## Hol onomic rank of A-hyper geometric <br> differential-differ ence equations

| 著者 | Takayana Nobuki, Ohar a Kat suyoshi |
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# Holonomic rank of $\mathcal{A}$-hypergeometric differential-difference equations 

Katsuyoshi Ohara ${ }^{\text {a }}$, Nobuki Takayama ${ }^{\text {b }}$<br>${ }^{\text {a }}$ Department of Computational Science, Kanazawa University<br>${ }^{\mathrm{b}}$ Department of Mathematics, Kobe University


#### Abstract

We introduce $\mathcal{A}$-hypergeometric differential-difference equation $\boldsymbol{H}_{A}$ and prove that its holonomic rank is equal to the normalized volume of $\mathcal{A}$ with the Gröbner basis theory and giving a set of convergent series solutions.


Key words: Ring of differential-difference operators, Hypergeometric functions, differential-differece equations
1991 MSC: 16S32, 33C70, 33C10, 39A99

## 1 Introduction

In this paper, we introduce $\mathcal{A}$-hypergeometric differential-difference equation $\boldsymbol{H}_{A}$ and study its series solutions and holonomic rank.

Let $A=\left(a_{i j}\right)_{i=1, \ldots, d, j=1, \ldots, n}$ be a $d \times n$-matrix whose elements are integers. We suppose that the set of the column vectors of $A$ spans $\mathbf{Z}^{d}$ and there is no zero column vector. Let $a_{i}$ be the $i$-th column vector of the matrix $A$ and $F(\beta, x)$ the integral

$$
F(\beta, x)=\int_{C} \exp \left(\sum_{i=1}^{n} x_{i} t^{a_{i}}\right) t^{-\beta-1} d t, \quad t=\left(t_{1}, \ldots, t_{d}\right), \beta=\left(\beta_{1}, \ldots, \beta_{d}\right)
$$

The integral $F(\beta, x)$ satisfies the $\mathcal{A}$-hypergeometric differential system associated to $A$ and $\beta$ "formally". We use the word "formally" because, there is no general and rigorous description about the cycle $C$ ([11, p.222]).

Email addresses: ohara@kanazawa-u.ac.jp (Katsuyoshi Ohara), takayama@math.kobe-u.ac.jp (Nobuki Takayama).

We will regard the parameters $\beta$ as variables. Then, the function $F(s, x)$ on the $(s, x)$ space satisfies differential-difference equations "formally", which will be our $\mathcal{A}$-hypergeometric differential-difference system defined in Section 3.

Rank theories of $\mathcal{A}$-hypergeometric differential system have been developed since Gel'fand, Zelevinsky and Kapranov [4]. In the end of 1980's, under the condition that the points lie on a same hyperplane, they proved that the rank of $\mathcal{A}$-hypergeometric differential system $H_{A}(\beta)$ agrees with the normalized volume of $A$ for any parameter $\beta \in \mathbf{C}^{d}$ if the toric ideal $I_{A}$ has the CohenMacaulay property. After their result had been gotten, many people have studied on conditions such that the rank equals the normalized volume. In particular, Matusevich, Miller and Walther proved that $I_{A}$ has the CohenMacaulay property if the rank of $H_{A}(\beta)$ agrees with the normalized volume of $A$ for any $\beta \in \mathbf{C}^{d}$ ([5]).

In this paper, we will introduce $\mathcal{A}$-hypergeometric differential-difference system, which can be regarded as a generalization of difference equation for the $\Gamma$-function, the Beta function, and the Gauss hypergeometric difference equations. As the first step on this differential-difference system, we will prove our main Theorem 3 (rank=volume) utilizing the Gröbner basis, theorems on $\mathcal{A}$-hypergeometric differential equations, construction of convergent series solutions with a homogenization technique, uniform convergence of series solutions, and Mutsumi Saito's results for contiguity relations [9], [10], [11, Chapter 4]. The existence theorem 2 on a fundamental set of convergent series solutions for $\mathcal{A}$-hypergeometric differential equation for generic $\beta$ is the second main theorem of our paper. Finally, we note that, for studying our $\mathcal{A}$-hypergeometric differential-difference system, we wrote a program "yang" (Yet another non-commutative Gröbner package) ([6], [8]) on a computer algebra system Risa/Asir and did several experiments on computers to conjecture and prove our theorems.

## 2 Holonomic rank

Let $\boldsymbol{D}$ be the ring of differential-difference operators

$$
\mathbf{C}\left\langle x_{1}, \ldots, x_{n}, s_{1}, \ldots, s_{d}, \partial_{1}, \ldots, \partial_{n}, S_{1}, \ldots, S_{d}, S_{1}^{-1}, \ldots, S_{d}^{-1}\right\rangle
$$

where the following (non-commutative) product rules are assumed

$$
S_{i} s_{i}=\left(s_{i}+1\right) S_{i}, \quad S_{i}^{-1} s_{i}=\left(s_{i}-1\right) S_{i}^{-1}, \quad \partial_{i} x_{i}=x_{i} \partial_{i}+1
$$

and the other types of the product of two generators commute.
Holonomic rank of a system of differential-difference equations will be defined
by using the following ring of differential-difference operators with rational function coefficients

$$
\mathbf{U}=\mathbf{C}\left(s_{1}, \ldots, s_{d}, x_{1}, \ldots, x_{n}\right)\left\langle S_{1}, \ldots, S_{d}, S_{1}^{-1}, \ldots, S_{d}^{-1}, \partial_{1}, \ldots, \partial_{n}\right\rangle
$$

It is a $\mathbf{C}$-algebra generated by rational functions in $s_{1}, \ldots, s_{d}, x_{1}, \ldots, x_{n}$ and differential operators $\partial_{1}, \ldots, \partial_{n}$ and difference operators $S_{1}, \ldots, S_{d}, S_{1}^{-1}, \ldots, S_{d}^{-1}$. The commutation relations are defined by $\partial_{i} c(s, x)=c(s, x) \partial_{i}+\frac{\partial c}{\partial x_{i}}, S_{i} c(s, x)=$ $c\left(s_{1}, \ldots, s_{i}+1, \ldots, s_{d}, x\right) S_{i}, S_{i}^{-1} c(s, x)=c\left(s_{1}, \ldots, s_{i}-1, \ldots, s_{d}, x\right) S_{i}^{-1}$.

Let $I$ be a left ideal in $\boldsymbol{D}$. The holonomic rank of $I$ is the number

$$
\operatorname{rank}(I)=\operatorname{dim}_{\mathbf{C}(s, x)} \mathbf{U} /(\mathbf{U} I) .
$$

In case of the ring of differential operators $(d=0)$, the definition of the holonomic rank agrees with the standard definition of holonomic rank in the ring of differential operators.

For a given left ideal $I$, the holonomic rank can be evaluated by a Gröbner basis computation in $\mathbf{U}$.

## $3 \mathcal{A}$-hypergeometric differential-difference equations

Let $A=\left(a_{i j}\right)_{i=1, \ldots, d, j=1, \ldots, n}$ be an integer $d \times n$ matrix of rank $d$. We assume that the column vectors $\left\{a_{i}\right\}$ of $A$ generates $\mathbf{Z}^{d}$ and there is no zero vector. The $\mathcal{A}$-hypergeometric differential-difference system $\boldsymbol{H}_{A}$ is the following system of differential-difference equations

$$
\begin{aligned}
\left(\sum_{j=1}^{n} a_{i j} x_{j} \partial_{j}-s_{i}\right) \bullet f=0 & \text { for } i=1, \ldots, d \quad \text { and } \\
\left(\partial_{j}-\prod_{i=1}^{n} S_{i}^{-a_{i j}}\right) \bullet f=0 & \text { for } j=1, \ldots, n
\end{aligned}
$$

Note that $\boldsymbol{H}_{A}$ contains the toric ideal $I_{A}$. (use [12, Algorithm 4.5] to prove it.)

Definition 1 Define the unit volume in $\mathbf{R}^{d}$ as the volume of the unit simplex $\left\{0, e_{1}, \ldots, e_{d}\right\}$. For a given set of points $\mathcal{A}=\left\{a_{1}, \ldots, a_{n}\right\}$ in $\mathbf{R}^{d}$, the normalized volume $\operatorname{vol}(\mathcal{A})$ is the volume of the convex hull of the origin and $\mathcal{A}$.

Theorem $1 \mathcal{A}$-hypergeometric differential-difference system $\boldsymbol{H}_{A}$ has linearly independent $\operatorname{vol}(A)$ series solutions.

The proof of this theorem is divided into two parts. The matrix $A$ is called homogeneous when it contains a row of the form $(1, \ldots, 1)$. If $A$ is homogeneous, then the associated toric ideal $I_{A}$ is homogeneous ideal [12]. The first part is the case that $A$ is homogeneous. The second part is the case that $A$ is not homogeneous.

Proof. ( $A$ is homogeneous.) We will prove the theorem with the homogeneity assumption of $A$. In other words, we suppose that $A$ is written as follows:

$$
A=\binom{1 \cdots}{\cdots}
$$

Gel'fand, Kapranov, Zelevinski gave a method to construct $m=\operatorname{vol}(A)$ linearly independent solutions of $H_{A}(\beta)$ with the homogeneity condition of $A$ ([4]). They suppose that $\beta$ is fixed as a generic $\mathbf{C}$-vector. Let us denote their series solutions by $f_{1}(\beta ; x), \ldots, f_{m}(\beta ; x)$. It is easy to see that the functions $f_{i}(s ; x)$ are solutions of the differential-difference equations $\boldsymbol{H}_{A}$. We can show, by carefully checking the estimates of their convergence proof, that there exists an open set in the $(s, x)$ space such that $f_{i}(s ; x)$ is locally uniformly convergent with respect to $s$ and $x$. Let us sketch their proof to see that their series converge as solutions of $\boldsymbol{H}_{A}$. The discussion is given in [4], but we need to rediscuss it in a suitable form to apply it to the case of inhomogeneous $A$.

Let $B$ be a matrix of which the set of column vectors is a basis of $\operatorname{Ker}(A$ : $\mathbf{Q}^{n} \rightarrow \mathbf{Q}^{d}$ ) and is normalized as follows:

$$
B=\left(\begin{array}{lll}
1 & & \\
& \ddots & \\
& & 1 \\
& & \\
& *
\end{array}\right) \in M(n, n-d, \mathbf{Q})
$$

We denote by $b^{(i)}$ the $i$-th column vector of $B$ and by $b_{i j}$ the $j$-th element of $b^{(i)}$. Then the homogeneity of $A$ implies

$$
\sum_{j=1}^{n} b_{i j}=0
$$

Let us fix a regular triangulation $\Delta$ of $\mathcal{A}=\left\{a_{1}, \ldots, a_{n}\right\}$ following the construction by Gel'fand, Kapranov, Zelevinsky. Take a $d$-simplex $\tau$ in the triangulation $\Delta$. If $\lambda \in \mathbf{C}^{n}$ is admissible for a $d$-simplex $\tau$ of $\{1,2, \ldots, n\}$ (admissible $\Leftrightarrow$ for all $j \notin \tau, \lambda_{j} \in \mathbf{Z}$ ), and $A \lambda=s$ holds, then $\boldsymbol{H}_{A}$ has a formal series
solution

$$
\phi_{\tau}(\lambda ; x)=\sum_{l \in L} \frac{x^{\lambda+l}}{\Gamma(\lambda+l+1)},
$$

where $L=\operatorname{Ker}\left(A: \mathbf{Z}^{n} \rightarrow \mathbf{Z}^{d}\right)$ and $\Gamma(\lambda+l+1)=\prod_{i=1}^{n} \Gamma\left(\lambda_{i}+l_{i}+1\right)$ and when a factor of the denominator of a term in the sum, we regard the term is zero. Put $\# \tau=n^{\prime}$. Note that there exists an open set $U$ in the $s$ space such that $\lambda_{i}, i \in \tau$ lie in a compact set in $\mathbf{C}^{n^{\prime}} \backslash \mathbf{Z}^{n^{\prime}}$. Moreover, this open set $U$ can be taken as a common open set for all $d$-simplices in the triangulation $\Delta$ and the associated admissible $\lambda$ 's when the integral values $\lambda_{j}(j \notin \tau)$ are fixed for all $\tau \in \Delta$.

Put $L^{\prime}=\left\{\left(k_{1}, \ldots, k_{n-d}\right) \in \mathbf{Z}^{n-d} \mid \sum_{i=1}^{n-d} k_{i} b^{(i)} \in \mathbf{Z}^{n}\right\}$. Then, $L^{\prime}$ is $\mathbf{Z}$-submodule of $\mathbf{Z}^{n-d}$ and $L=\left\{\sum_{i=1}^{n-d} k_{i} b^{(i)} \mid k \in L^{\prime}\right\}$. In other words, $L$ can be parametrized with $L^{\prime}$. Without loss of the generality, we may suppose that $\tau=\{n-d+$ $1, \ldots, n\}$. Then, we have

$$
\phi_{\tau}(\lambda ; x)=\sum_{l \in L} \frac{x^{\lambda+l}}{\Gamma(\lambda+l+1)}=\sum_{k \in L^{\prime}} \frac{x^{\lambda+\sum_{i=1}^{n-d} k_{i} b^{(i)}}}{\Gamma\left(\lambda+\sum_{i=1}^{n-d} k_{i} b^{(i)}+1\right)}
$$

Note that the first $n-d$ rows of $B$ are normalized. Then, we have

$$
\lambda_{j}+\sum_{i=1}^{n-d} k_{i} b_{i j}+1=\lambda_{j}+k_{j}+1 \in \mathbf{Z} \quad(j=1, \ldots, n-d)
$$

Since $1 / \Gamma(0)=1 / \Gamma(-1)=1 / \Gamma(-2)=\cdots=0$, the sum can be written as

$$
\phi_{\tau}(\lambda ; x)=\sum_{\substack{k \in L^{\prime} \\ \lambda_{j}+k_{j}+1 \in \mathbf{Z}_{>0} \\(j=1, \ldots, n-d)}} \frac{x^{\lambda+\sum_{i=1}^{n-d} k_{i} b^{(i)}}}{\Gamma\left(\lambda+\sum_{i=1}^{n-d} k_{i} b^{(i)}+1\right)}
$$

Moreover, when we put

$$
\begin{aligned}
k_{j}^{\prime} & =\lambda_{j}+k_{j}, \quad(j=1, \ldots, n-d) \\
\lambda^{\prime} & =\lambda-\sum_{i=1}^{n-d} \lambda_{i} b^{(i)} \\
\hat{\lambda} & =\left(\lambda_{1}, \ldots, \lambda_{n-d}\right)
\end{aligned}
$$

we have

$$
\sum_{i=1}^{n-d} k_{i} b^{(i)}=-\sum_{i=1}^{n-d} \lambda_{i} b^{(i)}+\sum_{i=1}^{n-d} k_{i}^{\prime} b^{(i)}
$$

Hence, the sum $\phi_{\tau}(\lambda ; x)$ can be written as

$$
\begin{aligned}
\phi_{\tau}(\lambda ; x) & =\sum_{\substack{k^{\prime} \in L^{\prime}+\hat{\lambda} \\
k^{\prime} \in \mathbf{Z}_{\geq 0}^{n-d}}} \frac{x^{\lambda-\sum_{i=1}^{n-d} \lambda_{i} b^{(i)}} \cdot x^{\sum_{i=1}^{n-d} k_{i}^{\prime} b^{(i)}}}{\Gamma\left(\lambda-\sum_{i=1}^{n-d} \lambda_{i} b^{(i)}+\sum_{i=1}^{n-d} k_{i}^{\prime} b^{(i)}+1\right)} \\
& =x^{\lambda^{\prime}} \sum_{\substack{k^{\prime} \in L^{\prime}+\hat{\lambda} \\
k^{\prime} \in \mathbf{Z}_{\geq 0}^{n-d}}} \frac{\left(x^{b^{(1)}}\right)^{k_{1}^{\prime}} \cdots\left(x^{b^{(n-d)}}\right)^{k_{n-d}^{\prime}}}{\Gamma\left(\lambda^{\prime}+\sum_{i=1}^{n-d} k_{i}^{\prime} b^{(i)}+1\right)}
\end{aligned}
$$

Note that our series with the coefficients in terms of Gamma functions agree with those in $[11, \S 3.4]$, which do not contain Gamma functions, by multiplying suitable constants. Hence we will apply some results on series solutions in [11] to our discussions in the sequel.

Lemma 1 Let $\left(k_{i}\right) \in\left(\mathbf{Z}_{\geq 0}\right)^{m}$ and $\left(b_{i j}\right) \in M(m, n, \mathbf{Q})$. Suppose that

$$
\sum_{i=1}^{m} k_{i} b_{i j} \in \mathbf{Z}, \quad \sum_{j=1}^{n} b_{i j}=0
$$

and parameters $\lambda=\left(\lambda_{1}, \ldots, \lambda_{n}\right)$ belongs to a compact set $K$. Then there exists a positive number $r$, which is independent of $\lambda$, such that the power series

$$
\sum_{\substack{k^{\prime} \in L^{\prime}+\hat{\lambda} \\ k^{\prime} \in \mathbf{Z}_{\geq 0}^{n-d}}} \frac{\left(x^{b^{(1)}}\right)^{k_{1}^{\prime}} \cdots\left(x^{b^{(n-d)}}\right)^{k_{n-d}^{\prime}}}{\Gamma\left(\lambda^{\prime}+\sum_{i=1}^{n-d} k_{i}^{\prime} b^{(i)}+1\right)}
$$

is convergent in $\left|x^{b^{(1)}}\right|, \cdots,\left|x^{b^{(n-d)}}\right|<r$.
The proof of this lemma can be done by elementary estimates of $\Gamma$ functions. See [7, pp.18-21] if readers are interested in the details. Since

$$
k^{\prime} \in L^{\prime}+\hat{\lambda} \Longleftrightarrow \sum_{i=1}^{n-d} k_{i}^{\prime} b^{(i)} \in \mathbf{Z}^{n}
$$

it follows from Lemma 1 that there exists a positive constant $r$ such that the series converge in

$$
\begin{equation*}
\left|x^{b^{(1)}}\right|, \cdots,\left|x^{b^{(n-d)}}\right|<r \tag{3.1}
\end{equation*}
$$

for any $s$ in the open set $U$. We may suppose $r<1$. Take the $\log$ of (3.1). Then we have

$$
\begin{equation*}
b^{(k)} \cdot\left(\log \left|x_{1}\right|, \ldots, \log \left|x_{n}\right|\right)<\log |r|<0 \quad \forall k \in\{1, \ldots, n-d\} \tag{3.2}
\end{equation*}
$$

Following [4], for the simplex $\tau$ and $r$, we define the set $C(A, \tau, r)$ as follows.

$$
C(A, \tau, r)=\left\{\psi \in \mathbf{R}^{n} \mid \exists \varphi \in \mathbf{R}^{d}, \quad \psi_{i}-\left(\varphi, a_{i}\right)\left\{\begin{array}{ll}
>-\log |r|, & i \notin \tau \\
=0, & i \in \tau
\end{array}\right\}\right.
$$

The condition (3.2) and $\left(-\log \left|x_{1}\right|, \ldots,-\log \left|x_{n}\right|\right) \in C(A, \tau, r)$ is equivalent (see $[3$, section 4$]$ as to the proof).

Since $\Delta$ is a regular triangulation of $A, \bigcap_{\tau \in \Delta} C(A, \tau, r)$ is an open set. Therefore, when $s$ lies in the open set $U$ and $-\log |x|$ lies in the above open set, the $\operatorname{vol}(A)$ linearly independent solutions converge.

Let us proceed on the proof for the inhomogeneous case. We suppose that $A$ is not homogeneous and has only non-zero column vectors. We define the homogenized matrix as

$$
\tilde{A}=\left(\begin{array}{cccc}
1 & \cdots & 1 & 1 \\
a_{11} & \cdots & a_{1 n} & 0 \\
\vdots & & \vdots & \vdots \\
a_{d 1} & \cdots & a_{d n} & 0
\end{array}\right) \in M(d+1, n+1, \mathbf{Z})
$$

For $s=\left(s_{1}, \ldots, s_{n}\right) \in \mathbf{C}^{d}$ and a generic complex number $s_{0}$, we put $\tilde{s}=$ $\left(s_{0}, s_{1}, \ldots, s_{d}\right)$. We suppose that $\tau=\{n-d+1, \ldots, d, d+1\}$ is a $(d+1)$-simplex. Let us take an admissible $\lambda$ for $\tau$ such that $\tilde{A} \tilde{\lambda}=\tilde{s}$ and $\tilde{\lambda}=\left(\lambda_{1}, \ldots, \lambda_{n+1}\right) \in$ $\mathbf{R}^{n+1}$ as in the proof of the homogeneous case. Put $\lambda=\left(\lambda_{1}, \ldots, \lambda_{n}\right)$. Consider the solution of the hypergeometric system for $\tilde{A}$

$$
\tilde{\phi}_{\tau}(\tilde{\lambda} ; \tilde{x})=\sum_{k^{\prime} \in L^{\prime} \cap S} \frac{\tilde{x}^{\lambda+\sum_{i=1}^{n-d} k_{i}^{\prime} b^{(i)}}}{\Gamma\left(\lambda+\sum_{i=1}^{n-d} k_{i}^{\prime} b^{(i)}+1\right)}
$$

and the series

$$
\phi_{\tau}(\lambda ; x)=\sum_{k^{\prime} \in L^{\prime} \cap S} \frac{\prod_{j=1}^{n} x_{j}^{\lambda+\sum_{i=1}^{n-d} k_{i}^{\prime} b_{i j}}}{\prod_{j=1}^{n} \Gamma\left(\lambda_{j}+\sum_{i=1}^{n-d} k_{i}^{\prime} b_{i j}+1\right)}
$$

$\left(\tilde{x}=\left(x_{1}, \ldots, x_{n+1}\right), x=\left(x_{1}, \ldots, x_{n}\right)\right)$. Here, the set $S$ is a subset of $L^{\prime}$ such that an integer in $\mathbf{Z}_{\leq 0}$ does not appear in the arguments of the Gamma functions in the denominator. We note that $L^{\prime}$ for $\tilde{A}$ and $L^{\prime}$ for $A$ agree, which can be proved as follows. Let $\left(k_{1}, \ldots, k_{n+1}\right)$ be in the kernel of $\tilde{A}$ in $\mathbf{Q}^{n+1}$. Since $\tilde{A}$ contains the row of the form $(1, \ldots, 1)$, then $\left(k_{1}, \ldots, k_{n}\right) \in \mathbf{Z}^{n}$ implies that $k_{n+1}$ is an integer. The conclusion follows from the definition of $L^{\prime}$.

Definition 2 We call $\phi_{\tau}(\lambda ; x)$ the dehomogenization of $\tilde{\phi}_{\tau}(\tilde{\lambda} ; \tilde{x})$.
Intuitively speaking, the dehomogenization is defined by "forgetting" the last variable $x_{n+1}$ associated $\Gamma$ factors. See Example 1.

Formal series solutions for the hypergeometric system for inhomogeneous $A$ do not converge in general. However, we can construct $\operatorname{vol}(A)$ convergent series
solutions as the dehomogenization of a set of series solutions for $\tilde{A}$ hypergeometric system associated to a regular triangulation on $\tilde{\mathcal{A}}$ induced by a "nice" weight vector $\tilde{w}(\varepsilon)$, which we will define. Put $\tilde{w}=(1, \ldots, 1,0) \in \mathbf{R}^{n+1}$. Since the Gröbner fan for the toric variety $I_{\tilde{A}}$ is a polyhedral fan, the following fact holds.

Lemma 2 For any $\varepsilon>0$, there exists $\tilde{v} \in \mathbf{R}^{n+1}$ such that $\tilde{w}(\varepsilon):=\tilde{w}+\varepsilon \tilde{v}$ lies in the interior of a maximal dimensional Gröbner cone of $I_{\tilde{A}}$. We may also suppose $\tilde{v}_{n+1}=0$.

Proof. Let us prove the lemma. The first part is a consequence of an elementary property of the fan. When $I$ is a homogeneous ideal in the ring of polynomials of $n+1$ variables, we have

$$
\begin{equation*}
\operatorname{in}_{\tilde{u}}(I)=\operatorname{in}_{\tilde{u}+t(1, \cdots, 1)}(I) \tag{3.3}
\end{equation*}
$$

for any $t$ and any weight vector $\tilde{u}$. In other words, $\tilde{u}$ and $\tilde{u}+t(1, \ldots, 1)$ lie in the interior of the same Gröbner cone.

When the weight vector $\tilde{w}(\varepsilon)=\tilde{w}+\varepsilon \tilde{v}$ lies in the interior of the Gröbner cone, we define a new $\tilde{v}$ by $\tilde{v}-\tilde{v}_{n+1}(1, \ldots, 1)$. Since the initial ideal does not change with this change of weight, we may assume that $\tilde{v}_{n+1}=0$ for the new $\tilde{v}$.

Since the Gröbner fan is a refinement of the secondary fan and hence $\tilde{w}(\varepsilon)$ is an interior point of a maximal dimensional secondary cone, it induces a regular triangulation ([12] p.71, Proposition 8.15). We denote by $\Delta$ the regular triangulation on $\tilde{A}$ induced by $\tilde{w}(\varepsilon)$. For a $d$-simplex $\tau \in \Delta$, we define $b^{(i)}$ as in the proof of the homogeneous case. Since the weight for $\tilde{a}_{n+1}$ is the lowest, $n+1 \in \tau$ holds. We can change indices of $\tilde{a}_{1}, \ldots, \tilde{a}_{n}$ so that $\tau=$ $\{n-d+1, \ldots, n+1\}$ without loss of generality.

Let us prove that the dehomogenized series $\phi_{\tau}(\lambda ; x)$ converge. It follows from a characterization of the support of the series [11, Theorem 3.4.2] that we have

$$
\tilde{w}(\varepsilon) \cdot\left(\sum_{i=1}^{n-d} k_{i}^{\prime} b^{(i)}+\lambda\right) \geq \tilde{w}(\varepsilon) \cdot \lambda, \quad \forall k^{\prime} \in L^{\prime} \cap S
$$

Here, $S$ is a set such that $\mathbf{Z}_{\leq 0}$ does not appear in the denominator of the $\Gamma$ factors. Take the limit $\varepsilon \rightarrow 0$ and we have

$$
\tilde{w} \cdot \sum_{i=1}^{n-d} k_{i}^{\prime} b^{(i)} \geq 0, \quad \forall k^{\prime} \in L^{\prime} \cap S
$$

From Lemma 2, $\tilde{w}(\varepsilon) \in C(\tilde{A}, \tau)$ holds and then

$$
\tilde{w}(\varepsilon) \cdot b^{(i)} \geq 0
$$

Similarly, by taking the limit $\varepsilon \rightarrow 0$, we have

$$
\tilde{w} \cdot b^{(i)}=\sum_{j=1}^{n} b_{i j} \geq 0 .
$$

Therefore, we have $\sum_{j=1}^{n+1} b_{i j}=0$, the inequality $b_{i, n+1} \leq 0$ holds for all $i$.
Since $k_{1}^{\prime} \geq-\lambda_{1}, \ldots, k_{n-d}^{\prime} \geq-\lambda_{n-d}$, we have

$$
\sum_{i=1}^{n-d} k_{i}^{\prime} b_{i, n+1} \leq-\sum_{i=1}^{n-d} \lambda_{i} b_{i, n+1}
$$

Note that the right hand side is a non-negative number. Suppose that $\lambda_{n+1}$ is negative. In terms of the Pochhammer symbol we have $\Gamma\left(\lambda_{n+1}-m\right)=$ $\Gamma\left(\lambda_{n+1}\right)\left(-\lambda_{n+1}+1 ; m\right)^{-1}(-1)^{m}$, then we can estimate the $(n+1)$-th gamma factors as

$$
\begin{align*}
\left|\Gamma\left(\lambda_{n+1}+\sum_{i=1}^{n-d} k_{i}^{\prime} b_{i, n+1}+1\right)\right| & =\left|\Gamma\left(\lambda_{n+1}+1\right)\right| \cdot\left|\left(-\lambda_{n+1} ;-\sum_{i=1}^{n-d} k_{i}^{\prime} b_{i, n+1}\right)\right|^{-1} \\
& \leq c^{\prime}\left|\Gamma\left(\lambda_{n+1}+1\right)\right| \cdot\left|\left(-\lambda_{n+1} ;-\sum_{i=1}^{n-d} \lambda_{i} b_{i, n+1}\right)\right|^{-1} \\
& =c \tag{3.4}
\end{align*}
$$

Here, $c^{\prime}$ and $c$ are suitable constants.
When $\lambda_{n+1} \geq 0$, there exists only finite set of values such that $\lambda_{n+1}+$ $\sum_{i=1}^{n-d} k_{i}^{\prime} b_{i, n+1} \geq 0$. Then, we can show the inequality (3.4) in an analogous way.

Now, by (3.4), we have

$$
\left|\frac{1}{\prod_{j=1}^{n} \Gamma\left(\lambda_{j}+\sum_{i=1}^{n-d} k_{i}^{\prime} b_{i j}+1\right)}\right| \leq c\left|\frac{1}{\Gamma\left(\lambda+\sum_{i=1}^{n-d} k_{i}^{\prime} b^{(i)}+1\right)}\right|
$$

We note that the right hand side is the coefficient of the series solution for the homogeneous system for $\tilde{A}$ and the series converge for $\left(-\log \left|x_{1}\right|, \ldots,-\log \left|x_{n+1}\right|\right) \in$ $C(\tilde{A}, \tau, r)(r<1)$ uniformly with respect to $\tilde{s}$ in an open set.

Put $x_{n+1}=1$. Since $-\log \left|x_{n+1}\right|=0$ and $\tilde{w}(\varepsilon) \in\left\{y \mid y_{n+1}=0\right\}$, we can see that

$$
\bigcap_{\tau \in \Delta} C(\tilde{A}, \tau, r) \cap\left\{y \mid y_{n+1}=0\right\}
$$

is a non-empty open set of $\mathbf{R}^{n}$. Therefore the dehomogenized series $\phi_{\tau}(\lambda ; x)$ converge in an open set in the $(s, x)$ space.

Theorem 2 The dehomogenized series $\phi_{\tau}(\lambda ; x)$ satisfies the hypergeometric differential-difference system $\boldsymbol{H}_{A}$ and they are linearly independent convergent solutions of $\boldsymbol{H}_{A}$ when $\lambda$ runs over admissible exponents associated to the initial system induced by the weight vector $\tilde{w}(\varepsilon)$.

Proof. Since $A \lambda=s$, it is easy to show that they are formal solutions of the differential-difference system $\boldsymbol{H}_{A}$. We will prove that we can construct $m$ linearly independent solutions. We note that the weight vector $\tilde{w}(\varepsilon)=$ $(1, \ldots, 1,0)+\varepsilon v \in \mathbf{R}^{n+1}$ is in the neighborhood of $(1, \ldots, 1,0) \in \mathbf{R}^{n+1}$ and in the interior of a maximal dimensional Gröbner cone of $I_{\tilde{A}}$.

It follows from [11, p.119] that the minimal generating set of $\operatorname{in}_{(1, \ldots, 1,0)} I_{\tilde{A}}$ does not contain $\partial_{n+1}$. Since

$$
\mathrm{in}_{\tilde{w}(\varepsilon)} I_{\tilde{A}}=\operatorname{in}_{v}\left(\mathrm{in}_{(1, \ldots, 1,0)} I_{\tilde{A}}\right)
$$

does not contain $\partial_{n+1}$, we have

$$
M=\left\langle\operatorname{in}_{\tilde{w}(\varepsilon)} I_{\tilde{A}}\right\rangle=\left\langle\mathrm{in}_{w(\varepsilon)} I_{A}\right\rangle \quad \text { in } \mathbf{C}\left[\partial_{1}, \ldots, \partial_{n+1}\right] .
$$

Here, we define $w(\varepsilon)$ with $\tilde{w}(\varepsilon)=(w(\varepsilon), 0)$. Put $\tilde{\theta}=\left(\theta_{1}, \ldots, \theta_{n+1}\right)$. From [11, Theorem 3.1.3], for generic $\tilde{\beta}=\left(\beta_{0}, \beta\right), \beta \in \mathbf{C}^{d}$, the initial ideal $\operatorname{in}_{(-\tilde{w}(\varepsilon), \tilde{w}(\varepsilon))} H_{\tilde{A}}(\tilde{\beta})$ is generated by $\operatorname{in}_{\tilde{w}(\varepsilon)}\left(I_{\tilde{A}}\right)$ and $\tilde{A} \tilde{\theta}-\tilde{\beta}$. Let us denote by $T(M)$ the standard pairs of $M$. From [11, Theorem 3.2.10], the initial ideal

$$
\begin{equation*}
\left\langle\mathrm{in}_{\tilde{w}(\varepsilon)} I_{\tilde{A}}, \tilde{A} \tilde{\theta}-\tilde{\beta}\right\rangle \tag{3.5}
\end{equation*}
$$

has $\# T(M)=\operatorname{vol}(\tilde{A})$ linearly independent solutions of the form

$$
\left\{\tilde{x}^{\tilde{}} \mid\left(\partial^{a}, T\right) \in T(M)\right\}
$$

Here, $\tilde{\lambda}$ is defined by $\tilde{\lambda}_{i}=a_{i} \in \mathbf{Z}_{>0}, \forall i \notin T$ and $\tilde{A} \tilde{\lambda}=\tilde{\beta}$. Note that $\tilde{\lambda}$ is admissible for the $d$-simplex $T$.

Since we have

$$
\left\langle\mathrm{in}_{\tilde{w}(\varepsilon)} I_{\tilde{A}}\right\rangle=\left\langle\mathrm{in}_{w(\varepsilon)} I_{A}\right\rangle
$$

the difference between

$$
\begin{equation*}
\left\langle\mathrm{in}_{w(\varepsilon)} I_{A}, A \theta-\beta\right\rangle \tag{3.6}
\end{equation*}
$$

and (3.5) is only

$$
\theta_{1}+\cdots+\theta_{n}+\theta_{n+1}-\beta_{0}
$$

and other equations do not contain $x_{n+1}, \partial_{n+1}$.
For any $\left(\partial^{a}, T\right) \in T(M)$, we have $n+1 \in T$. Therefore, the two solution spaces (3.6) and (3.5) are isomorphic under the correspondence

$$
\begin{equation*}
x^{\lambda} \mapsto \tilde{x}^{\tilde{\lambda}} \tag{3.7}
\end{equation*}
$$

Here, we put $\tilde{\lambda}=\left(\lambda, \lambda_{n+1}\right)$ and $\lambda_{n+1}$ is defined by

$$
\sum_{i=1}^{n} \lambda_{i}+\lambda_{n+1}-\beta_{0}=0
$$

It follows from [11, Theorem 2.3.11 and Theorem 3.2.10] that

$$
\left\{\tilde{x}^{\tilde{\lambda}} \mid\left(\partial^{a}, T\right) \in T(M)\right\}
$$

are C-linearly independent. Therefore, from the correspondence (3.7), the functions

$$
\left\{x^{\lambda} \mid\left(\partial^{a}, T\right) \in T(M)\right\},
$$

of which cardinality is $\operatorname{vol}(A)$, are C-linearly independent. Hence, series solutions with the initial terms

$$
\left\{\left.\frac{x^{\lambda}}{\Gamma(\lambda+1)} \right\rvert\,\left(\partial^{a}, T\right) \in T(M)\right\}
$$

are $\mathbf{C}$ linearly independent, which implies the linear independence of series solutions with these starting terms [11]. We have completed the proof of the theorem and also that of Theorem 1.

Theorem 3 The holonomic rank of $\boldsymbol{H}_{A}$ is equal to the normalized volume of A.

Proof. First we will prove $\operatorname{rank}\left(\boldsymbol{H}_{A}\right) \leq \operatorname{vol}(A)$. It follows from the Adolphson's theorem ([1]) that the holonomic rank of $\mathcal{A}$-hypergeometric system $H_{A}(\beta)$ is equal to the normalized volume of $A$ for generic parameters $\beta$. It implies that the standard monomials for a Gröbner basis of the $\mathcal{A}$-hypergeometric system $H_{A}(s)$ in $\mathbf{C}(s, x)\left\langle\partial_{1}, \ldots, \partial_{n}\right\rangle$ consists of $\operatorname{vol}(A)$ elements. We note that elements in the Gröbner basis can be regarded as an element in the ring of differentialdifference operators with rational function coefficients $\mathbf{U}$. We denote by $\partial_{j}$ and $r_{j}$ the creation and annihilation operators. The existence of them are proved in [10, Chapter 4]. Then, we have

$$
H_{j}=\partial_{j}-\prod_{i=1}^{n} S_{i}^{-a_{i j}} \in \boldsymbol{H}_{A}
$$

and

$$
B_{j}=r_{j}-\prod_{i=1}^{n} S_{i}^{a_{i j}} \in \boldsymbol{H}_{A}, \quad r_{j} \in \mathbf{C}(s, x)\left\langle\partial_{1}, \ldots, \partial_{n}\right\rangle .
$$

Since the column vectors of $A$ generate the lattice $\mathbf{Z}^{d}$, we obtain from $B_{j}$ 's and $H_{j}$ 's elements of the form $S_{i}-p(s, x, \partial), S_{i}^{-1}-q(s, x, \partial) \in \boldsymbol{H}_{A}$. It implies the number of standard monomials of a Gröbner basis of $\boldsymbol{H}_{A}$ with respect to a block order such that $S_{1}, \ldots, S_{n}>S_{1}^{-1}, \ldots, S_{n}^{-1}>\partial_{1}, \ldots, \partial_{n}$ is less than or equal to $\operatorname{vol}(A)$.

Second, we will prove $\operatorname{rank}\left(\boldsymbol{H}_{A}\right) \geq \operatorname{vol}(A)$. We suppose that $\operatorname{rank}\left(\boldsymbol{H}_{A}\right)<$ $\operatorname{vol}(A)$ and will induce a contradiction. For the block order $S_{1}, \cdots, S_{d}>$ $S_{1}^{-1}, \cdots, S_{d}^{-1}>\partial_{1}, \cdots, \partial_{n}$, we can show that the standard monomials $T$ of a Gröbner basis of $\boldsymbol{H}_{A}$ in $\mathbf{U}$ contains only differential terms and $\# T<\operatorname{vol}(A)$ by the assumption. Let $T^{\prime}$ be the standard monomials of Gröbner basis $G(s)$ of $H_{A}(s)$ in the ring of differential operators with rational function coefficients $D(s)$. Note that $\# T^{\prime}=\operatorname{vol}(A)$. Then $T$ is a proper subset of the set $T^{\prime}$. For $r \in T^{\prime} \backslash T$, it follows that

$$
\partial^{r} \equiv \sum_{\alpha \in T} c_{\alpha}(x, s) \partial^{\alpha} \quad \bmod \boldsymbol{H}_{A}
$$

From Theorem 2, we have convergent series solutions $f_{1}(s, x), \cdots, f_{m}(s, x)$ of $\boldsymbol{H}_{A}$, where $m=\operatorname{vol}(A)$. So,

$$
\begin{equation*}
\partial^{r} \bullet f_{i}=\sum_{\alpha \in T} c_{\alpha}(x, s) \partial^{\alpha} \bullet f_{i} \tag{3.8}
\end{equation*}
$$

Since $f_{1}(s, x), \ldots, f_{m}(s, x)$ are linearly independent, the Wronskian standing for $T^{\prime}$

$$
W\left(T^{\prime} ; f\right)(x, s)=\left|\begin{array}{ccc}
f_{1}(s ; x) & \cdots & f_{m}(\beta ; x) \\
\partial^{\delta} f_{1}(s ; x) & \cdots & \partial^{\delta} f_{m}(\beta ; x) \\
\vdots & \cdots & \vdots
\end{array}\right| \quad\left(\partial^{\delta} \in T^{\prime} \backslash\{1\}\right)
$$

is non-zero for generic number $s$. However $r \in T^{\prime}$ and (3.8) induce the Wronskian $W\left(T^{\prime} ; f\right)(s, x)$ is equal to zero.

Finally, by $\operatorname{rank}\left(\boldsymbol{H}_{A}\right) \leq \operatorname{vol}(A)$ and $\operatorname{rank}\left(\boldsymbol{H}_{A}\right) \geq \operatorname{vol}(A)$, the theorem is proved.

Example 1 Put $A=\left(\begin{array}{lll}1 & 2 & 3\end{array}\right)$ and $\tilde{A}=\left(\begin{array}{lll}1 & 1 & 1\end{array} 1\right.$ $\operatorname{gral}([11, \mathrm{p} .223])$.

The matrix $\tilde{A}$ is homogeneous. For $\tilde{w}(\varepsilon)=(1,1,1,0)+\frac{1}{100}(1,0,0,0)$, the initial ideal $\operatorname{in}_{\tilde{w}(\varepsilon)}\left(I_{\tilde{A}}\right)$ is generated by $\partial_{1}^{2}, \partial_{1} \partial_{2}, \partial_{1} \partial_{3}, \partial_{2}^{3}$. Note that the initial ideal does not contain $\partial_{4}$. We solve the initial system $(\tilde{A} \tilde{\theta}-\tilde{s}) \bullet g=0,\left(\operatorname{in}_{\tilde{w}(\varepsilon)}\left(I_{\tilde{A}}\right)\right) \bullet$ $g=0$. The standard pairs $\left(\partial^{a}, T\right)$ for $\operatorname{in}_{\tilde{w}(\varepsilon)}\left(I_{\tilde{A}}\right)$ are $\left(\partial_{1}^{0} \partial_{2}^{1},\{3,4\}\right),\left(\partial_{1}^{0} \partial_{2}^{0},\{3,4\}\right)$, $\left(\partial_{1}^{0} \partial_{2}^{2},\{3,4\}\right)$. Hence, the solutions for the initial system are
$x_{1}^{0} x_{2}^{1} x_{3}^{\left(s_{1}-2\right) / 3} x_{4}^{s_{0}-1-\left(s_{1}-2\right) / 3}, x_{1}^{0} x_{2}^{0} x_{3}^{s_{1} / 3} x_{4}^{a_{0}-s_{1} / 3}, x_{1}^{0} x_{2}^{2} x_{3}^{\left(s_{1}-4\right) / 3} x_{4}^{s_{0}-2-\left(s_{1}-4\right) / 3}([11])$. Therefore, the $\mathcal{A}$-hypergeometric differential-difference system $\boldsymbol{H}_{\tilde{A}}$ has the following series solutions.

$$
\begin{aligned}
\tilde{\phi}_{1}(\tilde{\lambda}, \tilde{x})= & x_{4}^{s_{0}}\left(\frac{x_{2}}{x_{4}}\right)\left(\frac{x_{3}}{x_{4}}\right)^{\frac{s_{1}-2}{3}} \\
& \cdot \sum_{\substack{k_{1} \geq 0, k_{2} \geq-1 \\
\left(k_{1}, k_{2}\right) \in L^{\prime}}} \frac{\left(x_{1} x_{3}^{-1 / 3} x_{4}^{-2 / 3}\right)^{k_{1}}\left(x_{2} x_{3}^{-2 / 3} x_{4}^{-1 / 3}\right)^{k_{2}}}{k_{1}!\left(k_{2}+1\right)!\Gamma\left(\frac{s_{1}-k_{1}-2 k_{2}+1}{3}\right) \Gamma\left(\frac{3 s_{0}-s_{1}-2 k_{1}-k_{2}+2}{3}\right)} \\
\tilde{\phi}_{2}(\tilde{\lambda}, \tilde{x})= & x_{4}^{s_{0}\left(\frac{x_{3}}{x_{4}}\right)^{\frac{s_{1}}{3}}} \\
& \cdot \sum_{\substack{k_{1} \geq 0, k_{2} \geq 0 \\
\left(k_{1}, k_{2}\right) \in L^{\prime}}} \frac{\left(x_{1} x_{3}^{-1 / 3} x_{4}^{-2 / 3}\right)^{k_{1}}\left(x_{2} x_{3}^{-2 / 3} x_{4}^{-1 / 3}\right)^{k_{2}}}{k_{1}!k_{2}!\Gamma\left(\frac{s_{1}-k_{1}-2 k_{2}+3}{3}\right) \Gamma\left(\frac{3 s_{0}-s_{1}-2 k_{1}-k_{2}+3}{3}\right)} \\
\tilde{\phi}_{3}(\tilde{\lambda}, \tilde{x})= & x_{4}^{s_{0}\left(\frac{x_{2}}{x_{4}}\right)^{2}\left(\frac{x_{3}}{x_{4}}\right)^{\frac{s_{1}-4}{3}}} \begin{array}{l}
\sum_{\substack{k_{1} \geq 0, k_{2} \geq-2 \\
\left(k_{1}, k_{2}\right) \in L^{\prime}}} \frac{\left(x_{1} x_{3}^{-1 / 3} x_{4}^{-2 / 3}\right)^{k_{1}}\left(x_{2} x_{3}^{-2 / 3} x_{4}^{-1 / 3}\right)^{k_{2}}}{k_{1}!\left(k_{2}+2\right)!\Gamma\left(\frac{s_{1}-k_{1}-2 k_{2}-1}{3}\right) \Gamma\left(\frac{3 s_{0}-s_{1}-2 k_{1}-k_{2}+1}{3}\right)}
\end{array}
\end{aligned}
$$

Here,
$L^{\prime}=\left\{\left(k_{1}, k_{2}\right) \mid k_{1} \equiv 0 \bmod 3, k_{2} \equiv 0 \bmod 3\right\} \cup\left\{\left(k_{1}, k_{2}\right) \mid k_{1} \equiv 1 \bmod 3, k_{2} \equiv 1 \bmod 3\right\}$.

The matrix $A$ is not homogeneous and by dehomogenizing the series solution for $\tilde{A}$ we obtain the following series solutions for the $\mathcal{A}$-hypergeometric differential-difference system $\boldsymbol{H}_{A}$.

$$
\begin{aligned}
& \phi_{1}(\lambda, x)=x_{2} x_{3}^{\frac{s_{1}-2}{3}} \sum_{\substack{k_{1} \geq 0, k_{2} \geq-1 \\
\left(k_{1}, k_{2}\right) \in L^{\prime}}} \frac{\left(x_{1} x_{3}^{-1 / 3}\right)^{k_{1}}\left(x_{2} x_{3}^{-2 / 3}\right)^{k_{2}}}{k_{1}!\left(k_{2}+1\right)!\Gamma\left(\frac{s_{1}-k_{1}-2 k_{2}+1}{3}\right)} \\
& \phi_{2}(\lambda, x)=x_{3}^{\frac{s_{1}}{3}} \sum_{\substack{k_{1} \geq 0, k_{2} \geq 0 \\
\left(k_{1}, k_{2}\right) \in L^{\prime}}} \frac{\left(x_{1} x_{3}^{-1 / 3}\right)^{k_{1}}\left(x_{2} x_{3}^{-2 / 3}\right)^{k_{2}}}{k_{1}!k_{2}!\Gamma\left(\frac{s_{1}-k_{1}-2 k_{2}+3}{3}\right)} \\
& \phi_{3}(\lambda, x)=x_{2}^{2} x_{3}^{x_{1}-4} 3 \\
& \sum_{\substack{k_{1} \geq 0, k_{2} \geq-2 \\
\left(k_{1}, k_{2}\right) \in L^{\prime}}} \frac{\left(x_{1} x_{3}^{-1 / 3}\right)^{k_{1}}\left(x_{2} x_{3}^{-2 / 3}\right)^{k_{2}}}{k_{1}!\left(k_{2}+2\right)!\Gamma\left(\frac{s_{1}-k_{1}-2 k_{2}-1}{3}\right)}
\end{aligned}
$$

Here $\phi_{k}(x)$ is the dehomogenization of $\tilde{\phi}_{k}(x)$.

Finally, let us present a difference Pfaffian system for $A$. It can be derived by
using Gröbner bases of $\boldsymbol{H}_{A}$ and has the following form:

$$
S_{1}\left(\begin{array}{c}
f \\
x_{3} \partial_{3} \bullet f \\
S_{1} \bullet f
\end{array}\right)=\left(\begin{array}{ccc}
0 & 0 & 1 \\
-\frac{s_{1} x_{1}}{6 x_{2}} & \frac{3 x_{1} x_{3}-4 x_{2}^{2}}{6 x_{2} x_{3}} & \frac{2\left(s_{1}-1\right) x_{2}+x_{1}^{2}}{6 x_{2}} \\
\frac{s_{1}}{2 x_{2}} & -\frac{3}{2 x_{2}} & -\frac{x_{1}}{2 x_{2}}
\end{array}\right)\left(\begin{array}{c}
f \\
x_{3} \partial_{3} \bullet f \\
S_{1} \bullet f
\end{array}\right) .
$$

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