# EXPONENTIATIONS OVER THE UNIVERSAL ENVELOPING ALGEBRA OF $s l_{2}(\mathbb{C})$ 

SONIA L'INNOCENTE ${ }^{1}$, ANGUS MACINTYRE, AND FRANÇOISE POINT ${ }^{2}$

Abstract. We construct, by model-theoretic methods, several exponentiations on the universal enveloping algebra $U$ of the Lie algebra $s l_{2}(\mathbb{C})$.

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## 1. Introduction

Consider the Lie algebra $M_{2}(\mathbb{C})$ and the Lie (sub)algebra $s l_{2}(\mathbb{C})$ of all $2 \times 2$ trace zero matrices with complex entries. Recall that a standard basis of $s l_{2}(\mathbb{C})$ (as $\mathbb{C}$-vectorspace) is given by: $x=\left(\begin{array}{ll}0 & 1 \\ 0 & 0\end{array}\right), y=\left(\begin{array}{ll}0 & 0 \\ 1 & 0\end{array}\right), h=\left(\begin{array}{cc}1 & 0 \\ 0 & -1\end{array}\right)=\operatorname{diag}(1,-1)$. The generators $x, y, h$ satisfy the relations: $[h, x]=2 x,[h, y]=-2 y,[x, y]=h$, where $[u, v]$ is the usual commutator of $u$ and $v$.

For each positive integer $\lambda$, we consider the finite-dimensional simple $s l_{2}(\mathbb{C})$-module $V_{\lambda}$ of dimension $\lambda+1$ and the (matrix) Lie algebra $M_{\lambda+1}(\mathbb{C})$ (the endomorphism ring of $V_{\lambda}$, viewed as a $\mathbb{C}$-vectorspace) and take the exponential maps from $M_{\lambda+1}(\mathbb{C})$ into the linear group $G L_{\lambda+1}(\mathbb{C})$. (In section 3, we recall some properties of these exponential maps.)

We connect these exponentials to the universal enveloping algebra $U$ of $s l_{2}(\mathbb{C})$ (whose definition and algebraic properties are described in Section 4). We will use some basic facts on the representation theory of this associative algebra (and its analogue over any algebraically closed field of characteristic 0 ). It has been studied from a model theoretical point of view by [11] and then by $[12,16,17]$.

Using on one hand the concrete exponential maps defined on the matrix rings $M_{\lambda+1}(\mathbb{C})$ and on the other hand the universal property of $U$, we define a sequence of exponential maps indexed by $\lambda$ from $U$ to $G L_{\lambda+1}(\mathbb{C})$. We describe some of the properties of these maps, which we have formalized (in section 2) by defining the notion of a non-commutative exponential ring (generalizing the commutative case) and we explicitly calculate elements lying in their kernels (respectively images). Then, we show that $U$ embeds into any nonprincipal ultraproduct $\prod_{\mathcal{V}} M_{\lambda+1}(\mathbb{C})$ and we define an exponential map $E X P$ from $U$ to any non-principal ultraproduct $\prod_{\mathcal{V}} G L_{\lambda+1}(\mathbb{C})$ of the groups $G L_{\lambda+1}(\mathbb{C})$, where $\mathcal{V}$ is a non principal ultrafilter on $\omega$. We show that $(U, E X P)$ is a non-commutative exponential ring, and we explicitly calculate a part of the kernel of $E X P$. Note that a formal exponential map exp was previously defined in the completion $\hat{U}$ of $U([22])$, on the ideal on $\hat{U}$ generated by the generators of $U$; in section 7 , we will compare the two approaches.

[^0][^1]We go on to endow $U$ with a topology using a norm in $\prod_{\mathcal{V}} M_{\lambda+1}(\mathbb{C})$ which takes its values in a non-standard ultrapower of $\mathbb{R}$, and we show that the exponential map $E X P$ is continuous and that the subgroup generated by $E X P(U)$ is a topological group.

Finally, by considering another norm on each $M_{\lambda+1}(\mathbb{C})$, and the asymptotic cone relative to these norms and a non-principal ultrafilter $\mathcal{V}$ on $\omega$, we embed $U$ in a complete metric space and show that $U$ has a faithful continuous action on that space.

## 2. Preliminaries on formalism

Let us set up the languages we need.
Let $\mathcal{L}_{g}:=\{\cdot, 1\}$ be the language of groups. Let $\mathcal{L}:=\{+,-, \cdot, 0,1\}$ be the language of (associative) rings, and let $\mathcal{L}_{l}:=\{+,-,[\cdot, \cdot], 0\}$ be the language of Lie rings. For a ring $R$, let $\mathcal{L}_{m, R}:=\{+,-, 0, \cdot r ; r \in R\}$ be the language of right $R$-modules.

For the language of $R$-algebras, where $R$ is a commutative ring, we will choose the expansion $\mathcal{L}_{A l g}$ of $\mathcal{L}$, a two-sorted language with a sort for a ring $R$, a sort for an algebra $A$ (associative or not) and a scalar multiplication map from $A \times R$ to $A$, where $A$ is either a $\mathcal{L}$-structure or a $\mathcal{L}_{l}$-structure and $R$ is a $\mathcal{L}$-structure.

For the language of Lie $K$-algebras, where $K$ is a field, we will choose either $\mathcal{L}_{\text {Lie }}:=$ $\mathcal{L}_{l} \cup \mathcal{L}_{m, K}$ or the two-sorted language $\mathcal{L}_{\text {Alg }}$. Note that for the former we omit reference to $K$ when it is understood. We will assume that $K$ is a field of characteristic 0 and is complete with respect to a nontrivial absolute value.

Let $T_{K}$ be the theory of $K$-vector spaces in $\mathcal{L}_{m, K}$.
Let $T_{\boldsymbol{L}}$ be the theory of Lie $K$-algebras in the language $\mathcal{L}_{L i e}$, namely
(1) $T_{K}$,
(2) $[\cdot, \cdot]$ is a $K$-bilinear map from $L \times L$ to $L$,
(3) $\forall x_{1}\left[x_{1}, x_{1}\right]=0$,
(4) $\forall x_{1}, \forall x_{2} \forall x_{3}\left(\left[\left[x_{1}, x_{2}\right], x_{3}\right]+\left[\left[x_{2}, x_{3}\right], x_{1}\right]+\left[\left[x_{3}, x_{1}\right], x_{2}\right]=0\right)$.

## 3. Axioms for semisimple Lie algebras

We will translate, into model-theoretic terms, the basic results on existence and uniqueness of a semi-simple Lie $K$-algebra with a given reduced abstract root system $\Phi$ ([13], chapter 18.2). This is not essential for the present paper, but may be of interest in future generalizations.

Recall that $\Phi$ is a subset of an Euclidean space $E$ endowed with a positive definite symmetric bilinear form (.,.). Denote by $<\beta, \alpha>:=2 \cdot \frac{(\beta, \alpha)}{(\alpha, \alpha)}$. For a root system $\Phi$ these values are integers.

For $x \in L$, let $a d x$ be the linear transformation of $L$ sending $y \in L$ to $[x, y]$.
Proposition 3.1. The theory of any semi-simple Lie algebra $L$ with given reduced root system $\Phi$ ( and inner product on it) is axiomatisable in $\mathcal{L}_{\text {Alg }}$ by the set $T_{\Phi}$ of axioms below. Moreover each $T_{\Phi}$ is $\aleph_{1}$-categorical.
(1) $T_{\text {Alg }}$ the theory of $K$-algebras in $\mathcal{L}_{\text {Alg }}$ over some field $K$;
(2) The scheme of axioms expressing that $K$ is an algebraically closed field of characteristic 0;
(3) (the $\alpha_{j}$ are the elements of the root system) $\exists h_{1} \cdots \exists h_{\ell} \exists e_{1} \cdots \exists e_{\ell} \exists e_{-1} \cdots \exists e_{-\ell}$

$$
\begin{gathered}
{\left[\bigwedge_{1 \leq i, j \leq \ell}\left[h_{i}, h_{j}\right]=0\right.} \\
\bigwedge_{1 \leq i \leq \ell}\left[e_{i}, e_{-i}\right]=h_{i} \& \bigwedge_{1 \leq i \neq j \leq \ell}\left[e_{i}, e_{j}\right]=0 \\
\bigwedge_{1 \leq i, j \leq \ell}\left[h_{i}, e_{j}\right]=<\alpha_{j}, \alpha_{i}>\cdot e_{j} \& \bigwedge_{1 \leq i, j \leq \ell}\left[h_{i}, e_{-j}\right]=-<\alpha_{j}, \alpha_{i}>\cdot e_{-j} \\
\bigwedge_{1 \leq i \neq j \leq \ell}\left(a d e_{i}\right)^{-<\alpha_{j}, \alpha_{i}>+1}\left(e_{j}\right)=0 \\
\bigwedge_{1 \leq i \neq j \leq \ell}\left(a d e_{-i}\right)^{-<\alpha_{j}, \alpha_{i}>+1}\left(e_{-j}\right)=0
\end{gathered}
$$

$$
\left.\& \forall x \exists k_{1} \in K \cdots \exists k_{3 \ell} \in K x=\sum_{1 \leq i \leq \ell} k_{i} . h_{i}+\sum_{1 \leq i \leq \ell} k_{\ell+j} \cdot e_{j}+\sum_{1 \leq j \leq \ell} k_{2 \ell+j} \cdot e_{-j}\right] .
$$

Proof: Serre's work tells us that given a root system $\Phi$ and a field of characteristic 0 , there exists a unique Lie algebra $L$ that can be presented by these relations and that it is semisimple. The second statement follows from the fact that if $L$ is a model of these axioms of cardinality $\aleph_{1}$, then it is a Lie algebra over an algebraically closed field $F$ of characteristic 0 of cardinality $\aleph_{1}$.

Question 3.1. Is the theory of any semi-simple Lie $\mathbb{C}$-algebra $L$ with given root system $\Phi$ finitely axiomatisable in $\mathcal{L}_{\text {Lie }}$ modulo $T_{\mathbb{C}}$ ?

Let Axiom (3') be got from Axiom (3) by deleting the last part where we quantify over $x$. Let $L$ be a model of (1),(2) and (3'), and let $L_{0}$ be the Lie subalgebra generated by the elements $h_{i}, e_{j}, e_{-j}$ satisfying the above relations. Then, Serre's theorem tells us that $L_{0}$ is a semi-simple finite-dimensional Lie algebra with root system $\Phi$ and Cartan subalgebra generated by $h_{i}, 1 \leq i \leq \ell$. Then can we add an axiom that forces $L$ to be equal to $L_{0}$ ?

We could have also worked in the language $\mathcal{L}_{l}$ since any semi-simple Lie algebra has a basis with integral structure constants (a Chevalley basis). We do not pursue this matter here.

## 4. Exponential Rings

Let $\mathcal{L}_{E}:=\mathcal{L} \cup\{E\}$ (respectively $\mathcal{L}_{A l g, E}:=\mathcal{L}_{A l g} \cup\{E\}$ ) where $E$ is a unary function symbol. We will introduce the notion of (non commutative) exponential ring generalizing the commutative case (see for instance [7]).

Definition 4.1. Let $(R, E, G)$ be a two-sorted structure with $R$ an $\mathcal{L}$-structure, $G$ a $\mathcal{L}_{g^{-}}$ structure and $E$ a map from $R$ to $G$. We will say that $(R, E, G)$ is an exponential ring if $R$ is an associative ring with $1, G$ a (multiplicative) group and if $E: R \rightarrow G$ satisfies the following axioms:
(1) $E(0)=1$,
(2) $\forall x \quad E(x) \cdot E(-x)=1$,
(3) $\forall x \forall y(x \cdot y=y \cdot x \rightarrow E(x+y)=E(x) \cdot E(y))$.

If in addition $R$ is a $K$-algebra, then $(R, K, E, G)$ is an exponential $K$-algebra if $(R, K)$ is a $\mathcal{L}_{\text {Alg }}$-structure such that the reduct $(R, E, G)$ is an exponential ring, the $\mathcal{L}_{\text {Alg }}$-reduct $(R, K)$ a $K$-algebra and

$$
\forall k_{1}, k_{2} \in K \quad \forall x \in R \quad E\left(k_{1} \cdot x\right) \cdot E\left(k_{2} \cdot x\right)=E\left(\left(k_{1}+k_{2}\right) \cdot x\right)
$$

Note that this last axiom together with (1) implies (2) above.
One recovers the classical case by taking $G$ the group of units of $R$, by assuming that $R$ is a commutative ring and then we revert to the one-sorted $\mathcal{L}_{E}$-structure $(R, E)$. In the case we deal with an exponential $K$-algebra, we will get that $(K, E)$ is an exponential field.

## 5. A Natural exponential map over $M_{\lambda+1}(\mathbb{C})$

Consider the field $\mathbb{C}$ of complex numbers and, for a fixed natural number $\lambda$, the associative $\mathbb{C}$-algebra $M_{\lambda+1}(\mathbb{C})$ of all $(\lambda+1) \times(\lambda+1)$ matrices with coefficients in $\mathbb{C}$ (with the matrix multiplication - as the underlying operation). It is also a Lie $\mathbb{C}$-algebra with the bracket $[A, B]:=A \cdot B-B \cdot A$ (see $[8,13])$. For $A \in M_{\lambda+1}(\mathbb{C})$, denote by $A^{*}$ the conjugate of the transpose of $A$, by $\operatorname{tr}(A)$ the trace of $A$, and finally by $\operatorname{det}(A)$ its determinant.

We will denote by $\operatorname{Diag}_{\lambda+1}(\mathbb{C})\left(\right.$ respectively $\left.U T_{\lambda+1}(\mathbb{C})\right)$ the subset of all diagonal matrices (respectively upper triangular matrices) in $M_{\lambda+1}(\mathbb{C})$.
Recall that on the Lie algebra $M_{\lambda+1}(\mathbb{C})$, we have a Hermitian sesquilinear form $(\cdot, \cdot)_{\lambda+1}$ defined by $(A, B)_{\lambda+1}:=\operatorname{tr}\left(B^{*} \cdot A\right)=\sum_{i, j} A_{i j} \cdot \bar{B}_{i j}$, where $A, B \in M_{\lambda+1}(\mathbb{C})$, its values are in $\mathbb{C}([21]$ page 9$)$. The Frobenius norm (denoted by $F$-norm) associated with it, is defined as follows: $\|A\|_{F, \lambda+1}^{2}:=(A, A)_{\lambda+1}$. We use this norm systematically later.

In addition to the triangle inequality and submultiplicativity (from which multiplication is continuous for the norm topology) the $F$ - norm satisfies the Cauchy-Schwartz inequality ([21] page 10). Note that for diagonizable matrices, the $F$ - norm is the square root of the sum of the squares of the norms of the eigenvalues of the matrix.

There are many norms on $\mathbb{C}^{\lambda+1}$, all giving the same topology. For example, on $\mathbb{C}$ we have the usual norm $|\cdot|$, inducing on $\mathbb{C}^{\lambda+1}$ the norm ("the 2 -norm") whose value is the square root of the sum of the squares of the absolute values of the entries. This norm, and the Frobenius norm, are both instances of Schatten 2-norms. When we refer later to norms, it will be to such norms, unless we deal explicitly with operator norms.

We consider the elements of $M_{\lambda+1}(\mathbb{C})$ as linear operators $\phi$ from $\left(\mathbb{C}^{\lambda+1},\|\cdot\|_{1}\right)$ to $\left(\mathbb{C}^{\lambda+1}, \| \cdot\right.$ $\|_{2}$ ). Then, for any ordered pair of norms on $\mathbb{C}^{\lambda+1}$ there is a corresponding operator norm on $M_{\lambda+1}(\mathbb{C})$. Later, when we consider ultraproducts of the $M_{\lambda+1}(\mathbb{C})$ we will return to discussion of such norms. We will use operator norms only with reference to Schatten 2 -norms.

From now on, we will assume that $M_{\lambda+1}(\mathbb{C})$ is equipped with a fixed norm $\|\cdot\|$ satisfying the Cauchy-Schwartz inequality. The topology on $M_{\lambda+1}(\mathbb{C})$ is independent of the norm, but in discussing convergence of series we will appeal to the fixed norm.

If $A$ is any matrix in $M_{\lambda+1}(\mathbb{C})$, one defines ([21] 1.1) the matrix exponential of $A$, denoted by $\exp (A)$, as the power series:

$$
\begin{equation*}
\exp (A)=\sum_{n=0}^{\infty} \frac{A^{n}}{n!}=I_{\lambda+1}+A+\frac{A^{2}}{2}+\frac{A^{3}}{3!}+\ldots \tag{1}
\end{equation*}
$$

where $I_{\lambda+1}$ denotes the $(\lambda+1) \times(\lambda+1)$ identity matrix. This exponential series converges in norm for all matrices, so the exponential of $A$ is well-defined. If $A$ is a $1 \times 1$ matrix, that is, a scalar $a$ of the field $\mathbb{C}$ ), then $\exp (A)=e^{a}$ where $e^{a}$ denotes the ordinary exponential of the element $a \in \mathbb{C}$.
Recall that the matrix exponential satisfies the following properties:
Proposition 5.1. Let $A, B \in M_{\lambda+1}(\mathbb{C})$ and $a, b \in \mathbb{C}$ we have:
(i) $\exp \left(0_{\lambda}\right)=I_{\lambda+1}$, where $0_{\lambda+1}$ denotes the zero matrix in $M_{\lambda+1}(\mathbb{C})$;
(ii) $\exp (a A) \cdot \exp (b A)=\exp ((a+b) A)$;
(iii) $\exp (A) \cdot \exp (-A)=I_{\lambda+1}$;
(iv) for $A$ and $B$ commuting, $\exp (A+B)=\exp (A) \cdot \exp (B)$;
(v) for an invertible matrix $B$, $\exp \left(B A B^{-1}\right)=B \exp (A) B^{-1}$;
(vi) $\operatorname{det}(\exp (A))=\exp (\operatorname{tr}(A))$.
(vii) If no two eigenvalues of $A$ have a difference belonging to $2 \pi . i . \mathbb{Z}$, then there exists a neighbourhood of $A$ on which exp is injective. The exponential is injective on the ball of radius $\log (2)$ around the origin, for the Frobenius norm.

Proof. See [21] Proposition 1 (b), (c), (d), Proposition 3 in section 1.1 and Proposition 7 in section 1.2. See also [1, Chapter 3]. It remains to prove (vi).

By (v), in order to prove (vi)one may assume without loss of generality that $A$ is in Jordan normal form. So $A$ can be written as a sum of a diagonal matrix $D$ and a nilpotent matrix $N$ with $N$ and $D$ commuting. So, by (iv) $\exp (D+N)=\exp (D) \cdot \exp (N)$. Therefore, $\operatorname{det}(\exp (A))=\operatorname{det}(\exp (D)) \cdot \operatorname{det}(\exp (N))=\operatorname{det}(\exp (D))$.

For non commuting matrices $A$ and $B$, the equality $\exp (A+B)=\exp (A) \cdot \exp (B)$ need not hold. In that case, the Baker-Campbell-Hausdorff formula can be used to express $\exp (A) \cdot \exp (B)($ see $[21$, Section 1.3]).

Now, using the matrix exponential, one defines the exponential map

$$
\exp : M_{\lambda+1}(\mathbb{C}) \rightarrow G L_{\lambda+1}(\mathbb{C}): A \mapsto \exp (A)
$$

and rephrasing, (part of) Proposition 5.1, we get that $\left(M_{\lambda+1}(\mathbb{C}), \exp , G L_{\lambda+1}(\mathbb{C})\right)$ is an exponential $\mathbb{C}$-algebra. Moreover, the map exp is surjective from $M_{\lambda+1}(\mathbb{C})$ to $G L_{\lambda+1}(\mathbb{C})$. (Every invertible matrix can be written as the exponential of some other matrix ([21] page 21).)

For future use, we recall some methods for explicitly calculating matrix exponentials.
Diagonalizable case. If a matrix $A \in M_{\lambda+1}(\mathbb{C})$ is diagonal $A=\operatorname{diag}\left(a_{1}, a_{2}, \ldots a_{\lambda+1}\right)$, then its exponential can be obtained by just exponentiating every entry on the diagonal: $\exp (A)=\operatorname{diag}\left(e^{a_{1}}, e^{a_{2}}, \ldots, e^{a_{\lambda+1}}\right)$.
This also allows one to exponentiate any diagonalizable (so-called semisimple) matrix $S \in M_{\lambda+1}(\mathbb{C})$. If $S=B D B^{-1}$ where $B$ is invertible and $D$ is diagonal, then, according to the property (v) in Proposition 5.1, we have that $\exp (S)=B \exp (D) B^{-1}$ and the exponential of the matrix $D$ is calculated as above.
Nilpotent case. Recall that a matrix $N \in M_{\lambda+1}(\mathbb{C})$ is nilpotent if $N^{q}=0$ for some positive integer $q$ ( without loss of generality $\leq \lambda+1$ ).
In this case, the matrix exponential $\exp (N)$ can be computed directly from the series expansion (expressed by (1)), as the series terminates after a finite number of terms: $\exp (N)=I_{\lambda+1}+N+\frac{N^{2}}{2}+\ldots+\frac{N^{q-1}}{(q-1)!}$.
General Case. Since any matrix $A \in M_{\lambda+1}(\mathbb{C})$ can be expressed uniquely as a sum $A=S+N$ where $S$ is diagonalizable, $N$ is nilpotent and $S \cdot N=N \cdot S$, then the exponential of $A$ can be computed by using the property (iv) of Proposition 5.1 and by reducing to the previous two cases, so:

$$
\exp (A)=\exp (S+N)=\exp (S) \cdot \exp (N)
$$

Note that this uniqueness easily translates, via quantifier elimination for algebraically closed fields, into a constructible version in the sense of algebraic geometry.

We will need a more thorough description of $\operatorname{Ker}(\exp )$. It is easy to see that the map exp is not injective. For instance, consider a nonzero diagonal matrix $I_{\lambda+1} \neq D \in M_{\lambda+1}(\mathbb{C})$, $D=\operatorname{diag}\left(d_{1}, d_{2}, \ldots d_{\lambda+1}\right)$, with its matrix exponential, $\exp (D)=\operatorname{diag}\left(e^{d_{1}}, e^{d_{2}}, \ldots, e^{d_{\lambda+1}}\right)$. Then, $\exp (D) \in \operatorname{Ker}(\exp )$ if and only if the entries of $D$ belong to the kernel of the standard complex exponential map, so if and only if $d_{1}, d_{2}, \ldots d_{\lambda+1} \in 2 \pi i \mathbb{Z}$.

Lemma 5.2. If the matrix $A \in M_{\lambda+1}(\mathbb{C})$ belongs to the kernel of exp, then $A$ is diagonalisable and its eigenvalues lie in the kernel of the exponential function $e$ in $\mathbb{C}$.
Proof: In order to determine whether $A \in \operatorname{Ker}(\exp )$, by Proposition 5.1 (v), since $\exp \left(B^{-1}\right.$. $A \cdot B)=B^{-1} \cdot \exp (A) \cdot B$, we have that $\exp A=I_{\lambda+1}$ if and only if $\exp \left(B^{-1} \cdot A \cdot B\right)=I_{\lambda+1}$, for any invertible matrix $B \in M_{\lambda+1}(\mathbb{C})$. Since $\mathbb{C}$ is algebraically closed, $A$ is conjugated to a matrix in Jordan normal form, which can be written as a sum of a diagonal matrix $D$ and a nilpotent matrix $N$ with $N$ and $D$ commuting. So, by Proposition 5.1 (iv), $\exp (D+N)=\exp (D) \cdot \exp (N)$.Now $\exp (N)$ is unipotent. So if $\exp (D+N)=I_{\lambda+1}$, then $\exp (D)=\exp (-N)$, so the diagonal matrix $\exp (D)=I_{\lambda+1}$. Thus the eigenvalues of $D$ are periods of the exponential on the complex numbers. Also, $\exp (N)$ must be $I_{\lambda+1}$. Finally, a simple calculation with the polynomial $\exp (N)$ (in $N$ ) gives $N=0_{\lambda+1}$. We conclude that the kernel of exp consists of the diagonalizable matrices with complex periods as eigenvalues.

Proposition 5.3. Each associative Lie algebra $\left(M_{\lambda+1}(\mathbb{C})\right.$, exp $)$ viewed as an $\mathcal{L}_{E}$ - structure is bi-interpretable with the exponential field $\left(\mathbb{C}, e^{x}\right)$.

Proof: (See [18]). We embed $M_{\lambda+1}(\mathbb{C})$ in the direct product $\mathbb{C}^{(\lambda+1)^{2}}$. Then since $\mathbb{C}$ is algebraically closed, any matrix is conjugate to a matrix in Jordan normal form, namely $D+N$ where $D$ and $N$ commute, $D$ is a diagonal matrix and $N$ a nilpotent matrix with $N^{\lambda+1}=0$. By Proposition 5.1, $\exp (D+N)=\exp (D) \cdot \exp (N)$. Furthermore, $\exp (N)=1+N+\cdots+N^{\lambda}$. For the other direction of interpretability, see [18].

However, note that the class of rings $\left\{M_{\lambda+1}(\mathbb{C}) ; \lambda \in \omega\right\}$ is undecidable since the class of their invertible elements $\left\{G L_{\lambda+1}(\mathbb{C}) ; \lambda \in \omega\right\}$ is undecidable (one interprets uniformly in $\lambda$ a class of finite models whose theory is undecidable) and this implies that the theory of any non-principal ultraproduct of the $M_{\lambda+1}(\mathbb{C})$ is undecidable. Moreover, note that one may replace $\mathbb{C}$ by an arbitrary field, and the group $G L_{\lambda+1}(\mathbb{C})$ by other algebraic groups like $S L_{\lambda+1}(\mathbb{C})$ (see [9]).

## 6. The universal enveloping algebra of $s l_{2}(\mathbb{C})$

Recall that the universal enveloping algebra $U$ of $s l_{2}(\mathbb{C})$ is an associative $\mathbb{C}$-algebra (hence, equipped by a Lie algebra structure) together with a canonical mapping $\sigma$ which is a Lie algebra homomorphism $\sigma: s l_{2}(\mathbb{C}) \rightarrow U$ such that, if $R$ is any associative $\mathbb{C}$-algebra and $f: s l_{2}(\mathbb{C}) \rightarrow R$ is a Lie algebra homomorphism, then there exists a unique algebra homomorphism $\Theta: U \rightarrow R$ sending 1 to 1 and such $f=\Theta \circ \sigma$ (see [6] chapter 2, sections $1,2)$.

Diagram 6.1. Let us choose as $R$ the Lie algebra $M_{2}(\mathbb{C})$ and as $f$ the Lie algebra homomorphism $f_{1}: \operatorname{sl}_{2}(\mathbb{C}) \rightarrow M_{2}(\mathbb{C})$, so there exists a unique algebra homomorphism $\Theta_{1}: U \rightarrow$ $M_{2}(\mathbb{C})$ such that (according to what just said above) the following diagram commutes.


Since the canonical mapping $\sigma$ of $s l_{2}(\mathbb{C})$ into $U$ is injective ([6] Proposition 2.1.9), from now on we will identify every element of $s l_{2}(\mathbb{C})$ to its canonical image in $U$.

By using this universal property of $U$, we can construct an exponential map over $U$. Let us define the exponential map from $U$ to $G L_{2}(\mathbb{C})$ as follows:

$$
\begin{gathered}
\operatorname{EXP}_{1}: \quad U \xrightarrow{\Theta_{1}} M_{2}(\mathbb{C}) \xrightarrow{\exp } G L_{2}(\mathbb{C}) \\
\operatorname{EXP}_{1}(\alpha)=\exp \left(\Theta_{1}(\alpha)\right) \quad \forall \alpha \in U .
\end{gathered}
$$

So, the values of $\operatorname{EXP}_{1}(U)$ are in $G L_{2}(\mathbb{C})$ and the restriction of $\operatorname{EXP}_{1}$ to $s l_{2}(\mathbb{C})$ coincides with the exponential map $\exp : s l_{2}(\mathbb{C}) \rightarrow G L_{2}(\mathbb{C})$ (viewing $s l_{2}(\mathbb{C}) \subset M_{2}(\mathbb{C})$ ), previously defined (see (1)). Note that the image of the restriction of exp to $s l_{2}(\mathbb{C})$ is included in $S L_{2}(\mathbb{C})$ (see Proposition $5.1(\mathrm{vi})$ ). Clearly $\left(U, E X P_{1}, G L_{2}(\mathbb{C})\right.$ ) is an exponential algebra.

Let $c=2 x \cdot y+2 y \cdot x+h^{2}$ be the Casimir element of $U$, where $x, y, h$ are the generators of $s l_{2}(\mathbb{C}) . c$ generates the center of $U$. Let us calculate $\operatorname{EXP}_{1}(c)$. First, let

$$
\begin{aligned}
\Theta_{1}(c) & =\Theta_{1}\left(2 x \cdot y+2 y \cdot x+h^{2}\right)=2\left(\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right) \cdot\left(\begin{array}{ll}
0 & 0 \\
1 & 0
\end{array}\right)+2\left(\begin{array}{ll}
0 & 0 \\
1 & 0
\end{array}\right) \cdot\left(\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right)+ \\
& +(\operatorname{diag}(1,-1))^{2}=2\left(\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right)+2\left(\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right)+\operatorname{diag}(1,1)=\operatorname{diag}(3,3) .
\end{aligned}
$$

By using the universal property of $U$, we have that $\operatorname{EXP}_{1}(c)=\exp \left(\Theta_{1}(c)\right)=\operatorname{diag}\left(e^{3}, e^{3}\right)$.
Now we want to describe the map $E X P_{1}$ on $U$. Recall that $U$ is a $\mathbb{Z}$-graded algebra with grading $g r(x)=1, \operatorname{gr}(y):=-1$ and so $g r(h)=0$; a $m$-homogeneous element $u \in U$ is an element such that $\operatorname{gr}(u)=m, m \in \mathbb{Z}$. So $U$ decomposes as a direct sum of $m$-homogeneous components $U_{m}$ consisting of $m$-homogeneous elements, $m \in \mathbb{Z}$,

$$
U=\oplus_{m \in \mathbb{Z}} U_{m} .
$$

Furthermore, every $m$-homogeneous component satisfies the following relation, depending on whether $m$ is positive or negative:

$$
\begin{aligned}
& U_{m}=x^{m} U_{0}=U_{0} x^{m} \\
& U_{m}=y^{|m|} U_{0}=\text { for every positive integer number } y^{|m|} \\
& \text { for every negative integer number } m
\end{aligned}
$$

where, as well-described in [11], the 0 -homogeneous component $U_{0}$ coincides with the ring of polynomials $\mathbb{C}[h, c]$ with coefficients in $\mathbb{C}$ and variables the diagonal matrix $h$ and the Casimir element $c$. Let $u_{m}$ be an element in $U_{m}$ for $m$ a positive integer. So, $u_{m}=$ $u_{o} \cdot x^{m}=x^{m} \cdot v_{0}$, for some $u_{0}, v_{0} \in U_{0}$, and $u_{m}^{2}=\left(u_{0} \cdot x^{m}\right) \cdot\left(x^{m} \cdot v_{0}\right)=u_{0} \cdot x^{2 m} \cdot v_{0}$. Applying $\Theta_{1}$ to $u_{m}$, we can see that $\Theta_{1}\left(u_{m}^{2}\right)=\Theta_{1}\left(u_{0} \cdot x^{2 m} \cdot v_{0}\right)=\Theta_{1}\left(u_{0}\right) \cdot \Theta_{1}(x)^{2 . m}$. $\Theta_{1}\left(v_{0}\right)=0$, (because $\Theta_{1}(x)^{2}=0$ ). By similar calculations, we can see that, $\forall u, v \in U$ with every degree different from $-1,0,1, \Theta_{1}(u \cdot v)=0$. Now, we focus on $U_{0}$, so pick an element $p=p(c, h)$ and calculate the corresponding value of $\operatorname{EXP}_{1}$. Since $\Theta_{1}(p(c, h))=$ $p\left(\Theta_{1}(c), \Theta_{1}(h)\right)=p(\operatorname{diag}(3,3), \operatorname{diag}(1,-1))$, we can deduce that $\forall p \in U_{0}, \Theta_{1}(p(c, h)=0)$ if and only if $p(3,1)=0$ and $p(3,-1)=0$. Note that the corresponding ideal is not prime. Anyway, $\Theta_{1}(p(c, h))$ is a diagonal matrix with its eigenvalues $p(3,1)$ and $p(3,-1)$, and the matrix $\operatorname{EXP}_{1}(p)=\operatorname{diag}\left(e^{p(3,1)}, e^{p(3,-1)}\right)$ with determinant equal to $e^{p(3,1)+p(3,-1)}$.
By what sketched above, $\Theta_{1}$ acts as zero on $U_{ \pm 2}, U_{ \pm 3}, \ldots$ So, we restrict our attention to $U_{-1}, U_{0}, U_{1}$. Let us pick up in $U_{-1} \oplus U_{0} \oplus U_{1}$ an element $\gamma=y p_{-1}(c, h)+p_{0}(c, h)+x p_{1}(c, h)$
where the polynomials $p(c, h), p_{0}(c, h), p_{1}(c, h)$ belong to $U_{0}$. We want to calculate the exponential value of $\gamma$, as follows:

$$
\begin{aligned}
\operatorname{EXP}_{1}(\gamma) & =\operatorname{EXP}_{1}\left(\left(y p_{-1}(c, h)\right)+\left(p_{0}(c, h)\right)+\left(x p_{0}(c, h)\right)\right)= \\
& =\exp \left(\Theta_{1}\left(y p_{-1}(h, c)\right)+\Theta_{1}\left(p_{0}(c, h)\right)+\Theta_{1}\left(x p_{1}(c, h)\right)\right)= \\
& \left.=\exp \left(\Theta_{1}(y) \Theta_{1}\left(p_{-1}(c, h)\right)+\Theta_{1}\left(p_{0}(c, h)\right)\right)+\Theta_{1}(x) \Theta_{1}\left(p_{1}(c, h)\right)\right) .
\end{aligned}
$$

Since the value of $\Theta_{1}$ calculated on any element in $U_{0}$ is represented by a diagonal matrix, so $\Theta_{1}\left(y p_{-1}(c, h)\right), \Theta_{1}\left(p_{0}(c, h)\right), \Theta_{1}\left(x p_{1}(c, h)\right)$ can be respectively represented by the diagonal matrices $\operatorname{diag}\left(a_{-1}, b_{-1}\right), \operatorname{diag}\left(a_{0}, b_{0}\right), \operatorname{diag}\left(a_{1}, b_{1}\right)$, where $a_{i}, b_{i} \in \mathbb{C}$, with $i=-1,0,1$. So, we have

$$
\begin{aligned}
\operatorname{EXP}_{1}(\gamma) & =\exp \left(\left(\begin{array}{cc}
0 & 0 \\
1 & 0
\end{array}\right) \cdot \operatorname{diag}\left(a_{-1}, b_{-1}\right)+\operatorname{diag}\left(a_{0}, b_{0}\right)+\left(\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right) \cdot \operatorname{diag}\left(a_{1}, b_{1}\right)\right)= \\
& =\exp \left(\left(\begin{array}{cc}
0 & 0 \\
a_{-1} & 0
\end{array}\right)+\operatorname{diag}\left(a_{0}, b_{0}\right)+\left(\begin{array}{cc}
0 & b_{1} \\
0 & 0
\end{array}\right)\right)= \\
& =\exp \left(\begin{array}{cc}
a_{0} & b_{1} \\
a_{-1} & b_{0}
\end{array}\right) .
\end{aligned}
$$

Thanks to these calculations, we can easily find the $\operatorname{EXP}_{1}$ of $x p_{1}(h, c)$ : indeed, $\operatorname{EXP}_{1}\left(x p_{1}(h, c)\right)=$ $\exp \left(\Theta_{1}\left(x p_{1}(h, c)\right)\right)=\exp \left(\begin{array}{cc}0 & b_{1} \\ 0 & 0\end{array}\right)=I_{2}+\left(\begin{array}{cc}0 & b_{1} \\ 0 & 0\end{array}\right)=\left(\begin{array}{cc}1 & b_{1} \\ 0 & 1\end{array}\right)$, (because the square of the matrix $x$, so of $x p_{1}(c, h)$ is null). So, $\operatorname{EXP}_{1}\left(x p_{1}(h, c)\right)=I_{2}+\Theta_{1}\left(x p_{1}(h, c)\right)$. A similar property holds for $y p_{-1}(c, h)$, in fact, $\operatorname{EXP}_{1}\left(y p_{-1}(c, h)\right)=\exp \left(\Theta_{1}\left(y p_{-1}(c, h)\right)\right)=$ $\exp \left(\begin{array}{cc}0 & 0 \\ a_{-1} & 0\end{array}\right)=\left(\begin{array}{cc}1 & 0 \\ a_{-1} & 1\end{array}\right)$.

## 7. Other exponential maps

In this section, we define other exponential maps over $U$ by using finite dimensional representations of $s l_{2}(\mathbb{C})$, that is, finite dimensional $s l_{2}(\mathbb{C})$-modules ([6] 1.2). (All our modules will be left modules). First, recall that by Weyl's theorem, any finite dimensional representation of $s l_{2}(\mathbb{C})$ can be decomposed as a direct sum of simple $s l_{2}(\mathbb{C})$-modules $([6]$ 1.8.5). For every positive integer $\lambda$, there exists a unique (up to isomorphism) simple $s l_{2}(\mathbb{C})$ module $V_{\lambda}$ of dimension $\lambda+1 ; V_{\lambda}$ can be described as the $\mathbb{C}$-vectorspace of all homogeneous polynomials of degree $\lambda$ with coefficients in $\mathbb{C}$ and variables $X$ and $Y$ (see [8, Chapter 5]). We decompose $V_{\lambda}$ with respect to the basis of monomials $X^{\lambda}, X^{\lambda-1} Y, \ldots, X Y^{\lambda-1}, Y^{\lambda}$, $V_{\lambda}=\oplus_{i=0}^{\lambda} \mathbb{C}\left[X^{\lambda-i} Y^{i}\right]$. The representation $f_{\lambda}$ of $s l_{2}(\mathbb{C})$ can be described as follows:

$$
\begin{aligned}
& x \text { acts as } X \frac{\partial}{\partial Y} \\
& y \text { acts as } Y \frac{\partial}{\partial X}, \\
& h \text { acts as } X \frac{\partial}{\partial X}-Y \frac{\partial}{\partial Y} .
\end{aligned}
$$

So, for $0<i \leq \lambda$, the basis element $X^{\lambda-i} Y^{i}$ is shifted to the left, by the action of $x$, sent to $i \cdot X^{(\lambda-i)+1} Y^{i-1}$ and for $i=0, X^{\lambda}$ is sent by $x$ to 0 . For $0 \leq i<\lambda$, the basis element $X^{\lambda-i} Y^{i}$ is shifted to the right, by the action of $y$, sent to $(\lambda-i) X^{(\lambda-i)-1} Y^{i+1}$ and for $i=\lambda$, $Y^{\lambda}$ is sent by $y$ to 0 . Each subspace generated by $X^{\lambda-i} Y^{i}$ is left invariant by the action
of $h$ : $X^{\lambda-i} Y^{i}$ is mapped to $(\lambda-2) X^{\lambda-i} Y^{i}$ (so the corresponding eigenvalue is equal to $\lambda-2 . i)$.

The $\mathbb{C}$-vectorspace $\operatorname{End}\left(V_{\lambda}\right)$ coincides with the $\mathbb{C}$-vectorspace $M_{\lambda+1}(\mathbb{C})$ of all $(\lambda+1) \times$ $(\lambda+1)$ matrices written with respect to a basis of eigenvectors for $h$.

More precisely, through the representation $f_{\lambda}$, the actions of $x, y$ and $h$ are represented respectively the following three $(\lambda+1) \times(\lambda+1)$ matrices $X_{\lambda+1}, Y_{\lambda+1}, H_{\lambda+1}, \lambda \in \omega-\{0\}$ :

$$
\begin{gather*}
X_{\lambda+1}=\left(\begin{array}{cccc}
0 & 1 & 0 \ldots & 0 \\
0 & 0 & 2 \ldots & 0 \\
\vdots & \vdots & & \lambda \\
0 & 0 & 0 \ldots & 0
\end{array}\right), \quad Y_{\lambda+1}=\left(\begin{array}{cccc}
0 & 0 & \ldots & 0 \\
\lambda & 0 & \ldots & 0 \\
0 & \lambda-1 & & 0 \\
\vdots & \vdots & & \\
0 & 0 & 1 & 0
\end{array}\right),  \tag{2}\\
H_{\lambda+1}=\operatorname{diag}(\lambda, \lambda-2, \ldots,-\lambda+2,-\lambda) .
\end{gather*}
$$

Note that the operator norm of $X_{\lambda+1}$ (respectively $Y_{\lambda+1}$ ) is equal to $|\lambda|$, as is the operator norm of $H_{\lambda}+1$. The norm of $\Theta_{\lambda}(c)$ is equal to $\left|\lambda^{2}+2 \lambda\right|$.

On the other hand the $F$-norm of $X_{\lambda+1}$ is equal to $\lambda(\lambda+1) / 2$.
For every positive integer $\lambda$, we have the following diagram.
Diagram 7.1. For any simple representation $V_{\lambda}$ of $\operatorname{sl}_{2}(\mathbb{C})$ of dimension $\lambda+1$ (with $\lambda \in$ $\omega-\{0\}$ ), let us consider the representation map $f_{\lambda}: s l_{2}(\mathbb{C}) \rightarrow M_{\lambda+1}(\mathbb{C})$, and the following (commutative) diagram determined by the universal property of $U$ :

where $\sigma$ is the canonical mapping (which is a Lie-algebra homomorphism) from $\operatorname{sl}_{2}(\mathbb{C})$ to $U$ and $\Theta_{\lambda}$ is the (unique) algebra homomorphism from $U$ to $M_{\lambda+1}(\mathbb{C})$ sending 1 to 1 making the diagram commutes.

Using the commutativity of the above diagram, we obtain that the images of $x, y, h$ by the representation map $\Theta_{\lambda}: U \rightarrow M_{\lambda+1}(\mathbb{C})$ coincide with their images by the representation $\operatorname{map} f_{\lambda}$, and so are equal to the matrices $X_{\lambda+1}, Y_{\lambda+1}, H_{\lambda+1}$, (see (2)).

The image by $\Theta_{\lambda}$ of the Casimir element $c$ in $U$ is given by the following calculation.

$$
\begin{aligned}
\Theta_{\lambda}(c) & =\Theta_{\lambda}\left(2 x \cdot y+2 y \cdot x+h^{2}\right)=2 \Theta_{\lambda}(x) \cdot \Theta_{\lambda}(y)+2 \Theta_{\lambda}(y) \cdot \Theta_{\lambda}(x)+\left(\Theta_{\lambda}(h)\right)^{2}= \\
& =2 X_{\lambda+1} \cdot Y_{\lambda+1}+2 Y_{\lambda+1} \cdot X_{\lambda+1}+H_{\lambda+1}^{2}= \\
& =\operatorname{diag}\left(\lambda^{2}+2 \lambda, \ldots, \lambda^{2}+2 \lambda\right)
\end{aligned}
$$

By the technique used for defining the exponential map $\operatorname{EXP}_{1}$ from $U$ to $G L_{2}(\mathbb{C})$, we can define the exponential map $\mathrm{EXP}_{\lambda}$ for every positive integer $\lambda$, as follows.

Definition 7.1. Let $\lambda \in \omega-\{0\}$. The exponential map $\operatorname{EXP}_{\lambda}$ over $U$ is obtained by composing $\Theta_{\lambda}$ with the natural exponential map $\exp$ from $M_{\lambda+1}(\mathbb{C})$ to $G L_{\lambda+1}(\mathbb{C})$ (see section 5):

$$
\operatorname{EXP}_{\lambda}(u)=\exp \left(\Theta_{\lambda}(u)\right) \quad \forall u \in U
$$

Proposition 7.2. $\forall \lambda \in \mathbb{N}-\{0\}$, the map $E X P_{\lambda}$ is surjective.
Proof. Since $\exp$ is surjective from $M_{\lambda+1}(\mathbb{C})$ to $G L_{\lambda+1}(\mathbb{C})$, it suffices to prove that $\Theta_{\lambda}: U \rightarrow M_{\lambda+1}(\mathbb{C})$ is surjective. The latter is deduced directly by Jacobson density theorem [14, Section 2.2]. For convenience of the reader, we indicate below the proof.

Let $V_{\lambda}$ be the irreducible representation of $s l_{2}(\mathbb{C})$ of dimension $\lambda+1$. As representation of $U$, we know by Schur's lemma, that $\operatorname{End}_{U}\left(V_{\lambda}\right) \cong \mathbb{C}$. Consider $\phi \in \operatorname{End}_{\mathbb{C}}\left(V_{\lambda}\right)(=$ $\left.M_{\lambda+1}(\mathbb{C})\right)$. Then by Jacobson density theorem we get that, for each finite subset of elements $v_{1}, \ldots, v_{\lambda+1} \in V_{\lambda}$, that there exists $u \in U$ such that $\bigwedge_{i=1}^{m}\left(\phi\left(v_{i}\right)=\Theta_{\lambda}(u) . v_{i}\right)$.
We can easily calculate (as matrices in $G L_{\lambda+1}(\mathbb{C})$ ) the values by $\operatorname{EXP}_{\lambda}$ of $x, y, h$, $c$, using on one hand that $\Theta_{\lambda}(x), \Theta_{\lambda}(y)$ are nilpotent matrices (in $\left.M_{\lambda+1}(\mathbb{C})\right)$ ), and on the other hand that $\Theta_{\lambda}(h), \Theta_{\lambda}(c)$ are diagonal matrices.

$$
\begin{aligned}
& \operatorname{EXP}_{\lambda}(x)=\exp \left(\Theta_{\lambda}(x)\right)=\exp \left(X_{\lambda+1}\right)=I_{\lambda+1}+X_{\lambda+1}+\frac{X_{\lambda+1}^{2}}{2}+\ldots+\frac{X_{\lambda+1}^{\lambda}}{\lambda!} \\
& \operatorname{EXP}_{\lambda}(y)=\exp \left(\Theta_{\lambda}(y)\right)=\exp \left(Y_{\lambda+1}\right)=I_{\lambda+1}+Y_{\lambda+1}+\frac{Y_{\lambda+1}^{2}}{2}+\ldots+\frac{Y_{\lambda+1}^{\lambda}}{\lambda!} \\
& \operatorname{EXP}_{\lambda}(h)=\exp \left(\Theta_{\lambda}(h)\right)=\exp \left(H_{\lambda+1}\right)=\operatorname{diag}\left(e^{\lambda}, e^{\lambda-2}, \ldots, e^{-\lambda+2}, e^{-\lambda}\right) \\
& \operatorname{EXP}_{\lambda}(c)=\exp \left(\Theta_{\lambda}(c)\right)=\exp \left(\operatorname{diag}\left(\lambda^{2}+2 \lambda, \ldots, \lambda^{2}+2 \lambda\right)\right)=\operatorname{diag}\left(e^{\lambda^{2}+2 \lambda}, \ldots, e^{\lambda^{2}+2 \lambda}\right)
\end{aligned}
$$

Furthermore, we easily see that $E X P_{\lambda}$ satisfies the properties properties of the matrix exponential exp described by Proposition 5.1.
Proposition 7.3. Let $\lambda \in \mathbb{N}-\{0\}$. Then $\left(U, E X P_{\lambda}, G L_{\lambda+1}(\mathbb{C})\right)$ is an exponential $\mathbb{C}$ algebra. More precisely, we have the following properties. Let $u, v \in U$ and let $a, b \in \mathbb{C}$, then :
(i) $E X P_{\lambda}\left(0_{U}\right)=I_{\lambda+1}$, where $0_{U}$ denotes the identity element (with respect to the addition) in $U$.
(ii) $E X P_{\lambda}(a \cdot u) \cdot E X P_{\lambda}(b \cdot u)=E X P_{\lambda}((a+b) \cdot u)$;
(iii) $E X P_{\lambda}(u) \cdot E X P_{\lambda}(-u)=I_{\lambda+1}$;
(iv) for $u$ and $v$ commuting, $E X P_{\lambda}(u+v)=E X P_{\lambda}(u) \cdot \exp (v)$;
(v) for an invertible element $v$ in $U, E X P_{\lambda}\left(v u v^{-1}\right)=\Theta_{\lambda}(v) E X P_{\lambda}(u) \Theta_{\lambda}(v)^{-1}$;

Proof. (i) By definition of $E X P_{\lambda}, \operatorname{EXP}_{\lambda}\left(0_{U}\right)=\exp \left(\Theta_{\lambda}\left(0_{U}\right)\right)=\exp \left(0_{\lambda}\right)=I_{\lambda+1}$ (see Proposition 5.1(i)).
(ii) $\operatorname{EXP}_{\lambda}(a u) \cdot \operatorname{EXP}_{\lambda}(b u)=\exp \left(\Theta_{\lambda}(a u)\right) \cdot \exp \left(\Theta_{\lambda}(b u)\right)=\exp \left(a \Theta_{\lambda}(u)\right) \cdot \exp \left(b \Theta_{\lambda}(u)\right)$. Since $\Theta_{\lambda}(u) \in M_{\lambda+1}(\mathbb{C})$ and Proposition 5.1, (ii) can be applied, then $\exp \left(a \Theta_{\lambda}(u)\right)$. $\exp \left(b \Theta_{\lambda}(u)\right)=\exp \left((a+b) \Theta_{\lambda}(u)\right)=\exp \left(\Theta_{\lambda}((a+b) u)\right)=\operatorname{EXP}_{\lambda}((a+b) u)$.
(iii) This follows immediately from the corresponding property for the matrix exponential.
(iv) First, note that if $u$ and $v$ commute in $U$, then $\Theta_{\lambda}(u)$ and $\Theta_{\lambda}(v)$ commute also (for $\Theta_{\lambda}$ is a homomorphism from $U$ to $M_{\lambda+1}(\mathbb{C})$ for every $\lambda$ ). Thus, by using Proposition 5.1 (iv) and the fact that $\Theta_{\lambda}$ is a homomorphism, we have: $\operatorname{EXP}_{\lambda}(u) \cdot \operatorname{EXP}_{\lambda}(v)=$ $\exp \left(\Theta_{\lambda}(u)\right) \cdot \exp \left(\Theta_{\lambda}(v)\right)=\exp \left(\Theta_{\lambda}(u)+\Theta_{\lambda}(v)\right)=\exp \left(\Theta_{\lambda}(u+v)\right)=\operatorname{EXP}_{\lambda}(u+v)$.
(v) The map $\Theta_{\lambda}$ is a morphism of associative rings, so if an element $v \in U$ is invertible, then so is $\Theta_{\lambda}(v)$. The result follows immediately by the corresponding property for the matrix exponential.

Note that since the Casimir element is central in $U$, its image $\Theta_{\lambda}(c)$ is central in $\Theta_{\lambda}(U) \subseteq M_{\lambda+1}(\mathbb{C})$, so for any $u \in U$, we get by Proposition 5.1 that $\exp \left(\Theta_{\lambda}(c)+\Theta_{\lambda}(u)\right)=$ $\exp \left(\Theta_{\lambda}(c)\right) \cdot \exp \left(\Theta_{\lambda}(u)\right)$. So, $\operatorname{EXP}_{\lambda}(c+u)=\operatorname{EXP}_{\lambda}(c) \cdot \operatorname{EXP}_{\lambda}(u)$.

As a direct consequence of the definition of the map $\operatorname{EXP}_{\lambda}$, we observe that $u \in$ $\operatorname{Ker}\left(\operatorname{EXP}_{\lambda}\right)$ if and only if $\Theta_{\lambda}(u) \in \operatorname{Ker}(\exp )$. So in order to describe $\operatorname{Ker}\left(\operatorname{EXP}_{\lambda}\right)$, we should say as much as possible about $\Theta_{\lambda}(u)$ for $u \in U$.

Proposition 7.4. Decompose $U=\oplus_{m \in \mathbb{Z}} U_{m}$. The representation map $\Theta_{\lambda}$ sends:
(i) an element $u_{0}$ of $U_{0}$ onto a diagonal matrix,
(ii) an element $u_{m} \in U_{m}, m>0$, onto an upper triangular matrix if $0<m \leq \lambda$, otherwise (when $m \geq \lambda+1) \Theta_{\lambda}\left(u_{m}\right)=0_{\lambda+1}$.
(iii) an element $u_{m} \in U_{m}, m<0$, is mapped to a lower triangular matrix, if $-\lambda \leq m \leq$ -1 and, otherwise, for $m \leq-\lambda-1$, to the zero matrix $0_{\lambda+1}$.

Proof. (i) Let $u_{0} \in U_{0}-\{0\}$; so $u_{0}$ is of the form $p(c, h)$ with $p\left(x_{1}, x_{2}\right) \in \mathbb{C}\left[\mathrm{x}_{1}, \mathrm{x}_{2}\right]$, where $x_{1}$ and $x_{2}$ are two commuting variables. We know that $\Theta_{\lambda}(p(c, h))=p\left(\Theta_{\lambda}(c), \Theta_{\lambda}(h)\right)$, where $\Theta_{\lambda}(c)$ and $\Theta_{\lambda}(h)$ are the diagonal matrices described respectively by 3 and 2 . Since any algebraic operation on diagonal matrices concerns just their diagonal entries, then for any polynomial $p\left(x_{1}, x_{2}\right) \in \mathbb{C}\left[x_{1}, x_{2}\right]$, we have that:

$$
\begin{equation*}
\Theta_{\lambda}(p(c, h))=\operatorname{diag}\left(p\left(\lambda^{2}+2 \lambda, \lambda\right), \ldots, p\left(\lambda^{2}+2 \lambda,-\lambda\right)\right) \quad\left(\in M_{\lambda+1}(\mathbb{C})\right) \tag{4}
\end{equation*}
$$

(ii) For the positive integer $m$, let $u_{m}$ be an element in $U_{m}$ of the form $u_{m}=x^{m} \cdot u_{0}$ where the 0 -component $u_{0}=p(c, h)$ as above. On one hand, suppose that $m \leq \lambda$. By using the fact that $\Theta_{\lambda}$ is a homomorphism and the values of $\Theta_{\lambda}(x)$ and $\Theta_{\lambda}(p(c, h))$ (described by (2) and (4) respectively) we have that $\Theta_{\lambda}\left(u_{m}\right)=\Theta_{\lambda}\left(x^{m} \cdot u_{0}\right)=\Theta_{\lambda}(x)^{m} \cdot \Theta_{\lambda}\left(u_{0}\right)=$ $X_{\lambda+1}^{m} \cdot \operatorname{diag}\left(p\left(\lambda^{2}+2 \lambda, \lambda\right), \ldots, p\left(\lambda^{2}+2 \lambda,-\lambda\right)\right)$, so $\Theta\left(u_{m}\right)$ is represented by the strictly upper triangular matrix with $\star_{l} \in \mathbb{C}, 1 \leq l \leq(\lambda+1)-m$.

$$
\left(\begin{array}{ccccc}
0 & 0 & \star_{1} & 0 \ldots & 0  \tag{5}\\
0 & 0 & 0 & \star_{2} \ldots & 0 \\
\vdots & \vdots & 0 & & \star_{m} \\
\vdots & \vdots & & 0 & 0 \\
0 & 0 & \ldots & 0 & 0
\end{array}\right)
$$

On the other hand, assume that $m \geq \lambda+1$. Since $\Theta_{\lambda}(x)$ is a nilpotent matrix, we can easily see: $\Theta_{\lambda}\left(u_{m}\right)=\Theta_{\lambda}(x)^{m} \cdot \Theta_{\lambda}\left(u_{0}\right)=0$.
(iii) Similarly, we can repeat the same argument for any element $u_{m}$, with $m<0$, of the form $y^{m} \cdot u_{0}$. So, for $-\lambda \leq m \leq-1$ the image by $\Theta_{\lambda}$ of $u_{m}$, is a lower triangular matrix of the form

$$
\left(\begin{array}{ccccc}
0 & 0 & 0 & 0 \ldots & 0  \tag{6}\\
0 & 0 & 0 & 0 \ldots & 0 \\
\star_{1} & 0 & 0 & 0 \ldots & 0 \\
0 & \star_{2} & 0 & 0 & \cdots \\
\vdots & \vdots & & 0 & \\
0 & \cdots & 0 & \star_{-m} \cdots &
\end{array}\right)
$$

If $m \leq-(\lambda+1)$. we have $\Theta_{\lambda}\left(u_{m}\right)=\Theta_{\lambda}\left(y^{m} \cdot \alpha_{0}\right)=0$.

Remark 1. An element $u_{0} \in U_{0}$ belongs to the kernel of $E X P_{\lambda}$ if and only if

$$
\begin{equation*}
\bigwedge_{0 \leq j \leq \lambda} p\left(\lambda^{2}+2 \lambda, \lambda-2 j\right) \in 2 \pi i \mathbb{Z} \tag{7}
\end{equation*}
$$

In fact, for $u_{0}=p(c, h)$ the diagonal matrix $\Theta_{\lambda} p(c, h)$ belongs to $\operatorname{Ker}(e x p)$ if and only if their diagonal entries described by (4) belongs to $\operatorname{Ker}(e)=2 \pi i \mathbb{Z}$.

Proposition 7.5. $E X P_{\lambda}$ maps an element $u$ of $U$ into $S L_{\lambda+1}(\mathbb{C})$ whenever

$$
\begin{equation*}
\operatorname{tr}\left(\Theta_{\lambda}(u)\right) \in 2 \pi i \mathbb{Z} \tag{8}
\end{equation*}
$$

In particular, if $u \in \oplus_{m \neq 0} U_{m}$, then its image by $E X P_{\lambda}$ lies always in $S L_{\lambda+1}(\mathbb{C})$.
Proof. For the first statement, it is enough to apply property (vi) of Proposition 5.1, so for any $u \in U$, the determinant of $\exp \left(\Theta_{\lambda}(u)\right)$ equals 1 if the trace of $\Theta_{\lambda}(u)$ belongs to $\operatorname{Ker}(e)=2 \pi i \mathbb{Z}$.
As to the second claim, first we can note that the map $E X P_{\lambda}$ maps $x, y$ and their powers into $S L_{\lambda+1}(\mathbb{C})$, because their images by $\Theta_{\lambda}$ are matrices of trace 0 . We get the same results with $x^{m}$ (respectively $y^{m}$ ). Since the subalgebra $U_{0}$ is sent to the subalgebra of diagonal matrices in $M_{\lambda+1}(\mathbb{C})$, the image of an element $\alpha_{m}=x^{m} \cdot \alpha_{0}$ in $U_{m}$ by $\Theta_{\lambda}$ is a matrix of trace 0 (as illustrated by (5)) and so its matrix exponential has determinant 1. The same argument holds where $\alpha_{m}=y^{m} \cdot \alpha_{0}$ (for negative $m$ ). Since the sum of matrices of trace 0 has trace 0 , an element of $\oplus_{m \neq 0} U_{m}$ is sent by $E X P_{\lambda}$ to $S L_{\lambda+1}(\mathbb{C})$.

When we restrict Proposition 7.5 to any element $u_{0}$ of $U_{0}$, where $u_{0}=p(c, h)$ (for some polynomial $\left.p\left(x_{1}, x_{2}\right) \in \mathbb{C}\left[x_{1}, x_{2}\right]\right)$ the condition (8) means that the sum of eigenvalues of $\Theta_{\lambda}(u), \sum_{0 \leq j \leq \lambda} p\left(\lambda^{2}+2 \lambda, \lambda-2 j\right)$, has to belong to $2 \pi i \mathbb{Z}$.
Put $p\left(x_{1}, x_{2}\right)=\sum_{l=0}^{d} q_{l}\left(x_{1}\right) x_{2}^{l}$, then $\sum_{0 \leq j \leq \lambda} \sum_{l=0}^{\lfloor d / 2\rfloor} q_{l}\left(\lambda^{2}+2 \lambda\right)(\lambda-2 j)^{2 l}=$ $\sum_{l=0}^{\lfloor d / 2\rfloor} q_{j}\left(\lambda^{2}+2 \lambda\right)\left[\sum_{0 \leq j \leq \lambda}(\lambda-2 j)^{2 l}\right]$.

Now, let us assume that $u_{0}$ is in the kernel of $\bigcap_{\lambda \in \mathbb{Z} ; \lambda>\lambda_{0}} E X P_{\lambda}$, for some $\lambda_{0}$. Then, $p\left(\lambda^{2}+2 \lambda, \lambda-j\right) \in 2 \pi i \mathbb{Z}$, for all $|\lambda|>\lambda_{0}$ and $0 \leq j \leq \lambda$.

In the remainder of this section, we will give a partial answer to the question of which elements $u$ of $U$ are such that $\Theta_{\lambda}(u) \in s u_{\lambda+1}$

Recall that $s u_{\lambda+1}:=\left\{A \in M_{\lambda+1}(\mathbb{C}): A^{*}=-A, \operatorname{tr}(A)=0\right\}$, and $S U_{\lambda+1}:=\{X \in$ $\left.G L_{\lambda+1}(\mathbb{C}): X \cdot X^{*}=I_{\lambda+1}, \operatorname{det}(X)=1\right\}$, where $X^{*}$ denotes the conjugate transpose of $X$; it is a compact Lie group.

Coming back first to the case $\lambda=1$, it is well-known that the exponential map exp (defined in $M_{2}(\mathbb{C})$ ) restricted to $s l_{2}(\mathbb{C})$ does not map it surjectively to its Lie group $S L_{2}(\mathbb{C})$ ([21] page 38). However if we restrict to the $\mathbb{R}$-subalgebra $s u_{2}$, exp is surjective onto the (compact) Lie group $S U_{2}(\mathbb{C})$ (see Lemma 2.a in section 2 of [21]). We have the following decomposition: $S L_{2}(\mathbb{C})=S U_{2}(\mathbb{C}) . B$, where $B$ is the subgroup of triangular matrices with determinant 1 and positive real diagonal entries ([21] page 39).

The surjectivity property of exp holds if one replaces $s u_{2}$ with $s u_{\lambda+1}$ and $S U_{2}$ by $S U_{\lambda+1}$ (see Corollary 2 in [21]).

Let $u \in U_{0}$, so $u=p(c, h)$. So, $\Theta_{\lambda}(u) \in s u_{\lambda+1}$, if $\sum_{j} p\left(\lambda^{2}+\lambda, \lambda-2 j\right)=0$ and for all $-\lambda \leq j \leq \lambda, p\left(\lambda^{2}+\lambda, \lambda-2 j\right)=-\bar{p}\left(\lambda^{2}+\lambda, \lambda-2 j\right)$. The last condition occurs, for instance if $p\left(x_{1}, x_{2}\right)$ is the multiple by the complex number $i$ of a polynomial with real coefficients.

Now consider elements $u \in \oplus_{m \neq 0} U_{m}$, namely $u=\sum_{\ell>0}\left(p_{\ell}(c, h) \cdot x^{\ell}+y^{\ell} \cdot q_{\ell}(c, h)\right)$ with $p_{\ell}, q_{\ell} \in \mathbb{C}[h, c]$. Then the condition under which $\Theta_{\lambda}(u) \in s u_{\lambda+1}$ is that $(\lambda-j) \cdot q_{\ell}\left(\lambda^{2}+\right.$
$\lambda, \lambda-2 j)=(-\lambda+j) \cdot p_{\ell}\left(\lambda^{2}+\lambda, \lambda-2 j\right)$, for all $-\lambda \leq j \leq \lambda$. Given a polynomial $p_{\ell}$, we can always find a polynomial $q_{\ell}$ (of degree $\leq \lambda-1$ ) meeting these $\lambda$ conditions, using Lagrange interpolation theorem.

So, given $u \in \oplus_{m>0} U_{m}$, there exists $u^{\prime} \in \oplus_{m<0} U_{m}$ such that $\Theta_{\lambda}\left(u+u^{\prime}\right) \in s u_{\lambda+1}$.

## 8. Exponentiations and ultraproducts

We will be considering a non principal ultraproduct of the Lie algebras $M_{\lambda+1}(\mathbb{C}), \lambda \in \omega$. Namely, let $\mathcal{V}$ be a non-principal ultrafilter on $\omega$ and consider the corresponding ultraproducts $\prod_{\mathcal{V}} M_{\lambda+1}(\mathbb{C})$ and $\prod_{\mathcal{V}} G L_{\lambda+1}(\mathbb{C})$.

By Łos's theorem, the structure $\left(\prod_{\mathcal{V}} M_{\lambda+1}(\mathbb{C}),+,-, 0,[\cdot, \cdot]\right)$ is a Lie algebra over $\mathbb{C}$ or over $\mathbb{C}^{*}:=\prod_{\mathcal{V}} \mathbb{C}$, which is infinite-dimensional.

We first observe the following.
Proposition 8.1. (i) If $u_{0}$ is any element of $U_{0}-\{0\}$, then there exists $\lambda_{0}$ such that for all $\lambda \geq \lambda_{0}$, we have $\Theta_{\lambda}\left(u_{0}\right) \neq 0$.
(ii) For any $u \in U-\{0\}$, there exists $\lambda_{0}$ such that for all $\lambda \geq \lambda_{0}$ we have $\Theta_{\lambda}(u) \neq 0$.

Proof. (i) Let $u_{0} \in U_{0}-\{0\}$; so $u_{0}$ is of the form $p(c, h)$ with $p\left(x_{1}, x_{2}\right) \in \mathbb{C}\left[x_{1}, x_{2}\right]$, where $x_{1}$ and $x_{2}$ are two commuting variables.

The claim can be deduced directly from [11, Lemma 19]. For convenience of the reader, we repeat the argument here. We argue by contradiction.

Assume that $\bigwedge_{0 \leq j \leq \lambda} p\left(\lambda^{2}+2 \lambda, \lambda-2 j\right)=0$. First, we choose $\lambda$ such that $p\left(\lambda^{2}+2 \lambda, x_{2}\right) \neq 0$, so as a polynomial in $j, p\left(\lambda^{2}+2 \lambda, \lambda-2 j\right)$ is non-trivial of degree $k$ and so the number of roots is bounded by $k$. So, if we choose $\lambda$ big enough, we will always find $j$ such that $p\left(\lambda^{2}+2 \lambda, \lambda-2 j\right) \neq 0$. Therefore, $\Theta_{\lambda}(p(c, h)) \neq 0$ for some $\lambda$.
(ii) Let $u \in U-\{0\}$, then there exists $m \in \mathbb{Z}$ such that its $m^{\text {th }}$ component $u_{m} \neq 0$. Assume that $m \geq 0$ and that $m$ is minimal such. Let $u_{m}=x^{m} u_{0}$, where $u_{0} \in U_{0}$. Let $p\left(x_{1}, x_{2}\right) \in \mathbb{C}\left[x_{1}, x_{2}\right]$ be such that $u_{0}=p(c, h)$. Write $p\left(x_{1}, x_{2}\right)=\sum_{l=0}^{d} q_{i}\left(x_{1}\right) x_{2}^{l}$. We can find (explicitly) an interval $[-r ; r]$ in $\mathbb{R}$ such that all the roots of the polynomial $q_{d}\left(x_{1}\right)$ are in that interval. Let $r^{\prime}=\max \{r, d\}$. Then if $\lambda>r^{\prime}$, then $q_{d}\left(\lambda^{2}+2 \lambda\right) \neq 0$ and so the polynomial $\sum_{l=0}^{d} q_{l}\left(\lambda^{2}+2 \lambda\right) x_{2}^{l}$ has less than $d$ roots and among the $\lambda+1$ elements of the form $(\lambda-2 j)$ where $0 \leq j \leq \lambda$, we have such $j$ with the property that $p\left(\lambda^{2}+2 \lambda,(\lambda-2 j)\right) \neq 0$.

Since the images of any homogeneous components $U_{m}$ with $-\lambda \leq m \leq \lambda$ are in direct sum and $\Theta_{\lambda}\left(u_{m}\right) \neq 0$, then we have $\Theta_{\lambda}(u) \neq 0$.

Define the obvious $\Theta:=\left[\Theta_{\lambda}\right]$ from $U$ to the ultraproduct of the $M_{\lambda+1}(\mathbb{C})$, over any non-principal ultrafilter $\mathcal{V}$ on $\omega$. By Proposition 8.1, the map $\Theta$ is an associative ring monomorphism. So, we get the following corollary.
Corollary 8.2. For any non-principal ultrafilter $\mathcal{V}$ on $\omega, U$ embeds in the associative Lie algebra $\prod_{\mathcal{V}} M_{\lambda+1}(\mathbb{C})$.

Recall that $U$ is a left and right Ore domain, so it has a left and right field of fractions which embeds in the ring $U^{\prime}$ of definable scalars of $U$. This ring $U^{\prime}$ has been shown to be von Neumann regular by I. Herzog ([11]), equivalently every left (right) principal ideal is generated by an idempotent. Moreover, since any $V_{\lambda}$ is also a $U^{\prime}$-module, we can send $r \in U^{\prime}$ in the direct product $\prod_{\lambda \in \omega} M_{\lambda+1}(\mathbb{C})$ by sending it in each factor to the element of $M_{\lambda+1}(\mathbb{C})$, representing its action on each $V_{\lambda}$.

Then, [16] explicitly identifies certain idempotents of $U^{\prime}$ of the form $e_{u}, u \in U$, corresponding to the projections on $\operatorname{ker}\left(\Theta_{\lambda}(u)\right)$ on $V_{\lambda}, \lambda \in \omega$. For instance $e_{x}$ is the projection on the highest weight space of $V_{\lambda}$. When $u \in U_{0}$, so of the form $p(c, h)$, with
$p\left(x_{1}, x_{2}\right) \in \mathbb{C}\left[x_{1}, x_{2}\right]$, they call $p$ standard if there are only finitely many $\lambda$ such that $p\left(\lambda^{2}+2 \lambda, \lambda-2 . j\right)=0$ for some $0 \leq j \leq \lambda$ (and non-standard otherwise). Note that if $u=p(c, h)$ with $p$ standard, then $\left[\Theta_{\lambda}(u)\right]_{\mathcal{V}}$ is invertible in $\prod_{\mathcal{V}} M_{\lambda+1}(\mathbb{C})$. (Note that the converse holds if $\left[\Theta_{\lambda}(u)\right]_{\mathcal{U}}$ is invertible with respect to any non-principal ultrafilter $\left.\mathcal{U}\right)$.

Let now $u=p(c, h) \in U_{0}$ be such that $p$ is non-standard, so for some non-principal ultrafilter $\mathcal{V}$ the action of $e_{u}$ in the ultraproduct $\prod_{\mathcal{V}} V_{\lambda}$ will be a non invertible element of the form $[(\operatorname{diag}(0, \ldots, 1, \ldots, 0,1, \ldots, 0)] \neq 0$, where the number of possible 0 's is bounded by the degree of $p$ with respect to the second variable.

We know that $\Theta$ is a surjection from $\prod_{\mathcal{V}} U$ to $\prod_{\mathcal{V}} M_{\lambda+1}(\mathbb{C})$ (see the proof of Proposition 7.2). Then, we will compose with the map

$$
\operatorname{Exp}: \prod_{\mathcal{V}} M_{\lambda+1}(\mathbb{C}) \rightarrow \prod_{\mathcal{V}} G L_{\lambda+1}(\mathbb{C}):\left[\mathrm{A}_{\lambda}\right]_{\mathcal{V}} \rightarrow\left[\exp \left(\mathrm{A}_{\lambda}\right)\right]_{\mathcal{V}}
$$

So, by composing with $\left[\Theta_{\lambda}\right]_{\mathcal{V}}$, we get a map $E X P^{*}=\operatorname{Exp}\left[\Theta_{\lambda}\right]_{\mathcal{V}}$ from $\prod_{\mathcal{V}} U$ to $\prod_{\mathcal{V}} G L_{\lambda+1}(\mathbb{C})$, which is surjective. The kernel of that map is in bijection with the kernel of Exp on $\prod_{\mathcal{V}} M_{\lambda+1}(\mathbb{C})$.
Definition 8.1. Let $E X P$ from $U$ to $\prod_{\mathcal{V}} G L_{\lambda+1}(\mathbb{C})$ be defined as follows:

$$
\operatorname{EXP}: U \rightarrow \prod_{\mathcal{V}} G L_{\lambda+1}(\mathbb{C}): \mathrm{u} \rightarrow\left[\operatorname{EXP}_{\lambda}(\mathrm{u})\right]_{\mathcal{V}}
$$

Proposition 8.3. Both $\left(U, E X P, \prod_{\mathcal{V}} G L_{\lambda+1}(\mathbb{C})\right)$ and $\left(U_{0}, E X P, \prod_{\mathcal{V}}\right.$ Diag $\left._{\lambda+1}(\mathbb{C})\right)$ are exponential $\mathbb{C}$-algebras. Moreover we have that $\operatorname{EXP}\left(\oplus_{m \neq 0} U_{m}\right) \subset \prod_{\mathcal{V}} S L_{\lambda+1}(\mathbb{C}), E X P\left(\oplus_{m \geq 0} U_{m}\right) \subset$ $\prod_{\mathcal{V}} U T_{\lambda+1}(\mathbb{C})$, and $E X P\left(U_{0}\right) \subset \prod_{\mathcal{V}} \operatorname{Diag}_{\lambda+1}(\mathbb{C})$.

Proof: A direct application of Łos Theorem shows that $E X P$ satisfies the properties stated for each $E X P_{\lambda}$ in Proposition 7.3.

Note that the above properties are independent of the non-principal ultrafilter $\mathcal{V}$ on $\omega$.
Question 8.1. What is the kernel of EXP?
It is the set of elements $u$ such that for a subset of $\lambda$ belonging to $\mathcal{V}, \exp \left(\Theta_{\lambda}(u)\right)=1$. So, the eigenvalues of $\Theta_{\lambda}(u)$ belong to $2 \pi i \cdot \mathbb{Z}$; does it translate into an independently interesting property of $u \in U$ ? For $u_{0} \in U_{0}$, we have the following answer. Let $p\left(x_{1}, x_{2}\right) \in \mathbb{C}\left[x_{1}, x_{2}\right]$ such that $u_{0}=p(c, h)$. Then, for almost all $\lambda$ and all $0 \leq j \leq \lambda$, we have $p\left(\lambda^{2}+\lambda, \lambda-2 j\right) \in$ $2 \pi i \cdot \mathbb{Z}$.

Proposition 8.4. Let $u:=p(c, h) \in U_{0}$, with $p\left(x_{1}, x_{2}\right) \in \mathbb{C}\left[x_{1}, x_{2}\right]$. Write $p\left(x_{1}, x_{2}\right)$ in the form $2 \pi i \cdot q\left(x_{1}, x_{2}\right)$. Then, if $u \in \operatorname{ker}(E X P)$, then $q\left(x_{1}, x_{2}\right) \in \mathbb{Q}\left[x_{1}, x_{2}\right]$.

Proof: Let $q\left(x_{1}, x_{2}\right)=\sum_{k=0}^{d} q_{k}\left(x_{1}\right) \cdot x_{2}^{k}$ and assume that $q(c, h) \in \operatorname{ker}(E X P)$. Then, the set $\left\{\lambda \in \omega: \bigwedge_{0 \leq j \leq \lambda} q\left(\lambda^{2}+2 \lambda, \lambda-2 j\right) \in 2 \pi i \cdot \mathbb{Z}\right\} \in \mathcal{V} \quad(\star)$.

Set $c_{k}:=q_{k}\left(\lambda^{2}+2 \lambda\right)$ and consider the following system of linear equations, with $z_{\ell} \in \mathbb{Z}$, $0 \leq \ell \leq n$ :
$\left(\begin{array}{ccccc}1 & y_{0} & y_{0}^{2} & \cdots & y_{0}^{d} \\ 1 & y_{1} & y_{1}^{2} & \cdots & y_{1}^{d} \\ & \vdots & & & \\ 1 & y_{n} & y_{n}^{2} & \cdots & y_{n}^{d}\end{array}\right) \cdot\left(\begin{array}{c}c_{0} \\ c_{1} \\ \vdots \\ c_{n}\end{array}\right)=\left(\begin{array}{c}z_{0} \\ z_{1} \\ \vdots \\ z_{n}\end{array}\right)$.

When $n=d$, the determinant of the (square) matrix $\left(\begin{array}{ccccc}1 & y_{0} & y_{0}^{2} & \cdots & y_{0}^{d} \\ 1 & y_{1} & y_{1}^{2} & \cdots & y_{1}^{d} \\ & \vdots & & & \\ 1 & y_{d} & y_{d}^{2} & \cdots & y_{d}^{d}\end{array}\right)$
is equal to $\prod_{0 \leq n_{1}<n_{2} \leq d}\left(y_{n_{1}}-y_{n_{2}}\right)$. So it is a non zero integer whenever the $y_{i}$ 's are $d$ pairwise distinct integers and so in that case, the coefficients $c_{k}$ are rational numbers.

So, it suffices to express hypothesis ( $\star$ ) for $\lambda>d$.
Now, write each $q_{k}\left(x_{1}^{2}+2 x_{1}\right)$ as $q_{k}^{\prime}\left(x_{1}\right)=\sum_{h=0}^{d_{k}} f_{h} \cdot x_{1}^{h}$ and again write the system of equations expressing that each $q_{k}\left(\lambda^{2}+2 \lambda\right) \in \mathbb{Q}$, for $\lambda \in \omega$. Let $q_{j} \in \mathbb{Q}, 0 \leq j \leq n$.

$$
\left(\begin{array}{ccccc}
1 & x_{0} & x_{0}^{2} & \cdots & x_{0}^{d_{k}} \\
1 & x_{1} & x_{1}^{2} & \cdots & x_{1}^{d_{k}} \\
& \vdots & & & \\
1 & x_{n} & x_{n}^{2} & \cdots & x_{n}^{d_{k}}
\end{array}\right) \cdot\left(\begin{array}{c}
f_{0} \\
f_{1} \\
\vdots \\
f_{n}
\end{array}\right)=\left(\begin{array}{c}
q_{0} \\
q_{1} \\
\vdots \\
q_{n}
\end{array}\right)
$$

Then, again when $n=d_{k}$, the determinant of the (square) matrix

$$
\mathrm{x}\left(\begin{array}{ccccc}
1 & x_{0} & x_{0}^{2} & \cdots & x_{0}^{d_{k}} \\
1 & x_{1} & x_{1}^{2} & \cdots & x_{1}^{d_{k}} \\
& \vdots & & & \\
1 & x_{d_{k}} & x_{d_{k}}^{2} & \cdots & x_{d_{k}}^{d_{k}}
\end{array}\right)
$$

is equal to $\prod_{0 \leq n_{1}<n_{2} \leq d_{k}}\left(x_{n_{1}}-x_{n_{2}}\right)$. So it is a non zero integer whenever the $x_{i}$ 's are $d_{k}$ pairwise distinct integers and so in that case, the coefficients $f_{k}$ are rational numbers. So, it suffices to express hypothesis $(\star)$ for $d_{k}+1$ values of $\lambda$ 's, as soon as $\lambda>d$.

Remark 2. We have a partial converse to the above proposition. Namely, let $q\left(x_{1}, x_{2}\right)=$ $\sum_{k=0}^{d} q_{k}\left(x_{1}\right) \cdot x_{2}^{k}$, where each $q_{k}\left(x_{1}\right) \in \mathbb{Q}\left[x_{1}\right]$, so can be written as $1 / n_{k} \cdot \sum_{h=1}^{d_{k}} z_{h} \cdot x_{1}^{h}+q_{0, k}$, where $n_{k} \in \mathbb{N}-\{0\}, \quad z_{h} \in \mathbb{Z}$ and $q_{0, k} \in \mathbb{Q}$.

If, we assume in addition that each $q_{0, k} \in \mathbb{Z}$, then for some ultrafilter $\mathcal{V}, 2 \pi i \cdot q(c, h) \in$ $\operatorname{ker}(\mathrm{EXP})$. Indeed, let $n=\operatorname{lcm}\left\{n_{k}: 0 \leq k \leq d\right\}$. Then we choose an ultrafilter $\mathcal{V}$ containing $2 n \cdot \omega$.

So, if $\lambda=2 n \cdot m$, for some $m \in \omega, q_{k}\left(\lambda^{2}+2 \lambda\right)=n / n_{k} \cdot \sum_{h=1}^{d_{k}} z_{h} \cdot\left(2 n \cdot m^{2}+2 m\right)^{h}+q_{0, k}$, then $q_{k}\left(\lambda^{2}+2 \lambda\right) \in \mathbb{Z}$ and so $\left\{\lambda \in \omega: \bigwedge_{0 \leq j \leq \lambda} q\left(\lambda^{2}+2 \lambda, \lambda-2 j\right) \in 2 \pi i \cdot \mathbb{Z}\right\} \in \mathcal{V}$.

Corollary 8.5. Let $u:=p(c, h) \in U_{0}$, with $p\left(x_{1}, x_{2}\right) \in \mathbb{C}\left[x_{1}, x_{2}\right]$. Write $p\left(x_{1}, x_{2}\right)$ in the form $2 \pi i \cdot q\left(x_{1}, x_{2}\right)$. Write $q\left(x_{1}, x_{2}\right)=\sum_{k=0}^{d} q_{k}\left(x_{1}\right) \cdot x_{2}^{k}$, with $q_{k}(x) \in \mathbb{Q}\left[x_{1}\right]$.

Then, $u \in \operatorname{ker}(E X P)$ for all non-principal ultrafilters on $\omega$, if and only if $q\left(x_{1}, x_{2}\right) \in$ $\mathbb{Q}\left[x_{1}, x_{2}\right]$ and for each $0 \leq k \leq d, q_{k}(0) \in \mathbb{Z}$.

Proposition 8.6. Let $u:=p(c, h) \in U_{0}$, with $p\left(x_{1}, x_{2}\right) \in \mathbb{C}\left[x_{1}, x_{2}\right]$. Write $p\left(x_{1}, x_{2}\right)$ in the form $2 \pi i \cdot q\left(x_{1}, x_{2}\right)$. Then, if $\operatorname{EXP}(u) \in \prod S L_{\lambda+1}(\mathbb{C})$, then $q\left(x_{1}, x_{2}\right) \in \mathbb{Q}\left[x_{1}, x_{2}\right]$.

Proof: Let $q\left(x_{1}, x_{2}\right)=\sum_{k=0}^{d} q_{k}\left(x_{1}\right) \cdot x_{2}^{k}$ and assume that the set $\left\{\lambda \in \omega: \operatorname{EXP}_{\lambda}(q(c, h)) \in\right.$ $\left.S L_{\lambda+1}(\mathbb{C})\right\} \in \mathcal{V}$. Equivalently, $\left\{\lambda \in \omega:\left[\sum_{\ell=0}^{\lfloor d / 2\rfloor} q_{\ell}\left(\lambda^{2}+2 \lambda\right) \cdot \sum_{0 \leq j \leq \lambda}(\lambda-2 j)^{2 \cdot \ell}\right] \in 2 \pi i \cdot \mathbb{Z}\right\} \in \mathcal{V}$ ( $\star$ ).

Set $c_{k}:=q_{k}\left(x_{1}\right)$ and consider the following system of linear equations, with $z_{\ell} \in \mathbb{Z}$, $0 \leq \ell \leq n$ :

$$
\left(\begin{array}{ccccc}
1 & y_{0} & y_{0}^{2} & \cdots & y_{0}^{d} \\
1 & y_{1} & y_{1}^{2} & \cdots & y_{1}^{d} \\
& \vdots & & & \\
1 & y_{n} & y_{n}^{2} & \cdots & y_{n}^{d}
\end{array}\right) \cdot\left(\begin{array}{c}
c_{0} \\
c_{1} \\
\vdots \\
c_{n}
\end{array}\right)=\left(\begin{array}{c}
z_{0} \\
z_{1} \\
\vdots \\
z_{n}
\end{array}\right) .
$$

When $n=d$, the determinant of the (square) matrix

$$
\left(\begin{array}{ccccc}
1 & y_{0} & y_{0}^{2} & \cdots & y_{0}^{d} \\
1 & y_{1} & y_{1}^{2} & \cdots & y_{1}^{d} \\
& \vdots & & & \\
1 & y_{d} & y_{d}^{2} & \cdots & y_{d}^{d}
\end{array}\right)
$$

is equal to $\prod_{0 \leq n_{1}<n_{2} \leq d}\left(y_{n_{1}}-y_{n_{2}}\right)$. So it is a non zero integer whenever the $y_{i}$ 's are $d$ pairwise distinct integers and so in that case, the coefficients $c_{k}$ are rational numbers.

So, it suffices to express hypothesis $(\star)$ for $\lambda>d$ and show that $\sum_{0 \leq j \leq \lambda}(\lambda-2 j)^{2 \cdot \ell}$ are pairwise distinct.

The rest of the proof is similar to the previous one.

## 9. Comparison with Serre's definition of an exponential map.

Recall that the completion $\hat{U}$ of $U([22])$ is defined as the infinite product $\Pi_{n=0}^{\infty} U^{n}$, where $U^{n}$ denotes the component of degree $n$ of $U$ (generated by all products of length $\leq n$ of generators $x, y$ of $U)$; an element $f \in \hat{U}$ can be represented as $\sum_{n=0}^{\infty} f_{n}$, where $f_{n} \in U^{n}$ (see [22] Part 1, chapter 4, paragraph 6). (Note that $U^{n}$ differs in general from $U_{n}$.)

Denote by $\mathcal{M}$ the ideal of $U$ generated by $x, y$ and let $\hat{\mathcal{M}}$ be the ideal of $\hat{U}$ generated by $\mathcal{M}$. For $f \in \hat{\mathcal{M}}$, J.-P. Serre defines $\exp _{S}$ by the usual formula $\exp _{S}(f):=\sum_{n} \frac{f^{n}}{n!}$. exp takes $\hat{\mathcal{M}}$ to $1+\hat{\mathcal{M}}$ (see [22] Part 1, chapter 4, paragraph 7). (Similarly, one can define $\log _{S}$ from $1+\hat{\mathcal{M}}$ to $\hat{\mathcal{M}}$ by $\log _{S}(1+x):=\sum_{n=1}^{\infty}(-1)^{n+1} \frac{x^{n}}{n}$, obtaining that $\exp _{S} \circ \log _{S}=1=$ $\log _{S} \circ \exp _{S}$ (see Theorem 7.2, Chapter 4, Part 1 in [22].)

Let $f:=\sum_{n=0}^{\infty} f_{n} \in \hat{U}$ and assume that $\sum_{n=0}^{\infty} \Theta_{\lambda}\left(f_{n}\right)$ belongs to $M_{\lambda+1}(\mathbb{C})$. Then, define $\hat{\Theta}(f):=\left[\sum_{n=0}^{\infty} \Theta_{\lambda}\left(f_{n}\right)\right]_{\mathcal{V}}$. Since, if $u \in U$, there exists a bound on the number of non-zero components, this map is always well-defined on the elements of $U$.

Proposition 9.1. For any $u \in \mathcal{M}, \hat{\Theta}\left(\exp _{S}(u)\right)=E X P(u)$.
Proof: Now, let $u \in \mathcal{M}$ with $u=\sum_{j=1}^{k} u_{j}$, where $u_{j} \in U^{j}$, then $u^{n}:=\left(\sum_{j=1}^{k} u_{j}\right)^{n}$. So, for each $m$, the $m$-component of $\exp _{S}(u)$ is a finite sum. Therefore $\hat{\Theta}\left(\exp _{S}(u)\right)$ is well-defined and $\hat{\Theta}\left(\exp _{S}(u)\right)=\left[\sum_{n} \Theta_{\lambda}\left(\frac{u^{n}}{n!}\right)\right]_{\mathcal{V}}=E X P(u)$.

## 10. A $\star$-NORM ON THE UNIVERSAL ENVELOPING ALGEBRA OF $s l_{2}(\mathbb{C})$.

Now, we would like to put a natural topology on $U$ in such a way that EXP is continuous. As in the previous section, we fix a non-principal ultrafilter $\mathcal{V}$ on $\omega$; let $\mathbb{C}^{*}:=\prod_{\mathcal{V}} \mathbb{C}$ be a non principal ultrapower of the field $(\mathbb{C},+, \cdot,-, 0)$. We equip $\mathbb{C}^{*}$ with the ultrapower of the standard complex conjugation, and in addition consider the ultraproduct of the various Frobenius norms. This takes values in the corresponding ultrapower of the reals, and satisfies the obvious modification of the norm axioms. By functoriality this norm comes formally from the ultraproduct of the hermitian sesquilinear forms.

Finally,by taking ultraproducts of normed algebras we get a natural notion of a $\star$-normed algebra, satisfying a natural version of the Cauchy-Schwartz inequality if the component algebras do. Since $\|\cdot\|_{\lambda+1}$ is a norm on each $M_{\lambda+1}(\mathbb{C})$, by the usual properties of an ultraproduct, we get a natural $\star$-norm $\|\cdot\|$ on $\prod_{\mathcal{V}} M_{\lambda+1}(\mathbb{C})$

This in turn, by 5.6 induces a star norm on $U$.
In the next proposition, we will give an estimate of the norm of $u \in U$ in terms of a polynomial in $\lambda$, with coefficients in $\mathbb{R}$.

Lemma 10.1. For each $u \in U$, there exist non-zero polynomials $q_{1}(),. q_{2}($.$) with coefficients$ in $\mathbb{R}$ such that for $\lambda$ sufficiently big, we have $q_{1}(\lambda) \leq\left\|\Theta_{\lambda}(u)\right\|_{F}^{2} \leq q_{2}(\lambda)$ and so $q_{1}\left([\lambda]_{\mathcal{V}}\right) \leq$ $\|u\| \leq q\left([\lambda]_{\mathcal{V}}\right)$.

Proof: Let us examine the norm of $\Theta_{\lambda}(u)$ for any element of $U$. Let $u=\sum_{m \in \mathbb{Z}} u_{m}$ (where $u_{m} \in U_{m}$ and $m \in \mathbb{Z}$ ). Moreover, for each $m \geq 0$, each $u_{m}=x^{m} \cdot p_{m}(c, h)$, and $u_{-m}=y^{m} \cdot p_{-m}(c, h)$, where $p_{m}\left(x_{1}, x_{2}\right), p_{-m}\left(x_{1}, x_{2}\right) \in \mathbb{C}\left[x_{1}, x_{2}\right]$. Assume that for some $k \in \mathbb{N}$, we have $u=\sum_{-k \leq m \leq k} u_{m}$, then we estimate $\left\|\Theta_{\lambda}(u)\right\|$ as follows. Assume $\lambda \geq k$, then

$$
\begin{aligned}
\left\|\Theta_{\lambda}(u)\right\|_{F}^{2}= & \left\|\Theta_{\lambda}\left(\sum_{m \in \mathbb{Z}} u_{m}\right)\right\|_{F}^{2}=\left\|\sum_{m \in \mathbb{Z}} \Theta_{\lambda}\left(u_{m}\right)\right\|_{F}^{2} \\
= & \left\|\sum_{m=-k}^{-1} \Theta_{\lambda}\left(u_{m}\right)+\Theta_{\lambda}\left(u_{0}\right)+\sum_{m=1}^{k} \Theta_{\lambda}\left(u_{m}\right)\right\|_{F}^{2} \\
= & \sum_{m=-k}^{-1}\left\|\Theta_{\lambda}\left(u_{m}\right)\right\|_{F}^{2}+\left\|\Theta_{\lambda}\left(u_{0}\right)\right\|_{F}^{2}+\sum_{m=1}^{k}\left\|\Theta_{\lambda}\left(u_{m}\right)\right\|_{F}^{2} \\
= & \sum_{m=-k}^{-1}\left\|\Theta_{\lambda}\left(y^{|m|} p_{-m}(c, h)\right)\right\|_{F}^{2}+\left\|\Theta_{\lambda}\left(p_{0}(c, h)\right)\right\|_{F}^{2}+\sum_{m=1}^{k}\left\|\Theta_{\lambda}\left(x^{m} p_{m}(c, h)\right)\right\|_{F}^{2} \\
= & \sum_{m=-k}^{-1}\left\|\Theta_{\lambda}(y)^{|m|} \cdot p_{-m}\left(\Theta_{\lambda}(c), \Theta_{\lambda}(h)\right)\right\|_{F}^{2}+\left\|p_{0}\left(\Theta_{\lambda}(c), \Theta_{\lambda}(h)\right)\right\|_{F}^{2} \\
& +\sum_{m=1}^{k}\left\|\Theta_{\lambda}(x)^{m} \cdot p_{m}\left(\Theta_{\lambda}(c), \Theta_{\lambda}(h)\right)\right\|_{F}^{2} .
\end{aligned}
$$

Then, we make the following estimate. Write $p_{m}\left(x_{1}, x_{2}\right)=\sum_{j=0}^{d_{m}} q_{j}\left(x_{1}\right) \cdot x_{2}^{j}$. Let $f_{j}\left(x_{1}\right)=$ $\frac{q_{j}\left(x_{1}\right)}{q_{d_{m}}\left(x_{1}\right)}$ and write the roots of $\sum_{j=0}^{d_{m}} f_{j}\left(x_{1}\right) \cdot x_{2}^{j}$ as $\alpha_{1}\left(x_{1}\right), \cdots, \alpha_{d_{m}}\left(x_{1}\right)$. Note that these roots are all in a ball of radius $M_{m}(\lambda):=1+\sum_{j=0}^{d_{m-1}}\left|f_{j}\left(\lambda^{2}+2 \lambda\right)\right|$; let $R_{m}(\lambda):=\sum_{j=0}^{d_{m}}\left|q_{j}\left(\lambda^{2}+2 \lambda\right)\right|$. Then $p_{m}\left(x_{1}, x_{2}\right)=q_{d_{m}}\left(x_{1}\right) . \prod_{j=1}^{d_{m}}\left(x_{2}-\alpha_{j}\left(x_{1}\right)\right)$. We have $\left|p_{m}\left(\lambda^{2}+2 . \lambda, \lambda-2 i\right)\right|=\mid q_{d_{m}}\left(\lambda^{2}+\right.$ $2 \lambda)\left|\cdot \prod_{j}\right|\left((\lambda-2 i)-\alpha_{j}\left(\lambda^{2}+2 . \lambda\right) \mid\right.$. Since the number of roots of $p_{m}\left(\lambda^{2}+2 \lambda, x_{2}\right)$ is at most $d_{m}$, there is at least one integer in the interval $[-\lambda ; \lambda]$ at distance bigger than $\left\lfloor\frac{\lambda}{d_{m}}\right\rfloor$ of all of these roots. So, $\left|q_{d_{m}}\left(\lambda^{2}+2 \lambda\right)\right|^{2} \cdot\left\lfloor\frac{\lambda}{d_{m}}\right\rfloor^{2 d_{m}} \leq \sum_{-\lambda \leq i \leq \lambda}\left|p_{m}\left(\lambda^{2}+2 \cdot \lambda, \lambda-2 i\right)\right|^{2} \leq$
$R_{m}(\lambda)^{2} \cdot\left(2 \cdot \lambda^{2 d_{m}+1}+1\right) \leq R_{m}(\lambda)^{2} \cdot\left(3 \lambda^{2 d_{m}+1}\right)$. So we get on one hand,

$$
\begin{aligned}
\left\|\Theta_{\lambda}(u)\right\|_{F}^{2} \leq & \sum_{m=-k}^{-1} \lambda^{2 \cdot|m|} \cdot \sum_{-\lambda \leq i \leq \lambda}\left|p_{-m}\left(\lambda^{2}+2 \lambda, \lambda-2 i\right)\right|^{2}+\sum_{-\lambda \leq i \leq \lambda}\left|p_{0}\left(\lambda^{2}+2 \lambda, \lambda-2 i\right)\right|^{2} \\
& +\sum_{m=1}^{k} \lambda^{2 m} \cdot \sum_{-\lambda \leq i \leq \lambda}\left|p_{m}\left(\lambda^{2}+2 \lambda, \lambda-2 i\right)\right|^{2} \\
\leq & \sum_{m=-k}^{-1} \lambda^{2 \cdot|m|} \cdot R_{m}(\lambda)^{2} \cdot\left(3 \lambda^{2 d_{m}+1}\right)+R_{0}(\lambda)^{2} \cdot\left(3 \lambda^{2 d_{0}+1}\right)+\sum_{m=1}^{k} \lambda^{2 m} \cdot R_{m}(\lambda)^{2} \cdot\left(3 \lambda^{2 d_{m}+1}\right)
\end{aligned}
$$

and on the other hand,

$$
\begin{aligned}
\left\|\Theta_{\lambda}(u)\right\|_{F}^{2} \geq & \sum_{m=-k}^{-1}(\lambda-k)^{2 \cdot|m|} \cdot \sum_{-\lambda \leq i \leq \lambda}\left|p_{-m}\left(\lambda^{2}+2 \lambda, \lambda-2 i\right)\right|^{2}+\sum_{-\lambda \leq i \leq \lambda}\left|p_{d_{0}}\left(\lambda^{2}+2 \lambda, \lambda-2 i\right)\right|^{2}+ \\
& \sum_{m=1}^{k}(\lambda-k)^{2 m} \cdot \sum_{-\lambda \leq i \leq \lambda}\left|p_{m}\left(\lambda^{2}+2 \lambda, \lambda-2 i\right)\right|^{2} \\
\geq & \sum_{m=-k}^{-1}(\lambda-k)^{2 \cdot|m|} \cdot\left|q_{d_{m}}\left(\lambda^{2}+2 \lambda\right)\right|^{2} \cdot\left\lfloor\frac{\lambda}{d_{m}}\right\rfloor^{2 d_{m}}+\left|q_{d_{0}}\left(\lambda^{2}+2 \lambda\right)\right|^{2} \cdot\left\lfloor\frac{\lambda}{d_{0}}\right\rfloor^{2 d_{0}}+ \\
& \sum_{m=1}^{k}(\lambda-k)^{2 m} \cdot\left|q_{d_{m}}\left(\lambda^{2}+2 \lambda\right)\right|^{2} \cdot\left\lfloor\frac{\lambda}{d_{m}}\right\rfloor^{2 d_{m}} .
\end{aligned}
$$

We can give an estimate of the degrees of $q_{1}$ and $q_{2}$. Namely, the degree of $q_{2}$ is equal to $\max _{-k \leq m \leq k}\left\{2 . \operatorname{deg}\left(R_{m}\right)+2|m|+2 . d_{m}+1\right\}$ and the degree $q_{1}$ is equal to $\max _{-k \leq m \leq k}\left\{4 . \operatorname{deg}\left(q_{d_{m}}\right)+\right.$ $\left.2|m|+2 . d_{m}\right\}$. (Note that $2 . \operatorname{deg}\left(q_{d_{m}}\right) \leq \operatorname{deg}\left(R_{m}\right)$.)

The ultraproduct of the norms induces a topology both on $\prod_{\mathcal{V}} M_{\lambda+1}(\mathbb{C})$ (under which + and . are continuous) and on the $U$. A basis of neighbourhoods $O_{\epsilon}$ of 0 (in $U$ ) is given by $O_{\epsilon}:=\{u \in U:\|u\| \leq \epsilon\}$, where $\epsilon \in \mathbb{R}^{*,+}-\{0\}$. When we just consider them as topological spaces, we will call them $\star$-normed spaces.

Then, we will consider the following topological subspaces $\prod_{\mathcal{V}} G L_{\lambda+1}(\mathbb{C})\left(\right.$ dense in $\left.\prod_{\mathcal{V}} M_{\lambda+1}(\mathbb{C})\right)$ and $\prod_{\mathcal{V}} S L_{\lambda+1}(\mathbb{C})$ which is a closed subspace of $\left.\prod_{\mathcal{V}} M_{\lambda+1}(\mathbb{C})\right)$.

Lemma 10.2. (See [19] Corollary 6.2.32.) Let $A, B \in M_{\lambda}(\mathbb{C})$, then $\| \exp (A+B)-$ $\exp (A)\left\|_{\lambda} \leq\right\| B \|_{\lambda} \exp \left(\|B\|_{\lambda}\right) \exp \left(\|A\|_{\lambda}\right)$. So, the exponential map is continuous on $M_{\lambda}(\mathbb{C})$ and Lipschitz continuous on each compact subset of $M_{\lambda}(\mathbb{C})$.

Proposition 10.3. Consider the $\star$-normed spaces $(U,\|\cdot\|)$ and $\left(\prod_{\mathcal{V}} M_{\lambda+1}(\mathbb{C}),\|\cdot\|_{\lambda+1}\right)$. The map EXP : $U \rightarrow \prod_{\mathcal{V}} G L_{\lambda+1}(\mathbb{C})$ is continuous and maps bounded sets to bounded sets. The image $\operatorname{EXP}\left(U_{0}\right)$ is an abelian subgroup of $\prod_{\mathcal{V}} G L_{\lambda+1}(\mathbb{C})$ and $\operatorname{EXP}\left(\oplus_{m \neq 0} U_{m}\right)$ is included in $\prod_{\mathcal{V}} S L_{\lambda+1}(\mathbb{C})$.
Proof: The continuity is clear from los and the preceding lemma.

Note that if the sequence $A_{\lambda+1} \in M_{\lambda+1}(\mathbb{C})$ is bounded, namely the sequence $\left\|\mathrm{A}_{\lambda+1}\right\|_{\lambda+1}$ is bounded, then the corresponding sequence $\left\|\exp \left(\mathrm{A}_{\lambda+1}\right)\right\|_{\lambda+1}$ is bounded. Indeed, by definition, $\exp \left(A_{\lambda+1}\right)=\sum_{k=0}^{\infty} \frac{A_{\lambda+1}^{k}}{(k)!}$, so the norm $\left\|\exp \left(\mathrm{A}_{\lambda+1}\right)\right\|_{\lambda+1} \leq e^{\left\|\mathrm{A}_{\lambda+1}\right\|_{\lambda+1}}$.

The last statement follows from Proposition 7.5.
Note that a priori, $\operatorname{EXP}(U)$ is not a subgroup of $\prod_{\mathcal{V}} G L_{\lambda+1}(\mathbb{C})$; we will denote by $<\operatorname{EXP}(U)>$ the subgroup generated by $\operatorname{EXP}(U)$ in $\prod_{\mathcal{V}} G L_{\lambda+1}(\mathbb{C})$. The Campbell-BakerHausdorff formula which expresses for two matrices $A, B, \exp (A) \cdot \exp (B)$ as $\exp (C)$ where $C$ is expressed as an infinite series in commutators in $A$ and $B$, can be translated back with $u$ and $v$ in place of $A$ and $B$ to express $\operatorname{EXP}(u) \cdot \operatorname{EXP}(v)$ in terms of an infinite series in $u$, $v$ ([21] section 1.3).

Does $<\operatorname{EXP}(U)>$ have finite width with respect to $\operatorname{EXP}(U)$, namely does there exist a finite number $k$ such that every element of $<\operatorname{EXP}(U)>$ can be written as a product of $k$ elements of $\operatorname{EXP}(U)$ ?

We consider the field $\mathcal{R}:=\left(\mathbb{R},+, ., 0,1, e^{x}\right)$, and we denote by $\mathcal{R}^{*}$ a non principal ultrapower of $\mathcal{R}$ with respect to the ultrafilter $\mathcal{V}$ on $\omega$. We will extend the exponential map EXP to $U \otimes \mathbb{R}^{*}$ as follows. Let $u \in U$ and $s:=\left[r_{\lambda}\right] \mathcal{V} \in \mathbb{R}^{*}$ with $r_{\lambda} \in \mathbb{R}$, then $\operatorname{EXP}(u \otimes s):=$ $\operatorname{Exp}\left[r_{\lambda} \cdot \theta_{\lambda}(u)\right]_{\mathcal{V}}=\operatorname{Exp}\left[\theta_{\lambda}\left(r_{\lambda} \cdot u\right)\right]_{\mathcal{V}}$ and $\operatorname{EXP}\left(\sum_{i} u_{i} \otimes s_{i}\right):=\operatorname{Exp}\left[\theta_{\lambda}\left(\sum_{i} u_{i} \cdot r_{i, \lambda}\right)\right] \mathcal{V}$, where $s_{i}:=\left[r_{i, \lambda}\right] \mathcal{V}$. Note that $\sum_{i} u_{i} \cdot r_{i, \lambda} \in U$. This is well defined.

We will say that a topological group $G$ is $\star$-path connected if given any two elements $h_{0}, h_{1} \in G$, there is a continuous map $g$ from $[0 ; 1]^{*}:=\mathbb{R}^{*} \cap[0 ; 1]$ to $G$ with $g(0)=h_{0}$ and $g(1)=h_{1}$.

Proposition 10.4. The subgroups $<E X P(U)>$ and $E X P\left(U_{0}\right)$ of $\prod_{\mathcal{V}} G L_{\lambda+1}(\mathbb{C})$ (respectively $<\operatorname{EXP}\left(U \otimes \mathbb{R}^{*}\right)>$ and $\operatorname{EXP}\left(U_{0} \otimes \mathbb{R}^{*}\right)$ are topological groups. Moreover, $<$ $\operatorname{EXP}\left(U \otimes \mathbb{R}^{*}\right)>$ and $\operatorname{EXP}\left(U_{0} \otimes \mathbb{R}^{*}\right)$ are $\star$-path connected.

Proof: First note that $\prod_{\mathcal{V}} G L_{\lambda+1}(\mathbb{C})$ is a topological group as an ultraproduct of topological groups. So, the subgroups $<\operatorname{EXP}(U)>, \operatorname{EXP}\left(U_{0}\right),<\operatorname{EXP}\left(U \otimes \mathbb{R}^{*}\right)>$ and $\operatorname{EXP}\left(U_{0} \otimes \mathbb{R}^{*}\right)$ are topological subgroups.

The groups $<\operatorname{EXP}\left(U \otimes \mathbb{R}^{*}\right)>$ and $\operatorname{EXP}\left(U_{0} \otimes \mathbb{R}^{*}\right)$ are $\star$-path connected. We only prove that $<\operatorname{EXP}\left(U \otimes \mathbb{R}^{*}\right)>$ is $\star$-path connected. Let $g_{0}, g_{1} \in<\operatorname{EXP}\left(U \otimes \mathbb{R}^{*}\right)>$. Then we can write $g_{1}=\operatorname{EXP}\left(u_{1}\right) \cdot \ldots \cdot \operatorname{EXP}\left(u_{n}\right)$ and $g_{0}=\operatorname{EXP}\left(v_{1}\right) \cdot \ldots \cdot \operatorname{EXP}\left(v_{m}\right)$, where $u_{1}, \cdots, u_{n}, v_{1}, \cdots, v_{m} \in U \otimes \mathbb{R}^{*}$. So, $g_{1}=g_{0} \cdot \operatorname{EXP}\left(y_{1}\right) \cdot \ldots \cdot \operatorname{EXP}\left(y_{k}\right)$, for some $y_{1}, \cdots, y_{k} \in U \otimes \mathbb{R}^{*}$. Let $t \in[0 ; 1]^{*}$ and set $g(t)=g_{0} \cdot \operatorname{EXP}\left(t \cdot y_{1}\right) \cdot \ldots \cdot \operatorname{EXP}\left(t \cdot y_{k}\right)$, so $g(0)=g_{0}$ and $g(1)=g_{1}$. Let us denote the set $\left\{g \in<\operatorname{EXP}(U)>: \exists t \in[0 ; 1]^{*} g=\right.$ $\left.E X P\left(t y_{1}\right) \cdot \ldots \cdot \operatorname{EXP}\left(t \cdot y_{k}\right)\right\}$ by $C_{g_{0}, g_{1}}$.

First, let us check that the map from $[0 ; 1]^{*}$ to $\operatorname{EXP}(U)$, sending $t$ to $\operatorname{EXP}(t u)$ is continuous at $t_{1} \in[0 ; 1]^{*}$.

Let $\epsilon \in[0 ; 1]^{*}$, then we have to find $\eta$ such that if $\left|t_{0}-t_{1}\right|<\eta$, then $\| \operatorname{EXP}\left(t_{0} \cdot u\right)-$ $\operatorname{EXP}\left(t_{1} \cdot u\right) \| \leq \epsilon$. We have $\operatorname{EXP}\left(t_{0} \cdot u\right)-\operatorname{EXP}\left(t_{1} \cdot u\right)=\operatorname{EXP}\left(t_{1} \cdot u\right) \cdot\left[\operatorname{EXP}\left(\left(t_{0}-t_{1}\right) \cdot u\right)-1\right]$. So, $\left.\left\|\operatorname{EXP}\left(t_{0} \cdot u\right)-\operatorname{EXP}\left(t_{1} \cdot u\right)\right\| \leq\left\|\operatorname{EXP}\left(t_{1} \cdot u\right)\right\| \cdot \| \operatorname{EXP}\left(\left(t_{0}-t_{1}\right) \cdot u\right)-1\right] \|$. Now, $\| \operatorname{EXP}\left(\left(t_{0}-\right.\right.$ $\left.\left.\left.t_{1}\right) \cdot u\right)-1\right]\left\|\leq\left|\left(t_{0}-t_{1}\right)\right| \cdot\right\| u \| \cdot e^{\left\|\left(\left(t_{0}-t_{1}\right) \cdot u\right)\right\|}$.

Then we use the fact that the product (possibly non commutative) of two continuous functions is continuous $(*)$. So, by induction on $n$, we may deduce that the map sending $t$ to $\operatorname{EXP}\left(t \cdot y_{1}\right) \cdot \operatorname{EXP}\left(t \cdot y_{2}\right) \cdot \ldots \cdot \operatorname{EXP}\left(t \cdot y_{k}\right)$ is also continuous.

Now suppose $<\operatorname{EXP}(U)>$ is the disjoint union of two open sets $U_{1}$ and $U_{2}$. Denote the intersection of $U_{1}$ (respectively $U_{2}$ ) with $C_{g_{0}, g_{1}}$ by $O_{1}$ (respectively $O_{2}$ ). The inverse image of $O_{1}$ and $O_{2}$ gives rise to a partition of $[0 ; 1]^{*}$, which is a contradiction.

For convenience of the reader, let us prove $(*)$. Let $f(t), g(t)$ be two continuous maps on the interval $[0 ; 1]^{*}$ and assume one of them is bounded. Then consider the map sending $t$ to the product $f(t) \cdot g(t)$; let us show it is continuous at $t_{1}$, assuming that $f$ is bounded. Estimate the difference: $f(t) \cdot g(t)-f\left(t_{1}\right) \cdot g\left(t_{1}\right)=\left(f(t)-f\left(t_{1}\right)\right) \cdot g\left(t_{1}\right)+f(t) \cdot\left(g(t)-g\left(t_{1}\right)\right)$. So, $\left\|f(t) \cdot g(t)-f\left(t_{1}\right) \cdot g\left(t_{1}\right)\right\| \leq\left\|\left(f(t)-f\left(t_{1}\right)\right)\right\| \cdot\left\|g\left(t_{1}\right)\right\|+\|f(t)\| \cdot\left\|\left(g(t)-g\left(t_{1}\right)\right)\right\|$. Note that the map sending $t$ to $E X P(t u)$ is bounded. Indeed, $\|E X P(t u)\| \leq e^{\|t \cdot u\|} \leq e^{|t| \cdot\|u\|} \leq e^{\|u\|}$.

## 11. The asymptotic cone

In the previous section, we embedded $U$ in a $\star$-normed space, namely $\prod_{\mathcal{V}} M_{\lambda+1}(\mathbb{C})$. Here, we will embed $U$ into a complete metric space (with an $\mathbb{R}$-valued metric) which will be the asymptotic cone associated with the family of normed algebras $M_{\lambda+1}(\mathbb{C}), \lambda \in \omega$, and a non-principal ultrafilter $\mathcal{V}$ on $\omega$. We will first endow each $M_{\lambda+1}(\mathbb{C})$ with a new norm scaled down by $\lambda$; this norm differs from the norms we previously introduced in the fact that the norms of $\Theta_{\lambda}(x), \Theta_{\lambda}(y), \Theta_{\lambda}(c), \Theta_{\lambda}(h)$ will be a multiple of $\lambda$ (see Proposition 11.2).

Even though they didn't name it asymptotic cone, it was introduced by L. van den Dries and A. Wilkie when they revisited Gromov's proof that a finitely generated group of polynomial growth is nilpotent-by-finite. Given a group of polynomial growth, M. Gromov associated a converging sequence of discrete metric spaces scaled down by a sequence of well-chosen natural numbers. Then, van den Dries and Wilkie associated with any finitely generated group $G$ a limited ultraproduct of discrete metric spaces quotiented out by infinitesimals. This space is usually denoted by $\operatorname{Cone}(X, \mathcal{V})$, where $X$ is a metric space associated with $G$ and $\mathcal{V}$ a non principal ultrafilter on $\omega$, note that $\operatorname{Cone}(X, \mathcal{V})$ may depend on $\mathcal{V}$ (see for instance [15], [5]). The advantage of using an ultraproduct construction is that one can easily transfer certain properties from the factors.

First, we introduce the map $\phi$ from $M_{\lambda+1}(\mathbb{C})$ to $\mathbb{N}$, sending $A \in M_{\lambda+1}(\mathbb{C})$ to the number of non-zero coefficients of $A$. Of course, $\phi(A)=0$ iff $A=0$.

Let us check that
(1) $\phi(A+B) \leq \phi(A)+\phi(B)$,
(2) $\phi(A \cdot B) \leq \phi(A) \cdot \phi(B)$

We denote the $i j$ coefficient of $A+B$ by $(A+B)_{i j}$. We have that if $(A+B)_{i j} \neq 0$, then either $A_{i j} \neq 0$ or $B_{i j} \neq 0$.

Let $C:=A \cdot B$, then $C_{i j}=\sum_{k} A_{i k} \cdot B_{k j}$ and so $C_{i j} \neq 0$ implies that for some $k, A_{i k} \neq 0$ and $B_{k j} \neq 0$. We prove the second claim by induction on the number $\phi(C)$. For $\phi(C)=1$, it is clear. By induction suppose that for any $1 \leq n \leq m$, if $\phi(C)=n$, then for some 2 -tuple ( $k_{1}, k_{2}$ ) with $k_{1} \geq 1, k_{2} \geq 1$, such that $\phi(A) \geq k_{1}$ and $\phi(B) \geq k_{2}$ and $n \leq k_{1} \cdot k_{2}$.

Assume now that $\phi(C)=m+1$, so there are $m+1$ tuples $(i, j)$ with $C_{i j} \neq 0$. For each of these tuples, there are two tuples $(i, k),(k, j)$ such that $A_{i k} \neq 0$ and $B_{k j} \neq 0$. By induction corresponding to the first $m$ non-zero tuples, we know that there are $k_{1}$ (respectively $k_{2}$ ) non-zero coefficients of the matrix $A$ (respectively of the matrix $B$ ) which are non-zero and such that $m \leq k_{1} . k_{2}$. Corresponding to the $m+1$ non-zero coefficient of $C$, there exists another non-zero coefficient of either $A$ or $B$ and so either $\phi(A) \geq k_{1}+1$, or $\phi(B) \geq k_{2}+1$, so $m+1 \leq \min \left\{\left(k_{1}+1\right) \cdot k_{2}, k_{1} \cdot\left(k_{2}+1\right)\right\}$.

So, this map $\phi$ defines a norm on $M_{\lambda+1}(\mathbb{C})$, that we will denote by $\|\cdot\|_{c, \lambda+1}$.
In the ultraproduct $\prod_{\mathcal{V}}\left(M_{\lambda+1}(\mathbb{C}), \frac{\|\cdot\|_{c, \lambda+1}}{\lambda}\right)$, we consider the set $\prod_{\mathcal{V}}^{*}\left(M_{\lambda+1}(\mathbb{C}), \frac{\|\cdot\|_{c, \lambda+1}}{\lambda}\right)$ of elements $\left[a_{\lambda}\right]$ such that for some natural number $N$, we have $\left\{\lambda \in \omega:\left\|a_{\lambda}\right\|_{c, \lambda} \leq N \cdot \lambda\right\} \in \mathcal{V}$. Then we quotient out this set by the equivalence relation $\sim$ defined by $\left[a_{\lambda}\right]_{\mathcal{V}} \sim\left[b_{\lambda}\right]_{\mathcal{V}}$ if
$\left(\frac{\left\|a_{\lambda}-b_{\lambda}\right\|_{c, \lambda}}{\lambda}\right) \rightarrow \mathcal{V} 0$. Let us denote the equivalence class of an element by $\left[a_{\lambda}\right]_{\sim}$ and by $s t$ the standard part of an element of $\prod_{\mathcal{V}} \mathbb{R}$ whose absolute value is bounded by some natural number.

On $X_{\mathcal{V}}:=\prod_{\mathcal{V}}^{*}\left(M_{\lambda+1}(\mathbb{C}), \frac{\|\cdot\|_{c, \lambda+1}}{\lambda}\right) / \sim$, we define the following distance with values in $\mathbb{R}^{\geq 0}$.

Let $a:=\left[a_{\lambda}\right]_{\sim}$ and $b:=\left[b_{\lambda}\right]_{\sim}$, then $d(a, b):=s t\left(\left[\frac{\left\|a_{\lambda}-b_{\lambda}\right\|_{c, \lambda}}{\lambda}\right]\right)$.
Lemma 11.1. The space $\left(X_{\mathcal{V}}(\mathbb{C}), d\right)$ is an infinite-dimensional complete metric space.
Proof: The only thing which remains to be checked is the completeness of the space, but this follows from the countable saturation of the ultraproduct.

We will say that $\left(X_{\mathcal{V}}(\mathbb{C}), d\right)$ is the asymptotic cone associated with $\left\{\left(M_{\lambda+1}(\mathbb{C}), \frac{\|\cdot\|_{c, \lambda+1}}{\lambda}\right) ; \lambda \in\right.$ $\mathbb{N}\}$ and $\mathcal{V}$.

Proposition 11.2. The universal enveloping Lie algebra $U$ of $s l_{2}(\mathbb{C})$ embeds in $\left(X_{\mathcal{V}}(\mathbb{C}), d\right)$ via its embedding in the ultraproduct of the matrix rings.

Proof: We proceed in two steps.
Firstly, we show that for any $u \in U,\left[\Theta_{\lambda}(u)\right]$ belongs to $\prod_{\mathcal{V}}^{*}\left(M_{\lambda+1}(\mathbb{C}), \frac{\|\cdot\|_{c, \lambda+1}}{\lambda}\right)$. This is direct by inspection of the proof of 5.4.

Secondly, let $u \in \oplus_{|j| \leq m} U_{j}$, then there exist $d_{0}, d_{1}, \ldots, d_{m}, d_{-1}, \ldots, d_{-m}$ such that for all $\lambda \in \mathbb{N}, \phi\left(\Theta_{\lambda}(u)\right)=\left(\lambda-d_{0}\right)+\sum_{j=1}^{m}\left((\lambda-i)-d_{i}\right)=\lambda . m-(m(m+1)) / 2-\sum_{i=1}^{m} d_{i}-\sum_{i=-1}^{-m} d_{-i}$. Again, this is seen by inspection of the proof of 5.4.

So if $\left[\Theta_{\lambda}(u)\right] \sim 0$, then $\left[\Theta_{\lambda}(u)\right]=0$.
We will denote the image of $U$ in $\left(X_{\mathcal{V}}(\mathbb{C}), d\right)$ by $U_{\sim}$.
Definition 11.1. A matrix $\left(a_{i j}\right)$ in $M_{\lambda+1}(\mathbb{C})$ is called a $m$-band matrix if there exists $m$ such that for any $1 \leq i, j \leq \lambda+1$, we have $a_{i j} \neq 0$ implies that $|i-j| \leq m$. (Namely, the non-zero entries of a $m$-band matrix are confined to a diagonal band comprising the main diagonal and the adjacent $m$ diagonals on either side. The band-width is equal to $2 . m+1$.

Proposition 11.3. Every element of $U$ acts by left multiplication on the image $U_{\sim}$ of $U$ in $\left(X_{\mathcal{V}}(\mathbb{C}), d\right)$, in a continuous way. More generally, any element $\left[a_{\lambda}\right]_{\mathcal{V}} \in \Pi_{\mathcal{V}} M_{\lambda+1}(\mathbb{C})$ acts by right multiplication on $\left(X_{\mathcal{V}}(\mathbb{C}), d\right)$, whenever there exists $m$ independently of $\lambda$ such that $a_{\lambda}$ is a m-band matrix.

Proof: Let $u, v \in U$. Let us show that $\left[\Theta_{\lambda}(u)\right] \cdot\left[\Theta_{\lambda}(v)\right]_{\sim}$ is well-defined. Namely, if $\left[\epsilon_{\lambda}\right] \in$ $\prod_{\mathcal{V}}\left(M_{\lambda+1}(\mathbb{C})\right.$ with $\left[\epsilon_{\lambda}\right] \sim 0$, then $\Theta_{\lambda}(u) \cdot \epsilon_{\lambda} \sim 0$. Assume that $u \in \oplus_{|j| \leq m} U_{j}$, so $\Theta_{\lambda}(u)$ is a band matrix of width $\leq m$. Namely $\Theta_{\lambda}(u)_{i j}=0$ unless $|i-j| \leq m$. So, if we denote by $c$ the matrix in $M_{\lambda+1}(\mathbb{C})$ which is the product $\Theta_{\lambda}(u) \cdot \epsilon_{\lambda}$, then $c_{i j}=\sum_{k=1}^{\lambda+1} \Theta_{\lambda}(u)_{i k} \cdot \epsilon_{k j}$. If we fix the matrix element $\epsilon_{k j}$, then there are at most $2 m$ indices $i$ such that $c_{i j} \neq 0$. Now since $\left[\epsilon_{\lambda}\right] \sim 0, \lim \frac{\phi\left(\epsilon_{\lambda}\right)}{\lambda}=0$. We have that $\phi(c) \leq \phi\left(\epsilon_{\lambda}\right) \cdot 2 m$, so $\lim \frac{\phi(c)}{\lambda} \leq \lim \frac{\phi\left(\epsilon_{\lambda}\right)}{\lambda} \cdot 2 m=0$.

This action is continuous. Let $\epsilon>0$, choose $\eta:=\frac{\epsilon \cdot \lambda}{\phi\left(\Theta_{\lambda}(u)\right)}$. Then for $v_{1}, v_{2} \in U_{\sim}$, if $d\left(v_{1}, v_{2}\right) \leq \eta$, then $d\left(u . v_{1}, u . v_{2}\right) \leq \epsilon$. Indeed, we have $\left\|\Theta_{\lambda}(u) \cdot \Theta_{\lambda}\left(v_{1}\right)-\Theta_{\lambda}(u) . \Theta_{\lambda}\left(v_{2}\right)\right\|_{c, \lambda}=$ $\phi\left(\Theta_{\lambda}(u) \cdot \Theta_{\lambda}\left(v_{1}\right)-\Theta_{\lambda}(u) \cdot \Theta_{\lambda}\left(v_{2}\right)\right)=\phi\left(\Theta_{\lambda}(u) \cdot\left(\Theta_{\lambda}\left(v_{1}\right)-\Theta_{\lambda}\left(v_{2}\right)\right)\right) \leq \phi\left(\Theta_{\lambda}(u)\right) \cdot \phi\left(\Theta_{\lambda}\left(v_{1}\right)-\right.$ $\left.\Theta_{\lambda}\left(v_{2}\right)\right)$.

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## References

[1] A. Baker, Matrix groups: An introduction to Lie Group Theory, Springer-Verlag, 2002.
[2] C.C. Chang, H.J. Keisler, Model theory, North-Holland, Amsterdam, 1973.
[3] R. Carter, Lie algebras of finite and affine type, Cambridge University Press, 2005.
[4] R. Carter, Simple groups of Lie type, John Wiley and Sons, 1988.
[5] Y. de Cornulier, Dimension of asymptotic cones of Lie Groups, Journal of Topology 1 (2008) 342-361.
[6] J. Dixmier, Enveloping Algebras, North Holland, 1977.
[7] L. van den Dries, Exponential rings, exponential polynomials and exponential functions, Pacific J. Math. 113 (1984) no. 1, 51-66.
[8] K. Erdmann, M. Wildon, Introduction to Lie algebras, Springer, SUMS series, 2006
[9] Ju.L. Ers̆ov, New examples of undecidable theories, Algebra i Logika Sem. 5 (1966) no. 5, 37-47.
[10] Goodearl K.R., Warfield R.B.Jr, An Introduction to Noncommutative Noetherian Rings, London Mathematical Society Student Texts 16, Cambridge University Press, 1989.
[11] I. Herzog, The pseudo-finite dimensional representations of $s l(2, k)$, Selecta Mathematica 7 (2001), 241-290.
[12] I. Herzog, S. L'Innocente, The nonstandard quantum plane, Annals of Pure and Applied Logic 156, Issue 1 (2008), 78-85.
[13] J. Humphreys, Introduction to Lie Algebras and Representation Theory, Graduate Texts in Mathematics 9, Springer, 1970.
[14] N. Jacobson, Structure of Rings. American Mathematical Society Colloquium Publications, 37, 1964.
[15] Kramer L., Tent K., Asymptotic cones and ultrapowers of Lie groups, The Bulletin of Symbolic Logic, vol. 10, number 2 (2004) 175-185.
[16] S. L'Innocente, A. Macintyre, Towards decidability of the theory of pseudo-finite dimensional representations of $s l_{2} k$; I, In: A. Ehrenfeucht, V.W. Marek, M. Srebrny. Andrzej Mostowski and Foundational Studies. IOS Press, 2007, 235-260.
[17] S. L'Innocente, M. Prest, Rings of definable scalars of Verma modules, Journal of Algebra and its applications, 6 (2007), no. 5, 779-787.
[18] A. Macintyre, Model theory of exponentials on Lie Algebras, Math. Struct. in Comp. Science (2008), vol. 18, 189-204.
[19] R.A. Horn, C.R. Johnson, Topics in Matrix Analysis, Cambridge University Press, 1994.
[20] M. Prest, Model theory and modules, Cambridge UP, 1988.
[21] W. Rossmann, Lie groups: An introduction through linear groups, Oxford University Press, 2002.
[22] J.P. Serre, Lie algebras and Lie groups: 1964 Lectures given at Harvard University, LNM 1500, Second Edition, 1992.
[23] J.S. Wilson, On pseudofinite simple groups, J. London Math. Soc. (2) 51 (1995), 471-490.
[24] M. Ziegler, Model theory of modules, Ann. Pure Appl. Logic 26 (1984), 149-213.
L'Innocente Sonia, Department of Mathematics and Computer Science, University of Camerino,
Via Madonna delle Carceri 9, 62032 Camerino (MC) Italy
E-mail address: sonia.linnocente@unicam.it
Angus MacIntyre, School of Mathematics, Queen Mary University of London, Mile End Road London E1 4NS, England

E-mail address: angus@dcs.qmul.ac.uk
Françoise Point, Institut de mathématique, Université de Mons, Le Pentagone, 20, place du Parc, B-7000 Mons, Belgium.

E-mail address: point@logique.jussieu.fr


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