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# Derivation and solution of the two-dimensional Toda lattice equations by use of the Iwasawa decomposition 

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#### Abstract

The Iwasawa decomposition is applied to the complexification of an arbitrary gauge group G over a two-dimensional space. The two-dimensional Toda lattice equations arise from the condition that the gauge field strength vanishes. The method is based on the generalisation of Yang's $R$-gauge for SU(2) to any semi-simple Lie group G. A parametrisation for the solutions to the equations governed by arbitrary classical groups is obtained.


## 1. Introduction

Recent years have seen the generation of much interest in the nonlinear lattice equations, discovered by Toda in 1966, which describe the motion of a one-dimensional system of particles. When the only interactions are between neighbouring particles and these have an exponential form, exact solutions of the equations have been found (Toda 1967). The surge of interest is due to the recognition that a rich mathematical structure, both algebraic and geometric, is associated with the equations $\ddagger$.

It is the association of the equations with the theory of groups and their corresponding Lie algebras which concerns us in this paper, so we shall briefly review this aspect. The equations of motion for a finite system of $(n+1)$ particles arranged in line are

$$
\begin{equation*}
\partial_{t}^{2} \rho_{i}=-\sum_{j=1}^{n} K_{i j} \exp \rho_{i}, \tag{1.1}
\end{equation*}
$$

where $\rho_{i}$ is the difference in the displacements of the $(i+1)$ th and the $i$ th particle. Toda's original model corresponds to the $n=\infty$ limit of this finite system. In (1.1), $K$ is an $n \times n$ matrix with the only non-zero entries as follows:

$$
\begin{array}{ll}
K_{i i}=2, & i=1,2, \ldots, n, \\
K_{i+1 i}=K_{i i+1}=-1, & i=1,2, \ldots, n-1 .
\end{array}
$$

[^0]However, this matrix arises in the theory of Lie groups as the Cartan matrix for $\mathrm{SU}(n+1)$. Furthermore, the Dynkin diagram for the corresponding Lie algebra pictorially resembles the system of particles interacting only with their neighbours. It is now apparent that the model may be generalised to include systems described by (1.1), but where ( $K_{i j}$ ) is the Cartan matrix for any semi-simple group. It has subsequently been realised that the integrability and solubility properties of the Toda lattice equations will also apply to the generalised equations.

Toda's model is not only of interest mathematically, since it also has some physical applications. There are similarities between the lattice and its continuum approximation, which is described by the Korteweg-de Vries equation and winch is studied in plasma physics. Other applications are found in gauge field theories. The simplest Toda lattice equation with $n=1$ and governed by $\mathrm{SU}(2)$ has the same form as the Liouville equation, which has appeared in $\operatorname{SU}(2)$ gauge theories. In Euclidean space, SU(2) instantons with cylindrical symmetry are solutions of this equation (Witten 1977), whilst in a broken $\operatorname{SU}(2)$ theory the only time-independent solution corresponds to the Prasad-Sommerfield monopole solution (Bais and Weldon 1978). Also in the string model, Omnés claims that the classical states correspond to solutions of the same equation (Omnés 1979).

The most recent development, and the one which has motivated our study, is the discovery that the Bogomolny-Prasad-Sommerfield equation for the spherically symmetric monopole for any group may also be written in the form of the Toda lattice equation (Olive 1980). The effective generalisation of the Toda model to the other groups has therefore suggested a way in which monopoles may be described in theories with an arbitrary semi-simple gauge group.

Our discussion thus far has been restricted to one-dimensional $S U(n+1)$ Toda lattices and also their generalisation to one-dimensional models with any semi-simple compact group. However, other models may be considered; for example, the onedimensional periodic Toda lattice represented by the extended Dynkin diagram and also two-dimensional models with one time and one space variable.

Leznov and Saveliev (1978, 1979a, b, 1980) have developed the theory of general lattices composed of a finite non-periodic chain in two time dimensions. The equations of motion for this system are given by

$$
\left(\partial_{i}^{2}+\partial_{s}^{2}\right) \rho_{i}=-4 \sum K_{i j} \exp \rho_{i} .
$$

Leznov and Saveliev have remarked that when $\left(K_{i j}\right)$ is a 2 nd-order generalised Cartan matrix for an infinite-dimensional contragredient Lie algebra, this becomes the sinhGordon equation.

Our particular study is concerned with similar non-periodic lattices and so with an arbitrary semi-simple group of finite rank. We specifically consider Euclidean space with the two variables ( $s, t$ ) and employ complex coordinates $u=s+\mathbf{i} t$ and $\bar{u}=s-\mathbf{i} t$. The generalised Toda lattice equations become

$$
\begin{equation*}
\partial_{u} \partial_{\bar{u}} \rho_{i}=-\sum_{i} K_{i j} \exp \rho_{i} \tag{1.2}
\end{equation*}
$$

or alternatively, defining

$$
\begin{equation*}
\rho_{i}=-\sum K_{i j} \psi_{i}, \tag{1.3}
\end{equation*}
$$

$$
\begin{equation*}
\partial_{u} \partial_{\bar{u}} \psi_{i}=\exp \left(-\sum_{i} K_{i j} \psi_{i}\right) . \tag{1.4}
\end{equation*}
$$

As the variables are complexified, the Lie group G governing the equation must also be complexified to

$$
\mathrm{G}_{\mathrm{c}}=\mathrm{G} \times \mathrm{i} \mathrm{G} .
$$

As a consequence of the ensuing non-compactness of $\mathrm{G}_{\mathrm{c}}$, we can apply the Iwasawa decomposition (Helgason 1978, Hermann 1966) to it in such a way that the origin of the Toda lattice variables will become apparent. It also leads to an easy means of parametrising the solutions of the equation.

Our method of deriving the Toda lattice equations is somewhat heuristic and so, to demonstrate its plausibility and neatness, we shall explain it here with reference to Yang's approach to the self-duality condition for $\mathrm{SU}(2)$ gauge fields. This approach is known as Yang's $R$-gauge method and is itself an application of the Iwasawa decomposition (Yang 1977, hereafter referred to as Y).

On four-dimensional Euclidean space with coordinates $\left\{x_{\mu}, \mu=1,2,3,4\right\}$ the self-duality condition for the gauge field strength $F_{\mu \nu}$ is

$$
\begin{equation*}
F_{\mu \nu}=\frac{1}{2} \varepsilon_{\mu \nu \rho \sigma} F^{\rho \sigma}, \tag{1.5}
\end{equation*}
$$

where $F_{\mu \nu}$ is written in terms of the gauge potentials $B_{\mu}$ by

$$
\begin{equation*}
F_{\mu \nu}=\partial_{\mu} B_{\nu}-\partial_{\nu} B_{\mu}+\left[B_{\mu}, B_{\nu}\right] . \tag{1.6}
\end{equation*}
$$

By complexifying the coordinates and defining four new variables

$$
\begin{array}{ll}
\sqrt{2} y=x_{1}+\mathrm{i} x_{2}, & \sqrt{2} \bar{y}=x_{1}-\mathrm{i} x_{2}, \\
\sqrt{2} z=x_{3}-\mathrm{i} x_{4}, & \sqrt{2} \bar{z}=x_{3}+\mathrm{i} x_{4}, \tag{1.7}
\end{array}
$$

and the corresponding components of the potentials, Yang is able to rewrite (1.5) in the form

$$
\begin{align*}
& F_{y z}=F_{\bar{y} \bar{z}}=0,  \tag{1.8}\\
& F_{y \bar{y}}+F_{z \bar{z}}=0 . \tag{1.9}
\end{align*}
$$

Equations (1.8) and (1.9) provide a means of determining $B_{\mu}$ : indeed, (1.8) imply that

$$
\begin{array}{ll}
B_{y}=R^{-1} \partial_{y} R, & B_{z}=R^{-1} \partial_{z} R, \\
B_{\bar{y}}=\bar{R}^{-1} \partial_{\bar{y}} \bar{R}, & B_{\bar{z}}=\bar{R}^{-1} \partial_{\bar{z}} \bar{R}, \tag{1.10}
\end{array}
$$

where $\operatorname{det} R=\operatorname{det} \bar{R}=1 . \bar{R}$ is defined such that in the real section

$$
\begin{equation*}
\bar{R}=\left(R^{\dagger}\right)^{-1} \tag{1.11}
\end{equation*}
$$

By making a gauge transformation, Yang argues that it is always possible to choose $R$ to have the lower triangular form

$$
R=\frac{1}{\sqrt{\phi}}\left(\begin{array}{ll}
1 & 0  \tag{1.12}\\
\rho & \phi
\end{array}\right)
$$

where $\phi$ is a real function and $\rho$ is a complex function on the real section. By using the relation (1.11), $\bar{R}$ is seen to have an upper triangular form. The resultant gauge is known as the $R$-gauge.

By substituting (1.10) into (1.9) when $R$ is given by the expression (1.12), the equation (1.9) simplifies (see equation (27) in Y) and the problem of determining $B_{\mu}$ is reduced to one of solving this simplified form for $\rho$ and $\phi$.

Our first observation is that in the sector $y=z \equiv u$, Yang's equation (1.9) (and so equation (27) in Y) becomes over the real section

$$
\begin{align*}
& \phi \partial_{u} \partial_{\bar{u}} \phi-\left(\partial_{u} \phi\right)\left(\partial_{\bar{u}} \phi\right)+\left(\partial_{u} \rho\right)\left(\partial_{\bar{u}} \bar{\rho}\right)=0,  \tag{1.13}\\
& \phi \partial_{\bar{u}}\left(\partial_{u} \rho\right)-2\left(\partial_{u} \rho\right) \partial_{\bar{u}} \phi=0,  \tag{1.14}\\
& \phi \partial_{u}\left(\partial_{\bar{u}} \bar{\rho}\right)-2\left(\partial_{\bar{u}} \bar{\rho}\right) \partial_{u} \phi=0 .
\end{align*}
$$

A notable fact is that (1.14) (hereafter called the subsidiary equations) are soluble. Indeed, if we let

$$
\begin{equation*}
\phi=\exp (-2 \psi) \tag{1.15}
\end{equation*}
$$

then the solutions are

$$
\begin{equation*}
\partial_{u} \rho=\varepsilon(u) \mathrm{e}^{-4 \psi}, \quad \partial_{\bar{u}} \bar{\rho}=\bar{\varepsilon}(\bar{u}) \mathrm{e}^{-4 \psi} . \tag{1.16}
\end{equation*}
$$

By putting (1.16) into (1.13), we obtain

$$
\begin{equation*}
\partial_{u} \partial_{\bar{u}} \psi=\varepsilon \bar{\varepsilon} \exp (-4 \psi) \tag{1.17}
\end{equation*}
$$

If we choose $\varepsilon \bar{\varepsilon}=1 \mathrm{in}$ (1.17), comparison of the resultant equation with (1.4) shows that we have obtained the simplest Toda lattice equation governed by $\mathrm{G}=\mathrm{SU}(2)$.

In the sector $y=z,(1.8)$ are automatically satisfied and (1.9) becomes

$$
\begin{equation*}
F_{u \bar{u}} \equiv \partial_{u} B_{\bar{u}}-\partial_{\bar{u}} B_{u}+\left[B_{u}, B_{\bar{u}}\right]=0 . \tag{1.18}
\end{equation*}
$$

We have shown that when $B_{u}$ and $B_{\bar{u}}$ are given by

$$
\begin{equation*}
B_{u}=R^{-1} \partial_{u} R, \quad B_{\bar{u}}=\bar{R}^{-1} \partial_{\bar{u}} \bar{R}, \tag{1.19}
\end{equation*}
$$

the equation (1.18) reduces to the Toda lattice equation.
Our derivation in $\S 2$ of the Toda lattice equations using the Iwasawa decomposition is based upon the example above. We show that the same argument may be applied even when the group governing the equation is extended to any semi-simple Lie group. The method also suggests a means of solving the equations governed by any classical group and this is presented in $\S 3$, together with some specific examples. The general parametrisation of the $\mathrm{SU}(n+1)$ solution is included in appendix 2 . In $\S 4$ we make some concluding remarks about the reality conditions on $B_{u}$ and $B_{\bar{u}}$ and also the similarity between the one-dimensional Toda model and our two-dimensional model.

## 2. Iwasawa decomposition and the Toda lattice equations

We first note that Yang's $R$-matrix (1.12), having complex entries and unit determinant, belongs to the group $\operatorname{SL}(2, \mathbb{C})$, which is the complexification of the original gauge group, $\operatorname{SU}(2)$. Furthermore $R$ has a lower triangular form with real diagonal entries. It is possible to show that such a form is obtained by applying the Iwasawa decomposition NAK to $\operatorname{SL}(2, \mathbb{C})$, where $K$ is the maximal compact subgroup and so is $\mathrm{SU}(2)$ in this case (Ardalan 1978). In particular, by choosing a gauge so that the element belonging to K in the decomposition is effectively unity, the explicit form (1.12) is obtained. Having observed these properties of the $R$-gauge, it is obvious that the method previously
outlined in the introduction of deriving the Toda lattice equations governed by $\operatorname{SU(2)}$ can be extended to derive the equations governed by $\mathrm{SU}(n+1)$. The Iwasawa decomposition applied to $\operatorname{SL}(n+1, \mathbb{C})$ specifies an $\operatorname{SU}(n+1)$ gauge in which the gauge potentials have triangular form with real diagonal elements. For this reason, Yang's $R$-gauge and its generalisation to $\mathrm{SU}(n+1)$ have been referred to as the triangular gauge (Brihaye et al 1978).

We follow Ardalan (1978) and obtain the analogue of the $R$-gauge not only when the gauge group $\mathrm{G}=\mathrm{SU}(n+1)$, but also when G is any compact semi-simple group. However, we choose representations such that $R$ is not always triangular. Our intention is to derive the Toda lattice equations governed by an arbitrary group and then find solutions of the equation for all the groups associated with the classical simple Lie algebras.

We must first complexify the group $G$ to $G_{c}$ and then invoke the corresponding complex-valued Lie algebra. The group $G$ may either be classical or exceptional and the former class may be divided into the following subclasses: $\mathrm{SU}(n+1), \mathrm{SO}(2 n+1)$, $\mathrm{Sp}(n)$ and $\mathrm{SO}(2 n) \dagger$. It should be noted that in ail cases the maximal compact subgroup of the complex group $G_{c}$ coincides with $G$ (Gilmore 1974). If $g$ is a semi-simple Lie algebra over $\mathbb{C}$, then it may be decomposed according to Cartan's root space decomposition (Helgason 1978) as

$$
\begin{equation*}
\mathfrak{g}=\mathfrak{h}+\sum_{\alpha \in \Delta^{+}}\left(\mathfrak{g}_{+\alpha}+\mathfrak{g}_{-\alpha}\right), \tag{2.1}
\end{equation*}
$$

where $\mathfrak{h}$ is a fixed Cartan subalgebra, $\Delta^{+}$is the positive root system of $g$ with respect to $\mathfrak{h}$ and $g_{ \pm \alpha}$ are the root spaces defined by

$$
\mathfrak{g}_{ \pm \alpha}=\{E \in \mathfrak{g}: \operatorname{ad}(H) E= \pm \alpha(H) E \text { for all } H \in \mathfrak{h}\}
$$

We shall use the Chevalley basis ( $H_{\alpha}, E_{ \pm \alpha}$ ) for $g$ (Humphreys 1972, Carter 1972). $H_{\alpha}$ is a basis for $\mathfrak{h}$ and $E_{ \pm \alpha}$ a basis for $\mathfrak{g}_{ \pm \alpha}$, such that

$$
\begin{align*}
& {\left[H_{\alpha}, H_{\beta}\right]=0,}  \tag{2.2a}\\
& {\left[H_{\alpha}, E_{ \pm \beta}\right]= \pm K_{\beta \alpha} E_{ \pm \beta},}  \tag{2.2b}\\
& {\left[E_{\alpha}, E_{-\alpha}\right]=H_{\alpha},}  \tag{2.2c}\\
& {\left[E_{\alpha}, E_{\beta}\right]= \begin{cases}0 & \alpha+\beta \notin \Delta, \alpha \neq \beta, \\
N_{\alpha, \beta} E_{\alpha+\beta} & \text { otherwise, } \alpha \neq \beta .\end{cases} } \tag{2.2d}
\end{align*}
$$

Only $n$, where $n=$ rank g , of the $H$ 's are linearly independent and the corresponding set of $n$ positive roots is denoted by $\pi^{+}$. Any root $\alpha \in \pi^{+}$is known as a simple root, that is, it cannot be written as the sum of other positive roots. As a consequence of this property, we can show that for $\alpha, \beta \in \pi^{+}$

$$
N_{\alpha, \beta}=N_{\alpha,-\beta}=0
$$

In addition for $\alpha, \beta$ simple, the matrix ( $K_{\beta \alpha}$ ) defined by the commutation relation ( $2.2 b$ ) is the ( $n \times n$ ) Cartan matrix associated with $g$.

Since g is a semi-simple Lie algebra over $\mathbb{C}$, we may consider it as a Lie algebra $\mathrm{g}^{R}$ over $\mathbb{R}$ with complex structure, and the associated Lie group may be denoted by $\mathrm{G}_{\mathrm{c}}$. By

[^1]combining the Cartan decomposition of $\mathfrak{g}$ and the root space decomposition of $g^{R}$, the Iwasawa decomposition for $G_{c}$ is obtained. It is given by
\[

$$
\begin{equation*}
\mathrm{G}_{\mathrm{c}}=\mathrm{NAK}, \tag{2.3}
\end{equation*}
$$

\]

where $N$ and $A$ are subgroups of $G_{c}$ corresponding to the subalgebras $g_{-\alpha}$ and $\mathfrak{h}$ respectively, and $K$ is the maximal compact subgroup. The decomposition NAK is quite general and may be applied to any non-compact, semi-simple Lie group. However, if $G_{c}$ is of the form

$$
\begin{equation*}
\mathrm{G}_{\mathrm{c}}=\mathrm{i} G \times \mathrm{G}, \tag{2.4}
\end{equation*}
$$

we may be more specific about $\mathrm{N}, \mathrm{A}$ and K . Indeed,

$$
\mathrm{K}=\mathrm{G}
$$

the original gauge group; also $A$ is generated by $\mathfrak{b}$ over $\mathbb{R}$ and $N$ by $\mathfrak{g}_{-\alpha}$ over $\mathbb{C}$. Hence, any group element $g \in G_{c}$ may be written as

$$
\begin{equation*}
g=n a k \tag{2.5}
\end{equation*}
$$

with $n \in \mathrm{~N}, a \in \mathrm{~A}, k \in \mathrm{G} . n$ and $a$ are given by

$$
\begin{equation*}
n=\exp \left(\sum_{\alpha \in \Delta^{+}} z_{\alpha} E_{-\alpha}\right), \quad a=\exp \left(\sum_{\alpha \in \pi^{+}} \psi_{\alpha} H_{\alpha}\right) \tag{2.6}
\end{equation*}
$$

where $z_{\alpha}$ are complex functions and $\psi_{\alpha}$ are real functions, all dependent on both coordinates, $u$ and $\bar{u}$,

$$
\begin{equation*}
u=s+\mathrm{i} t, \quad \bar{u}=s-\mathrm{i} t \tag{2.7}
\end{equation*}
$$

Note that the suffix $\alpha$ of $H_{\alpha}$ corresponds to the simple roots, so only the linearly independent basis elements are included.

Our discussion concerning the Iwasawa decomposition has involved the subgroup of $\mathrm{G}_{\mathrm{c}}$ corresponding to the nilpotent subalgebra $\mathrm{g}_{-\alpha}$, but not the subgroup corresponding to the equivalent subalgebra $\mathfrak{g}_{+\alpha}$, which is spanned by $E_{+\alpha}$. It is possible to include this other subgroup in another version of the Iwasawa decomposition in which any element $\tilde{g} \in \mathrm{G}_{\mathrm{c}}$ is given by

$$
\begin{equation*}
\tilde{g}=\tilde{n} \tilde{a} \tilde{k} . \tag{2.8}
\end{equation*}
$$

In (2.8) $\tilde{k} \in \mathrm{G}$ and the representations for $\tilde{n}$ and $\tilde{a}$ are

$$
\begin{equation*}
\tilde{n}=\exp \left(-\sum_{\alpha \in \Delta^{+}} \tilde{z}_{\alpha} E_{+\alpha}\right), \quad \tilde{a}=\exp \left(-\sum_{\alpha \in \pi^{+}} \tilde{\psi}_{\alpha} H_{\alpha}\right) \tag{2.9}
\end{equation*}
$$

To generalise the $R$-gauge method, we gauge away the elements $k, \tilde{k} \in \mathrm{G}$ in the decompositions (2.5) and (2.8) for $g$ and $\tilde{g}$ respectively. This depends on choosing the same gauge for both elements, $g$ and $\tilde{g}$, and so we restrict them by imposing the condition

$$
\begin{equation*}
k=\tilde{k} \tag{2.10}
\end{equation*}
$$

However, it should be noted that to derive the Toda lattice equations it is not necessary to assume any relationship between the parameters $z_{\alpha}, \psi_{\alpha}$ in the representation of $n a$ and $\tilde{z}_{a}, \tilde{\psi}_{\alpha}$ in $\tilde{n} \tilde{a}$.

The condition (2.10) allows us to choose the gauge in which $g$ and $\tilde{g}$ are given by $R$ and $\tilde{R}$, where

$$
\begin{equation*}
R=n a, \quad \tilde{R}=\tilde{n} \tilde{a} \tag{2.11}
\end{equation*}
$$

By analogy with (1.19) we define the left-invariant gauge potentials $B_{u}$ and $B_{u}$ by

$$
\begin{equation*}
B_{u}(u, \bar{u})=R^{-1} \partial_{u} R, \quad B_{\bar{u}}(u, \bar{u})=\tilde{R}^{-1} \partial_{\bar{u}} \tilde{R} \tag{2.12}
\end{equation*}
$$

After substituting the explicit forms (2.11) for $R$ and $\tilde{R}$ in (2.12), we obtain

$$
\begin{equation*}
B_{u}=a^{-1}\left(n^{-1} \partial_{u} n\right) a+a^{-1} \partial_{u} a, \quad B_{\bar{u}}=\tilde{a}^{-1}\left(\tilde{n}^{-1} \partial_{\bar{u}} \tilde{n}\right) \tilde{a}+\tilde{a}^{-1} \partial_{\bar{u}} \tilde{a} . \tag{2.13}
\end{equation*}
$$

We now adopt the following ansatz:

$$
\begin{equation*}
n^{-1} \partial_{u} n=\sum_{\alpha \in \pi^{+}} y_{\alpha} E_{-\alpha}, \quad \tilde{n}^{-1} \partial_{\bar{u}} \tilde{n}=-\sum_{\alpha \in \pi^{+}} \tilde{y}_{\alpha} E_{+\alpha} \tag{2.14}
\end{equation*}
$$

which is equivalent to the assumption that, when $\alpha$ is not a simple root,

$$
\begin{equation*}
y_{\alpha}=\tilde{y}_{\alpha}=0 \tag{2.15}
\end{equation*}
$$

that is, for $\alpha \notin \pi^{+}$. This restriction to include only the simple roots is in accordance with the Toda model developed from the theory of Lax pairs (Bogoyavlensky 1976) in which the non-simple roots do not play an important role.

By using the representations for $n$ and $\tilde{n}$ in (2.6) and (2.9) and the commutation relations ( $2.2 d$ ), we can show that for $\alpha \in \pi^{+}$

$$
\begin{equation*}
y_{\alpha}=\partial_{u} z_{\alpha}, \quad \tilde{y}_{\alpha}=\partial_{\bar{u}} \tilde{z}_{\alpha} . \tag{2.16}
\end{equation*}
$$

Also from the representations for $a$ and $\tilde{a}$, it is straightforward to obtain

$$
\begin{equation*}
a^{-1} \partial_{u} a=\sum_{\alpha \in \pi^{+}}\left(\partial_{u} \psi_{\alpha}\right) H_{\alpha}, \quad \tilde{a}^{-1} \partial_{\tilde{u}} \tilde{a}=-\sum_{\alpha \in \pi^{+}}\left(\partial_{\tilde{u}} \tilde{\psi}_{\alpha}\right) H_{\alpha} \tag{2.17}
\end{equation*}
$$

due to the commutation of the basis elements $H_{\alpha}$. Further,
$a^{-1} E_{-\alpha} a=\exp \left(\sum_{\beta \in \pi^{+}} K_{\alpha \beta} \psi_{\beta}\right) E_{-\alpha}, \quad \tilde{a}^{-1} E_{+\alpha} \tilde{a}=\exp \left(\sum_{\beta \in \pi^{+}} K_{\alpha \beta} \tilde{\psi_{\beta}}\right) E_{+\alpha}$
and so by combining (2.14)-(2.18) with (2.13), we have

$$
\begin{align*}
& B_{u}=\sum_{\alpha \in \pi^{+}}\left[y_{\alpha} \exp \left(\sum_{\beta \in \pi^{+}} K_{\alpha \beta} \psi_{\beta}\right) E_{-\alpha}+\left(\partial_{u} \psi_{\alpha}\right) H_{\alpha}\right],  \tag{2.19}\\
& B_{\bar{u}}=-\sum_{\alpha \in \pi^{+}}\left[\tilde{y}_{\alpha} \exp \left(\sum_{\beta \in \pi^{+}} K_{\alpha \beta} \tilde{\psi}_{\beta}\right) E_{+\alpha}+\left(\partial_{\bar{u}} \tilde{\psi}_{\alpha}\right) H_{\alpha}\right] .
\end{align*}
$$

In order to impose the condition (1.18) that

$$
F_{u \bar{u}}=\partial_{u} B_{\bar{u}}-\partial_{\bar{u}} B_{u}+\left[B_{u}, B_{\bar{u}}\right]
$$

vanishes, we calculate firstly

$$
\begin{align*}
{\left[B_{u}, B_{\bar{u}}\right]=- } & \sum_{\alpha, \gamma \in \pi^{+}}\left[y_{\alpha}\left(\partial_{\bar{u}} \tilde{\psi}_{\gamma}\right) \exp \left(\sum_{\beta \in \pi^{+}} K_{\alpha \beta} \psi_{\beta}\right) K_{\alpha \gamma} E_{-\alpha}\right. \\
& \left.+\tilde{y}_{\alpha}\left(\partial_{u} \psi_{\gamma}\right) \exp \left(\sum_{\beta \in \pi^{+}} K_{\alpha \beta} \tilde{\psi}_{\beta}\right) K_{\alpha \gamma} E_{+\alpha}\right] \\
& +\sum_{\alpha \in \pi^{+}} y_{\alpha} \tilde{y}_{\alpha} \exp \left(\sum_{\beta \in \pi^{+}} K_{\alpha \beta}\left(\psi_{\beta}+\tilde{\psi}_{\beta}\right)\right) H_{\alpha} \tag{2.20}
\end{align*}
$$

and secondly

$$
\begin{align*}
\partial_{\bar{u}} B_{u}-\partial_{u} B_{\bar{u}}= & \sum_{\alpha \in \pi^{+}}\left[\left(\partial_{\bar{u}} y_{\alpha}+y_{\alpha} \sum_{\beta \in \pi^{+}} K_{\alpha \beta} \partial_{\bar{u}} \psi_{\beta}\right) \exp \left(\sum_{\beta \in \pi^{+}} K_{\alpha \beta} \psi_{\beta}\right) E_{-\alpha}+\left(\partial_{\bar{u}} \partial_{u} \psi_{\alpha}\right) H_{\alpha}\right. \\
& \left.+\left(\partial_{u} \tilde{y}_{\alpha}+\tilde{y}_{\alpha} \sum_{\beta \in \pi^{+}} K_{\alpha \beta} \partial_{u} \tilde{\psi}_{\beta}\right) \exp \left(\sum_{\beta \in \pi^{+}} K_{\alpha \beta} \tilde{\psi}_{\beta}\right) E_{+\alpha}+\left(\partial_{u} \partial_{\bar{u}} \tilde{\psi}_{\alpha}\right) H_{\alpha}\right] . \tag{2.21}
\end{align*}
$$

Now by equating the coefficients of $H_{\alpha}$ and $E_{ \pm \alpha}$ in (2.20) and (2.21) the condition reduces to the following three equations:

$$
\begin{align*}
& y_{\alpha} \tilde{y}_{\kappa} \exp \left(\sum_{\beta \in \pi^{+}} K_{\alpha \beta}\left(\psi_{\beta}+\tilde{\psi}_{\beta}\right)\right)=\partial_{u} \partial_{\bar{u}}\left(\psi_{\alpha}+\tilde{\psi}_{\alpha}\right),  \tag{2.22}\\
& \partial_{\bar{u}} y_{\alpha}+y_{\alpha} \sum_{\beta \in \pi^{+}} K_{\alpha \beta} \partial_{\bar{u}} \psi_{\beta}=-y_{\alpha} \sum_{\beta \in \pi^{+}} K_{\alpha \beta}\left(\partial_{\tilde{u}} \tilde{\psi}_{\beta}\right),  \tag{2.23}\\
& \partial_{u} \tilde{y}_{\alpha}+\tilde{y}_{\alpha} \sum_{\beta ؟ \pi^{+}} K_{\alpha \beta} \partial_{u} \tilde{\psi}_{\beta}=-\tilde{y}_{\alpha} \sum_{\beta \in \pi^{+}} K_{\alpha \beta}\left(\partial_{u} \psi_{\beta}\right)
\end{align*}
$$

We note that for $G=S U(2)$, there is only one simple root and $K=2$, so in this case (2.22) and (2.23) correspond to (1.13) and (1.14) respectively with

$$
\begin{equation*}
y=\partial_{u} \rho, \quad \tilde{y}=g_{u} \tilde{\rho}, \quad \phi=\mathrm{e}^{-(\psi+\tilde{\psi}\}} \tag{2.24}
\end{equation*}
$$

Therefore, by analogy (2.22) is the main equation and (2.23) are the subsidiary equations. Fortunately, as in the $S U(2)$ case, we are able to solve the subsidiary equations to give

$$
\begin{equation*}
y_{0}=\varepsilon(u) \exp \left(-\sum_{\beta \in \pi} K_{\alpha \beta}\left(\psi_{\beta}+\tilde{\psi}_{\beta}\right)\right), \quad \tilde{y}_{\alpha}=\tilde{\varepsilon}(\tilde{u}) \exp \left(-\sum_{\beta \in \pi^{*}} \boldsymbol{K}_{\alpha \beta}\left(\psi_{\beta}+\tilde{\psi}_{\beta}\right)\right), \tag{2.25}
\end{equation*}
$$

where $\varepsilon(u)$ and $\tilde{\varepsilon}(\bar{u})$ are to be determined by the 'boundary conditions'. If we put the solutions (2.25) into (2.22), then the main equation becomes

$$
\begin{equation*}
\partial_{\mu} \partial_{\tilde{u}}\left(\psi_{\alpha}+\dot{\psi}_{\alpha}\right)=\varepsilon \tilde{\varepsilon} \exp \left(-\sum_{\beta \in \pi^{+}} K_{\alpha \beta}\left(\psi_{\beta}+\tilde{\psi}_{\beta}\right)\right) . \tag{2.26}
\end{equation*}
$$

When $\varepsilon \ddot{\varepsilon}=1$, the equation (2.26) is one particular form of the Toda lattice equation governed by an arbitrary group $G$ with Cartan matrix $K$. The form (1.2) of the Toda lattice equation may be derived by defining the new variables

$$
\begin{equation*}
\sigma_{\gamma}=-\sum_{\beta \in \pi^{-}} K_{\alpha \beta} \psi_{\beta}, \quad \dot{\sigma}_{\alpha}=-\sum_{\beta \in \pi^{+}} K_{\alpha \beta} \tilde{\psi}_{\beta} \tag{2.27a}
\end{equation*}
$$

and

$$
\begin{equation*}
\rho_{\alpha}=\sigma_{\alpha}+\tilde{\sigma}_{\alpha} \tag{2.27b}
\end{equation*}
$$

However, it should be said that the algebraic meaning of the variables $\rho_{\alpha}$ is somewhat vague, but in our form (2.26) we can trace the origin of the $\psi_{\alpha}$ and $\tilde{\psi}_{\alpha}$ back to the coefficients of the generators of the abelian subgroup in the Iwasawa decomposition.

Also we can give explicit expressions for the potentials $B_{u}$ and $B_{\bar{u}}$ by substituting the solutions (2.25) for $y_{\alpha}$ and $\tilde{y}_{\alpha}$ into (2.19). In particular, when $\varepsilon=\tilde{\varepsilon}=1$, this gives

$$
\begin{align*}
& B_{u}=\sum_{\alpha \in \pi^{+}}\left[\exp \left(-\sum_{\beta \in \pi^{+}} K_{\alpha \beta} \tilde{\psi}_{\beta}\right) E_{-\alpha}+\left(\partial_{u} \psi_{\alpha}\right) H_{\alpha}\right] \\
& B_{\bar{u}}=-\sum_{\alpha \in \pi^{+}}\left[\exp \left(-\sum_{\beta \in \pi^{+}} K_{\alpha \beta} \psi_{\beta}\right) E_{+\alpha}+\left(\partial_{\bar{u}} \tilde{\psi}_{\alpha}\right) H_{\alpha}\right] . \tag{2.28}
\end{align*}
$$

We close this section by remarking that indeed the Iwasawa decomposition has allowed us to generalise Yang's $R$-gauge method and hence obtain the Toda lattice equations governed by an arbitrary semi-simple Lie group. We reiterate that to derive these equations it is not necessary to impose any condition relating the parameters of $R$ and $\tilde{R}$. This differs from the $R$-gauge where from (2.24) we see that in the real section $\tilde{y}=\bar{y}$, where the bar denotes complex conjugation, and, for consistency with (1.15), $\tilde{\psi}=\psi$. We discuss this point in greater detail in a later section.

## 3. Parametrisation of the solutions

The purpose of this section is to demonstrate a method of obtaining the solutions to the Toda lattice equations (2.26) when the governing group $G$ is classical. The solutions are parametrised in the form of determinants. As in § 2, it is the Iwasawa decomposition of $R$ and $\tilde{R}$ given by ( 2.11 ) which plays an important role.

Since we have factored out the gauge terms in $g$ and $\hat{g}$, the resultant elements $R$ and $\tilde{R}$ must belong to the symmetric space $\mathrm{G}_{\mathrm{c}} / \mathrm{G}$. In addition $B_{u}$ and $B_{\bar{u}}$ given by (2.19) or (2.28) compose a matrix connection one-form, taking values in the algebra corresponding to $G_{c}$, which is defined by

$$
\begin{equation*}
\omega=B_{u} \mathrm{~d} u+B_{\bar{u}} \mathrm{~d} \bar{u} \tag{3.1}
\end{equation*}
$$

Then the condition (1.18) implies that the curvature two-form must vanish, that is,

$$
\begin{equation*}
\Omega \equiv \mathrm{d} \omega+\frac{1}{2}[\omega, \omega]=0 \tag{3.2}
\end{equation*}
$$

Conversely, if $\Omega=0$, then the connection $\omega$ should be a left-invariant one-form given by

$$
\begin{equation*}
\omega=g_{0}^{-1} \mathrm{~d} g_{0} \tag{3.3}
\end{equation*}
$$

where $g_{0} \in \mathrm{G}_{\mathrm{c}}$. Consequently $B_{u}$ and $B_{\bar{u}}$ are given by

$$
\begin{equation*}
B_{u}=g_{0}^{-1} \partial_{u} g_{0}, \quad B_{\bar{u}}=g_{0}^{-1} \partial_{\tilde{u}} g_{0} \tag{3.4}
\end{equation*}
$$

However, we have already capitalised on the definitions of $B_{u}$ and $B_{\bar{u}}$ given by (2.12). To show that the two sets of definitions (3.4) and (2.12) are non-contradictory, we make explicit the property alluded to in $\S 2$ that the gauge potentials are left-invariant. In fact, $B_{u}$ and $B_{\bar{u}}$ given by (3.4) are unchanged when $g_{0}$ is transformed by the respective left-translations

$$
\begin{equation*}
g_{0} \rightarrow \tilde{r}^{-1}(\ddot{u}) g_{0}, \quad g_{0} \rightarrow r^{-1}(u) g_{0} \tag{3.5}
\end{equation*}
$$

Hence, if we are permitted to make the identifications

$$
\begin{equation*}
R=\tilde{r}^{-1}(\bar{u}) g_{0}, \quad \tilde{R}=r^{-1}(u) g_{0} \tag{3.6}
\end{equation*}
$$

at least up to a gauge transformation, then the two definitions are consistent. In general
$r$ and $\tilde{r}$ may belong to $\mathrm{G}_{\mathrm{c}}$, but are restricted so that $r$ and $\tilde{r}$ depend only on $u$ and $\bar{u}$ respectively. We use (2.3) to write $\tilde{r}$ and $r$ in decomposed form as

$$
\tilde{r}=\tilde{n}_{0} \tilde{a}_{0} \tilde{k}_{0}, \quad r=n_{0} a_{0} k_{0} .
$$

By use of these equations, (3.6) can be rearranged to give

$$
\begin{equation*}
g_{0}=\tilde{n}_{0} \tilde{a}_{0} \tilde{k}_{0} R, \quad g_{0}=n_{0} a_{0} k_{0} \tilde{R} . \tag{3.7}
\end{equation*}
$$

It is known that for all $k^{\prime \prime} \in \mathrm{G}, \operatorname{Ad} k^{\prime \prime}$ leaves $\mathrm{G}_{\mathrm{c}}$ invariant. In particular (Hermann 1966, Kostant 1973)

$$
k^{\prime \prime} n a k^{\prime \prime-1}=n^{\prime} a^{\prime} k^{\prime}
$$

where $n^{\prime} \in \mathbf{N}, a^{\prime} \in \mathrm{A}, k^{\prime} \in \mathrm{K}$, and this may be rewritten as

$$
\begin{equation*}
k^{\prime \prime} n a=n^{\prime} a^{\prime} k^{\prime} k^{\prime \prime} \tag{3.8}
\end{equation*}
$$

We use (3.8) to rearrange equations (3.7) to give

$$
\begin{align*}
& g_{0}=\tilde{n}_{0} \tilde{a}_{0} R^{\prime} k^{\prime} \tilde{k_{0}},  \tag{3.9a}\\
& g_{0}=n_{0} a_{0} \tilde{R}^{\prime} \tilde{k^{\prime}} k_{0} \tag{3.9b}
\end{align*}
$$

where

$$
R^{\prime}=n^{\prime} a^{\prime}, \quad \tilde{R}^{\prime}=\tilde{n}^{\prime} \tilde{a}^{\prime}
$$

We now require that $R$ and $\tilde{R}$ of $\S 2$ be replaced by $R^{\prime}$ and $\tilde{R}^{\prime}$. If we define $B_{u}^{\prime}$ and $B_{\bar{u}}^{\prime}$ from $R^{\prime}$ and $\tilde{R}^{\prime}$ respectively by analogy with (2.12), then these are just gauge transformations of the potentials $B_{u}$ and $B_{\bar{u}}$ given by (3.4). For consistency, $B_{u}^{\prime}$ and $B_{\bar{u}}^{\prime}$ must be defined in the same gauge, and hence in (3.9) we put

$$
\begin{equation*}
k^{\prime} \tilde{k}_{0}=\tilde{k}^{\prime} k_{0} \tag{3.10}
\end{equation*}
$$

By equating the right-hand sides of (3.9) and imposing the condition (3.10), we obtain

$$
\begin{equation*}
\tilde{a}_{0} R \tilde{R}^{-1} a_{0}^{-1}=\tilde{n}_{0}^{-1} n_{0}, \tag{3.11}
\end{equation*}
$$

where we have dropped the primed notation. Hence (3.11) may be obtained effectively from (3.6) when $\tilde{r}$ and $r$ are given by $\dagger$

$$
\begin{equation*}
\tilde{r}=\tilde{n}_{0}(\bar{u}) \tilde{a}_{0}(\bar{u}), \quad r=n_{0}(u) a_{0}(u) . \tag{3.12}
\end{equation*}
$$

We use the following representations for the decompositions (3.12):

$$
\begin{array}{ll}
a_{0}=\exp \left(\sum_{\alpha \in \pi^{+}} \mu_{\alpha}(u) H_{\alpha}\right), & \tilde{a}_{0}=\exp \left(-\sum_{\alpha \in \pi^{+}} \tilde{\mu}_{\alpha}(\bar{u}) H_{\alpha}\right) \\
n_{0}=\exp \left(\sum_{\alpha \in \Delta^{+}} \nu_{\alpha}(u) E_{-\alpha}\right), & \tilde{n}_{0}=\exp \left(-\sum_{\alpha \in \Delta^{+}} \tilde{\nu}_{\alpha}(\bar{u}) E_{+\alpha}\right) \tag{3.13}
\end{array}
$$

Our method of solution depends on taking determinants of submatrices of the left- and right-hand side of (3.11) and then equating. Firstly, we consider the left-hand side,

[^2]which is explicitly given by
$L=\exp \left(-\tilde{\mu}_{\alpha} H_{\alpha}\right) \exp \left(z_{\beta} E_{-\beta}\right) \exp \left[\left(\psi_{\alpha}+\tilde{\psi}_{\alpha}\right) H_{\alpha}\right] \exp \left(\tilde{z}_{\beta} E_{+\beta}\right) \exp \left(-\mu_{\alpha} H_{\alpha}\right)$,
where in the exponents the index $\beta$ is summed over all positive roots, while the index $\alpha$ is summed only over the simple roots. Since for all the classical algebras the $H_{\alpha}$ are diagonal, the properties of $L$ are essentially those of the product of $\exp \left(z_{\beta} E_{-\beta}\right)$ and $\exp \left(\tilde{z}_{\beta} E_{+\beta}\right)$. Now if $L^{\prime}$ is an $(n \times n)$ lower triangular matrix and $L^{\prime \prime}$ is an ( $n \times n$ ) upper triangular matrix with
$$
L=L^{\prime} L^{\prime \prime},
$$
then the only contributions to the $(m \times m)$ submatrix $\left(L_{k l}\right)_{1 \leqslant k, l \leqslant m}$ are from the ( $m \times m$ ) submatrices $\left(L_{k l}^{\prime}\right)_{1 \leqslant k, l \leqslant m}$ and $\left(L_{k l}^{\prime \prime}\right)_{1 \leqslant k, l \leqslant m}$. Hence, for $1 \leqslant k, l \leqslant m$,
\[

$$
\begin{equation*}
\operatorname{det}\left(L_{k l}\right)=\operatorname{det}\left(L_{k l}^{\prime}\right) \operatorname{det}\left(L_{k l}^{\prime \prime}\right) \tag{3.15}
\end{equation*}
$$

\]

In appendix 1 for each classical Lie algebra we specify some set of submatrices $\left\{\left(\left(E_{-\alpha}\right)_{l m}\right): l \leqslant m\right\}$ in which each member has a lower triangular matrix lying in a block where all the other entries are zero. We restrict our attention to these submatrices, and then from (3.14), we can show that by using (3.15)

$$
\begin{aligned}
\operatorname{det}\left(L_{k l}\right)=\operatorname{det}[ & \left.\left(\exp \left(-\bar{\mu}_{\alpha} H_{\alpha}\right)\right)_{k l}\right] \operatorname{det}\left[\left(\exp \left(z_{\alpha} E_{-\alpha}\right)\right)_{k l}\right] \operatorname{det}\left[\left(\exp \left[\left(\psi_{\alpha}+\bar{\psi}_{\alpha}\right) H_{\alpha}\right]\right)_{k l}\right] \\
& \times \operatorname{det}\left[\left(\exp \left(\tilde{z}_{\alpha} E_{+\alpha}\right)\right)_{k l}\right] \operatorname{det}\left[\left(\exp \left(-\mu_{\alpha} H_{\alpha}\right)\right)_{k l}\right]
\end{aligned}
$$

where $1 \leqslant k, l \leqslant m \leqslant n$ for $\mathfrak{a}_{n}, \mathbf{c}_{n}, \boldsymbol{d}_{n}, 2 \leqslant k, l \leqslant m \leqslant n+1$ for $\mathfrak{b}_{n}$.
Since for each algebra in the appropriate range

$$
\operatorname{det}\left[\left(\exp \left(z_{\alpha} E_{-\alpha}\right)\right)_{k l}\right]=1
$$

and since the $H_{\alpha}$ are diagonal, this becomes
$\operatorname{det}\left(L_{k l}\right)=\exp \left\{\operatorname{Tr}\left[\left(-\tilde{\mu}_{\alpha} H_{\alpha}\right)_{k l}\right]\right\} \exp \left\{\operatorname{Tr}\left[\left(\left(\psi_{\alpha}-\tilde{\psi}_{\alpha}\right) H_{\alpha}\right)_{k l}\right]\right\} \exp \left\{\operatorname{Tr}\left[\left(-\mu_{\alpha} H_{\alpha}\right)_{k l}\right]\right\}$.
Therefore, (3.13) may be written as
$\exp \left\{\operatorname{Tr}\left[\left(\left(\psi_{\alpha}+\tilde{\psi}_{\alpha}\right) H_{\alpha}\right)_{k l}\right]\right\}=\exp \left\{\operatorname{Tr}\left[\left(\left(\mu_{\alpha}+\tilde{\mu}_{\alpha}\right) H_{\alpha}\right)_{k l}\right]\right\} \operatorname{det}\left[\left(\tilde{n}_{0}^{-1} n_{0}\right)_{k l}\right]$.
Before reducing this expression further, we shall consider the representations (3.12) for $a_{0}, \tilde{a}_{0}, n_{0}$ and $\tilde{n}_{0}$. (3.11) provides a relationship between $R$ and $\tilde{R}$,

$$
R=\tilde{r}^{-1} r \tilde{R}
$$

and so $B_{u}$, for instance, is given by

$$
\begin{equation*}
B_{u}=\tilde{R}^{-1} r^{-1}\left(\partial_{u} r\right) \tilde{R}+\tilde{R}^{-1} \partial_{u} \tilde{R}, \tag{3.17}
\end{equation*}
$$

as well as by (2.12). In (2.13) $B_{u}$ has a specific form depending only on $H_{\alpha}$ and $E_{-\alpha}$ and, moreover, as a consequence of the ansatz (2.14), the basis elements correspond only to the simple roots. Hence there must be some condition on the parameters in $r$ so that $B_{u}$ given by (2.17) also has these properties. In particular, we are concerned with the coefficient of $E_{-\alpha}$ in (3.17). By using (3.12),

$$
r^{-1} \partial_{u} r=a_{0}^{-1}\left(n_{0}^{-1} \partial_{u} n_{0}\right) a_{0}+a_{0}^{-1} \partial_{u} a_{0}
$$

and so, by referring to the commutation relations (2.2), we see that the only term containing $E_{-\alpha}$ is

$$
\begin{equation*}
\tilde{R}^{-1} a_{0}^{-1}\left(n_{0}^{-1} \partial_{u} n_{0}\right) a_{0} \tilde{R} \tag{3.18}
\end{equation*}
$$

Firstly, by analogy with the ansatz (2.14), we assume that

$$
\begin{equation*}
n_{0}^{-1} \partial_{u} n_{0}=\sum_{\alpha \in \pi^{+}} \phi_{\alpha} E_{-\alpha} \tag{3.19a}
\end{equation*}
$$

that is, we impose the condition that, for $\alpha \notin \pi^{+}$,

$$
\begin{equation*}
\phi_{\alpha}=0, \tag{3.20a}
\end{equation*}
$$

and then (3.18) becomes

$$
\begin{gathered}
\tilde{R}^{-1} a_{0}^{-1} \sum_{\alpha \leqslant \pi^{+}} \phi_{\alpha} E_{-\alpha} a_{0} \tilde{R}=\tilde{R}^{-1} \sum_{\alpha \in \pi^{+}} \exp \left(\sum_{\beta \in \pi^{+}} K_{\alpha \beta} \mu_{\beta}\right) \phi_{\alpha} E_{-\alpha} \tilde{R} \\
=\exp \left(\tilde{\psi}_{\alpha} H_{\alpha}\right) \sum_{\alpha \in \pi^{+}} \exp \left(\sum_{\beta \in \pi^{+}} K_{\alpha \beta} \mu_{\beta}\right) \phi_{\alpha} E_{-\alpha} \exp \left(-\tilde{\psi}_{\alpha} H_{\alpha}\right) \\
=\sum_{\alpha \in \pi^{+}} \exp \left(\sum_{\beta \in \pi^{+}} K_{\alpha \beta}\left(\mu_{\beta}-\tilde{\psi}_{\beta}\right)\right) \phi_{\alpha} E_{-\alpha}
\end{gathered}
$$

This last term should be compared with the corresponding term in $B_{u}=R^{-1} \partial_{u} R$; that is, from (2.28)

$$
\sum_{\alpha \in \pi^{+}} \exp \left(-\sum_{\beta \in \pi^{+}} K_{\alpha \beta} \tilde{\psi}_{\beta}\right) E_{--\alpha}
$$

Hence

$$
\begin{equation*}
\phi_{\alpha}=\exp \left(-\sum_{\beta \in \pi^{,}} K_{\alpha \beta} \mu_{\beta}\right) . \tag{3.21a}
\end{equation*}
$$

Similarly, by considering $B_{\tilde{u}}$ in the same way, we may show that

$$
\begin{equation*}
\tilde{\phi}_{\alpha}=\exp \left(-\sum_{\beta \in \pi^{+}} K_{\alpha \beta} \tilde{\mu}_{\beta}\right), \tag{3.21b}
\end{equation*}
$$

where when $\alpha$ is a simple root $\tilde{\phi}_{\alpha}$ is defined by

$$
\begin{equation*}
\tilde{n}_{0}^{-1} \partial_{\bar{u}} \tilde{n}_{0}=-\sum_{\alpha \in \pi^{+}} \tilde{\phi}_{\alpha} E_{+\alpha} \tag{3.19b}
\end{equation*}
$$

and for $\alpha \notin \pi^{+}$

$$
\begin{equation*}
\dot{\phi}_{\alpha}=0 . \tag{3.20b}
\end{equation*}
$$

However, if we use the representations for $n_{0}$ and $\tilde{n}_{0}$ in (3.19), it is possible to show that, for $\alpha \in \pi^{+}$,

$$
\begin{equation*}
\phi_{\alpha}=\partial_{u} \nu_{\alpha}, \quad \tilde{\phi}_{\alpha}=\partial_{\tilde{u}} \tilde{\nu}_{\alpha} \tag{3.22}
\end{equation*}
$$

Hence (3.21) and (3.22) provide a relationship between the parameters $\nu_{\alpha}$ and $\mu_{\alpha}$ appearing on the right-hand side of (3.16) and, bearing this in mind, we shall return to the further simplification of that equation. By considering the explicit form of the
matrices $H_{\alpha}$ for each algebra in appendix 1, we can reduce (3.16) to

$$
\left.\begin{array}{lr}
\exp \left(\psi_{m}+\tilde{\psi}_{m}\right)=\exp \left(\mu_{m}+\tilde{\mu}_{m}\right) \operatorname{det}\left(\tilde{n}_{0}^{-1} n_{0}\right)_{m}, & 1 \leqslant m \leqslant n \text { for } \mathfrak{a}_{n}, \mathfrak{c}_{n}, \\
\exp \left(\psi_{m-1}+\tilde{\psi}_{m-1}\right)=\exp \left(\mu_{m-1}+\tilde{\mu}_{m-1}\right) \operatorname{det}\left(\tilde{n}_{0}^{-1} n_{0}\right)_{m}, & 1 \leqslant m \leqslant n-2 \text { for } \mathfrak{D}_{n},  \tag{3.23}\\
\exp \left(\psi_{n-1}+\tilde{\psi}_{n-1}\right)=\exp \left(\mu_{n-1}+\tilde{\mu}_{n-1}\right)\left[\operatorname{det}\left(\tilde{n}_{0}^{-1} n_{0}\right)_{n-1} / \operatorname{det}\left(\tilde{n}_{0}^{-1} n_{0}\right)_{n}\right]^{1 / 2}, & \text { for } \mathfrak{D}_{n} \\
\exp \left(\psi_{n}+\tilde{\psi}_{n}\right)=\exp \left(\mu_{n}+\tilde{\mu}_{n}\right)\left[\operatorname{det}\left(\tilde{n}_{0}^{-1} n_{0}\right)_{n-1} \operatorname{det}\left(\tilde{n}_{0}^{-1} n_{0}\right)_{n}\right]^{1 / 2}, &
\end{array}\right\}
$$

where the abbreviation $\operatorname{det}\left(\tilde{n}_{0}^{-1} n_{0}\right)_{m}$ refers to the determinant of the submatrix [ $\left(\tilde{n}_{0}^{-1} n_{0}\right)_{k l}$ ] with $k$ and $l$ in the appropriate range.

The parameters in the exponent on the left-hand side of the equations (3.23) are just those which satisfy the Toda lattice equations (2.26). In fact, by taking logarithms of each side of (3.23) and then substituting for $\tilde{n}_{0}$ and $n_{0}$ the representations (3.12) subjected to the conditions (3.19)-(3.22), we obtain a parametrisation for the solutions of the Toda lattice equation. These comments shall be exemplified below by referring to some specific groups.

In all the examples below, we shall use the representations for $E_{-\alpha}$ and $E_{+\alpha}$ given in appendix 1. For those involving an $\mathrm{SU}(r+1)$ group the notation will be that of appendix 2 .
(i) $S U(2)$. There is only one root and so, with

$$
n_{0}=\exp \left(\nu_{21} e_{21}\right), \quad \tilde{n}_{0}=\exp \left(-\tilde{\nu}_{12} e_{12}\right)
$$

( $\tilde{n}_{0}^{-1} n_{0}$ ) is simply given by

$$
\left(\begin{array}{cc}
1+\tilde{\nu}_{12} \nu_{21} & \tilde{\nu}_{12} \\
\nu_{21} & 1
\end{array}\right)
$$

Also

$$
\phi=\partial_{u} \nu_{21}=\mathrm{e}^{-2 \mu}, \quad \tilde{\phi}=\partial_{\bar{u}} \tilde{\nu}_{12}=\mathrm{e}^{-2 \tilde{\mu}}
$$

so (3.23) gives

$$
\begin{equation*}
\mathrm{e}^{(\psi+\tilde{\psi})}=\left[\left(\partial_{u} \nu_{21}\right)\left(\partial_{\tilde{u}} \tilde{\nu}_{12}\right)\right]^{-1 / 2}\left(1+\tilde{\nu}_{12} \nu_{21}\right) . \tag{3.24}
\end{equation*}
$$

If we use the variable

$$
\rho=-(\psi+\tilde{\psi})
$$

then (3.24) is

$$
\rho=-\ln \left(1+\tilde{\nu}_{12} \nu_{21}\right)+\frac{1}{2} \ln \left(\partial_{u} \nu_{21} \partial_{\bar{u}} \tilde{\nu}_{12}\right)
$$

reproducing Witten's $S U(2)$ solution (Witten 1977).
(ii) $\operatorname{SU}(3)$.

$$
n_{0}=\exp \left(\nu_{i j} e_{i j}\right)=1+\nu_{i j} e_{i j}+\frac{1}{2} \delta_{m i} e_{k j} \nu_{i j} \nu_{k m}
$$

and so

$$
\phi_{1}=\partial_{u} \nu_{21}=\exp \left(-2 \mu_{1}+\mu_{2}\right), \quad \phi_{2}=\partial_{u} \nu_{32}=\exp \left(\mu_{1}-2 \mu_{2}\right) .
$$

There is also the additional constraint

$$
\begin{equation*}
0=\phi_{3}=\partial_{u} \nu_{31}+\frac{1}{2}\left(\partial_{u} \nu_{21}\right) \nu_{32}-\frac{1}{2} \nu_{21}\left(\partial_{u} \nu_{32}\right) . \tag{3.25a}
\end{equation*}
$$

Similarly,

$$
\begin{align*}
& \tilde{\phi}_{1}=\partial_{\tilde{u}} \tilde{\nu}_{12}=\exp \left(-2 \tilde{\mu}_{1}+\tilde{\mu}_{2}\right), \quad \tilde{\phi}_{2}=\partial_{\bar{u}} \tilde{\nu}_{23}=\exp \left(\tilde{\mu}_{1}-2 \tilde{\mu}_{2}\right), \\
& 0=\tilde{\phi}_{3}=\partial_{\bar{u}} \tilde{\nu}_{13}-\frac{1}{2}\left(\partial_{\tilde{u}} \tilde{\nu}_{12}\right) \tilde{\nu}_{23}+\frac{1}{2} \tilde{\nu}_{12}\left(\partial_{\bar{u}} \tilde{\nu}_{23}\right) . \tag{3.25b}
\end{align*}
$$

If we denote $\tilde{n}_{0}^{-1} n_{0}$ by $\mathcal{N}$ then

$$
\begin{aligned}
& \mathcal{N}_{11}=1+\tilde{\nu}_{12} \nu_{21}+\left(\tilde{\nu}_{13}+\frac{1}{2} \tilde{\nu}_{12} \tilde{\nu}_{23}\right)\left(\nu_{31}+\frac{1}{2} \nu_{32} \nu_{21}\right), \\
& \mathcal{N}_{22}=1+\tilde{\nu}_{23} \nu_{32} \\
& \mathcal{N}_{12}=\tilde{\nu}_{12}+\nu_{32}\left(\tilde{\nu}_{13}+\frac{1}{2} \tilde{\nu}_{12} \tilde{\nu}_{23}\right) \\
& \mathcal{N}_{21}=\nu_{21}+\tilde{\nu}_{23}\left(\nu_{31}+\frac{1}{2} \nu_{32} \nu_{21}\right)
\end{aligned}
$$

So using (3.23)

$$
\begin{align*}
\exp \left(\psi_{1}+\tilde{\psi}_{1}\right)= & {\left[\left(\partial_{u} \nu_{21}\right)\left(\partial_{u} \tilde{\nu}_{12}\right)\right]^{-2 / 3}\left[\left(\partial_{u} \nu_{32}\right)\left(\partial_{\tilde{u}} \tilde{\nu}_{23}\right)\right]^{-1 / 3} } \\
& \times\left[1+\tilde{\nu}_{12} \nu_{21}+\left(\tilde{\nu}_{13}+\frac{1}{2} \tilde{\nu}_{12} \tilde{\nu}_{23}\right)\left(\nu_{31}+\frac{1}{2} \nu_{32} \nu_{21}\right)\right],  \tag{3.26a}\\
\exp \left(\psi_{2}+\tilde{\psi}_{2}\right)= & {\left[\left(\partial_{u} \nu_{21}\right)\left(\partial_{\bar{u}} \tilde{\nu}_{12}\right)\right]^{-1 / 3}\left[\left(\partial_{u} \nu_{32}\right)\left(\partial_{\tilde{u}} \tilde{\nu}_{23}\right)\right]^{-2 / 3} } \\
& \times\left[1+\tilde{\nu}_{23} \nu_{32}+\left(\tilde{\nu}_{13}-\frac{1}{2} \tilde{\nu}_{12} \tilde{\nu}_{23}\right)\left(\nu_{31}-\frac{1}{2} \nu_{32} \nu_{21}\right)\right] . \tag{3.26b}
\end{align*}
$$

(iii) $S p(2)$. The first example with $R$ and $\tilde{R}$ non-triangular.

From appendix 1,
$E_{-1}=e_{21}-e_{34}, \quad E_{-2}=e_{42}, \quad E_{-3}=e_{31}, \quad E_{-4}=e_{32}+e_{41}$.
Again denoting

$$
\tilde{n}_{0}^{-1} n_{0}=\exp \left(\tilde{\nu}_{\alpha} E_{\alpha}\right) \exp \left(\nu_{\alpha} E_{-\alpha}\right)
$$

by $\mathcal{N}$, we have
$\mathcal{N}_{11}=1+\nu_{1} \tilde{\nu}_{1}+V_{3} \tilde{V}_{3}+V_{4}^{+} \tilde{V}_{4}^{+}, \quad \mathcal{N}_{12}=\tilde{\nu}_{1}+\tilde{V}_{3} V_{4}^{-}+\nu_{2} \tilde{V}_{4}^{+}$,
$\mathscr{N}_{22}=1+V_{4}^{-} \tilde{V}_{4}^{-}+\nu_{2} \tilde{\nu}_{2}, \quad \mathcal{N}_{21}=\nu_{1}+V_{3} \tilde{V}_{4}^{-}+\tilde{\nu}_{2} V_{4}^{+} ;$
the other elements of the matrix need not be specified and we have used the abbreviations

$$
\begin{array}{ll}
V_{4}^{ \pm}=\nu_{4} \pm \frac{1}{2} \nu_{1} \nu_{2}, & \tilde{V}_{4}^{ \pm}=\tilde{\nu}_{4} \pm \frac{1}{2} \tilde{\nu}_{1} \tilde{\nu}_{2}, \\
V_{3}=\nu_{3}-\frac{1}{6} \nu_{1}^{2} \nu_{2}, & \tilde{V}_{3}=\tilde{\nu}_{3}-\frac{1}{6} \tilde{\nu}_{1}^{2} \tilde{\nu}_{2} .
\end{array}
$$

Also the parameters are related by the equations
$\phi_{1}=\partial_{u} \nu_{1}=\exp \left(-2 \mu_{1}+\mu_{2}\right), \quad \phi_{2}=\partial_{u} \nu_{2}=\exp \left(-2 \mu_{2}+2 \mu_{1}\right)$,
together with similar equations for $\tilde{\phi}_{1}$ and $\tilde{\phi}_{2}$. The constraints $\phi_{3}=\phi_{4}=0, \tilde{\phi}_{3}=\tilde{\phi}_{4}=0$ reduce to respectively

$$
\begin{array}{ll}
\partial_{u} V_{3}-V_{4}^{-}\left(\partial_{u} \nu_{1}\right)=0, & \partial_{u} V_{4}^{+}-\nu_{2} \partial_{u} \nu_{1}=0 \\
\partial_{u} \tilde{V}_{3}-\tilde{V}_{4}^{-}\left(\partial_{u} \tilde{\nu}_{1}\right)=0, & \partial_{u} \tilde{V}_{4}^{+}-\tilde{\nu}_{2} \partial_{u} \tilde{\nu}_{1}=0 \tag{3.28b}
\end{array}
$$

By straightforward substitution of $\operatorname{det} \mathcal{N}$ and the relations (3.27) into (3.23), the solutions are given by
$\exp \left(\psi_{1}+\tilde{\psi}_{1}\right)=\left[\left(\partial_{\mu} \nu_{1}\right)\left(\partial_{\bar{u}} \tilde{\nu}_{1}\right)\right]^{-1}\left[\left(\partial_{u} \nu_{2}\right)\left(\partial_{\bar{u}} \tilde{\nu}_{2}\right)\right]^{-1 / 2}\left[1+\nu_{1} \tilde{\nu}_{1}+V_{3} \tilde{V}_{3}+V_{4}^{+} \tilde{V}_{4}^{+}\right]$,

$$
\begin{align*}
\exp \left(\psi_{2}+\tilde{\psi}_{2}\right)= & {\left[\left(\partial_{u} \nu_{1}\right)\left(\partial_{\tilde{u}} \tilde{\nu}_{1}\right)\right]^{-1}\left[\left(\partial_{u} \nu_{2}\right)\left(\partial_{\tilde{u}} \tilde{\nu}_{2}\right)\right]^{-1} } \\
& \times\left[\left(1+\nu_{1} \tilde{\nu}_{1}+V_{3} \tilde{V}_{3}+V_{4}^{+} \tilde{V}_{4}^{+}\right)\left(1+\nu_{2} \tilde{\nu}_{2}+V_{4}^{-} \tilde{V}_{4}^{-}\right)\right. \\
& \left.-\left(\nu_{1}+\tilde{\nu}_{2} V_{4}^{+}+V_{3} \tilde{V}_{4}^{-}\right)\left(\tilde{\nu}_{1}+\nu_{2} \tilde{V}_{4}^{+}+\tilde{V}_{3} V_{4}^{-}\right)\right] . \tag{3.29b}
\end{align*}
$$

Of course, the solutions are more accurately expressed by introducing the explicit forms for $V_{3}, \tilde{V}_{3}, V_{4}^{ \pm}$and $\tilde{V}_{4}^{ \pm}$. However, the form (3.29) is often more convenient to manipulate; for example, to prove that (3.29) really do satisfy the Toda lattice equation. This is also where the constraints ( 3.28 ) come into play. To demonstrate this, we calculate $\partial_{u} \partial_{\tilde{u}}\left(\psi_{1}+\tilde{\psi}_{1}\right)$ by taking logarithms of (3.29a) and then differentiating:

$$
\begin{align*}
\partial_{u} \partial_{\bar{u}}\left(\psi_{1}+\tilde{\psi}_{1}\right) & =\Psi\left(1+\nu_{1} \tilde{\nu}_{1}+V_{3} \tilde{V}_{3}+V_{4}^{+} \tilde{V}_{4}^{+}\right)^{-2} \\
= & \Psi\left[\left(\partial_{u} \nu_{1}\right)\left(\partial_{\tilde{u}} \tilde{\nu}_{1}\right)\right]^{-2}\left[\left(\partial_{u} \nu_{2}\right)\left(\partial_{u} \tilde{\nu}_{2}\right)\right]^{-1} \exp \left[-2\left(\psi_{1}+\tilde{\psi}_{1}\right)\right] . \tag{3.30}
\end{align*}
$$

$\Psi$ is used to denote

$$
\begin{aligned}
&\left(1+\nu_{1} \tilde{\nu}_{1}+V_{3} \tilde{V}_{3}+V_{4}^{+} \tilde{V}_{4}^{+}\right)\left(\partial_{u} \nu_{1} \partial_{\bar{u}} \tilde{\nu}_{1}+\partial_{u} V_{3} \partial_{\tilde{u}} \tilde{V}_{3}+\partial_{u} V_{4}^{+} \partial_{\tilde{u}} \tilde{V}_{4}^{+}\right) \\
&-\left(\nu_{1} \partial_{\tilde{u}} \tilde{\nu}_{1}+V_{3} \partial_{\tilde{u}} \tilde{V}_{3}+V_{4}^{+} \partial_{\bar{u}} \tilde{V}_{4}^{+}\right)\left(\tilde{\nu}_{1} \partial_{u} \nu_{1}+\tilde{V}_{3} \partial_{u} V_{3}+\tilde{V}_{4}^{+} \partial_{u} V_{4}^{+}\right) \\
&=\left(\partial_{u} \nu_{1}\right)\left(\partial_{\tilde{u}} \tilde{\nu}_{1}\right)\left[\left(1+\nu_{1} \tilde{\nu}_{1}+V_{3} \tilde{V}_{3}+V_{4}^{+} \tilde{V}_{4}^{+}\right)\left(1+V_{4}^{-} \tilde{V}_{4}^{-}+\nu_{2} \tilde{\nu}_{2}\right)\right. \\
&\left.-\left(\nu_{1}+\tilde{V}_{4}^{-} V_{3}+V_{4}^{+} \tilde{\nu}_{2}\right)\left(\tilde{\nu}_{1}+\tilde{V}_{3} V_{4}^{-}+\tilde{V}_{4}^{+} \nu_{2}\right)\right] \\
&= {\left[\left(\partial_{u} \nu_{1}\right)\left(\partial_{\tilde{u}} \tilde{\nu}_{1}\right)\right]^{2}\left[\left(\partial_{u} \nu_{2}\right)\left(\partial_{\tilde{u}} \tilde{\nu}_{2}\right)\right] \exp \left(\psi_{2}+\tilde{\psi}_{2}\right), }
\end{aligned}
$$

using (3.28) and (3.29b). So (3.30) becomes

$$
\partial_{\bar{u}} \partial_{u}\left(\psi_{1}+\tilde{\psi}_{1}\right)=\exp \left[\left(\psi_{2}+\tilde{\psi}_{2}\right)-2\left(\psi_{1}+\tilde{\psi}_{1}\right)\right]
$$

and comparison of this with (2.26) shows that $\psi_{1}+\tilde{\psi}_{1}$ given by (3.29a) indeed satisfies the Toda lattice equation.
(iv) $S U(4)$. We shall briefly consider the main features of the solutions. It is illustrative of the general $\mathrm{SU}(n+1)$ solution, since unlike $\mathrm{SU}(2)$ and $\mathrm{SU}(3)$, but similar to $\mathrm{Sp}(2)$, the original form of the constraints (3.20) is rearranged. In this way the derivatives of coefficients of the $E_{-\alpha}$ in $n_{0}$ corresponding to non-simple roots are expressed as the derivatives of coefficients of those for simple roots. Here, if

$$
n_{0}=\exp \left(\nu_{i j} e_{i j}\right),
$$

then the form (3.20a) of the constraints are

$$
\begin{aligned}
& \partial_{u} V_{31}-\nu_{32}\left(\partial_{u} \nu_{21}\right)=0, \quad \partial_{u} V_{42}-\nu_{43}\left(\partial_{u} \nu_{32}\right)=0, \\
& \partial_{u} V_{41}-\nu_{43}\left(\partial_{u} V_{31}\right)-V_{42} \partial_{u} \nu_{21}=0,
\end{aligned}
$$

where

$$
\begin{aligned}
& V_{31}=\nu_{31}+\frac{1}{2} \nu_{32} \nu_{21}, \quad V_{42}=\nu_{42}+\frac{1}{2} \nu_{43} \nu_{32}, \\
& V_{41}=\nu_{41}+\frac{1}{2}\left(\nu_{42} \nu_{21}+\nu_{43} \nu_{31}\right) .
\end{aligned}
$$

The first two constraint equations have a convenient form and, using the first, the third can be reduced to

$$
\partial_{u} V_{41}-V_{42} \partial_{u} \nu_{21}=0
$$

as required. The value of this rearrangement was shown for $S p(2)$ since it provides for easy manipulation of $\Psi$ in (3.30) to give a common factor $\left(\partial_{u} \nu_{1}\right)\left(\partial_{\bar{u}} \tilde{\nu}_{1}\right)$. There are similar conditions for the tilde parameters.

The general parametrisation for the $\operatorname{SU}(n+1)$ solution is given in appendix 2 , so we shall not consider further the $\mathrm{SU}(4)$ case.
(v) $S O(5)$ and $S O$ (6). Since the algebras $b_{2}$ and $c_{2}$ are isomorphic, the solutions of the Toda lattice equations governed by $\mathrm{SO}(5)$ are the same as those for $\mathrm{Sp}(2)$. The Cartan matrix for $\mathrm{SO}(5)$ is the transpose of that for $\mathrm{Sp}(2)$, so if we exchange the parameters $\psi_{1}$ and $\psi_{2}$ in (3.29) we have solutions (3.23) for $\mathrm{SO}(5)$. Also, $\mathfrak{a}_{3}$ and $\mathrm{D}_{3}$ are isomorphic and similar remarks hold for their Cartan matrices. So the solutions for $\operatorname{SO}(6)$ are given by those for $\operatorname{SU}(4)$ with $\psi_{3}$ unchanged and $\psi_{1}$ and $\psi_{2}$ interchanged.

## 4. Discussion and concluding remarks

In $\S \S 2$ and 3 we have applied the Iwasawa decomposition to the complexification of some semi-simple Lie group $G$. By this method we have been able not only to derive the Toda lattice equations governed by G, but also to parametrise the solutions to the equations.

The derivation of the equations is based on a generalisation of Yang's $R$-gauge method. However, unlike Yang, we have not imposed any relation between the parameters of $R$ and $\tilde{R}$. This is possible since we are always able to solve the subsidiary equations for $z_{\alpha}$ and $\tilde{z}_{\alpha}$ in terms of $\psi_{\alpha}$ and $\tilde{\psi}_{\alpha}$, and the latter pair of parameters always appear in the combination

$$
\psi_{\alpha}+\tilde{\psi}_{\alpha}=-\left(K^{-1}\right)_{\alpha \beta} \rho_{\beta}
$$

However, if we wish to take a more physical viewpoint, then we should be able to recover the original gauge theory when we transform back to the real coordinates, $s$ and $t$. In this case we must demand that $B_{u}$ and $B_{\bar{u}}$ have values in the Lie algebra of $G$ and not $G_{c}$. Consequently, we have to impose some condition on $B_{u}$ and $B_{\bar{u}}$ and hence impose some relationship between $R$ and $\tilde{R}$. The condition is constructed by invoking the differences between the Lie algebras of $G$ and $G_{c}$. We show that a sufficient but not necessary condition is given by

$$
\begin{equation*}
\tilde{R}=\left(R^{+}\right)^{-1}, \tag{4.1}
\end{equation*}
$$

not only for $\mathrm{G}=\mathrm{SU}(n+1)$ (Yates 1978), but also when G is any classical group.
The complex extensions $G_{c}$ must first be divided into two categories to state explicitly the differences between $G$ and $G_{c}$. Category (i) contains $\operatorname{SL}(n+1, \mathbb{C})$ and $\operatorname{Sp}(n, \mathbb{C})$, whilst (ii) is comprised of $\mathrm{SO}(2 n+1, \mathbb{C}), \mathrm{SO}(2 n, \mathbb{C})$. For those in (i) the elements of $G$ are unitary and in (ii) the elements of $G$ are real. These properties are not exhibited by the corresponding $G_{c}$ and so the respective conditions on the potentials are

$$
\begin{align*}
& \text { (i) } B_{u}^{*}=-B_{\bar{u}},  \tag{4.2a}\\
& \text { (ii) } B_{u}^{*}=B_{\bar{u}} . \tag{4.2b}
\end{align*}
$$

These conditions are satisfied if (4.1) holds for elements $R, \tilde{R} \in \mathrm{G}_{\mathrm{c}} / \mathrm{G}$ in both categories. Since

$$
B_{\bar{u}}=R^{\dagger}\left(\partial_{\bar{u}}\left(R^{\dagger}\right)^{-1}\right)=-\left[R^{-1} \partial_{u} R\right]^{\dagger}=-B_{u}^{\dagger}
$$

(4.2a) is satisfied. Also for category (ii) the elements of $\mathrm{G}_{\mathrm{c}}$ and G are orthogonal, so

$$
\left(R^{\dagger}\right)^{-1}=\left(R^{T}\right)^{\dagger}=R^{*}
$$

and (4.1) becomes

$$
\tilde{R}=R^{*}
$$

which ensures that ( $4.2 b$ ) is satisfied.
Since the potentials are left-invariant the condition (4.1) is only sufficient for (4.2) to be satisfied. We may left translate $R$ and $\tilde{R}$, so that the transformed elements are no longer related by (4.1), but (4.2) still remains true.

Using the representations (2.6) and (2.9), the condition (4.1) reduces to a relationship between the parameters of $R$ and $\tilde{R}$, namely

$$
\begin{equation*}
\tilde{z}_{\alpha}=\bar{z}_{\alpha}, \quad \tilde{\psi}_{\alpha}=\psi_{\alpha} \tag{4.3}
\end{equation*}
$$

A second notable feature of the method here is the introduction of the factors $r(u)$ and $\tilde{r}(\bar{u})$ to obtain explicit expressions for the solutions. This feature is not paralleled by the one-dimensional model, obviously since the functions are only dependent on the one coordinate, $t$. However, the algebraic properties of the one-dimensional model of Olshanetsky and Perelomov (1979) do have some similarity to the two-dimensional case. If we wish to compare our method with that of Olshanetsky and Perelomov (hereafter referred to as OP), we must relate $r$ and $\tilde{r}$ to some factor in the onedimensional model.

To achieve this, we first reduce $u$ and $\bar{u}$ to $t$, so that (1.18) becomes

$$
\begin{equation*}
\dot{A}-\dot{B}=[A, B] \tag{4.4}
\end{equation*}
$$

with $A=R^{-1} \dot{R}$ and $B=\bar{R}^{-1} \hat{R}$, the dot referring to differentiation with respect to $t$. If we introduce

$$
\begin{equation*}
L=B-A \tag{4.5}
\end{equation*}
$$

then (4.4) has the form of a Lax pair

$$
\begin{equation*}
\dot{L}=[L, A] \tag{4.6}
\end{equation*}
$$

Drawing on the theory of the Lax pair, we can say that (4.6) is equivalent to

$$
\begin{equation*}
L(t)=R^{-1} L(0) R, \quad A=R^{-1} \dot{R} \tag{4.7}
\end{equation*}
$$

By rearranging (4.7) and then using (4.5), we obtain

$$
L(0)=R L(t) R^{-1}=R \tilde{R}^{-1} \dot{R} R^{-1}-\dot{R} R^{-1}=-\partial_{t}\left(R \tilde{R}^{-1}\right)\left(R \tilde{R}^{-1}\right)^{-1},
$$

which becomes

$$
\begin{equation*}
L(0)=-\dot{x} x^{-1} \tag{4.8}
\end{equation*}
$$

when we define

$$
\begin{equation*}
x=R \tilde{R}^{-1} \tag{4.9}
\end{equation*}
$$

Substitution of the decompositions (2.11) for $R$ and $\tilde{R}$ in (4.9) gives

$$
\begin{equation*}
x=n a \tilde{a}^{-1} \tilde{n}^{-1} \equiv n h \tilde{n}^{-1} \tag{4.10}
\end{equation*}
$$

where $h=a \tilde{a}^{-1}$ is a diagonal matrix. (4.10) exactly reproduces the results in OP. $x$ is
described by the initial data since it follows from (4.8) that

$$
\begin{equation*}
x(t)=x(0) \exp [-L(0) t] \tag{4.11}
\end{equation*}
$$

From (4.9) and (3.13) we observe that $x(t)$ corresponds to $\tilde{r}_{0}^{-1} r_{0}$ and also in (4.10) $h$ contains the Toda lattice variables. So the matrix manipulations on $x(t)$ used in OP to solve the equations are equivalent to those we apply to (3.13), also to solve the equations.

There may be a fundamental reason mathematically for splitting up $x$ into the product of functions (4.9). Such an idea deserves more thought, particularly with reference to the theory of the symmetric space $G_{\mathrm{c}} / \mathrm{G}$ to which $R$ and $\tilde{R}$ belong.

In § 2 the Toda lattice equations are derived from the condition that the gauge field strength $F_{u \bar{u}}$ should be sourceless, and in $\S 3$ this condition was shown to imply that the curvature two-form (3.2) should vanish. Now (3.2) may be considered as the integrability condition for the set of linear equations

$$
\begin{equation*}
\mathrm{d} \theta=\theta \omega \tag{4.12}
\end{equation*}
$$

where $\theta$ is a row of 0 -forms. The set of equations (4.12) should be investigated for our model, since it is anticipated that they might provide some insight into the Bäcklund transformations of the Toda lattice equations (Farwell and Minami 1982).

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## Appendix 1. A matrix representation for the Chevalley basis

We shall use $e_{i j}, 1 \leqslant i, j \leqslant n$, to denote the $n \times n$ matrix in which the only non-zero entry is 1 , where the $i$ th row and $j$ th column intersect (Humphreys 1972). From this definition it follows that

$$
\begin{equation*}
e_{i j} e_{k l}=\delta_{i k} e_{i l} \tag{A1.1}
\end{equation*}
$$

and so

$$
\begin{equation*}
\left[e_{i j}, e_{k l}\right]=\delta_{j k} e_{i l}-\delta_{l i} e_{k j} \tag{A1.2}
\end{equation*}
$$

In each case below the suffix $i$ on the basis elements $H_{i}$ of $\mathfrak{h}$ refers to some specific ordering of the simple roots comprising $\pi^{+}$. By $E_{+i}$ and $E_{-i}$, we mean the corresponding elements of the bases of $g_{+\alpha}$ and $g_{-\alpha}$ respectively; that is, those which satisfy the commutation relations ( 2.2 C ), namely

$$
\left[E_{+i}, E_{--i}\right]=H_{i} .
$$

The required $E_{+i}, E_{-i}$ may be selected from the sets $\left\{E_{+\alpha}\right\},\left\{E_{-\alpha}\right\}$ by inspection of the commutation relations derived using (A1.2).

Note that in the following representations for each class of Lie algebra, the matrix $E_{+i}$ is taken to be the transpose of $E_{-i}$. Such a specification conforms with the
commutation relations (2.2).
(i) $a_{n}-(n+1) \times(n+1)$ matrix representation

Basis for $\mathfrak{b}$

| $\mathfrak{g}_{-\alpha, \alpha \in \Delta^{+}}$ | $e_{i j}$ | $(1 \leqslant j<i \leqslant n+1)$ |
| :--- | :--- | :--- |
| $\mathfrak{g}_{+\alpha, \alpha \in \Delta^{+}}$ | $e_{j i}=e_{i j}^{\mathrm{T}}$ | $(1 \leqslant j<i \leqslant n+1)$. |

(ii) $\mathbf{b}_{n}-(2 n+1) \times(2 n+1)$ matrix representation

Basis for $\mathfrak{b}$

$$
\begin{aligned}
& H_{i-1}=e_{i i}-e_{i+1 i+1}-e_{n+i n+i}+e_{n+i+1 n+i+1} \quad(2 \leqslant i \leqslant n) \\
& H_{n}=e_{n+1 n+1}-e_{2 n+12 n+1}
\end{aligned}
$$

| $\mathfrak{g}_{-\alpha, \alpha \in \Delta^{+}}$ | $e_{i+1 j+1}-e_{j+n+1 i+n+1}$ | $(1 \leqslant j<i \leqslant n)$ |
| :--- | :--- | :--- |
|  | $e_{i+1+n 1}-e_{1 i+1}$ | $(1 \leqslant i \leqslant n)$ |
|  | $e_{n+i+1 j+1}-e_{n+j+1 i+1}$ | $(1 \leqslant j<i \leqslant n)$ |
| $\mathbf{g}_{+\alpha, \alpha \in \Delta^{+}}$ | $e_{j+1 i+1}-e_{i+n+1 j+n+1}$ | $(1 \leqslant j<i \leqslant n)$ |
|  | $e_{1 i+1+n}-e_{i+11}$ | $(1 \leqslant i \leqslant n)$. |
|  | $e_{i+1 n+i+1}-e_{i+1 n+j+1}$ | $(1 \leqslant j<i \leqslant n)$. |

(iii) $\mathrm{c}_{n}-2 n \times 2 n$ matrix representation

Basis for $\mathfrak{h}$
$H_{i}=e_{i i}-e_{i+1 i+1}-e_{n+i n+i}+e_{n+i+1 n+i+1} \quad(1 \leqslant i \leqslant n-1)$
$H_{n}=e_{n n}-e_{2 n 2 n}$

| $\mathbf{g}_{-\alpha, \alpha \in \Delta^{+}}$ | $e_{i j}-e_{j+n i+n}$ | $(1 \leqslant j<i \leqslant n)$ |
| :--- | :--- | :--- |
|  | $e_{n+i i}$ | $(1 \leqslant i \leqslant n)$ |
|  | $e_{n+j i}+e_{n+i j}$ | $(1 \leqslant j<i \leqslant n)$ |
| $\mathbf{g}_{+\alpha, \alpha \in \Delta^{+}}$ | $e_{j i}-e_{n+i j+n}$ | $(1 \leqslant j<i \leqslant n)$ |
|  | $e_{i n+i}$ | $(1 \leqslant i \leqslant n)$ |
|  | $e_{i n+j}+e_{j n+i}$ | $(1 \leqslant j<i \leqslant n)$. |

(iv) $\mathrm{D}_{n}-2 n \times 2 n$ matrix representation

Basis for $\mathfrak{h}$
$H_{i}=e_{i i}-e_{i+1 i+1}-e_{n+i n+i}+e_{n+i+1 n+i+1} \quad(1 \leqslant i \leqslant n-1)$
$H_{n}=e_{n-1 n-1}+e_{n n}-e_{2 n-12 n-1}-e_{2 n 2 n}$

| $\mathbf{g}_{-\alpha, \alpha \in \Delta^{+}}$ | $e_{i j}-e_{j+n i+n}$ | $(1 \leqslant j<i \leqslant n)$ |
| :--- | :--- | :--- |
|  | $e_{n+i j}-e_{n+i j}$ | $(1 \leqslant j<i \leqslant n)$ |
| $\mathbf{g}_{+\alpha, \alpha \in \Delta^{+}}$ | $e_{j i}-e_{i+n j+n}$ | $(1 \leqslant j<i \leqslant n)$ |
|  | $e_{j n+i}-e_{n+j i}$ | $(1 \leqslant j<i \leqslant n)$. |

It is important to note various properties of the bases for each class. In all cases, the $H_{i}$ are traceless diagonal matrices. Furthermore, for $a_{n}$ the set $\left\{E_{-\alpha}\right\}$ is composed of lower triangular matrices with zeros on the diagonal, and so the transposed set $\left\{E_{+\alpha}\right\}$ consists of upper triangular matrices with no diagonal entries. Although the sets $\left\{E_{-\alpha}\right\}$
are not triangular for the other three classes, it is possible to select submatrices which do have this property. Let $\left[\left(E_{-\alpha}\right)_{l m}\right]$ denote a submatrix of the basis element $E_{-\alpha}$ where the range of $l$ and $m$ is specified. Then
(a) for $\mathfrak{b}_{n}$ : the set $\left\{\left[\left(E_{-\alpha}\right)_{l m}\right]_{2 \leqslant l, m \leqslant n+1}\right\}$ are lower triangular with zero diagonal entries and the set $\left\{\left[\left(E_{-\alpha}\right)_{l m}\right]_{2 \leqslant 1 \leqslant n+1, m=1 \text { or } n+2 \leqslant m \leqslant 2 n+1}\right\}$ have all zero entries;
(b) for $\mathrm{c}_{n}$ and $\mathcal{D}_{n}$ : the set $\left\{\left[\left(E_{-\alpha}\right)_{l m}\right]_{1 \leq l, m \leq n}\right\}$ are lower triangular with zero diagonal entries and the set $\left\{\left[\left(E_{-\alpha}\right)_{l m}\right]_{1 \leqslant l \leqslant n, n+1 \leqslant m \leqslant 2 n}\right\}$ have all zero entries.

Corresponding results concerning upper diagonal and zero submatrices of $E_{+\alpha}$ may be deduced using the operation of transposition.

## Appendix 2. Parametrisation of the $\mathbf{S U}(n+1)$ solution

In $\S 3$ we show that the solution of the Toda lattice equation governed by $\mathrm{SU}(n+1)$ is given by

$$
\begin{equation*}
\psi_{m}+\tilde{\psi}_{m}=\mu_{m}+\tilde{\mu}_{m}+\ln \operatorname{det}\left(\tilde{n}_{0}^{-1} n_{0}\right)_{m}, \quad 1 \leqslant m \leqslant n . \tag{A2.1}
\end{equation*}
$$

In general, we use the representations (3.12) for $n_{0}$ and $\tilde{n}_{0}$. However, since in appendix 1 , for $\mathfrak{a}_{n}$ each $E_{-\alpha}$ has the simple form $e_{i j}, 1 \leqslant j<i \leqslant n+1$, we adopt the double index notation for $\nu$ and $\tilde{\nu}$ and write

$$
\begin{align*}
n_{0} & =\exp \left(\sum_{i>j} \nu_{i j} e_{i j}\right) \\
& =1+\sum_{1 \leqslant j<i \leqslant n+1} \nu_{i j} e_{i j}+\sum_{1 \leqslant j<k<i \leqslant n+1} \nu_{i k} \nu_{k j} e_{i j}  \tag{A2.2}\\
\tilde{n}_{0}^{-1} & =\exp \left(\sum_{i>i} \tilde{\nu}_{j i} e_{j i}\right) \\
& =1+\sum_{1 \leqslant i<i \leqslant n+1} \tilde{\nu}_{i j} e_{j i}+\sum_{1 \leqslant i<k<i \leqslant n+1} \tilde{\nu}_{j k} \tilde{\nu}_{k i} e_{j i} \tag{A2.3}
\end{align*}
$$

We use the abbreviations

$$
\begin{equation*}
V_{i j}^{ \pm}=\nu_{i j} \pm \sum_{k=j+1}^{i-1} \nu_{i k} \nu_{k j}, \quad \tilde{V}_{i j}^{ \pm}=\tilde{\nu}_{j i} \pm \sum_{k=j+1}^{i-1} \tilde{\nu}_{j k} \tilde{\nu}_{k i} . \tag{A2.4}
\end{equation*}
$$

The range of summation in the second terms ensures that corresponding to the simple roots, $i=j+1$,

$$
V_{i j}^{ \pm} \equiv \nu_{i j}, \quad \tilde{V}_{i i}^{ \pm} \equiv \tilde{\nu}_{j i}
$$

Then (A2.2), (A2.3) become

$$
\begin{equation*}
n_{0}=1+\sum_{1 \leqslant j<i \leqslant n+1} V_{i j}^{+} e_{i j}, \quad \tilde{n}_{0}^{-1}=1+\sum_{1 \leqslant j<i \leqslant n+1} \tilde{V}_{i i}^{+} e_{i j} . \tag{A2.5}
\end{equation*}
$$

Conditions on $V_{i j}$ and $\tilde{V}_{j i}$ arise as a consequence of the constraints (3.20):

$$
\begin{array}{ll}
\partial_{u} V_{i j}^{+}-V_{i k}^{-}\left(\partial_{u} V_{k j}^{+}\right)=0 & \text { for } i>j+1, \\
\partial_{\tilde{u}} \tilde{V}_{i i}^{-}+\tilde{V}_{j k}^{+}\left(\partial_{\bar{u}} \tilde{V}_{k i}^{-}\right)=0 & \text { for } i>j+1 . \tag{A2.7}
\end{array}
$$

Now, from (A2.5)

$$
\begin{gathered}
\tilde{n}_{0}^{-1} n_{0}=1+\sum_{1 \leqslant i<k \leqslant n+1} \tilde{V}_{i k}^{+} V_{k i}^{+} e_{i i}+\sum_{1 \leqslant j<i \leqslant n+1}\left(\tilde{V}_{j i}^{+}+\sum_{k=j+1}^{i-1} \tilde{V}_{i k}^{+} V_{k i}^{+}\right) e_{j i} \\
+\sum_{1 \leqslant j<i \leqslant n+1}\left(V_{i j}^{+}+\sum_{k=j+1}^{i-1} V_{i k}^{+} \tilde{V}_{k j}^{+}\right) e_{i j}
\end{gathered}
$$

and so the submatrix

$$
\begin{align*}
&\left(\left(\tilde{n}_{0}^{-1} n_{0}\right)_{p q}\right)_{1 \leqslant p, q \leqslant m<n}=I_{m}+\sum_{i=1}^{m} \sum_{k=i+1}^{n+1} \tilde{V}_{i k}^{+} V_{k i}^{+} e_{i i} \\
&+\sum_{1 \leqslant j<i \leqslant m}\left(\tilde{V}_{j i}^{+}+\sum_{k=j+1}^{i-1} \tilde{V}_{j k}^{+} V_{k i}^{+}\right) e_{i j}+\sum_{1 \leqslant j<i \leqslant m}\left(V_{i j}^{+}+\sum_{k=j+1}^{i-1} V_{i k}^{+} \tilde{V}_{k i}^{+}\right) e_{i j} \tag{A2.8}
\end{align*}
$$

The relations (3.21) and (3.22) become for each $j=1,2, \ldots, n$

$$
\begin{equation*}
\left(\partial_{u} V_{i+1 j}\right)\left(\partial_{\bar{u}} \tilde{V}_{i j+1}\right)=\exp \left(-\sum_{m=1}^{n} K_{i m}\left(\mu_{m}+\tilde{\mu}_{m}\right)\right) \tag{A2.9}
\end{equation*}
$$

where $\left(K_{j m}\right): j, m=1,2, \ldots, n$ is the Cartan matrix for $\operatorname{SU}(n+1)$.
If we define the variable

$$
\begin{equation*}
\rho_{j}=-K_{j m}\left(\psi_{m}+\tilde{\psi}_{m}\right), \tag{A2.10}
\end{equation*}
$$

then the solution (A2.1) becomes for $j=1,2, \ldots, n$,

$$
\begin{equation*}
\rho_{i}=\ln \left[\left(\partial_{u} \nu_{j+1 j}\right)\left(\partial_{u} \tilde{\nu}_{j+1 j}\right)\right]+\ln \left(\frac{\operatorname{det}\left(\tilde{n}_{0}^{-1} n_{0}\right)_{i+1} \operatorname{det}\left(\tilde{n}_{0}^{-1} n_{0}\right)_{i-1}}{\left[\operatorname{det}\left(\tilde{n}_{0}^{-1} n_{0}\right)_{j}\right]^{2}}\right) \tag{A2.11}
\end{equation*}
$$

with $\left(\tilde{n}_{0}^{-1} n_{0}\right)_{m}$ given by (A2.8), with the proviso that the parameters are subjected to (A2.6) and (A2.7).

## References

Ardalan F 1978 Phys. Rev. 18 1960-8
Bais F A and Weldon H A 1978 Phys. Rev. Lett. 41 601-4
Bogoyavlensky O I 1976 Commun. Math. Phys. 51 201-9
Brihaye Y, Fairlie D B, Nuyts J and Yates R G 1978 J. Math. Phys. 19 2528-32
Carter R W 1972 Simple Groups of Lie Type (London: Wiley)
Farwell R S and Minami M 1982 J. Phys. A: Math. Gen. 15 in press
Flaschka H 1974a Phys. Rev. B 9 1924-5
—— 1974b Prog. Theor. Phys. 51 703-16
Gilmore R 1974 Lie Groups, Lie Algebras and Some of Their Applications (New York: Wiley)
Helgason S 1978 Differential Geometry, Lie Groups, and Symmetric Spaces (New York: Academic)
Hermann R 1966 Lie Groups for Physicists (New York: Benjamin)
Humphreys J E 1972 Introduction to Lie Algebras and Representation Theory (New York: Springer)
Kostant B 1973 Ann. Sci. Éc. Norm. Sup. $4{ }^{e}$ Serie 6 413-55

- 1979 Adv. Math. 34 195-338

Lax P D 1968 Comm. Pure Appl. Math. 21 467-91
Leznov A N and Saveliev M V 1978 Phys. Lett. B 79 294-6

- 1979a Phys. Lett. B 83 314-6
_-1979b Lett. Math. Phys. 3 489-94
_- 1980 Commun. Math. Phys. 74 111-8

Manakov S V 1974 Sov. Phys.-JETP 40 269-74
Moser J 1975 Lecture Notes in Physics (Dynamical systems; theory and applications) ed J Moser 38 467-97
Olive D 1980 Imperial College Preprint ICTP/80/81-1
Olshanetsky M A and Perelomov A M 1979 Invent. Math. 54 261-9
Omnés R 1979 Nucl. Phys. B 149 269-84
Toda M 1967 J. Phys. Soc. Japan 22 431-6
Witten E 1977 Phys. Rev. Lett. 38 121-4
Yang C N 1977 Phys. Rev. Lett. 38 1377-9
Yates R G 1978 Some Classical Solutions of Yang-Mills Equations (PhD Thesis, Durham)


[^0]:    $\ddagger$ On leave of absence from the Research Institute for Mathematical Sciences, Kyoto University, Kyoto, $\ddagger$ Some of the more seminal contributions to the mathematical structure were: discovery of Lax pair; derivation of equations for finite systems and their solution; analogy with group structure (Lax 1968, Flaschka 1974a, b, Manakov 1974, Bogoyavlensky 1976, etc). Subsequently, Leznov and Saveliev (1978) in considering the two-dimensional model used the connection with the Cartan matrix; Olshanetsky and Perelomov (1979) couched the solution to the one-dimensional model in algebraic terminology, and independently Kostant (1979) found solutions to more general integrable systems.

[^1]:    $\dagger$ Complex extensions of $\operatorname{SU}(n+1), \mathrm{SO}(2 n+1), \mathrm{Sp}(n)$ and $\mathrm{SO}(2 n)$ are denoted $\mathrm{SL}(n+1, \mathbb{C}), \mathrm{SO}(2 n+1, \mathbb{C})$, $\operatorname{Sp}(n, \mathbb{C})$ and $\operatorname{SO}(2 n, \mathbb{C})$ and their algebras $a_{n}, b_{n}, c_{n}$ and $b_{n}$ respectively. It should be remarked that $\operatorname{Sp}(n)$ is not $\operatorname{Sp}(n, R)$, the latter not being compact.

[^2]:    $\dagger$ David Olive has made us aware that these restricted forms for $r$ and $\tilde{r}$ may be ascribed to the path independency of $g_{0}=T \exp \int_{(0,0)}^{(u, \overline{0})}\left(B_{u} \mathrm{~d} u+B_{\bar{u}} \mathrm{~d} \bar{u}\right)$, Since $R(u, \tilde{u})=T \exp \int_{(0, \tilde{u}}^{(u, \bar{u})} B_{u} \mathrm{~d} u$ and $\tilde{R}(u, \bar{u})=$ $T \exp \int_{(u, 0)}^{(u, \bar{u})} \boldsymbol{B}_{\bar{u}}$ d $\bar{u}$, we can write $g_{0}$ in two ways: $g_{0}=R(u, 0) \tilde{R}(u, \bar{u})=\tilde{R}(0, \bar{u}) \boldsymbol{R}(u, \bar{u})$. We thank him for this communication.

