# On the class of attainable multidimensional NMR spectra 

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#### Abstract

It is shown that a large class of two-dimensional NMR spectra is characterized by a matrix algebra and an invariant subspace. Both the matrix algebra and the invariant subspace are determined by the system matrices of the bilinear system which describes the NMR experiments.


## 1. Introduction

In an earlier paper ([4]) the following system theoretic setup was introduced to describe (multidimensional) NMR experiments. (For a general introduction to NMR experimentation see, e.g., [1,2].) The basic relationship between inputs $u_{1}$ and $u_{2}$ to the system, i.e., the excitation signals or, in particular, the radiofrequency pulses, and the measured output $y$, i.e., the measured induced magnetization, is given described by a bilinear system,

$$
\begin{aligned}
\dot{x}(t) & =A x(t)+u_{1}(t) N_{1} x(t)+u_{2}(t) N_{2} x(t)+b_{1} u_{1}(t)+b_{2} u_{2}(t), \quad x\left(t_{0}\right)=x_{0} \\
y(t) & =c x(t)
\end{aligned}
$$

where $x$ is a state vector, $A, N_{1}, N_{2}$ are square matrices, $b_{1}$ and $b_{2}$ are column vectors and $c$ is a row vector.

A key step in analyzing NMR experiments is the analysis of the effects of a typical input, i.e., a pulse. In [4] it was shown that, for example, in the case of a weakly coupled two spin system, the effects of a radiofrequency input can be described after a suitable state-space transformation by a bilinear system with constant input. More precisely, an NMR system is usually such that if the inputs are given by

$$
u_{1}(t)=B_{1} \cos \left(\omega_{p}\left(t-t_{0}+\Delta t\right)\right), \quad u_{2}(t)=B_{1} \sin \left(\omega_{p}\left(t-t_{0}+\Delta t\right)\right), \quad t \geqslant t_{0}
$$

the input-output relationship can be described by a bilinear system of the form

$$
\begin{aligned}
\dot{x}_{r}(t) & =A_{r}\left(\omega_{p}\right) x_{r}(t)+N_{r}\left(\omega_{p}, \Delta t\right) x_{r}(t)+b_{r}\left(\omega_{p}, \Delta t\right) u_{r}, \quad x_{r}\left(t_{0}\right)=x_{0} \\
y(t) & =\mathrm{e}^{\mathrm{i}\left(t-t_{0}\right) \omega_{p}} c x_{r}(t)
\end{aligned}
$$

with $u_{r}=B_{1}$. The important fact is that the input $B_{1}=u_{r}$ is constant. If $A_{p}\left(\omega_{p}, \Delta t, B_{1}\right):=A_{r}\left(\omega_{p}, \Delta t\right)+u_{r} N_{r}\left(\omega_{p}, \Delta t\right)$ is invertible, this equation can be solved with the solution given by

$$
\begin{aligned}
& x_{r}(t)=\mathrm{e}^{\left(t-t_{0}\right) A_{p}\left(\omega_{p}, \Delta t, B_{1}\right)} x_{0}+\left(\mathrm{e}^{\left(t-t_{0}\right) A_{p}\left(\omega_{p}, \Delta t, B_{1}\right)}-I\right) A_{p}\left(\omega_{p}, \Delta t, B_{1}\right)^{-1} b_{r}\left(\omega_{p}, \Delta t\right) B_{1} \\
& \quad t \geqslant t_{0}
\end{aligned}
$$

For the remainder of the paper we will assume that $A_{p}\left(\omega, \Delta t, B_{1}\right)$ is invertible whenever we write $\left(A_{p}\left(\omega, \Delta t, B_{1}\right)\right)^{-1}$. If a system is in state $x_{0}$ at time $t_{0}$ just before a pulse is applied, the pulse will move the system to the state

$$
x_{1}=P_{1} x_{0}+z_{1}
$$

at time $T>t_{0}$, where

$$
\begin{aligned}
P_{1} & :=P_{1}\left(T, t_{0}, \omega_{p, 1}, \Delta_{1} t, B_{1}\right):=\mathrm{e}^{\left(T-t_{0}\right) A_{p}\left(\omega_{p, 1}, \Delta_{1} t, B_{1}\right)} \\
z_{1} & :=\left(\mathrm{e}^{\left(T-t_{0}\right) A_{p}\left(\omega_{p, 1}, \Delta_{1} t, B_{1}\right)}-I\right) A_{p}\left(\omega_{p, 1}, \Delta_{1} t, B_{1}\right)^{-1} b_{r}\left(\omega_{p, 1}, \Delta_{1} t\right) B_{1}
\end{aligned}
$$

Note that this representation also includes the case in which no pulse has been applied. In this case $z_{1}=0$ and $P_{1}=\mathrm{e}^{\left(T-t_{0}\right) A_{r}}$. Hence the effects of a sequence of $k$ pulses on an initial state is given by

$$
\begin{aligned}
x_{k} & =P_{k}\left(P_{k-1}\left(\cdots\left(P_{2}\left(P_{1} x_{0}+z_{1}\right)+z_{2}\right) \cdots\right)+z_{k-1}\right)+z_{k} \\
& =P_{k} P_{k-1} \cdots P_{1} x_{0}+P_{k} P_{k-1} \cdots P_{2} z_{1}+P_{k} P_{k-1} \cdots P_{3} z_{2}+\cdots+P_{k} z_{k-1}+z_{k}
\end{aligned}
$$

Therefore we can write

$$
x_{k}:=T_{1} x_{0}+e_{1}
$$

where

$$
T_{1}:=P_{k} P_{k-1} \cdots P_{1}
$$

and

$$
e_{1}:=P_{k} P_{k-1} \cdots P_{1} x_{0}+P_{k} P_{k-1} \cdots P_{2} z_{1}+P_{k} P_{k-1} \cdots P_{3} z_{2}+\cdots+P_{k} z_{k-1}+z_{k}
$$

The three blocks of pulses that often characterize a two-dimensional experiment are therefore determined by three matrices $T_{1}, T_{2}$ and $T_{3}$, and three vectors $e_{1}, e_{2}$ and $e_{3}$. Here the notation is such that the pair $\left(T_{1}, e_{1}\right)$ describes the preparation block of pulses, $\left(T_{2}, e_{2}\right)$ represents any possible pulses in the middle of the evolution period, and $\left(T_{3}, e_{3}\right)$ describes the pulses during the mixing period. Note that any pulse within a block of pulses during which no input signal is applied is also considered to be a pulse, i.e., a pulse with zero level input. Hence ([4]) the free induction decay of such a system is given by

$$
\begin{aligned}
s\left(t_{1}, t_{2}\right)= & c \mathrm{e}^{t_{2} A} T_{3} \mathrm{e}^{\left(t_{1} / 2\right) A} T_{2} \mathrm{e}^{\left(t_{1} / 2\right) A} T_{1} x_{0}+c \mathrm{e}^{t_{2} A} T_{3} \mathrm{e}^{\left(t_{1} / 2\right) A} T_{2} \mathrm{e}^{\left(t_{1} / 2\right) A} e_{1} \\
& +c \mathrm{e}^{t_{2} A} T_{3} \mathrm{e}^{\left(t_{1} / 2\right) A} e_{2}+c \mathrm{e}^{t_{2} A} e_{3}, \quad t_{1}, t_{2} \geqslant 0 .
\end{aligned}
$$

As usual, in the above expression $t_{1}$ stands for the measured time and $t_{2}$ for the length of the evolution period.

The spectrum of a two-dimensional experiment is then given by
$G\left(\omega_{1}, \omega_{2}\right)=c\left(2 \pi \mathrm{i} \omega_{1} I-A\right)^{-1}\left[T_{3} P\left(\omega_{1}\right)\left(T_{1} x_{0}+e_{1}\right)+T_{3}\left(2 \pi \mathrm{i} \omega_{1} I-\frac{1}{2} A\right)^{-1} e_{2}+\delta_{0}\left(\omega_{1}\right) e_{3}\right]$, where

$$
P\left(\omega_{1}\right):=\int_{0}^{\infty} \mathrm{e}^{\left(t_{1} / 2\right) A} T_{2} \mathrm{e}^{\left(t_{1} / 2\right) A} \mathrm{e}^{-2 \pi \mathrm{i} \omega_{1} t_{1}} \mathrm{~d} t_{1}, \quad \omega_{1} \in \Re
$$

and $\delta_{0}\left(\omega_{1}\right)$ stands for the delta function with mass concentrated at 0 .
In this paper we will only consider the case, where $T_{2}=I$ and $e_{2}=0$, i.e., the case when no pulses are applied in the center of the evolution period. Moreover, we shall assume that before each scan the system is in equilibrium, i.e., $x_{0}=0$. Then the spectrum is given by

$$
G\left(\omega_{1}, \omega_{2}\right)=c\left(2 \pi \mathrm{i} \omega_{1} I-A\right)^{-1} T_{3}\left(2 \pi \mathrm{i} \omega_{2} I-A\right)^{-1} e_{1}+\delta_{0}\left(\omega_{1}\right) e_{3}, \quad \omega_{1}, \omega_{2} \in \Re
$$

The term $\delta_{0}\left(\omega_{1}\right) e_{3}$ arises from the term $c \mathrm{e}^{t_{2} A}$ in the time domain data. Note that since it is independent of $t_{1}$, it is a constant in the $t_{1}$ time direction. In the analysis of experimental data this term would be removed before applying the Fourier transform. We can therefore assume that the spectrum is given by

$$
G\left(\omega_{1}, \omega_{2}\right)=c\left(2 \pi \mathrm{i} \omega_{1} I-A\right)^{-1} T_{3}\left(2 \pi \mathrm{i} \omega_{2} I-A\right)^{-1} e_{1}, \quad \omega_{1}, \omega_{2} \in \Re
$$

Thus far, we have not taken into consideration techniques such as phase-cycling (see, e.g., [1,2]). Phase-cycling can be seen to be a special case of adding together free induction decays corresponding to different experiments. We will therefore study experiments which are obtained by adding free induction decays together with the same evolution period $t_{2}$ and for which only one of the blocks of pulses is changed. Using such addition schemes we therefore obtain spectra of the form

$$
\begin{aligned}
G\left(\omega_{1}, \omega_{2}\right) & =c\left(2 \pi \mathrm{i} \omega_{1} I-A\right)^{-1}\left(\sum_{j=1}^{k_{3}} \lambda_{j} T_{3, j}\right)\left(2 \pi \mathrm{i} \omega_{2} I-A\right)^{-1}\left(\sum_{l=1}^{k_{1}} \mu_{l} e_{1, l}\right), \\
\omega_{1}, \omega_{2} & \in \Re .
\end{aligned}
$$

Here $T_{3,1}, T_{3,2}, \ldots, T_{3, k_{3}}$ determine $k_{3}$ different pulse blocks in the mixing period, and $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{k_{3}}$ are complex constants. The $k_{1}$ vectors $e_{1,1}, e_{1,2}, \ldots, e_{1, k_{1}}$ are determined by the $k_{1}$ different preparation pulse blocks, and $\mu_{1}, \ldots, \mu_{k_{1}}$ are complex constants.

In this paper we analyze the class of spectra that can be obtained by applying different pulse sequences and addition schemes. It is clear that such an analysis of the mixing period necessitates a careful study of all the possible combinations $\sum_{j=1}^{k_{3}} \lambda_{j} T_{3, j}$. This is the topic of section 2 . In section 3 we will examine the term $\sum_{l=1}^{k_{1}} \mu_{l} e_{1, l}$, which characterizes the preparation period of the experiment.

We will only study 2-D experiments in this paper. However, the extension of the methods presented here to higher dimensional experiments is obvious.

## 2. The mixing period and matrix algebras

In this section we are going to study the set $\mathcal{T}$ given by all possible matrix expressions of the form $\sum_{j=1}^{k_{3}} \lambda_{j} T_{3, j}$. As pointed out in section 1 the mixing period of a two-dimensional NMR experiment is characterized by such an expression. We aim to derive representations of this vector space which are more amenable to computations and analysis than this general description. In fact, we are going to show that this set is a matrix algebra and we are going to determine generators of it.

It follows from the discussion in section 1 that each matrix $T_{3}$ has a representation of the form

$$
T_{3}=\prod_{j=1}^{k} \mathrm{e}^{\left(t_{j}-t_{j-1}\right)\left(A_{r}\left(\omega_{p, j}\right)+B_{1, j} N_{r}\left(\omega_{p, j}, \Delta_{j} t\right)\right)}
$$

for $0<t_{0}<t_{1}<\cdots<t_{k}$. Note that this representation immediately implies that the matrix product of two elements in $\mathcal{T}$ is again in $\mathcal{T}$, i.e., that $\mathcal{T}$ forms an algebra. (See, e.g., [3] for the definition of an algebra.)

We now show that for each $\omega \geqslant 0, B_{1} \geqslant 0, \Delta t \geqslant 0, A_{r}\left(\omega_{p}\right)+B_{1} N_{r}(\omega, \Delta t) \in \mathcal{T}$. To see this first note that $\mathcal{T}$ is closed since it is a finite dimensional vector space. As $\mathrm{e}^{t\left(A_{r}\left(\omega_{p}\right)+B_{1} N_{r}(\omega, \Delta t)\right)} \in \mathcal{T}$ for each $t \geqslant 0$ we also have that the derivative $\left(A_{r}\left(\omega_{p}\right)+\right.$ $\left.B_{1} N_{r}(\omega, \Delta t)\right) \mathrm{e}^{t\left(A_{r}\left(\omega_{p}\right)+B_{1} N_{r}(\omega, \Delta t)\right)} \in \mathcal{T}$ for $t \geqslant 0$. Therefore the derivative evaluated at 0 , i.e., $\left(A_{r}\left(\omega_{p}\right)+B_{1} N_{r}(\omega, \Delta t)\right)$ is also in $\mathcal{T}$. Since $\left(A_{r}\left(\omega_{p}\right)+B_{1} N_{r}(\omega, \Delta t)\right) \in \mathcal{T}$ for all $B_{1} \geqslant 0$ we also have that $A_{r}\left(\omega_{p}\right) \in \mathcal{T}$. Since $\mathcal{T}$ is a vector space this also implies that $N_{r}\left(\omega_{p}, \Delta t\right) \in \mathcal{T}$.

We now show the following theorem.
Theorem 1. (1) $\mathcal{T}$ is a matrix algebra (over the complex field) which is generated by the matrices $A_{r}\left(\omega_{p}\right), N_{r}\left(\omega_{p}, \Delta t\right), \omega_{p} \geqslant 0, \Delta t \geqslant 0$, i.e., each element in $\mathcal{T}$ can be obtained by forming linear combinations of arbitrary products of these matrices.
(2) If $A_{r}\left(\omega_{p}\right)=A_{r}^{1}+\omega_{p} A_{r}^{2}$ and if $N_{r}\left(\omega_{p}, \Delta t\right)=\mathrm{e}^{\mathrm{i} \omega_{p} \Delta t} N_{r}^{1}+\mathrm{e}^{-\mathrm{i} \omega_{p} \Delta t} N_{r}^{2}$, then $\mathcal{T}$ is generated by $A_{r}^{1}, A_{r}^{2}, N_{r}^{1}$ and $N_{r}^{2}$. Here $A_{r}^{1}, A_{r}^{2}, N_{r}^{1}$ and $N_{r}^{2}$ are assumed to be constant matrices independent of $\omega_{p}$ and $\Delta t$.

Proof. (1) As was argued above $\mathcal{T}$ is an algebra. It follows from the remarks preceding the statement of the theorem that the algebra $\mathcal{T}^{\prime}$ generated by $A_{r}\left(\omega_{p}\right), N_{r}\left(\omega_{p}, \Delta t\right)$, $\omega_{p} \geqslant 0, \Delta t \geqslant 0$, is a subalgebra of $\mathcal{T}$.

To show the other inclusion, i.e., that $\mathcal{T} \subseteq \mathcal{T}^{\prime}$ recall that by the Cayley-Hamilton theorem each matrix exponential can be written as the linear combination of a finite number of powers of the matrix. This implies that each element in $\mathcal{T}$ can be generated by the matrices $A_{r}\left(\omega_{p}\right), N_{r}\left(\omega_{p}, \Delta t\right), \omega_{p} \geqslant 0, \Delta t \geqslant 0$.
(2) If $A_{r}\left(\omega_{p}\right)=A_{r}^{1}+\omega_{p} A_{r}^{2}, \omega_{p} \geqslant 0$, then clearly $A_{r}^{1}$ and $A_{r}^{2}$ are in $\mathcal{T}$. If $N_{r}\left(\omega_{p}, \Delta t\right)=\mathrm{e}^{\mathrm{i} \omega_{p} \Delta t} N_{r}^{1}+\mathrm{e}^{-\mathrm{i} \omega_{p} \Delta t} N_{r}^{2}$, then

$$
N_{r}\left(\omega_{p}, \Delta t\right)=\cos \left(\omega_{p} \Delta t\right)\left(N_{r}^{1}+N_{r}^{2}\right)+\mathrm{i} \sin \left(\omega_{p} \Delta t\right)\left(N_{r}^{1}-N_{r}^{2}\right), \quad \omega_{p} \geqslant 0
$$

From this it follows immediately that $N_{r}^{1}+N_{r}^{2}$ and $N_{r}^{1}-N_{r}^{2}$ are in $\mathcal{T}$. Hence, $N_{r}^{1}$ and $N_{r}^{2}$ are also in $\mathcal{T}$.

Given $A_{r}^{1}, A_{r}^{2}, N_{r}^{1}$ and $N_{r}^{2}$, the generating matrices of (1) can be obtained in the obvious way. Hence, $A_{r}^{1}, A_{r}^{2}, N_{r}^{1}$ and $N_{r}^{2}$ indeed generated $\mathcal{T}$.

The special situation which is treated in part (2) of the theorem is encountered in NMR, for example, for the case of two weakly coupled spins (see [4]). In this case it is also straightforward to obtain spanning vectors for the vector space $\mathcal{T}$, which is useful for computational purposes.

Corollary 1. With the assumption and notation of part (2) of the theorem let

$$
M_{0}:=\left[\begin{array}{llll}
A_{r}^{1} & A_{r}^{2} & N_{r}^{1} & N_{r}^{2}
\end{array}\right]
$$

and recursively define $M_{k}$ by setting

$$
M_{k}:=\left[A_{r}^{1} M_{k-1} A_{r}^{2} M_{k-1} N_{r}^{1} M_{k-1} N_{r}^{2} M_{k-1}\right], \quad k=1,2, \ldots
$$

Then the block entries of $M_{k}, k=1,2, \ldots, \operatorname{span} \mathcal{T}$.
As discussed in section 1 the spectrum of a 2-D experiment is given by

$$
G\left(\omega_{1}, \omega_{2}\right)=c\left(2 \pi \mathrm{i} \omega_{1} I-A\right)^{-1} T_{3}\left(2 \pi \mathrm{i} \omega_{2} I-A\right)^{-1} b
$$

where we have set $b:=e_{1}$. If we assume that $A$ is diagonal with eigenvalues $a_{1}, a_{2}, \ldots, a_{n}$, then the spectrum can be written as

$$
G\left(\omega_{1}, \omega_{2}\right)=\sum_{j=1}^{n} \sum_{l=1}^{n} \frac{c_{j}}{2 \pi \mathrm{i} \omega_{1}-a_{j}} t_{j l} \frac{b_{l}}{2 \pi \mathrm{i} \omega_{2}-a_{l}}
$$

where $c:=c\left(c_{1}, c_{2}, \ldots, c_{n}\right), b:=\left(b_{1}, b_{2}, \ldots, b_{n}\right)^{\mathrm{T}}, T_{3}:=\left(t_{i j}\right)_{1 \leqslant i, j \leqslant n}$.
Redundancies occur in this description of the spectrum if $A$ has eigenvalues with multiplicities. Assume that $a_{1}^{r}, a_{2}^{r}, \ldots, a_{k}^{r}$ are the eigenvalues of $A$ with all possible multiplicities removed. To remove the redundancies in the description of the spectrum, we construct the $k \times n$ matrix $Q$, whose entries are either zeros or ones. The first row of $Q$ is constructed as follows. In the first entry enter 1 corresponding to the eigenvalue $a_{1}$. Then enter a zero in the second entry if $a_{2} \neq a_{1}$, otherwise enter 1 . Similarly, in the $j$ th entry enter 0 if $a_{k} \neq a_{1}$ and 1 if $a_{k}=a_{1}$. The $r$ th row, $1<r \leqslant k$, of $Q$ is constructed as follows: Consider the first eigenvalue $a_{l}$ in the list $a_{1}, a_{2}, \ldots, a+l, \ldots, a_{n}$ which is not equal to any of the eigenvalues $a_{1}^{r}, a_{2}^{r}, \ldots, a_{k-1}^{r}$ and enter 1 at the corresponding entry of the row, i.e., its $l$ th entry. Enter 0 in the first $l-1$ entries of this row. Then enter a zero in the $(l+s)$ th entry of the row if $a_{l+s} \neq a_{l}$ and 1 otherwise, for $s=1, \ldots, n-k$.

Now let

$$
Q b:=b^{1}:=\left(\begin{array}{c}
b_{1}^{1} \\
b_{2}^{1} \\
\vdots \\
b_{k}^{1}
\end{array}\right),
$$

$$
\begin{aligned}
c Q^{*} & :=c^{1}:=\left(c_{1}^{1} c_{2}^{1} \ldots c_{k}^{1}\right) \\
Q T_{3} Q^{*} & :=T_{3}^{1}:=\left(t_{j k}^{1}\right)_{1 \leqslant j, l \leqslant k} .
\end{aligned}
$$

With this notation we then have a 'reduced' representation of the spectrum given by

$$
G\left(\omega_{1}, \omega_{2}\right)=\sum_{j=1}^{k} \sum_{l=1}^{k} \frac{c_{j}^{1}}{2 \pi \mathrm{i} \omega_{1}-a_{j}^{r}} t_{j l}^{1} \frac{b_{l}^{1}}{2 \pi \mathrm{i} \omega_{2}-a_{l}^{r}}
$$

Further redundancies can occur in the description of the spectrum if $b^{1}$ or $c^{1}$ have zero entries. If $c_{j}^{1}$, the $j$ th entry of $c^{1}$, is zero then the pole $1 /\left(2 \pi \mathrm{i} \omega_{1}-a_{j}^{r}\right)$ will not appear in the spectrum. Similarly, if the $l$ th entry of $b^{1}$ is zero, then the pole $1 /\left(2 \pi \mathrm{i} \omega_{2}-a_{l}^{r}\right)$ will not appear in the spectrum. Moreover, if $c_{j}^{1}$ is zero then the $j$ th row of $T_{3}^{1}$ is irrelevant for the appearance of the spectrum and similarly, if $b_{l}^{1}$ is zero then the $l$ th column of $T_{3}^{1}$ is irrelevant for the appearance of the spectrum. To study this effect on the algebra $\mathcal{T}$ let $\mathcal{C}$ be the 'projection matrix' that is obtained from the $k \times k$ identity matrix by deleting the $j$ th row if the $c_{j}^{1}=0$, for some $j \in\{1, \ldots, k\}$. Analogously, define $\mathcal{B}$ to be the projection matrix that is obtained from the $k \times k$ identity matrix by deleting the $l$ th column if $b_{l}^{1}=0$, for some $l \in\{1,2, \ldots, k\}$. Let now

$$
\begin{aligned}
T_{3}^{r} & :=\left(t_{j l}^{r}\right)_{\substack{1 \leqslant \leqslant \leqslant n_{b} \\
1 \\
1 \\
n_{b}}}:=\mathcal{C} T_{3}^{1} \mathcal{B}, \\
b^{r} & :=\left(\begin{array}{c}
b_{1}^{r} \\
b_{2}^{r} \\
\vdots \\
b_{n_{b}}^{r}
\end{array}\right):=\mathcal{B}^{*} b^{1}, \\
c^{r} & :=\left(c_{1}^{r} c_{2}^{r} \ldots c_{n_{c}}^{r}\right):=c^{1} \mathcal{C}^{*} .
\end{aligned}
$$

Then

$$
G\left(\omega_{1}, \omega_{2}\right)=\sum_{j=1}^{n_{c}} \sum_{l=1}^{n_{b}} \frac{c_{j}^{r}}{2 \pi \mathrm{i} \omega_{1}-a_{j}^{r}} t_{j l}^{r} \frac{b_{l}^{r}}{2 \pi \mathrm{i} \omega_{2}-a_{l}^{r}}
$$

Note that in general $n_{b} \neq n_{c}$, i.e., $T_{3}^{r}$ is not necessarily a square matrix.
In the 'reduced' representation of the spectrum which was just discussed, the set of all $T_{3}^{r}$ matrices is given by

$$
\mathcal{T}:=\mathcal{C} Q \mathcal{T} Q^{*} \mathcal{B}=\left\{\mathcal{C} A T_{3} Q^{*} B \mid T_{3} \in \mathcal{T}\right\}
$$

While in general $\mathcal{T}_{r}$ is not a matrix algebra it is in an obvious way a subspace of the vector space of $n_{c} \times n_{b}$ matrices with complex entries.

Of particular interest from an experimental point of view is whether an experiment can be devised such that a particular diagonal- or cross-peak can for example be suppressed or made to appear without unduly influencing other parts of the spectrum. Since a peak has the general form $c_{j}^{r} t_{j l}^{r} b_{l}^{r} /\left(\left(2 \pi \mathrm{i} \omega_{1}-a_{j}^{r}\right)\left(2 \pi \mathrm{i} \omega_{2}-b_{l}^{r}\right)\right)$, this problem
is related to the question whether or not experiments can be found that can set the $t_{j l}^{r}$-entry of $T_{3}^{r}$ to zero or to another value without changing the other entries of $T_{3}^{r}$. If this is possible, we call this a switchable 2-D pole. We say that 2-D pole is an $(a, b)$-pole if the pole is given by $v /\left(\left(2 \pi \mathrm{i} \omega_{1}-a\right)\left(2 \pi \mathrm{i} \omega_{2}-b\right)\right)$ for some constant $v$. It is easily seen that the $\left(a_{j}^{r}, b_{j}^{r}\right)$-pole is switchable if and only if $E_{j l} \in \mathcal{T}_{r}$, where $E_{j l}$ is the $n_{c} \times n_{b}$ matrix such that all entries are zero with the exception of the ( $j, l$ ) entry which is 1 .

Similarly, we call a collection $\left(a_{n_{1}}^{r}, b_{n_{1}}^{r}\right), \ldots,\left(a_{n_{s}}^{r}, b_{n_{s}}^{r}\right)$ of 2-D poles switchable if two experiments can found such that in the first experiment all the coefficients of the poles are non-zero and in the second experiment all coefficients are zero, but the coefficients for the other 2-D poles are the same for both experiments.

We have the following result.
Proposition 1. The collection $\left(a_{n_{1}}^{r}, b_{m_{1}}^{r}\right), \ldots,\left(a_{n_{s}}^{r}, b_{m_{s}}^{r}\right)$ of 2-D poles is switchable if and only if there exist nonzero coefficients $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{s}$ such that $\sum_{l=1}^{n} \lambda_{l} E_{n_{s}, m_{s}} \in$ $\mathcal{T}_{r}$.

## 3. The preparation period and invariant subspaces

It was shown in section 1 that the preparation period of a two-dimensional NMR experiment is determined by the expression $\sum_{l=1}^{k_{1}} \mu_{l} e_{1, l}$. It will be the topic of this section to analyze the set $\mathcal{E}$ of all such expressions. In particular, we are going to show that the set of all such expressions is an invariant subspace of $\mathcal{C}^{n}$. It was shown in section 1 that each $e_{1, l}$ is given by

$$
e_{1, l}:=P_{k, l} P_{k-1, l} \cdots P_{2, l} z_{1, l}+P_{k, l} P_{k-1, l} \cdots P_{3, l} z_{2, l}+\cdots+P_{k, l} z_{k-1, l}+z_{k, l}
$$

where

$$
z_{j, l}:=\left(\mathrm{e}^{\left(t_{j, l}-t_{j, l-1}\right)} A_{p}\left(\omega_{p, j, l}, \Delta_{j, l} t, B_{j, l}\right)-I\right) A_{p}\left(\omega_{p, j, l}, \Delta_{j, l} t, B_{j, l}\right)^{-1} b_{r}\left(\omega_{p, j, l}, \Delta_{j, l} t\right) B_{j, l}
$$

and

$$
P_{j, l}:=P_{j, l}\left(t_{j, l}, t_{j, l-1}, \omega_{p, j, l}, \Delta_{j, l} t, B_{j, l}\right):=\mathrm{e}^{\left(t_{j, l}-t_{j, l-1}\right) A_{p}\left(\omega_{p, j, l}, \Delta_{j, l}, B_{j, l}\right)},
$$

$j=1, \ldots, k, l=1,2, \ldots, k_{1}$. Here we have assumed as will be done throughout the remainder of this paper that $x_{0}=0$.

We have the following theorem.
Theorem 2. (1) We have

$$
\mathcal{E}=\operatorname{span}\left\{\mathcal{T} b_{r}\left(\omega_{p}, \Delta t\right) \mid \omega_{p} \geqslant 0, \Delta t \geqslant 0\right\},
$$

where $\operatorname{span}(V)$ denotes the linear span of the set of vectors $V$ and $\mathcal{T}$ is defined as in section 2.
(2) $\mathcal{E}$ is the smallest linear subspace of $\mathcal{C}^{n}$ that contains $b_{r}\left(\omega_{p}, \Delta t\right)$ for $\omega_{p} \geqslant 0$, $\Delta t \geqslant 0$, and is invariant under $A_{r}\left(\omega_{p}\right)$ and $N_{r}\left(\omega_{p}, \Delta t\right)$ for all $\omega_{p} \geqslant 0, \Delta t \geqslant 0$.
(3) If $A_{r}\left(\omega_{p}\right)=A_{r}^{1}+\omega_{p} A_{r}^{2}$, if $N_{r}\left(\omega_{p}, \Delta t\right)=\mathrm{e}^{\mathrm{i} \omega_{p} \Delta t} N_{r}^{1}+\mathrm{e}^{-\mathrm{i} \omega_{p} \Delta t} N_{r}^{2}$ and if $b_{r}\left(\omega_{p}, \Delta t\right)=\mathrm{e}^{\mathrm{i} \omega_{p} \Delta t} b_{r}^{1}+\mathrm{e}^{\mathrm{i} \omega_{p} \Delta t} b_{r}^{2}$, then $\mathcal{E}$ is the smallest subspace of $\mathcal{C}^{n}$ that contains $b_{r}^{1}$ and $b_{r}^{2}$ and is invariant under $A_{r}^{1}, A_{r}^{2}, N_{r}^{1}$ and $N_{r}^{2}$. Here $A_{r}^{1}, A_{r}^{2}, N_{r}^{1}, N_{r}^{2}, b_{r}^{1}$ and $b_{r}^{2}$ are assumed to be constant matrices independent of $\omega_{p}$ and $\Delta t$.

Proof. (1) We begin by showing that $\mathcal{E} \subseteq \operatorname{span}\left\{\mathcal{T} b_{r}\left(\omega_{p}, \Delta t\right) \mid \omega_{p} \geqslant 0, \Delta t \geqslant 0\right\}$. To do this we first show that a vector $z$ as defined above is in $\operatorname{span}\left\{\mathcal{T} b_{r}\left(\omega_{p}, \Delta t\right) \mid \omega_{p} \geqslant\right.$ $0, \Delta t \geqslant 0)\}$. Let $\omega_{p} \geqslant 0, \Delta t \geqslant 0, B_{1} \geqslant 0$ and set

$$
z:=\left(\mathrm{e}^{t A_{p}\left(\omega_{p}, \Delta t, B_{1}\right)}-I\right) A_{p}\left(\omega_{p}, \Delta t, B_{1}\right)^{-1} b_{r}\left(\omega_{p}, \Delta t\right) B_{1}
$$

Then

$$
\begin{aligned}
z & =\left(\mathrm{e}^{t A_{p}\left(\omega_{p}, \Delta t, B_{1}\right)}-I\right) A_{p}\left(\omega_{p}, \Delta t, B_{1}\right)^{-1} b_{r}\left(\omega_{p}, \Delta t\right) B_{1} \\
& =\left(\sum_{r=0}^{\infty} \frac{1}{r!}\left(t A_{p}\left(\omega_{p}, \Delta t, B_{1}\right)\right)^{r}-I\right) A_{p}\left(\omega_{p}, \Delta t, B_{1}\right)^{-1} b_{r}\left(\omega_{p}, \Delta t\right) B_{1} \\
& =\left(\sum_{r=1}^{\infty} \frac{1}{r!}\left(t A_{p}\left(\omega_{p}, \Delta t, B_{1}\right)\right)^{r}\right) b_{r}\left(\omega_{p}, \Delta t\right) B_{1}
\end{aligned}
$$

which shows that $\left.z \in \operatorname{span}\left\{\mathcal{T} b_{r}\left(\omega_{p}, \Delta t\right) \mid \omega_{p} \geqslant 0, \Delta t \geqslant 0\right)\right\}$.
We now show that if $y_{0} \in \operatorname{span}\left\{\mathcal{T} b_{r}\left(\omega_{p}, \Delta t\right) \mid \omega_{p} \geqslant 0, \Delta t \geqslant 0\right\}$, then for $\omega_{p} \geqslant 0$, $\Delta t \geqslant 0, t>t_{0}$ and $B_{1} \geqslant 0$, we have that

$$
P y_{0}+z
$$

is in $\operatorname{span}\left\{\mathcal{T} b_{r}\left(\omega_{p}, \Delta t\right) \mid \omega_{p} \geqslant 0, \Delta t \geqslant 0\right\}$, where

$$
\begin{aligned}
P & :=\mathrm{e}^{\left(t_{1}-t_{0}\right) A_{p}\left(\omega_{p}, \Delta t, B_{1}\right)} \\
z & :=\left(\mathrm{e}^{\left(t_{1}-t_{0}\right) A_{p}\left(\omega_{p}, \Delta t, B_{1}\right)}-I\right) A_{p}\left(\omega_{p}, \Delta t, B_{1}\right)^{-1} b_{r}\left(\omega_{p}, \Delta t\right) B_{1} .
\end{aligned}
$$

This is the case since $z \in \operatorname{span}\left\{\mathcal{T} b_{r}\left(\omega_{p}, \Delta t\right) \mid \omega_{p} \geqslant 0, \Delta t \geqslant 0\right\}$ by the above argument and since $P \in \mathcal{T}$. Since a general element in $\mathcal{E}$ is the linear combination of elements of the form (see section 1)

$$
P_{k}\left(P_{k-1}\left(\cdots\left(P_{2}\left(P_{1} x_{0}+z_{1}\right)+z_{2}\right) \cdots\right)+z_{k-1}\right)+z_{k}
$$

the inclusion $\mathcal{E} \subseteq \operatorname{span}\left\{\mathcal{T} b_{r}\left(\omega_{p}, \Delta t\right) \mid \omega_{p} \geqslant 0, \Delta t \geqslant 0\right\}$ follows immediately from what was just shown.

We now need to show that $\operatorname{span}\left\{\mathcal{T} b_{r}\left(\omega_{p}, \Delta t\right) \mid \omega_{p} \geqslant 0, \Delta t \geqslant 0\right\} \subseteq \mathcal{E}$. To do this let $v \in\left\{\mathcal{T} b_{r}\left(\omega_{p}, \Delta t\right) \mid \omega_{p} \geqslant 0, \Delta t \geqslant 0\right\}$, i.e.,

$$
v=P_{k} P_{k-1} \cdots P_{2} b_{r}\left(\omega_{p}, \Delta t\right)
$$

where $\Delta t \geqslant 0, \omega_{p} \geqslant 0$ and

$$
P_{j}=\mathrm{e}^{\left(t_{j}-t_{j-1}\right) A_{p}\left(\omega_{p, j}, \Delta_{j} t, B_{j}\right)}
$$

with $0<t_{1}<t_{2}<\cdots<t_{k}$ and $\omega_{p, j} \geqslant 0, \Delta_{j} t \geqslant 0, B_{j} \geqslant 0, j=2,3, \ldots, k$. In order to show that $v \in \mathcal{E}$, we first show that

$$
P_{k} P_{k-1} \cdots P_{2} z_{1}(t) \in \mathcal{E}
$$

where

$$
z_{1}(t):=\left(\mathrm{e}^{\left(t-t_{0}\right) A_{p}\left(\omega_{p}, \Delta t, B\right)}-I\right)\left(A_{p}\left(\omega_{p}, \Delta t, B\right)\right)^{-1} b_{r}\left(\omega_{p}, \Delta t\right), \quad 0<t<t_{1}, B>0
$$

Set

$$
z_{j}:=\left(\mathrm{e}^{\left(t_{j}-t_{j-1}\right) A_{p}\left(\omega_{p, j}, \Delta_{j} t, B_{j}\right)}-I\right)\left(A_{p}\left(\omega_{p, j}, \Delta_{j} t, B_{j}\right)\right)^{-1} b_{r}\left(\omega_{p, j}, \Delta_{j} t\right) B_{j}
$$

$j=2,3, \ldots, k$. By definition of $\mathcal{E}, z_{j} \in \mathcal{E}$ for $j=1,2, \ldots, k$, and $P_{k} z_{k-1}+z_{k} \in \mathcal{E}$. Since $\mathcal{E}$ is a vector space this implies that $P_{k} z_{k-1} \in \mathcal{E}$. Similarly, by definition of $\mathcal{E}$, we have that $P_{k} P_{k-1} z_{k-2}+P_{k} z_{k-1}+z_{k} \in \mathcal{E}$. Again using that $\mathcal{E}$ is a vector space, this implies that $P_{k} P_{k-1} z_{k-2} \in \mathcal{E}$. Proceeding recursively we show that

$$
P_{k} P_{k-1} \cdots P_{2} z_{1}(t) \in \mathcal{E}, \quad 0<t<t_{1}
$$

Let $T_{3}:=P_{k} P_{k-1} \cdots P_{2}$ and set $r(t):=T_{3} z_{1}(t), 0<t<t_{1}$. Since $\mathcal{E}$ is a finite dimensional subspace of $\mathcal{C}^{n}$, it is closed. Hence,

$$
\begin{aligned}
v & =P_{k} P_{k-1} \cdots P_{2} b_{r}\left(\omega_{p}, \Delta t\right)=T_{3} A_{p}\left(\omega_{p}, \Delta t, B\right)\left(A_{p}\left(\omega_{p}, \Delta t, B\right)\right)^{-1} b_{r}\left(\omega_{p}, \Delta t\right) \\
& =T_{3} \lim _{t \rightarrow 0} \frac{z_{1}(t)}{t}=\lim _{t \rightarrow 0} \frac{r(t)}{t} \in \mathcal{E} .
\end{aligned}
$$

(2) Clearly, $\mathcal{E}=\operatorname{span}\left\{\mathcal{T} b_{r}\left(\omega_{p}, \Delta t\right) \mid \omega_{p} \geqslant 0, \Delta t \geqslant 0\right\}$ is a linear subspace of $\mathcal{C}^{n}$ that contains $b_{r}\left(\omega_{p}, \Delta t\right)$ for $\omega_{p} \geqslant 0, \Delta t \geqslant 0$, and by the characterization of $\mathcal{T}$ in theorem 1 is invariant under $A_{r}\left(\omega_{p}\right)$ and $N_{r}\left(\omega_{p}, \Delta t\right)$ for all $\omega_{p} \geqslant 0, \Delta t \geqslant 0$.

Conversely, an invariant subspace of $\mathcal{C}^{n}$ that contains $b_{r}\left(\omega_{p}, \Delta t\right)$ for $\omega_{p} \geqslant 0$, $\Delta t \geqslant 0$, and is invariant under $A_{r}\left(\omega_{p}\right)$ and $N_{r}\left(\omega_{p}, \Delta t\right)$ for all $\omega_{p} \geqslant 0, \Delta t \geqslant 0$, contains the elements of $\left\{\mathcal{T} b_{r}\left(\omega_{p}, \Delta t\right) \mid \omega_{p} \geqslant 0, \Delta t \geqslant 0\right\}$.
(3) The statement follows from (2) in conjunction with theorem 1.

In [4] it was shown that the spectrum of a 1-D NMR experiment is given by

$$
G(\omega)=c(2 \pi \mathrm{i} I-A)^{-1} e_{0}, \quad \omega \in \Re
$$

for some vector $e_{0}$. The only parameter in this description which can be adjusted in an experiment is $e_{0}$. If addition schemes are introduced that are analogous to those introduced in section 1 for 2-D systems, then the set $\mathcal{E}$ describes all possible vectors $e_{0}$.

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