i^{\prime \prime}}}{\left(Q_{j}^{*}\right)^{2}} \\
& =e^{C} \sum_{j=1}^{n} \frac{P_{j}^{i^{\prime}} P_{j}^{i^{\prime \prime}}}{Q_{j}^{*}} \sum_{k=1}^{m_{1}} \frac{\lambda_{k}^{*} P_{j}^{k}}{Q_{j}^{*}} \\
& =-e^{C} D_{i^{\prime}, i^{\prime \prime} .}^{*} . \tag{234}
\end{align*}
$$

Therefore, part $\star 4$ becomes, at $\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}$,

$$
\begin{align*}
\frac{\partial^{2}}{\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}} & \left.\sum_{k=1}^{m} \lambda_{k} e^{D_{k}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}}=e^{D_{i^{\prime}}^{*}} D_{i^{\prime}, i^{\prime \prime}}^{*}+e^{D_{i^{\prime \prime}}^{*}} D_{i^{\prime}, i^{\prime \prime}}^{*} \\
& +e^{C} \sum_{k=1}^{m_{1}} \lambda_{k}^{*} D_{k, i^{\prime}}^{*} D_{k, i^{\prime \prime}}^{*}-e^{C} D_{i^{\prime}, i^{\prime \prime}}^{*} \tag{235}
\end{align*}
$$

Let us define $D_{i, i^{\prime}, i^{\prime \prime}}^{*} \equiv \partial^{2} D_{i} /\left.\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}}$ and $E_{i^{\prime}, i^{\prime \prime}} \equiv$ $\sum_{k=1}^{m_{1}} \lambda_{k}^{*} D_{k, i^{\prime}}^{*} D_{k, i^{\prime \prime}}^{*}$. Then, based on the above calculation, we have

$$
\begin{align*}
& \left.\frac{\partial^{2} F_{i}}{\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}} e^{C} \\
& \quad+\left\{e^{D_{i}^{*}-C}\left(\delta_{i i^{\prime}}+\lambda_{i}^{*} D_{i, i^{\prime}}^{*}\right)+\lambda_{i}^{*}\left(1-e^{D_{i^{\prime}}^{*}-C}\right)\right\} \\
& \quad \times\left(e^{D_{i^{\prime \prime}}^{*}}-e^{C}\right) \\
& \quad+\left\{e^{D_{i}^{*}-C}\left(\delta_{i i^{\prime \prime}}+\lambda_{i}^{*} D_{i, i^{\prime \prime}}^{*}\right)+\lambda_{i}^{*}\left(1-e^{D_{i^{\prime \prime}}^{*}-C}\right)\right\} \\
& \quad \times\left(e^{D_{i^{\prime}}^{*}}-e^{C}\right) \\
& \quad+F_{i}^{*}\left(e^{D_{i^{\prime}}^{*}} D_{i^{\prime}, i^{\prime \prime}}^{*}+e^{D_{i^{\prime \prime}}^{*}} D_{i^{\prime}, i^{\prime \prime}}^{*}+e^{C} E_{i^{\prime}, i^{\prime \prime}}-e^{C} D_{i^{\prime}, i^{\prime \prime}}^{*}\right) \\
& =\delta_{i i^{\prime}} e^{D_{i}^{*}} D_{i, i^{\prime \prime}}^{*}+\delta_{i i^{\prime \prime}} e^{D_{i}^{*}} D_{i, i^{\prime}}^{*}+\lambda_{i}^{*} e^{D_{i}^{*}} D_{i, i^{\prime}}^{*} D_{i, i^{\prime \prime}}^{*} \\
& \quad+\lambda_{i}^{*} e^{D_{i}^{*}} D_{i, i^{\prime}, i^{\prime \prime}}^{*} . \tag{236}
\end{align*}
$$

By rearranging this equation, we obtain, using (48) of Lemma 4,

$$
\begin{align*}
& \left.\frac{\partial^{2} F_{i}}{\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}} \\
& =e^{D_{i}^{*}-C}\left\{\delta_{i i^{\prime}} D_{i, i^{\prime \prime}}^{*}+\delta_{i i^{\prime \prime}} D_{i, i^{\prime}}^{*}+\lambda_{i}^{*}\left(D_{i, i^{\prime}}^{*} D_{i, i^{\prime \prime}}^{*}+D_{i, i^{\prime}, i^{\prime \prime}}^{*}\right)\right. \\
& \quad+\left(\delta_{i i^{\prime}}+\lambda_{i}^{*} D_{i, i^{\prime}}^{*}\right)\left(1-e^{D_{i^{\prime \prime}}^{*}-C}\right) \\
& \left.\quad+\left(\delta_{i i^{\prime \prime}}+\lambda_{i}^{*} D_{i, i^{\prime \prime}}^{*}\right)\left(1-e^{D_{i^{\prime}}^{*}-C}\right)\right\} \\
& \quad+2 \lambda_{i}^{*}\left(1-e^{D_{i^{\prime}}^{*}-C}\right)\left(1-e^{D_{i^{\prime \prime}}^{*}-C}\right) \\
& \quad-\lambda_{i}^{*}\left(e^{D_{i^{\prime}}^{*}-C} D_{i^{\prime}, i^{\prime \prime}}^{*}+e^{D_{i^{\prime \prime}}^{*}-C} D_{i^{\prime}, i^{\prime \prime}}^{*}+E_{i^{\prime}, i^{\prime \prime}}-D_{i^{\prime}, i^{\prime \prime}}^{*}\right) \tag{237}
\end{align*}
$$

## Appendix D

## Proof of Step2

For $i=m^{\prime}+1, \ldots, m$, let $H_{i, i^{\prime} i^{\prime \prime}}$ be the $\left(i^{\prime}, i^{\prime \prime}\right)$ component of the Hessian matrix $H_{i}\left(\boldsymbol{\lambda}^{*}\right)$. Then, by Theorem 6,

$$
\begin{align*}
& H_{i, i^{\prime} i^{\prime \prime}}=\left.\frac{\partial^{2} F_{i}}{\partial \lambda_{i^{\prime}} \partial \lambda_{i^{\prime \prime}}}\right|_{\boldsymbol{\lambda}=\boldsymbol{\lambda}^{*}} \\
& =\delta_{i i^{\prime \prime}}\left(1-e^{D_{i^{\prime}}^{*}-C}+D_{i, i^{\prime}}^{*}\right)+\delta_{i i^{\prime}}\left(1-e^{D_{i^{\prime \prime}}^{*}-C}+D_{i, i^{\prime \prime}}^{*}\right),  \tag{238}\\
& \quad i=m^{\prime}+1, \ldots, m, i^{\prime}, i^{\prime \prime}=1, \ldots, m . \tag{20}
\end{align*}
$$

Here, for the simplicity of symbols, define

$$
\begin{equation*}
S_{i i^{\prime}} \equiv 1-e^{D_{i^{\prime}}^{*}-C}+D_{i, i^{\prime}}^{*}, i^{\prime}=1, \ldots, m \tag{239}
\end{equation*}
$$

then (238) becomes

$$
\begin{equation*}
H_{i, i^{\prime} i^{\prime \prime}}=\delta_{i i^{\prime \prime}} S_{i i^{\prime}}+\delta_{i i^{\prime}} S_{i i^{\prime \prime}} \tag{240}
\end{equation*}
$$

We have $S_{i i^{\prime}}=D_{i, i^{\prime}}^{*}$ for $i^{\prime}=m^{\prime}+1, \ldots, m$, then by (240), the Hessian matrix $H_{i}\left(\boldsymbol{\lambda}^{*}\right)$ is given by (241)-(245).

Therefore, by (99),

$$
\begin{align*}
& \frac{1}{2} \overline{\boldsymbol{\mu}}^{N} H_{i}\left(\boldsymbol{\lambda}^{*}\right)^{t} \overline{\boldsymbol{\mu}}^{N}=\frac{1}{2}\left(\overline{\boldsymbol{\mu}}_{\mathrm{I}, \mathrm{III}}^{N}, \overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N}\right)\left(\begin{array}{cc}
O & H_{i}^{1} \\
{ }^{t} H_{i}^{1} & H_{i}^{2}
\end{array}\right)\binom{{ }^{t} \overline{\boldsymbol{\mu}}_{\mathrm{I}, \mathrm{III}}^{N}}{{ }^{t} \overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N}} \\
& =\frac{1}{2}\left(-\overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N} A_{2} A_{1}^{-1}, \overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N}\right)\left(\begin{array}{cc}
O & H_{i}^{1} \\
{ }^{t} H_{i}^{1} & H_{i}^{2}
\end{array}\right)\binom{-{ }^{t} A_{1}^{-1}{ }^{t} A_{2}{ }^{t} \overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N}}{{ }^{t} \overline{\boldsymbol{\mu}}_{\mathrm{II}}} \\
& =\frac{1}{2} \overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N}\left(-A_{2} A_{1}^{-1} H_{i}^{1}-{ }^{t} H_{i}^{1 t} A_{1}^{-1}{ }^{t} A_{2}+H_{i}^{2}\right)^{t} \overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N} . \tag{246}
\end{align*}
$$

$$
\begin{align*}
& H_{i}\left(\boldsymbol{\lambda}^{*}\right)=\left(\begin{array}{cc:ccccccc} 
& & 0 & \ldots & 0 & S_{i 1} & 0 & \ldots & 0 \\
& O & \vdots & & \vdots & \vdots & \vdots & & \vdots \\
& & & 0 & \ldots & 0 & S_{i m^{\prime}} & 0 & \ldots \\
0 & 0 & 0 & 0 & \ldots & 0 & D_{i, m^{\prime}+1}^{*} & 0 & \ldots \\
0 & \ldots & 0 & 0 \\
\vdots & & \vdots & \vdots & & \vdots & \vdots & \vdots & \\
0 & \ldots & 0 & 0 & \ldots & 0 & D_{i, i-1}^{*} & 0 & \ldots \\
S_{i 1} & \ldots & S_{i m^{\prime}} & D_{i, m^{\prime}+1}^{*} & \ldots & D_{i, i-1}^{*} & 2 D_{i, i}^{*} & D_{i, i+1}^{*} & \ldots \\
0 & \ldots & 0 & 0 & \ldots & 0 & D_{i,, m}^{*} \\
\vdots & & \vdots & \vdots & & \vdots & \vdots & \vdots & \\
0 & \ldots & 0 & 0 & \ldots & 0 & D_{i, m}^{*} & 0 & \ldots \\
0
\end{array}\right)  \tag{241}\\
& \equiv\left(\begin{array}{cc}
O & H_{i}^{1} \\
{ }^{t} H_{i}^{1} & H_{i}^{2}
\end{array}\right), \tag{242}
\end{align*}
$$

where $O \in \mathbb{R}^{m^{\prime} \times m^{\prime}}, H_{i}^{1} \in \mathbb{R}^{m^{\prime} \times m_{2}}, H_{i}^{2} \in \mathbb{R}^{m_{2} \times m_{2}}$ and

$$
\begin{align*}
H_{i}^{1} & \equiv\left(\delta_{i i^{\prime \prime}} S_{i i^{\prime}}\right), i^{\prime}=1, \ldots, m^{\prime}, i^{\prime \prime}=m^{\prime}+\ldots, m,  \tag{243}\\
{ }^{t} H_{i}^{1} & \equiv\left(\delta_{i i^{\prime}} S_{i i^{\prime \prime}}\right), i^{\prime}=m^{\prime}+1, \ldots, m, i^{\prime \prime}=1, \ldots, m^{\prime},  \tag{244}\\
H_{i}^{2} & \equiv\left(\delta_{i i^{\prime \prime}} D_{i i^{\prime}}^{*}+\delta_{i i^{\prime}} D_{i i^{\prime \prime}}^{*}\right), i^{\prime}, i^{\prime \prime}=m^{\prime}+1, \ldots, m . \tag{245}
\end{align*}
$$

Now, define

$$
\begin{equation*}
G_{i} \equiv-A_{2} A_{1}^{-1} H_{i}^{1} \in \mathbb{R}^{m_{2} \times m_{2}} \tag{247}
\end{equation*}
$$

and $A_{1}^{-1} \equiv\left(\zeta_{i^{\prime} i^{\prime \prime}}\right)$. Further, let $G_{i, i^{\prime} i^{\prime \prime}}$ be the $\left(i^{\prime}, i^{\prime \prime}\right)$ component of $G_{i}$. Then, by (247), we have

$$
\begin{align*}
G_{i, i^{\prime} i^{\prime \prime}} & =-\sum_{k, k^{\prime}=1}^{m^{\prime}} a_{i^{\prime} k} \zeta_{k k^{\prime}} H_{i, k^{\prime} i^{\prime \prime}}^{1}  \tag{248}\\
& =-\sum_{k, k^{\prime}=1}^{m^{\prime}} a_{i^{\prime} k} \zeta_{k k^{\prime}} \delta_{i i^{\prime \prime}} S_{i k^{\prime}}  \tag{249}\\
& =-\delta_{i i^{\prime \prime}} \sum_{k, k^{\prime}=1}^{m^{\prime}} a_{i^{\prime} k} \zeta_{k k^{\prime}} S_{i k^{\prime}}  \tag{250}\\
& i^{\prime}, i^{\prime \prime}=m^{\prime}+1, \ldots, m
\end{align*}
$$

Define

$$
\begin{equation*}
T_{i i^{\prime}} \equiv-\sum_{k, k^{\prime}=1}^{m^{\prime}} a_{i^{\prime} k} \zeta_{k k^{\prime}} S_{i k^{\prime}}, i^{\prime}=m^{\prime}+1, \ldots, m \tag{251}
\end{equation*}
$$

then (250) becomes

$$
\begin{equation*}
G_{i, i^{\prime} i^{\prime \prime}}=\delta_{i i^{\prime \prime}} T_{i i^{\prime}} \tag{252}
\end{equation*}
$$

Thus, we have

$$
\begin{align*}
& G_{i}=\left(\delta_{i i^{\prime \prime}} T_{i i^{\prime}}\right),{ }^{t} G_{i}=\left(\delta_{i i^{\prime}} T_{i i^{\prime \prime}}\right)  \tag{253}\\
& i^{\prime}, i^{\prime \prime}=m^{\prime}+1, \ldots, m
\end{align*}
$$

Hence, by (246),

$$
\begin{equation*}
\frac{1}{2} \overline{\boldsymbol{\mu}}^{N} H_{i}\left(\boldsymbol{\lambda}^{*}\right)^{t} \overline{\boldsymbol{\mu}}^{N}=\frac{1}{2} \overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N}\left(G_{i}+{ }^{t} G_{i}+H_{i}^{2}\right)^{t} \overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N} \tag{254}
\end{equation*}
$$

Define

$$
\begin{equation*}
\hat{H}_{i} \equiv G_{i}+{ }^{t} G_{i}+H_{i}^{2} \tag{255}
\end{equation*}
$$

and

$$
\begin{equation*}
r_{i i^{\prime}} \equiv T_{i i^{\prime}}+D_{i, i^{\prime}}^{*} \tag{256}
\end{equation*}
$$

and let $\hat{H}_{i, i^{\prime} i^{\prime \prime}}$ be the $\left(i^{\prime}, i^{\prime \prime}\right)$ component of $\hat{H}_{i}$. Then, by (255), (253), (245), and (256), we have

$$
\begin{align*}
\hat{H}_{i, i^{\prime} i^{\prime \prime}} & =\delta_{i i^{\prime \prime}} T_{i i^{\prime}}+\delta_{i i^{\prime}} T_{i i^{\prime \prime}}+\left(\delta_{i i^{\prime \prime}} D_{i, i^{\prime}}^{*}+\delta_{i i^{\prime}} D_{i, i^{\prime \prime}}^{*}\right)  \tag{257}\\
& =\delta_{i i^{\prime \prime}}\left(T_{i i^{\prime}}+D_{i, i^{\prime}}^{*}\right)+\delta_{i i^{\prime}}\left(T_{i i^{\prime \prime}}+D_{i, i^{\prime \prime}}^{*}\right)  \tag{258}\\
& =\delta_{i i^{\prime \prime}} r_{i i^{\prime}}+\delta_{i i^{\prime}} r_{i i^{\prime \prime}} \tag{259}
\end{align*}
$$

Therefore, (254) becomes

$$
\begin{align*}
& \frac{1}{2} \overline{\boldsymbol{\mu}}^{N} H_{i}\left(\boldsymbol{\lambda}^{*}\right)^{t} \overline{\boldsymbol{\mu}}^{N}=\frac{1}{2} \overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N} \hat{H}_{i}^{t} \overline{\boldsymbol{\mu}}_{\mathrm{II}}^{N}  \tag{260}\\
& =\frac{1}{2} \sum_{i^{\prime}, i^{\prime \prime}=m^{\prime}+1}^{m} \bar{\mu}_{i^{\prime}}^{N} \hat{H}_{i, i^{\prime} i^{\prime \prime}} \bar{\mu}_{i^{\prime \prime}}^{N}  \tag{261}\\
& =\frac{1}{2} \sum_{i^{\prime}, i^{\prime \prime}=m^{\prime}+1}^{m} \bar{\mu}_{i^{\prime}}^{N}\left(\delta_{i i^{\prime \prime}} r_{i i^{\prime}}+\delta_{i i^{\prime}} r_{i i^{\prime \prime}}\right) \bar{\mu}_{i^{\prime \prime}}^{N}  \tag{262}\\
& =\frac{1}{2}\left(\sum_{i^{\prime}=m^{\prime}+1}^{m} \bar{\mu}_{i^{\prime}}^{N} r_{i i^{\prime}} \bar{\mu}_{i}^{N}+\sum_{i^{\prime \prime}=m^{\prime}+1}^{m} \bar{\mu}_{i}^{N} r_{i i^{\prime \prime}} \bar{\mu}_{i^{\prime \prime}}^{N}\right)  \tag{263}\\
& =\bar{\mu}_{i}^{N} \sum_{i^{\prime}=m^{\prime}+1}^{m} r_{i i^{\prime}} \bar{\mu}_{i^{\prime}}^{N} \tag{264}
\end{align*}
$$

Summarizing the above, the recurrence formula satisfied by $\bar{\mu}_{i}^{N}, i=m^{\prime}+1, \ldots, m$, is

$$
\begin{aligned}
& \bar{\mu}_{i}^{N+1}=\bar{\mu}_{i}^{N}+\bar{\mu}_{i}^{N} \sum_{i^{\prime}=m^{\prime}+1}^{m} r_{i i^{\prime}} \bar{\mu}_{i^{\prime}}^{N} \\
& \quad i=m^{\prime}+1, \ldots, m
\end{aligned}
$$

## Appendix E <br> Proof of Theorem 10

We can prove Theorem 10 for any $m_{2} \geq 3$, but because the symbols become complicated, we will give a proof for the case of $m_{2}=3$. The proof for $m_{2}=3$ does not lose generality, hence the extension to $m_{2} \geq 3$ is easy. We should prove the following Theorem.

Theorem 10 (case of $m_{2}=3$ ): Let us consider a sequence $\left\{\xi_{i}^{N}\right\}, i=1,2,3, N=0,1, \ldots$ defined by

$$
\begin{align*}
\xi_{i}^{N+1} & =\xi_{i}^{N}-\xi_{i}^{N} \sum_{i^{\prime}=1}^{3} q_{i i^{\prime}} \xi_{i^{\prime}}^{N}, i=1,2,3  \tag{266}\\
0 & <\xi_{i}^{0} \leq 1 / 2, i=1,2,3 \tag{267}
\end{align*}
$$

where $\boldsymbol{q}_{i} \equiv\left(q_{i 1}, q_{i 2}, q_{i 3}\right)$ is a probability vector.
Further, we assume the diagonally dominant condition (122), i.e.,

$$
\begin{equation*}
q_{i i}>\sum_{i^{\prime}=1, i^{\prime} \neq i}^{3} q_{i i^{\prime}}, i=1,2,3 \tag{268}
\end{equation*}
$$

and the conjecture (123), i.e.,

$$
\begin{equation*}
\xi_{1}^{N} \geq \xi_{2}^{N} \geq \xi_{3}^{N}, N \geq N_{0} \tag{269}
\end{equation*}
$$

Then, we have

$$
\begin{equation*}
\lim _{N \rightarrow \infty} N \xi_{i}^{N}=1, i=1,2,3 \tag{270}
\end{equation*}
$$

Lemma 12: $0<\xi_{i}^{N} \leq 1 / 2, i=1,2,3$, holds for $N=$ $0,1, \ldots$.

Proof: We prove by mathematical induction. For $N=0$, the assertion holds by (267). Assuming that the assertion holds for $N$, by noting $1 / 2 \leq 1-\sum_{i^{\prime}=1} q_{i i^{\prime}} \xi_{i^{\prime}}^{N}<1$, we have $0<\xi_{i}^{N+1} \leq 1 / 2$.

Lemma 13: The sequence $\left\{\xi_{i}^{N}\right\}, N=0,1, \ldots$ is strictly decreasing.

Proof: Because $\xi_{i}^{N}-\xi_{i}^{N+1}=\xi_{i}^{N} \sum_{i^{\prime}=1}^{3} q_{i i^{\prime}} \xi_{i^{\prime}}^{N}>0$ holds by Lemma 12.

Lemma 14: $\lim _{N \rightarrow \infty} \xi_{i}^{N}=0, i=1,2,3$.
Proof: $\xi_{i}^{\infty} \equiv \lim _{N \rightarrow \infty} \xi_{i}^{N} \geq 0$ exists by Lemmas 12 and 13. Then, $\xi_{i}^{\infty}=\xi_{i}^{\infty}-\xi_{i}^{\infty} \sum_{i^{\prime}=1}^{3} q_{i i^{\prime}} \xi_{i^{\prime}}^{\infty}$ holds by (266), hence we have $\xi_{i}^{\infty}=0$.

Lemma 15: $\liminf _{N \rightarrow \infty} N \xi_{1}^{N} \geq 1$.
Proof: By (269), for $N \geq N_{0}$,

$$
\begin{align*}
\xi_{1}^{N+1} & =\xi_{1}^{N}-\xi_{1}^{N} \sum_{i^{\prime}=1}^{3} q_{1 i^{\prime}} \xi_{i^{\prime}}^{N}  \tag{271}\\
& \geq \xi_{1}^{N}-\xi_{1}^{N} \sum_{i^{\prime}=1}^{3} q_{1 i^{\prime}} \xi_{1}^{N}  \tag{272}\\
& =\xi_{1}^{N}-\left(\xi_{1}^{N}\right)^{2} . \tag{273}
\end{align*}
$$

Now, we define a sequence $\left\{\hat{\xi}_{1}^{N}\right\}, N=0,1, \ldots$ by the recurrence formula

$$
\begin{align*}
\hat{\xi}_{1}^{0} & =\xi_{1}^{N_{0}},  \tag{274}\\
\hat{\xi}_{1}^{N+1} & =\hat{\xi}_{1}^{N}-\left(\hat{\xi}_{1}^{N}\right)^{2}, N=0,1, \ldots \tag{275}
\end{align*}
$$

Then, we will prove

$$
\begin{equation*}
\xi_{1}^{N+N_{0}} \geq \hat{\xi}_{1}^{N}, N=0,1, \ldots \tag{276}
\end{equation*}
$$

by mathematical induction. For $N=0$, (276) holds by the assumption (274). Assume that (276) holds for $N$. Because the function $f(\xi)=\xi-\xi^{2}$ is monotonically increasing in $0<\xi \leq 1 / 2$, by (273), we have $\xi_{1}^{N+N_{0}+1} \geq \xi_{1}^{N+N_{0}}-$ $\left(\xi_{1}^{N+N_{0}}\right)^{2} \geq \hat{\xi}_{1}^{N}-\left(\hat{\xi}_{1}^{N}\right)^{2}=\hat{\xi}_{1}^{N+1}$; thus, (276) also holds for $N+1$. Therefore, $\liminf _{N \rightarrow \infty} N \xi_{1}^{N} \geq \liminf _{N \rightarrow \infty} N \hat{\xi}_{1}^{N}=$ $\lim _{N \rightarrow \infty} N \hat{\xi}_{1}^{N}=1$, where the last equality is because of Lemma 6.

Lemma 16: $\lim \sup _{N \rightarrow \infty} N \xi_{3}^{N} \leq 1$.
Proof: By (269), for $N \geq N_{0}$, we have $\xi_{3}^{N+1}=\xi_{3}^{N}-$ $\xi_{3}^{N} \sum_{i^{\prime}=1}^{3} q_{3 i^{\prime}} \xi_{i^{\prime}}^{N} \leq \xi_{3}^{N}-\xi_{3}^{N} \sum_{i^{\prime}=1}^{3} q_{3 i^{\prime}} \xi_{3}^{n}=\xi_{3}^{N}-\left(\xi_{3}^{N}\right)^{2}$. Now, we define a sequence $\left\{\hat{\xi}_{3}^{N}\right\}$ by

$$
\begin{align*}
\hat{\xi}_{3}^{0} & =\xi_{3}^{N_{0}}  \tag{277}\\
\hat{\xi}_{3}^{N+1} & =\hat{\xi}_{3}^{N}-\left(\hat{\xi}_{3}^{N}\right)^{2}, N=0,1, \ldots \tag{278}
\end{align*}
$$

Then, we can prove $\xi_{3}^{N+N_{0}} \leq \hat{\xi}_{3}^{N}, N=0,1, \ldots$ in a similar way as the proof of Lemma 15. Therefore, we have $\lim \sup _{N \rightarrow \infty} N \xi_{3}^{N} \leq \lim \sup _{N \rightarrow \infty} N \hat{\xi}_{3}^{N}=$ $\lim _{N \rightarrow \infty} N \hat{\xi}_{3}^{N}=1$, where the last equality is because of Lemma 6.

Lemma 17: Let $\tau_{1}^{N} \equiv \sum_{i^{\prime}=1}^{3} q_{1 i^{\prime}} \xi_{i^{\prime}}^{N}, \tau_{3}^{N} \equiv \sum_{i^{\prime}=1}^{3} q_{3 i^{\prime}} \xi_{i^{\prime}}^{N}$. Then, there exists a constant $K>0$ with

$$
\begin{equation*}
\tau_{1}^{N}-\tau_{3}^{N} \geq K\left(\xi_{1}^{N}-\xi_{3}^{N}\right), N \geq N_{0} \tag{279}
\end{equation*}
$$

Proof: We have

$$
\begin{align*}
& \tau_{1}^{N}-\tau_{3}^{N}=\sum_{i^{\prime}=1}^{3} q_{1 i^{\prime}} \xi_{i^{\prime}}^{N}-\sum_{i^{\prime}=1}^{3} q_{3 i^{\prime}} \xi_{i^{\prime}}^{N}  \tag{280}\\
& \geq q_{11} \xi_{1}^{N}+q_{12} \xi_{3}^{N}+q_{13} \xi_{3}^{N}-q_{31} \xi_{1}^{n}-q_{32} \xi_{1}^{N}-q_{33} \xi_{3}^{N}  \tag{281}\\
& =\left(q_{11}-q_{31}-q_{32}\right) \xi_{1}^{N}-\left(q_{33}-q_{12}-q_{13}\right) \xi_{3}^{N}, N \geq N_{0}, \tag{282}
\end{align*}
$$

and $q_{11}-q_{31}-q_{32}=q_{33}-q_{12}-q_{13}$ holds by $q_{11}+q_{12}+$ $q_{13}=q_{31}+q_{32}+q_{33}=1$. Defining $K \equiv q_{11}-q_{31}-q_{32}=$ $q_{33}-q_{12}-q_{13}$, we have $\tau_{1}^{N}-\tau_{3}^{N} \geq K\left(\xi_{1}^{N}-\xi_{3}^{N}\right)$. By the assumption (268), we have $q_{11}>1 / 2, q_{12}+q_{13}<1 / 2, q_{33}>$ $1 / 2, q_{31}+q_{32}<1 / 2$, thus $K>0$.

Lemma 18: $\sum_{N=0}^{\infty}\left(\xi_{1}^{N}-\xi_{3}^{N}\right)<\infty$.
Proof: We have

$$
\begin{align*}
\frac{\xi_{1}^{N}}{\xi_{3}^{N}}-\frac{\xi_{1}^{N+1}}{\xi_{3}^{N+1}} & =\frac{\xi_{1}^{N}}{\xi_{3}^{N}}-\frac{\xi_{1}^{N}\left(1-\tau_{1}^{N}\right)}{\xi_{3}^{N}\left(1-\tau_{3}^{N}\right)}  \tag{283}\\
& =\frac{\xi_{1}^{N}}{\xi_{3}^{N}} \cdot \frac{\tau_{1}^{N}-\tau_{3}^{N}}{1-\tau_{3}^{N}} \tag{284}
\end{align*}
$$

By (269), we have $\xi_{1}^{N} / \xi_{3}^{N} \geq 1, N \geq N_{0}$, and by Lemma 17, $\tau_{1}^{N}-\tau_{3}^{N} \geq K\left(\xi_{1}^{N}-\xi_{3}^{N}\right), K>0, N \geq N_{0}$. Further, by Lemma 12, we have $0<\tau_{3}^{N} \leq 1 / 2$, thus $1 /\left(1-\tau_{3}^{N}\right)>1$. Therefore, by (284),

$$
\begin{equation*}
\frac{\xi_{1}^{N_{0}}}{\xi_{3}^{N_{0}}}-\frac{\xi_{1}^{N+1}}{\xi_{3}^{N+1}}>K \sum_{l=N_{0}}^{N}\left(\xi_{1}^{l}-\xi_{3}^{l}\right), N \geq N_{0} \tag{285}
\end{equation*}
$$

hence $\sum_{l=N_{0}}^{N}\left(\xi_{1}^{l}-\xi_{3}^{l}\right)$ has an upper bound $K^{-1}\left(\xi_{1}^{N_{0}} / \xi_{3}^{N_{0}}\right)$ which is unrelated to $N$. Then, the sum is convergent.

Here, we cite the following theorem.
Theorem B: ([5], p.31) Let $\left\{a_{N}\right\}_{N=0,1, \ldots \text {, be a decreasing }}$ positive sequence. If $\sum_{N=0}^{\infty} a_{N}$ is convergent, then $N a_{N} \rightarrow$ $0, N \rightarrow \infty$.

Lemma 19: The sequence $\left\{\xi_{1}^{N}-\xi_{3}^{N}\right\}$ is decreasing for $N \geq$ $N_{0}$.

Proof: We have

$$
\begin{align*}
\xi_{1}^{N+1}-\xi_{3}^{N+1} & =\xi_{1}^{N}-\xi_{1}^{N} \sum_{i^{\prime}=1}^{3} q_{1 i^{\prime}} \xi_{i^{\prime}}^{N}-\xi_{3}^{N}+\xi_{3}^{N} \sum_{i^{\prime}=1}^{3} q_{3 i^{\prime}} \xi_{i^{\prime}}^{N} \\
& \leq \xi_{1}^{N}-\xi_{1}^{N} \sum_{i^{\prime}=1}^{3} q_{1 i^{\prime}} \xi_{3}^{N}-\xi_{3}^{N}+\xi_{3}^{N} \sum_{i^{\prime}=1}^{3} q_{3 i^{\prime}} \xi_{1}^{N}  \tag{286}\\
& =\xi_{1}^{N}-\xi_{1}^{N} \xi_{3}^{N}-\xi_{3}^{N}+\xi_{3}^{N} \xi_{1}^{N}  \tag{287}\\
& =\xi_{1}^{N}-\xi_{3}^{N}, N \geq N_{0} \tag{289}
\end{align*}
$$

hence the assertion holds.
Lemma 20: $\lim _{N \rightarrow \infty} N\left(\xi_{1}^{N}-\xi_{3}^{N}\right)=0$.
Proof: The assertion holds by Lemmas 18, 19, and Theorem B.

Summarizing the above, we have
Theorem 10 (case of $m_{2}=3$ ):

$$
\lim _{N \rightarrow \infty} N \xi_{i}^{N}=1, i=1,2,3
$$

Proof: By Lemmas 16 and 20, we have

$$
\begin{align*}
\limsup _{N \rightarrow \infty} N \xi_{1}^{N} & =\limsup _{N \rightarrow \infty}\left(N \xi_{1}^{N}-N \xi_{3}^{N}+N \xi_{3}^{N}\right)  \tag{290}\\
& \leq \limsup _{N \rightarrow \infty} N\left(\xi_{1}^{N}-\xi_{3}^{N}\right)+\limsup _{N \rightarrow \infty} N \xi_{3}^{N} \tag{291}
\end{align*}
$$

$$
\begin{equation*}
\leq 1 \tag{292}
\end{equation*}
$$

thus, together with Lemma 15, we have

$$
\begin{equation*}
\lim _{N \rightarrow \infty} N \xi_{1}^{N}=1 \tag{293}
\end{equation*}
$$

Further, we have

$$
\begin{align*}
\lim _{N \rightarrow \infty} N \xi_{3}^{N} & =\lim _{N \rightarrow \infty}\left(N \xi_{3}^{N}-N \xi_{1}^{N}+N \xi_{1}^{N}\right)  \tag{294}\\
& =-\lim _{N \rightarrow \infty} N\left(\xi_{1}^{N}-\xi_{3}^{N}\right)+\lim _{N \rightarrow \infty} N \xi_{1}^{N}  \tag{295}\\
& =1 \tag{296}
\end{align*}
$$

Finally, by (293), (296) and the squeeze theorem, we have

$$
\begin{equation*}
\lim _{N \rightarrow \infty} N \xi_{2}^{N}=1 \tag{297}
\end{equation*}
$$

## Appendix F

Proof of Lemma 9
Let $0=\theta_{1} \leq \ldots \leq \theta_{m-1} \leq \theta_{m}<1$ be the eigenvalues of $J\left(\boldsymbol{\lambda}^{*}\right)$. We have $\theta_{\max }=\theta_{m}, \theta_{\text {sec }}=\theta_{m-1}$. Let $\boldsymbol{b}_{i}, i=1, \ldots, m$, be the left eigenvector of $J\left(\boldsymbol{\lambda}^{*}\right)$ for $\theta_{i}, i=1, \ldots, m$. We have $\boldsymbol{b}_{\max }=\boldsymbol{b}_{m}$. Because of (75), we can assume that $\left\{\boldsymbol{b}_{i}\right\}_{i=1, \ldots, m}$ forms a basis of $\mathbb{R}^{m}$. Suppose $\boldsymbol{\mu}^{N} \in \boldsymbol{b}_{\max }^{\perp}$ for $N=0,1, \ldots$ Then, $\boldsymbol{\mu}^{N}$ is uniquely represented as

$$
\begin{equation*}
\boldsymbol{\mu}^{N}=\sum_{i=1}^{m-1} c_{i}^{N} \boldsymbol{b}_{i}, c_{i}^{N} \in \mathbb{R} \tag{298}
\end{equation*}
$$

in the $m-1$ dimensional subspace $\boldsymbol{b}_{\text {max }}^{\perp}$. By (298), we have

$$
\begin{align*}
\boldsymbol{\mu}^{N+1} & =\boldsymbol{\mu}^{N} J\left(\boldsymbol{\lambda}^{*}\right)  \tag{299}\\
& =\sum_{i=1}^{m-1} c_{i}^{N} \boldsymbol{b}_{i} J\left(\boldsymbol{\lambda}^{*}\right)  \tag{300}\\
& =\sum_{i=1}^{m-1} c_{i}^{N} \theta_{i} \boldsymbol{b}_{i} \tag{301}
\end{align*}
$$

Comparing the coefficients of $\boldsymbol{\mu}^{N+1}=\sum_{i=1}^{m-1} c_{i}^{N+1} \boldsymbol{b}_{i}$ and (301), we have $c_{i}^{N+1}=\theta_{i} c_{i}^{N}=\ldots=\left(\theta_{i}\right)^{N+1} c_{i}^{0}, i=$ $1, \ldots, m-1$, thus $\boldsymbol{\mu}^{N}=\sum_{i=1}^{m-1}\left(\theta_{i}\right)^{N} c_{i}^{0} \boldsymbol{b}_{i}$. Therefore,

$$
\begin{align*}
\left\|\boldsymbol{\mu}^{N}\right\| & \leq \sum_{i=1}^{m-1}\left(\theta_{i}\right)^{N}\left|c_{i}^{0}\right|\left\|\boldsymbol{b}_{i}\right\|  \tag{302}\\
& \leq K \cdot\left(\theta_{m-1}\right)^{N}, K>0  \tag{303}\\
& =K \cdot\left(\theta_{\mathrm{sec}}\right)^{N} \tag{304}
\end{align*}
$$

## Appendix G <br> Proof of Lemma 10

Suppose that ${ }^{t} \boldsymbol{b}_{\text {max }}$ is a right eigenvector for $\theta_{\text {max }}$. For any $\boldsymbol{\mu} \in \boldsymbol{b}_{\max }^{\perp}, \boldsymbol{\mu} J\left(\boldsymbol{\lambda}^{*}\right)^{t} \boldsymbol{b}_{\max }=\theta_{\max } \boldsymbol{\mu}^{t} \boldsymbol{b}_{\max }=0$ holds. Thus, we obtain $\boldsymbol{\mu} J\left(\boldsymbol{\lambda}^{*}\right) \in \boldsymbol{b}_{\text {max }}^{\perp}$.

Conversely, suppose $\boldsymbol{\mu} J\left(\boldsymbol{\lambda}^{*}\right) \in \boldsymbol{b}_{\text {max }}^{\perp}$ for any $\boldsymbol{\mu} \in \boldsymbol{b}_{\text {max }}^{\perp}$. Our goal is to show $J\left(\boldsymbol{\lambda}^{*}\right)^{t} \boldsymbol{b}_{\text {max }}=\theta_{\text {max }}{ }^{t} \boldsymbol{b}_{\text {max }}$, which is equivalent to

$$
\begin{equation*}
\boldsymbol{\mu} J\left(\boldsymbol{\lambda}^{*}\right)^{t} \boldsymbol{b}_{\max }=\theta_{\max } \boldsymbol{\mu}^{t} \boldsymbol{b}_{\max } \text { holds for any } \boldsymbol{\mu} \tag{305}
\end{equation*}
$$

We will prove (305). Since we can write $\boldsymbol{\mu}$ uniquely as $\boldsymbol{\mu}=$ $K \boldsymbol{b}_{\max }+\tilde{\boldsymbol{\mu}}$ with $K \in \mathbb{R}$ and $\tilde{\boldsymbol{\mu}} \in \boldsymbol{b}_{\max }^{\perp}$, we have

$$
\begin{align*}
& \boldsymbol{\mu} J\left(\boldsymbol{\lambda}^{*}\right)^{t} \boldsymbol{b}_{\max }=K \boldsymbol{b}_{\max } J\left(\boldsymbol{\lambda}^{*}\right)^{t} \boldsymbol{b}_{\max }+\tilde{\boldsymbol{\mu}} J\left(\boldsymbol{\lambda}^{*}\right)^{t} \boldsymbol{b}_{\max }  \tag{306}\\
& =K \theta_{\max } \boldsymbol{b}_{\max }{ }^{t} \boldsymbol{b}_{\max }+0 \quad \text { (by the assumption) }  \tag{307}\\
& =\theta_{\max } K \boldsymbol{b}_{\max }{ }^{t} \boldsymbol{b}_{\max }+\theta_{\max } \tilde{\boldsymbol{\mu}}^{t} \boldsymbol{b}_{\max } \quad\left(\text { by } \tilde{\boldsymbol{\mu}} \in \boldsymbol{b}_{\max }^{\perp}\right)  \tag{308}\\
& =\theta_{\max } \boldsymbol{\mu}^{t} \boldsymbol{b}_{\max }, \tag{309}
\end{align*}
$$

which proves (305).

## Appendix H <br> Proof of Theorem 12

Define $Q^{N} \equiv \boldsymbol{\lambda}^{N} \Phi$. Noting that $\sum_{i=1}^{m} \lambda_{i}^{*} D\left(P^{i} \| Q^{*}\right)=C$ holds by the Kuhn-Tucker condition (4), we have

$$
\begin{align*}
0 & <C-I\left(\boldsymbol{\lambda}^{N}, \Phi\right)=C-\sum_{i=1}^{m} \lambda_{i}^{N} \sum_{j=1}^{n} P_{j}^{i} \log \left(\frac{P_{j}^{i}}{Q_{j}^{*}} \cdot \frac{Q_{j}^{*}}{Q_{j}^{N}}\right)  \tag{310}\\
& =C-\sum_{i=1}^{m}\left(\lambda_{i}^{*}+\mu_{i}^{N}\right) D\left(P^{i} \| Q^{*}\right)+D\left(Q^{N} \| Q^{*}\right)  \tag{311}\\
& =-\sum_{i=1}^{m} \mu_{i}^{N} D\left(P^{i} \| Q^{*}\right)+D\left(Q^{N} \| Q^{*}\right) \tag{312}
\end{align*}
$$

We will evaluate $D\left(Q^{N} \| Q^{*}\right)$ in (312). Defining $R^{N} \equiv Q^{N}-$ $Q^{*}$ with $R_{j}^{N}=Q_{j}^{N}-Q_{j}^{*}, j=1, \ldots, n$, we have

$$
\begin{align*}
& D\left(Q^{N} \| Q^{*}\right)=\sum_{j=1}^{n}\left(Q_{j}^{*}+R_{j}^{N}\right) \log \left(1+\frac{R_{j}^{N}}{Q_{j}^{*}}\right)  \tag{313}\\
& =\sum_{j=1}^{n}\left(Q_{j}^{*}+R_{j}^{N}\right)\left\{\frac{R_{j}^{N}}{Q_{j}^{*}}-\frac{1}{2}\left(\frac{R_{j}^{N}}{Q_{j}^{*}}\right)^{2}+o\left(\left\|R^{N}\right\|^{2}\right)\right\} \tag{314}
\end{align*}
$$

$$
\begin{equation*}
=\sum_{j=1}^{n} R_{j}^{N}+\sum_{j=1}^{n} \frac{\left(R_{j}^{N}\right)^{2}}{Q_{j}^{*}}-\frac{1}{2} \sum_{j=1}^{n} \frac{\left(R_{j}^{N}\right)^{2}}{Q_{j}^{*}}+o\left(\left\|R^{N}\right\|^{2}\right) \tag{315}
\end{equation*}
$$

$$
\begin{equation*}
=0+\frac{1}{2} \sum_{j=1}^{n} \frac{1}{Q_{j}^{*}}\left(\sum_{i=1}^{m} \mu_{i}^{N} P_{j}^{i}\right)^{2}+o\left(\left\|\boldsymbol{\mu}^{N}\right\|^{2}\right) \tag{316}
\end{equation*}
$$

Let $\mathcal{I}_{\text {III }}=\emptyset$. Then, we have $D\left(P^{i} \| Q^{*}\right)=C, i=1, \ldots, m$. By $\sum_{i=1}^{m} \mu_{i}^{N}=0$, the first term of (312) is 0 . Therefore, if $\left\|\boldsymbol{\mu}^{N}\right\|<K_{1} \cdot(\theta)^{N}$, then we obtain (203) by (316). Meanwhile, if $\lim _{N \rightarrow \infty} N \mu_{i}^{N}=\sigma_{i}$, then we obtain (204) by (316).

Next, let $\mathcal{I}_{\text {III }} \neq \emptyset$. By (312), we have

$$
\begin{equation*}
C-I\left(\boldsymbol{\lambda}^{N}, \Phi\right) \leq \sum_{i=1}^{m}\left|\mu_{i}^{N}\right| D\left(P^{i} \| Q^{*}\right)+O\left(\left\|\boldsymbol{\mu}^{N}\right\|^{2}\right) . \tag{317}
\end{equation*}
$$

Therefore, if $\left\|\boldsymbol{\mu}^{N}\right\|<K_{1} \cdot(\theta)^{N}$, then we obtain (205) by (317). Meanwhile, if $\lim _{N \rightarrow \infty} N \mu_{i}^{N}=\sigma_{i}$, then we obtain (206) by (317).

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