Integrable scalar cosmologies
II. Can they fit into Gauged Extended Supergravity
or be encoded in $\mathcal{N} = 1$ superpotentials?

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Abstract

The question whether the integrable one-field cosmologies classified in a previous paper by Fré, Sagnotti and Sorin can be embedded as consistent one-field truncations into Extended Gauged Supergravity or in $\mathcal{N} = 1$ supergravity gauged by a superpotential without the use of $D$-terms is addressed in this paper. The answer is that such an embedding is very difficult and rare but not impossible. Indeed, we were able to find two examples of integrable models embedded in supergravity in this way. Both examples are fitted into $\mathcal{N} = 1$ supergravity by means of a very specific and interesting choice of the superpotential $W(z)$. The question whether there are examples of such an embedding in Extended Gauged Supergravity remains open.

In the present paper, relying on the embedding tensor formalism we classified all gaugings of the $\mathcal{N} = 2$ STU model, confirming, in the absence on hypermultiplets, the uniqueness of the stable de Sitter vacuum found several years ago by Fré, Trigiante and Van Proeyen and excluding the embedding of any integrable cosmological model. A detailed analysis of the space of exact solutions of the first supergravity-embedded integrable cosmological model revealed several new features worth an in-depth consideration. When the scalar potential has an extremum at a negative value, the Universe necessarily collapses into a Big Crunch notwithstanding its spatial flatness. The causal structure of these Universes is quite different from that of

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the closed, positive curved, Universe: indeed, in this case the particle and event horizons do not coincide and develop complicated patterns. The cosmological consequences of this unexpected mechanism deserve careful consideration.

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1. Introduction

In a recent paper [1] some of us have addressed the question of classifying integrable one-field cosmological models based on a slightly generalized ansatz for the spatially flat metric,

\[ ds^2 = e^{2B(t)} dt^2 - a^2(t) dx^i dx^j, \]  

(1.1)

and on a suitable choice of a potential \( V(\phi) \) for the unique scalar field \( \phi \), whose kinetic term is supposed to be canonical:

\[ \mathcal{L}_{\text{kin}}(\phi) = \frac{1}{2} \partial^{\mu} \phi \partial_{\mu} \phi \sqrt{-g}. \]  

(1.2)

The suitable potential functions \( V(\phi) \) that lead to exactly integrable Maxwell–Einstein field equations were searched within the family of linear combinations of exponential functions \( \exp \beta \phi \) or rational functions thereof. The motivations for such a choice were provided both with string and supergravity arguments and a rather remarkable bestiary of exact cosmological solutions was uncovered, endowed with quite interesting mathematical properties. Some of these solutions have also some appeal as candidate models of the inflationary scenario, capable of explaining the structure of the primordial power spectrum.

In [1] the classical Friedman equations

\[ H^2 = \frac{1}{3} \dot{\phi}^2 + \frac{2}{3} V(\phi), \]
\[ \dot{H} = -\dot{\phi}^2, \]
\[ \ddot{\phi} + 3H \dot{\phi} + V' = 0, \]  

(1.3)

where

\[ a(t) = e^{A(t)}, \quad H(t) \equiv \frac{\dot{a}(t)}{a(t)} = \dot{A}(t) \]  

(1.4)

are respectively the scale factor and the Hubble function, were revisited in the more general gauge with \( B \neq 0 \) which allows for the construction of exact solutions, whenever the effective two-dimensional dynamical system lying behind Eqs. (1.3) can be mapped, by means of a suitable canonical transformation, into an integrable dynamical model endowed with two Hamiltonians in involution. Such procedure produced the bestiary constructed and analyzed in [1].

After the change of perspective produced by the recent series of papers [2–12] and in particular after [8,9], we know that all positive definite members of the above mentioned bestiary can be embedded into \( N = 1 \) supergravity as \( D \)-terms produced by the gauging of an axial symmetry, provided the Kähler manifold to which we assign the Wess–Zumino multiplet of the inflaton is consistent with the chosen potential \( V(\phi) \), namely it has an axial symmetric Kähler potential defined in a precise way by \( V(\phi) \). In [13] which is published at the same time as the present paper, two of us have analyzed the mathematical algorithm lying behind this embedding mechanism which we have named the \( D \)-map. In the same paper a possible path toward the microscopic
interpretation of the peculiar axial symmetric Kähler manifolds requested by the $D$-type supergravity embedding of the integrable potentials is proposed and discussed. Such a microscopic interpretation is obligatory in order to give a sound physical meaning to the supergravity embedding.

The main theme that we are going to debate in the present paper is instead the following. Can integrable cosmologies be embedded into gauged extended supergravity or in $\mathcal{N}=1$ supergravity gauged by the $F$-terms that are produced by the choice of some suitable superpotential $W(z)$? In the present paper the choice of the Kähler geometry for the inflaton will not depend on the potential $V(\phi)$. The inflaton Wess–Zumino multiplet will always be assigned to a constant curvature Kähler manifold as it is the case in compactifications on torii, orbifolds or orientifolds.

Having clarified this fundamental distinction between the complementary approaches of the present paper and of the parallel paper [13], we continue our discussion of the Friedman system (1.3). Referring to the classical cosmic time formulation (1.3) of Friedman equations and to the very instructive hydrodynamical picture, we recall that the energy density and the pressure of the fluid describing the scalar matter can be identified with the two combinations

$$\rho = \frac{1}{4} \dot{\phi}^2 + \frac{1}{2} V(\phi),$$
$$p = \frac{1}{4} \dot{\phi}^2 - \frac{1}{2} V(\phi),$$

(1.5)
since, in this fashion, the first of Eqs. (1.3) translates into the familiar link between the Hubble constant and the energy density of the Universe,

$$H^2 = \frac{4}{3} \rho.$$  
(1.6)

A standard result in General Relativity (see for instance [39]) is that for a fluid whose equation of state is

$$p = w\rho, \quad w \in \mathbb{R},$$

(1.7)
the relation between the energy density and the scale factor takes the form

$$\frac{\rho}{\rho_0} = \left( \frac{a_0}{a} \right)^{3(1+w)},$$

(1.8)
where $\rho_0$ and $a_0$ are their values at some reference time $t_0$. Combining Eq. (1.7) with the first of Eqs. (1.3) one can then deduce that

$$a(t) \sim (t-t_i)^2w^{3(w+1)}$$

(1.9)
where $t_i$ is an initial cosmic time. All values $-1 \leq w \leq 1$ can be encompassed by Eqs. (1.5), including the two particularly important cases of a dust-filled Universe, for which $w=0$ and $a(t) \sim (t-t_i)^\frac{2}{3}$, and of a radiation-filled Universe, for which $w = \frac{1}{3}$ and $a(t) \sim (t-t_i)^\frac{1}{2}$. Moreover, when the potential energy $V(\phi)$ becomes negligible with respect to the kinetic energy in Eqs. (1.5), $w \approx 1$. On the other hand, when the potential energy $V(\phi)$ dominates $w \approx -1$, and Eq. (1.8) implies that the energy density is approximately constant (vacuum energy) $\rho = \rho_0$. The behavior of the scale factor is then exponential, since the Hubble function is also a constant $H_0$ on account of Eq. (1.6), and therefore

$$a(t) \sim \exp[H_0 t], \quad H_0 = \sqrt[3]{\frac{4}{3} \rho_0}.$$
The actual solutions of the bestiary described in [1] correspond to complicated equations of state whose index $w$ varies in time. Nonetheless they are qualitatively akin, at different epochs, to these simple types of behavior.

As we just stressed, the next question that constitutes the main issue of the present paper is whether the integrable potentials classified in [1] play a role in consistent one-field truncations of four-dimensional Gauged Supergravity. A striking and fascinating feature of supergravity is in fact that its scalar potentials are not completely free. Rather, they emerge from a well defined gauging procedure that becomes more and more restrictive as the number $\mathcal{N}$ of supercharges increases, so that the link between the integrable cosmologies of [1] and this structure is clearly of interest.

The first encouraging observation was already mentioned: in all integrable models found in [1] the potential $V(\phi)$ is a polynomial or rational function of exponentials $\exp[\beta \phi]$ of a field $\phi$ whose kinetic term is canonical. If we discard the rational cases and we retain only the polynomial ones that are the majority, this feature connects naturally such cosmological models to Gauged Supergravity with scalar fields belonging to non-compact, symmetric coset manifolds $G/H$. This wide class encompasses not only all $\mathcal{N} > 2$ theories, but also some $\mathcal{N} \leq 2$ models that are frequently considered in connection with Cosmology, Black Holes, Compactifications and other issues.

Since the coset manifolds $G/H$ relevant to supergravity are characterized by a numerator group $G$ that is a non-compact semisimple group, in these models one can always resort to a solvable parameterization of the scalar manifold [15], so that the scalar fields fall into two classes:

1. The Cartan fields $h^i$ associated with the Cartan generators of the Lie algebra $G$, whose number equals the rank $r$ of $G/H$. For instance, in models associated with toroidal or orbifold compactifications, fields of this type are generically interpreted as radii of the underlying multi-tori.
2. The axion fields $b^I$ associated with the roots of the Lie algebra $G$.

The kinetic terms of Cartan scalars have the canonical form

$$
\sum_i \frac{\alpha_i^2}{2} \partial_\mu h^i \partial^\mu h^i,
$$

up to constant coefficients, while for the axion scalars entering solvable coset representatives, the $\alpha_i^2$ factors leave way to exponential functions $\exp[\beta_i h^i]$ of Cartan fields. The scalar potentials of Gauged Supergravity are polynomial functions of the coset representatives, so that after the truncation to Cartan sectors, setting the axions to constant values, one is led naturally to combinations of exponentials of the type encountered in [1]. Yet the devil lies in the details, since the integrable potentials do result from exponential functions $\exp[\beta h]$, but with rigidly fixed ratios between the $\beta_i$ entering the exponents and the $\alpha_i$ entering the kinetic terms. The candidate potentials are displayed in Tables 1 and 2 following the notations and the nomenclature of [1]. As a result, the possible role of integrable potentials in Gauged Supergravity theories is not evident a priori, and actually, the required ratios are quite difficult to be obtained. Notwithstanding these difficulties we were able to identify a pair of examples, showing that although rare, supergravity
convenient to adopt a slightly different starting point which touches upon some of the fundamental model considered in [1] with the possible one-field truncations of supergravity, it is

2. The set up for comparison with supergravity

In this paper we focus on $D = 4$ supergravity models. In order to compare the effective dynamical model considered in [1] with the possible one-field truncations of supergravity, it is convenient to adopt a slightly different starting point which touches upon some of the fundamen-

\begin{table}[h]
\centering
\begin{tabular}{|c|c|}
\hline
Potential function & \\
I_1 & $C_{11}e^{\psi} + 2C_{12} + C_{22}e^{-\psi}$ \\
I_2 & $C_1e^{2\psi} + C_2e^{(\gamma+1)\psi}$ \\
I_3 & $C_1e^{2\psi} + C_2$ \\
I_7 & $C_1(\cosh \gamma \psi) \frac{2}{\gamma} - 2 + C_2(\sinh \gamma \psi) \frac{2}{\gamma} - 2$ \\
I_8 & $C_1(\cosh(2\gamma \psi)) \frac{2}{\gamma} - 1 \cos([\frac{1}{\gamma} - 1 \arccos(\tanh(2\gamma \psi) + C_2)])$ \\
I_9 & $C_1e^{2\psi} + C_2e^{\frac{2}{\gamma} \psi}$ \\
\hline
\end{tabular}
\caption{Potential function}
\end{table}

\begin{table}[h]
\centering
\begin{tabular}{|c|c|}
\hline
Sporadic integrable potentials & \\
\begin{align*}
\mathcal{V}_{\text{Ia}}(\psi) &= \frac{1}{2}[(a + b) \cosh(\frac{2}{3} \psi) + (3a - b) \cosh(\frac{2}{3} \psi)] \\
\mathcal{V}_{\text{Ib}}(\psi) &= \frac{1}{4}[(a + b) \sinh(\frac{6}{5} \psi) - (3a - b) \sinh(\frac{2}{5} \psi)] \\
&\text{where } \{a, b\} = \begin{pmatrix} 1 & -3 \\ 1 & -\frac{1}{2} \\ 1 & -\frac{3}{16} \end{pmatrix} \\
\mathcal{V}_{\text{II}}(\psi) &= \frac{1}{2}(3a + 3b - c + 4(a - b) \cosh(\frac{2}{3} \psi) + (a + b + c) \cosh(\frac{4}{5} \psi)), \\
&\text{where } \{a, b, c\} = \begin{pmatrix} 1 & 1 & -2 \\ 1 & 1 & -6 \\ 1 & 8 & -6 \\ 1 & 16 & -12 \\ 1 & \frac{8}{3} & -\frac{3}{4} \\ 1 & \frac{1}{16} & -\frac{3}{4} \end{pmatrix} \\
\mathcal{V}_{\text{IIIA}}(\psi) &= \frac{1}{16}[(1 - \frac{1}{3\sqrt{3}})e^{-6\psi/5} + (7 + \frac{1}{\sqrt{3}})e^{-2\psi/5} + (7 - \frac{1}{\sqrt{3}})e^{2\psi/5} + (1 + \frac{1}{3\sqrt{3}})e^{6\psi/5}] \\
\mathcal{V}_{\text{IIIB}}(\psi) &= \frac{1}{16}[(2 - 18\sqrt{3})e^{-6\psi/5} + (6 + 30\sqrt{3})e^{-2\psi/5} + (6 - 30\sqrt{3})e^{2\psi/5} + (2 + 18\sqrt{3})e^{6\psi/5}] \\
\end{align*}
\end{tabular}
\caption{Sporadic integrable potentials}
\end{table}

integrable cosmological models based on $G/S$ scalar manifolds do exist and might provide a very useful testing ground where exact calculations can be performed ab initio to the very end.

2. The set up for comparison with supergravity

In this paper we focus on $D = 4$ supergravity models. In order to compare the effective dynamical model considered in [1] with the possible one-field truncations of supergravity, it is convenient to adopt a slightly different starting point which touches upon some of the fundamen-

\footnote{The main consequence of the $D$-embedding of integrable potentials discussed in the parallel paper [13] is that the Kähler manifold hosting the inflaton is not a constant curvature coset manifold $G/H$.}
tal features of all supersymmetric extension of gravity. As we have already mentioned, differently from non-supersymmetric theories, where the kinetic and potential terms of the scalar fields are uncorrelated and disposable at will, the fascination of *sugras* is precisely that a close relation between these two terms here exists and is mandatory. Indeed the potential is just created by means of the gauging procedure. The explicit formulae for the potential always involve the metric of the target manifold which, on the other hand, determines the scalar field kinetic terms. Thus, in one-field truncations, the form of the kinetic term cannot be normalized at will but comes out differently, depending on the considered model and on the chosen truncation. A sufficiently ductile Lagrangian that encodes the various sugra-truncations discussed in this paper is the following one:

\[
L_{\text{eff}} = \text{const} \times e^{3A - B} \left[ -\frac{3}{2} \dot{A}^2 + \frac{q}{2} \dot{h}^2 - e^{2B} V(h) \right]. \tag{2.1}
\]

The field \(h\) is a residual dilaton field after all the other dilatonic and axionic fields have been fixed to their stationary values and \(q\) is a parameter, usually integer, that depends both on the chosen supergravity model and on the chosen truncation. The correspondence with the set up of [1] is simple: \(\phi = \sqrt{q} h\). Hence altogether the transformation formulae that correlate the general discussion of this paper with the bestiary of supergravity potentials, found in [1] and displayed in Tables 1 and 2 are the following ones:

\[
\begin{align*}
\dot{A}(t) &= 3H(t) = 3 \frac{d}{dt} \log[a(t)], \\
\ddagger \\
\dot{A}(t) &= 3A(t), \\
B(t) &= B(t), \\
\phi &= \sqrt{3q} h, \\
V(\phi) &= 3V(h) = 3V\left(\frac{\phi}{\sqrt{3q}}\right). \tag{2.2}
\end{align*}
\]

We will consider examples of \(\mathcal{N} = 2\) and \(\mathcal{N} = 1\) models trying to spot the crucial points that make it unexpectedly difficult to fit integrable cosmological models into the well established framework of *gauged supergravities*. Difficult but not impossible since we were able to identify at least one integrable \(\mathcal{N} = 1\) supergravity model based on the coupling of a single Wess–Zumino multiplet endowed with a very specific superpotential. While postponing to a further paper the classification of all the gaugings of the \(\mathcal{N} = 2\) models based on symmetric spaces [50] (see Table 3) and the analysis of their one-field truncation in the quest of possible matching with the integrable potentials, in the present paper we will consider in some detail another possible point of view. It was named the *minimalist approach* in the conclusions of [1]. Possibly no physically relevant cosmological model extracted from Gauged Supergravity is integrable, yet the solution of its field equations might be effectively simulated in their essential behavior by the exact solution of a neighboring integrable model. Relying on the classification of fixed points presented in [1] we advocate that if there is a one parameter family of potentials that includes both an integrable case and a case derived from supergravity and if the fixed point is the same for the integrable case and for the supergravity case, then the integrable model provides a viable substitute of the physical one and its solutions provide good approximations of the physical ones accessible only with numerical evaluations. We will illustrate this viewpoint with the a detailed analysis of one particularly relevant case.
Table 3
List of special Kähler homogeneous spaces in $D = 4$ with their $D = 3$ enlarged counterparts, obtained through Kaluza–Klein reduction. The number $n$ denotes the number of vector multiplets. The total number of vector fields is therefore $n_V = n + 1$.

<table>
<thead>
<tr>
<th>Coset $D = 4$</th>
<th>Coset $D = 3$</th>
<th>Susy $\mathcal{N}$</th>
</tr>
</thead>
<tbody>
<tr>
<td>SU(1,1) U(1)</td>
<td>G$_2$(2)</td>
<td>$\mathcal{N} = 2, n = 1$</td>
</tr>
<tr>
<td>Sp(6,R)</td>
<td>SU(2) × SU(2)</td>
<td></td>
</tr>
<tr>
<td>SU(3) U(1)</td>
<td>E$_6$(2)</td>
<td>$\mathcal{N} = 2, n = 6$</td>
</tr>
<tr>
<td>SU(3) × SU(3) U(1)</td>
<td>G$^*$ (12)</td>
<td></td>
</tr>
<tr>
<td>SU(6) U(1)</td>
<td>SO(12) × SU(2)</td>
<td>$\mathcal{N} = 2, n = 9$</td>
</tr>
<tr>
<td>E$_7(-25)$ × U(1)</td>
<td>E$_8(-24)$</td>
<td>$\mathcal{N} = 2, n = 15$</td>
</tr>
<tr>
<td>SL(2,2) × SO(2,2+p)</td>
<td>SO(4,4+p)</td>
<td>$\mathcal{N} = 2, n = 3 + p$</td>
</tr>
<tr>
<td>SU(p+1,1) U(1)</td>
<td>SU(p+2,2)</td>
<td>$\mathcal{N} = 2$</td>
</tr>
</tbody>
</table>

The obvious limitation of this approach is the absence of an algorithm to evaluate the error that separates the unknown physical solution from its integrable model clone. Yet a posteriori numerical experiments show that this error is rather small and that all essential features of the physical solution are captured by the solutions of the appropriate integrable member of the same family.

Certainly it would be very much rewarding if other integrable potentials could be derived from specific truncations of specially chosen supergravity gaugings. If such a case is realized, the particular choice of parameters that leads to integrability would certainly encode some profound physical significance.

3. $\mathcal{N} = 2$ models and stable de Sitter vacua

An issue of high relevance for a theoretical explanation of current cosmological data is the construction of stable de Sitter string vacua that break all supersymmetries [20], a question that is actually formulated at the level of the low-energy $N$-extended supergravity. As recently reviewed in [25], for $\mathcal{N} > 2$ no stable de Sitter vacua have ever been found and do not seem to be possible. In $\mathcal{N} = 1$ supergravity coupled only to chiral multiplets, stability criteria can be formulated in terms of sectional curvatures of the underlying Kähler manifold that are quite involved, so that their general solution has not been worked out to date.

In $\mathcal{N} = 2$ supergravity stable de Sitter vacua have been obtained, until very recently, only in a unique class of models [26] (later generalized in [27]) and, as stressed there, the mechanism that generates a scalar potential with the desired properties results from three equally essential ingredients:

1. The gauging of a non-compact, non-abelian group that in the models that were considered is $\text{so}(2,1)$.

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3 For a recent construction of meta-stable de Sitter vacua in abelian gaugings of $\mathcal{N} = 2$ supergravity, see [28].
2. The introduction of Fayet–Iliopoulos terms corresponding to the gauging of compact $u(1)$ factors.
3. The introduction of a Wagemans–de Roo angle that within special Kähler geometry rotates the directions associated to the non-compact gauge group with respect to those associated with the compact one.

The class of models constructed in [26] relies on the coupling of vector multiplets to supergravity as dictated by the special Kähler manifold

$$S\mathcal{K}_n = S\mathcal{T}[2, n] \equiv \frac{SU(1, 1)}{U(1)} \times \frac{SO(2, n)}{SO(2) \times SO(n)},$$

which accommodates the scalar fields and governs the entire structure of the Lagrangian. There are two interesting special cases: for $n = 1$ one obtains the ST model, which describes two vector multiplets, while for $n = 2$ one obtains the STU-model, which constitutes the core of most supergravity theories and is thus ubiquitous in the study of string compactifications at low energies. In this case, due to accidental Lie algebra automorphisms, the scalar manifold factorizes, since

$$S\mathcal{T}[2, 2] \equiv \frac{SU(1, 1)}{U(1)} \times \frac{SU(1, 1)}{U(1)} \times \frac{SU(1, 1)}{U(1)}.$$  

Starting from the Lagrangian of ungauged $\mathcal{N} = 2$ supergravity based on this special Kähler geometry, the scalar potential is generated gauging a subgroup $G_{\text{gauge}} \subset SU(1, 1) \times SO(2, n)$. The three models explicitly constructed in [26], whose scalar potential admits stable de Sitter extrema, are

- The STU model with 3 vector multiplets, in the manifold $S\mathcal{T}[2, 2]$, which, together with the graviphoton, are gauging $SO(2, 1) \times U(1)$, with a Fayet–Iliopoulos term for the $U(1)$ factor.
- A model with 5 vector multiplets, in the manifold $S\mathcal{T}[2, 4]$, which, together with the graviphoton, are gauging $SO(2, 1) \times SO(3)$, with a Fayet–Iliopoulos term for the $SO(3)$; and
- The last model extended with 2 hypermultiplets with 8 real scalars in the coset $\frac{SO(4, 2)}{SO(4) \times SO(2)}$.

The choice of the hypermultiplet sector for the third model is possible since the coset $\frac{SO(4, 2)}{SO(4) \times SO(2)}$ can be viewed as a factor in the special Kähler manifold $S\mathcal{T}[2, 4]$, or alternatively as a quaternionic-Kähler manifold by itself. The scalar potentials of the three models are qualitatively very similar, while the key ingredient behind the emergence of de Sitter extrema is the introduction of a non-trivial Wagemans–de Roo angle. For this reason we shall analyze only the first and simplest of these three models.

The explicit form of the scalar potential obtained in this gauging can be illustrated by introducing a parametrization of the scalar sector according to Special Geometry, and symplectic sections are the main ingredient to this effect. In the notation of [29], the holomorphic section reads

$$\Omega = \begin{pmatrix} X^A \\ F_\Sigma \end{pmatrix},$$

where

$$X^A(S, y) = \begin{pmatrix} \frac{1}{2}(1 + y^2) \\ \frac{1}{2}i(1 - y^2) \\ y^a \end{pmatrix}, \quad a = 1, \ldots, n,$$
\[ F_A(S, y) = \begin{pmatrix} \frac{1}{2}S(1 + y^2) \\ \frac{1}{2}iS(1 - y^2) \\ -Sy^a \end{pmatrix}, \quad y^2 = \sum_{a=1}^{n} (y^a)^2. \] (3.4)

The complex \( y^a \) fields are Calabi–Vesentini coordinates for the homogeneous manifold \( \text{SO}(2n) / \text{SO}(2) \times \text{SO}(n) \), while the complex field \( S \) parameterizes the homogeneous space \( \text{SU}(1, 1) / \text{U}(1) \), which is identified with the complex upper half-plane. With these conventions, the positivity domain of our Lagrangian is

\[ \text{Im} \ S > 0. \] (3.5)

The Kähler potential is by definition

\[ K = -\log(-i\langle \Omega | \Omega \rangle) = -\log[-i(\overline{X}^A F_A - \overline{F}_X X^X)], \] (3.6)

so that in this example the Kähler potential and the Kähler metric read

\[ K = K_1 + K_2, \quad K_1 = -\log[-i(S - \overline{S})], \quad K_2 = -\log\left[ \frac{1}{2} (1 - 2\overline{y}^a y^a + |y^a y^a|^2) \right], \]

\[ g_{S\overline{S}} = \frac{1}{(2 \text{Im} S)^2}, \quad g_{a\overline{b}} = \frac{\partial}{\partial y^a} \frac{\partial}{\partial \overline{y}^b} K_2. \] (3.7)

### 3.1. de Roo–Wagemans angles

As we have stressed, in the construction of [26], the de Roo–Wagemans angles are essential ingredients for the existence of de Sitter extrema. They were originally introduced [30,31] in \( \mathcal{N} = 4 \) supergravity with semisimple gaugings to characterize the relative embeddings of each simple factor \( G_k \) of the gauge group inside \( \text{Sp}(2(n + 2), \mathbb{R}) \), performing a symplectic rotation on the holomorphic section of the manifold prior to gauging. Different choices of the angles yield inequivalent gauged models with different properties. For \( n = 2 \), with \( \text{SO}(2, 1) \times \text{U}(1) \) gauging, there is just one de Roo–Wagemans’ angle and the corresponding rotation matrix reads

\[ R = \begin{pmatrix} A & B \\ -B & A \end{pmatrix}, \] (3.8)

where

\[ A = \begin{pmatrix} 1_{3 \times 3} & 0 \\ 0 & \cos(\theta) \end{pmatrix}, \quad B = \begin{pmatrix} 0_{3 \times 3} & 0 \\ 0 & \sin(\theta) \end{pmatrix}. \] (3.9)

The symplectic section is rotated as

\[ \Omega \rightarrow \Omega_R \equiv R \cdot \Omega, \] (3.10)

while the Kähler potential is clearly left invariant by the transformation. The de Roo–Wagemans’ angle appears explicitly in the scalar potential, which is determined by the symplectic section \( \Omega_R \) and by

\[ V_R \equiv \exp[K] \Omega_R \] (3.11)

and reads [26]
\[ V_{SO(2,1) \times U(1)} = V_2 + V_1 = \frac{1}{2 \text{Im } S} \left( e_1^2 \left| \cos \theta - S \sin \theta \right|^2 + e_0^2 \frac{P_2^+(y)}{P_2^-(y)} \right), \]  

(3.12)

where \( P_2^\pm(y) \) are polynomial functions in the Calabi–Vesentini variables of holomorphic degree specified by their index,

\[ P_2^-(y) = 1 - 2y_0\bar{y}_0 \pm 2y_1\bar{y}_1 + y^2\bar{y}^2, \]  

(3.13)

while \( e_{0,1} \) are the coupling constants for the \( \mathfrak{so}(2,1) \) and \( \mathfrak{u}(1) \) gauge algebras.

In order to study the properties of this potential one has to perform a coordinate transformation from the Calabi–Vesentini coordinates to the standard ones that provide a \textit{solvable parametrization} of the three Lobachevsky–Poincaré planes displayed in Eq. (3.2). With some care such a transformation can be worked out and reads

\begin{align*}
    y_1 &= -\frac{i(ib_1(b_2 + e^{h_2}) + e^{h_1+h_2} + i e^{h_1} b_2 - 1)}{(ib_1 + e^{h_1} + 1)(b_2 + e^{h_2} + 1)}, \\
    y_2 &= \frac{ib_1 + e^{h_1} - e^{h_2} - ib_2}{(ib_1 + e^{h_1} + 1)(b_2 + e^{h_2} + 1)}, \\
    S &= ie^h + b.
\end{align*}

(3.14)

After this coordinate change, the complete Kähler potential becomes

\[ K = -\log \left( \frac{16b_0 e^{h_1+h_2}}{(1 + e^{h_1})^2 + b_1^2)((1 + e^{h_2})^2 + b_2^2)} \right) \]  

(3.15)

so that the Kähler metric is

\[ ds_K^2 = \frac{1}{4} \left( e^{-2h} db^2 + dh^2 + e^{-2h_1} db_1^2 + e^{-2h_2} db_2^2 + dh_1^2 + dh_2^2 \right), \]  

(3.16)

while in the new coordinates the scalar potential takes the form

\[ V = -\frac{1}{8} e^{-h-h_1-h_2} \left[ 2e^{h_1+h_2} (-b^2 + 2 \sin(2\theta) b - e^{2h} + (b_1^2 + e^{2h} - 1) \cos(2\theta) - 1) e_1^2 \right. \]

\[ \left. -((e^{h_1} + e^{h_2})^2 + b_1^2 + b_2^2 - 2b_1 b_2) e_0^2 \right]. \]  

(3.17)

Let us now turn to exploring consistent truncation patterns to one-field models with standard kinetic terms for the residual scalars. To this effect, one can verify that the constant values

\[ \{b, b_1, b_2\} \Rightarrow \tilde{b}_0 = \left\{ \frac{-\sin(2\theta)}{\cos(2\theta) - 1}, \kappa, \kappa \right\} \]  

(3.18)

result in the vanishing of the derivatives of the potential with respect to the three axions, identically in the remaining fields, so that one can safely introduce these values (3.18) in the potential to arrive at the reduced form

\[ V = V_{\tilde{b} = \tilde{b}_0} = V_{\tilde{b} = \tilde{b}_0} = \frac{1}{4} e^{-h} e_0^2 + \frac{1}{8} e^{-h-h_1-h_2} e_0^2 + \frac{1}{8} e^{-h-h_1+h_2} e_0^2 + \frac{1}{2} e^h \sin^2(\theta) e_1^2. \]  

(3.19)

The last step of the reduction is performed setting the two fields \( h_{1,2} \) to a common constant value:

\[ h_{1,2} = \ell. \]  

(3.20)

Indeed, it can be simply verified that for these values the derivatives of \( \tilde{V} \) with respect to \( h_{1,2} \) vanish identically. Finally, redefining the field \( h \) by means of the constant shift
The one-field potential becomes

\[ V(h) = \sin(\theta)e_0e_1 \cosh \mu. \]  

(3.22)

This information suffices to determine the corresponding dynamical system. We start from the general form of the $\mathcal{N} = 2$ supergravity action truncated to the scalar sector which is the following:

\[ S^\mathcal{N}=2 = \int d^4x L^{\mathcal{N}=2}_{\text{SUGRA}}, \]

\[ L^{\mathcal{N}=2}_{\text{SUGRA}} = \sqrt{-g}\left[ R[g] + 2g^{SK}_{ij}\partial_\mu z^i \partial^{\mu} \bar{z}^j - 2V(z, \bar{z}) \right]. \]  

(3.23)

where $g^{SK}_{ij}$ is the special Kähler metric of the target manifold and $V(z, \bar{z})$ is the potential that we have been discussing. Reduced to the residual dynamical field content, after fixing the other fields to their extremal values, the above action becomes:

\[ S^\mathcal{N}=2 = \int d^4x \sqrt{-g}\left( R[g] + \frac{1}{2} \partial_\mu h \partial^{\mu} h - \mu^2 \cosh h \right), \]  

(3.24)

where we have redefined $\mu^2 = 2\bar{\mu}^2$. Hence the effective one-field dynamical system is described by the following Lagrangian

\[ L_{\text{eff}} = \exp[3A - B]\left( \frac{1}{2} \dot{h}^2 - \frac{3}{2} \dot{A}^2 - \exp[2B] \frac{\mu^2 \cosh h}{V(h)} \right) \]  

(3.25)

which agrees with the general form (2.1), introduced above.

In light of this, the effective dynamical model of the gauged STU model would be integrable if the potential

\[ V(\varphi) = 3\mu^2 \cosh\left( \frac{1}{\sqrt{3}} \varphi \right) \]  

(3.26)

could be identified with any of the integrable potentials listed in Tables 1 and 2. We show below that this is not the case. Nonetheless, the results of [1] provide a qualitative information on the behavior of the solutions of this supergravity model. As a special case, one can simply retrieve the de Sitter vacuum from this formulation in terms of a dynamical system. Choosing the gauge $B = 0$, the field equations associated with the Lagrangian (3.25) are solved by setting $h = 0$, which corresponds to the extremum of the potential, and

\[ A(t) = \exp[H_0t], \quad H_0 = \sqrt{\frac{2}{3}} \mu^2 = \sqrt{\frac{4}{3}} \frac{\sin(\theta)e_0e_1}{e_1}, \]  

(3.27)

which corresponds to the eternal exponential expansion of de Sitter space. This solution is an attractor for all the other solutions as shown in [1].
\[ \mathcal{L}_{\text{cosh}} = \exp[3A - B]\left(\frac{q}{2}h^2 - \frac{3}{2}A^2 - \exp[2B]\mu^2 \cosh[\mu h]\right) \] (3.28)

that depends on two parameters \( q \) and \( p \). Comparing with the list of integrable models one can see that there are just two integrable cases corresponding to the choices

\[ \frac{p}{\sqrt{3q}} = 1, \quad \frac{p}{\sqrt{3q}} = \frac{2}{3}. \] (3.29)

The first case \( \frac{p}{\sqrt{3q}} = 1 \), corresponding to the potential \( V(\varphi) \sim \cosh[\varphi] \) can be mapped into three different integrable series among those displayed in Table 1. The first embedding is into the series \( I_1 \) by choosing \( C_{11} = C_{22} \neq 0 \) and \( C_{12} = 0 \). The second embedding is into series \( I_2 \), by choosing \( \gamma = \frac{1}{2} \) and \( C_1 = C_2 \). The third embedding is into model \( I_7 \), by choosing once again \( \gamma = \frac{1}{2} \) and \( C_1 = C_2 \). The second case \( \frac{p}{\sqrt{3q}} = \frac{2}{3} \) corresponding to the potential \( V(\varphi) \sim \cosh[\frac{2}{3}\varphi] \) can be mapped into series \( I_2 \) of Table 1, by choosing \( \gamma = -\frac{1}{3} \) and \( C_1 = C_2 \). It can also be mapped into the series \( I_7 \) by choosing \( \gamma = \frac{1}{3} \) and \( C_1 = -C_2 \). Unfortunately, none of these solutions correspond to the Lagrangian (3.25), where

\[ p = 1, \quad q = 1 \] (3.30)

so that the one-field cosmology emerging from the non-compact non-abelian \( \mathfrak{so}(2, 1) \) gauging of the STU model is indeed not integrable! This analysis emphasizes that embedding an integrable model into the gauging of an extended supergravity theory is a difficult task.

In Section 6 we will consider in more detail the cosh-model defined by Eq. (3.28). There we will show that it can be reduced to a normal form depending only on one parameter that we name the index:

\[ \omega = \frac{p}{\sqrt{3q}} \] (3.31)

and we will compare its behavior for various values of the index \( \omega \). The two integrable cases mentioned above correspond respectively to the following critical indices,

\[ \omega^f_c = \sqrt{3}, \quad \omega^n_c = \frac{2}{\sqrt{3}}. \] (3.32)

The first critical index has been denoted with the superscript \( f \) since, in the language of [1] the fixed point of the corresponding dynamical system is of focus type. Similarly, the second critical index has been given the superscript \( n \) since the fixed point of the corresponding dynamical system is of the node type. In these two cases we are able to integrate the field equations explicitly. For the other values of \( \omega \) we are confined to numerical integration. Such a numerical study reveals that when the initial conditions are identical, the solutions of the non-integrable models have a behavior very similar to that of the exact solutions of the integrable model, as long as the type of fixed point defined by the extremum of the potential is the same. Hence the behavior of the one-field cosmology emerging from the \( \mathfrak{so}(2, 1) \) gauging of the STU model can be approximated by the exact analytic solutions of the cosh-model with index \( \omega^n_c \).

It remains a fact that the value of \( \omega \) selected by the Gauged Supergravity model is \( \omega = 1 \) rather than the integrable one, a conclusion that will be reinforced by a study of Fayet–Iliopoulos gaugings in the \( S^3 \) model [33].

Considering instead the integrable series \( I_2 \) of Table 1, we will show in Section 5 that there is just one case there that can be fitted into a Gauged Supergravity model. It corresponds to the
Fig. 1. The numerical solution shows that the de Sitter solution is indeed an attractor.

value $\gamma = \frac{2}{3}$, which can be realized in $\mathcal{N} = 1$ supergravity by an acceptable and well defined superpotential. After a wide inspection, this seems to be one of the very few integrable supersymmetric models so far available. A second one will be identified in Section 7. As we shall emphasize, the superpotential underlying both instances of supersymmetric integrable models is strictly $\mathcal{N} = 1$ and does not arise from a Fayet–Iliopoulos gauging of a corresponding $\mathcal{N} = 2$ model.

3.2. Behavior of the solutions in the $\mathcal{N} = 2$ STU model with so(1,2)-gauging

Although the $\mathcal{N} = 2$ model that we have been considering is not integrable, its Friedman equations can be integrated numerically providing a qualitative understanding of the nature of the solutions.

In Fig. 1 we show the behavior of both the scale factor and the scalar field for any regular initial conditions. The plot clearly shows that the de Sitter solution corresponding to an indefinite exponential expansion is an attractor as predicted by the fixed point analysis of the differential system. Indeed the numerical integration reveals a slow-roll phase that works in reversed order with respect to the standard inflationary scenario [21]. When the scalar field is high up and descends rapidly, the expansion of the Universe proceeds rather slowly, then the scalar field reaches the bottom of the potential and rolls slowly toward its minimum, while the Universe expands exponentially becoming asymptotically de Sitter.

3.3. A more systematic approach: The embedding tensor formalism

In the previous section we have reviewed and analyzed, from a cosmological perspective, a special class of $\mathcal{N} = 2$ models which exhibit stable de Sitter vacua. A complete analysis of one-field cosmological models emerging from $\mathcal{N} = 2$ supergravities (or even $\mathcal{N} > 2$ theories) is a considerably more ambitious project, which requires a systematic study of the possible gaugings of extended supergravities. A precious tool in this respect is the embedding tensor formulation of gauged extended supergravities [34]. This approach consists in writing the gauged theory as a deformation of an ungauged one (with the same field content and supersymmetry) in which the additional terms in the Lagrangian (minimal couplings, fermion mass terms and scalar potential) and in the supersymmetry transformation laws, which are needed in order to make the theory locally invariant with respect to the chosen gauge group $G$ while keeping the original
supersymmetry unbroken, are all expressed in terms of a single matrix of coupling constants (the embedding tensor) which can be described as a covariant tensor with respect to the global symmetry group $G$ of the original ungauged model. If we denote by $\{t_\alpha\}$ the generators of the Lie algebra $\mathfrak{g}$ of $G$ and by $X_A$ the gauge generators, to be gauged by the vector fields $A^A_\mu$ of the model, since the gauge group must be contained in $G$, $X_A$ must be a linear combination of the $t_\alpha$:

$$X_A = \Theta_A^\alpha t_\alpha.$$  \hfill (3.33)

The matrix $\Theta_A^\alpha$ is the embedding tensor and defines all the information about the embedding of the gauge algebra inside $\mathfrak{g}$. A formulation of the gauging which is independent of the symplectic frame of the original ungauged theory, was given in [35] and extends the definition of the embedding tensor by including, besides the electric components defined above, also magnetic ones:

$$\Theta_M^\alpha = \{\Theta_A^\alpha, \Theta^A_\alpha\}.$$  \hfill (3.34)

The index $M$ is now associated with the symplectic duality-representation $W$ of $G$ in which the electric field strengths and their magnetic duals transform, so that $\Theta_M^\alpha$ formally belongs to the product $W \otimes \text{Adj}(G)$, namely is a $G$-covariant tensor. Since all the deformations of the original ungauged action, implied by the gauging procedure, are written in terms of $\Theta_M^\alpha$ in a $G$-covariant way, the gauged equations of motion and the Bianchi identities formally retain the original global $G$-invariance, provided $\Theta_M^\alpha$ is transformed as well. Since, however, the action of $G$, at the gauged level, affects the coupling constants of the theory, encoded in the embedding tensor, it should be viewed as an equivalence between theories rather than a symmetry, and gauged models whose embedding tensors are related by $G$-transformations, share the same physics. Thus gauged extended supergravities obtained from the same ungauged model, can be classified in universality classes defined by the orbits of the embedding tensor under the action of $G$. Classifying such classes is a rather non-trivial task. In simple models like the STU one, this can be done thoroughly. In the following we perform this analysis and analyze the possible one-field cosmological models for each class, leaving its extension to more general $\mathcal{N} = 2$ gauged models to a future investigation [33].

To set the stage, let us consider an $\mathcal{N} = 2$ theory with $n_V$ vector fields and a global symmetry group of the form:

$$G = \text{USK} \times \text{GQK},$$  \hfill (3.35)

where $\text{USK}, \text{GQK}$ are the isometry groups of the Special Kähler and Quaternionic Kähler manifolds (in the absence of hypermultiplets $\text{GQK} = \text{SO}(3)$). Let $\mathfrak{g}, \mathfrak{g}_{\text{SK}}, \mathfrak{g}_{\text{QK}}$ denote the Lie algebras of $G, \text{USK}, \text{GQK}$ and $\{t_\alpha\}, \{t_A\}, \{t_a\}, \alpha = 1, \ldots, \text{dim}(G), A = 1, \ldots, \text{dim}(\text{USK}), a = 1, \ldots, \text{dim}(\text{GQK})$, a set of corresponding bases. Only the group $\text{USK}$ has a symplectic duality action on the $2n_V$-dimensional vector $F^M_{\mu\nu}, M = 1, \ldots, 2n_V$, consisting of the electric field strengths and their duals

$$F^M_{\mu\nu} \equiv \begin{pmatrix} F^A_{\mu\nu} \\ G^A_{\mu\nu} \end{pmatrix}$$  \hfill (3.36)

namely:

$$\forall u \in \text{USK}: F^M_{\mu\nu} \rightarrow F'^M_{\mu\nu} = u^M_N F^N_{\mu\nu}.$$  \hfill (3.37)

We have denoted by $W$ the corresponding $2n_V$-dimensional, symplectic representation of $\text{USK}$. 
For reader’s convenience we summarize the index conventions in Table 4.

The embedding tensor has the general form:

\[
\{ \Theta^{\alpha}_M \} = \{ \Theta^{A}_M, \Theta^{a}_M \},
\]

and defines the embedding of the gauge algebra \( g_{\text{gauge}} = \{ X_M \} \) inside \( g \):

\[
X_M = \Theta^{A}_M t_A + \Theta^{a}_M t_a.
\]

In the absence of hypermultiplets, the components \( \Theta^{a}_M, a = 1, 2, 3 \), running over the adjoint representation of the global symmetry \( \text{SO}(3) \), are the Fayet–Iliopoulos terms. The generators \( t_A \) of \( g_{\text{SK}} \) have a non-trivial \( W \)-representation: \( t_A = (t_{AMN})^{N} \) while the generators \( t_a \) do not. Thus we can define the following tensor:

\[
X^{MN} = \Theta^{A}_M t_{AN}^{P}, \quad X_{MN}^P = X_{MN}^{Q} C^{QP},
\]

where \( C \) is the \( 2n_V \times 2n_V \) skew-symmetric, invariant \( \text{Sp}(2n_V, \mathbb{R}) \)-matrix.

Gauge-invariance and supersymmetry of the action impose linear and quadratic constraints on \( \Theta \):

- The linear constraints are:
  \[
  X^{MN} = 0, \quad X_{(MN)P} = 0.
  \]

- The quadratic constraints originate from the condition that \( X_M \) close an algebra inside \( g \) with structure constants given in terms of \( X^{MN} \), and from the condition that the symplectic vectors \( \Theta^{a}_M \), labeled by \( \alpha \), be mutually local:
  \[
  [X_M, X_N] = -X^{MN} P X_P, \quad C^{MN} \Theta^{a}_M \Theta^{b}_N = 0 \Leftrightarrow \Theta^{A\alpha} \Theta^{A\beta} = 0.
  \]

the former can be rewritten as the following set of two equations:

\[
\Theta^{A\alpha} \Theta^{B\beta} f_{AB}^{C} + \Theta^{A\alpha} t_{AN}^{P} \Theta^{P\beta} = 0,
\]

where \( f_{AB}^{C}, f_{ab}^{c} \) are the structure constants of \( g_{\text{SK}} \) and \( g_{\text{QK}} \), respectively. It can be shown that Eqs. (3.41) and (3.44) imply \( \Theta^{A\alpha} \Theta^{A\beta} = 0 \), which is the part of (3.43) corresponding to \( \alpha = A, \beta = B \).

Let us now denote by \( k^{i}_A, k^{i}_a \) the Killing vectors on the Special Kähler manifold corresponding to the isometry generator \( t_A \), \( k^{i}_a \) the Killing vectors on the Quaternionic Kähler manifold corresponding to the isometry generator \( t_a \), and \( P^{i}_a, x = 1, 2, 3 \), the corresponding momentum map.
(note that these quantities are, by definition, associated only with the geometry of the scalar manifold and thus independent of the gauging). The scalar potential can be written in the following way [29]

\[
\mathcal{V} = \mathcal{V}_{\text{hyperino}} + \mathcal{V}_{\text{gaugino}, 1} + \mathcal{V}_{\text{gaugino}, 2} + \mathcal{V}_{\text{gravitino}},
\]

\[
\mathcal{V}_{\text{hyperino}} = 4 V^M V^N \Theta_M^a \Theta_N^b k^i_a k^j_b g_{i j},
\]

\[
\mathcal{V}_{\text{gaugino}, 1} = V^M V^N \Theta_M^A \Theta_N^B \bar{k}^i_A \bar{k}^j_B g_{i j},
\]

\[
\mathcal{V}_{\text{gaugino}, 2} = g^{i j} D_i V^M D_j V^N \Theta_M^a \Theta_N^b p^x_a p^x_b,
\]

\[
\mathcal{V}_{\text{gravitino}} = -3 V^M V^N \Theta_M^a \Theta_N^b p^x_a p^x_b,
\]

where \( V^M \) is the covariantly holomorphic symplectic section of the Special Kähler manifold under consideration: \( (V^M) = (L^A, M_A) \).

### 3.4. Scan of the STU model gaugings and their duality orbits

Consider the STU model with no hypermultiplets. This corresponds to the sixth item in Table 3 for \( p = 0 \). The global symmetry group is \( G = \text{SL}(2, \mathbb{R})^3 \times \text{SO}(3) \), the latter factor being the form of \( G_{\text{OK}} \) in the absence of hypermultiplets. It is relevant to our discussion only in the case we want to add FI terms, i.e. when we introduce non-vanishing components \( \Theta_M^a, a = 1, 2, 3 \).

The symplectic \( \mathbf{W} \)-representation of the electric–magnetic charges is the \((\frac{1}{2}, \frac{1}{2}, \frac{1}{2})\) of \( \text{USK} = \text{SL}(2, \mathbb{R})^3 \). Let us use the indices \( i, j, k = 1, 2 \) to label the fundamental representation of \( \text{SL}(2, \mathbb{R}) \). As \( \mathfrak{sl}(2) \)-generators in this spinor representation of \( \mathfrak{so}(2, 1) \sim \mathfrak{sl}(2) \), we make the following choice: \( \{s_x\} = \{\sigma_1, i\sigma_2, \sigma_3\} \), \( \sigma_x \) being the Pauli matrices. The index \( M \) can be written as \( M = (i_1, i_2, i_3) \) and the embedding tensor \( \Theta_M^A \) takes the following form:

\[
\Theta_M^A = \left\{ \Theta_{(i_1, i_2, i_3)}^{x_1}, \Theta_{(i_1, i_2, i_3)}^{x_2}, \Theta_{(i_1, i_2, i_3)}^{x_3} \right\}
\]

\[
\in \left( \frac{1}{2}, \frac{1}{2}, \frac{1}{2} \right) \times \left[ (1, 0, 0) + (0, 1, 0) + (0, 0, 1) \right],
\]

where \( x_i \) run over the adjoint (vector)-representations of the three \( \mathfrak{sl}(2) \) algebras. Since:

\[
\left( \frac{1}{2}, \frac{1}{2}, \frac{1}{2} \right) \times \left[ (1, 0, 0) + (0, 1, 0) + (0, 0, 1) \right]
\]

\[
= 3 \times \left( \frac{1}{2}, \frac{1}{2}, \frac{1}{2} \right) + \left( \frac{3}{2}, \frac{1}{2}, \frac{1}{2} \right) + \left( \frac{1}{2}, \frac{3}{2}, \frac{1}{2} \right) + \left( \frac{1}{2}, \frac{1}{2}, \frac{3}{2} \right),
\]

each component of the embedding tensor can be split into its \((\frac{1}{2}, \frac{1}{2}, \frac{1}{2})\) and \((\frac{3}{2}, \frac{1}{2}, \frac{1}{2})\) irreducible parts as follows:

\[
\Theta_{(i_1, i_2, i_3)}^{x_1} = (s^{x_1})_{i_1} j \xi_{i_1 j}^{(1)} + \mathfrak{S}_{i_1, i_2, i_3}^{x_1},
\]

\[
\Theta_{(i_1, i_2, i_3)}^{x_2} = (s^{x_2})_{i_2} j \xi_{i_2 j}^{(1)} + \mathfrak{S}_{i_1, i_2, i_3}^{x_2},
\]

\[
\Theta_{(i_1, i_2, i_3)}^{x_3} = (s^{x_3})_{i_3} j \xi_{i_3 j}^{(1)} + \mathfrak{S}_{i_1, i_2, i_3}^{x_3}.
\]

The irreducible \((\frac{3}{2}, \frac{1}{2}, \frac{1}{2})\) tensors \( \mathfrak{S}_{i_1, i_2, i_3}^{x_i} \) are defined by the vanishing of the appropriate gamma-trace, namely:

\[
\mathfrak{S}_{i_1, i_2, i_3}^{x_i} (s^{x_1})_{i_1} j = \mathfrak{S}_{i_1, i_2, i_3}^{x_2} (s^{x_2})_{i_2} j = \mathfrak{S}_{i_1, i_2, i_3}^{x_3} (s^{x_3})_{i_3} j = 0.
\]
Let us now define embedded gauge generators $X_{MN}^P$:

$$X_{(i_1,i_2,i_3),(j_1,j_2,j_3)}^{(k_1,k_2,k_3)} = \Theta_{(i_1,i_2,i_3)}^{x_1} (s_{x_1})_{i_1}^{j_1} \delta_{j_2}^{k_2} \delta_{j_3}^{k_3} + \Theta_{(i_1,i_2,i_3)}^{x_2} (s_{x_2})_{j_2}^{k_2} \delta_{j_1}^{k_1} \delta_{j_3}^{k_3} + \Theta_{(i_1,i_2,i_3)}^{x_3} (s_{x_3})_{j_3}^{k_3} \delta_{j_1}^{k_1} \delta_{j_2}^{k_2},$$

(3.51)

The linear constraints (3.41) become:

$$X_{(i_1,i_2,i_3),(j_1,j_2,j_3)}^{(i_1,i_2,i_3)} = 0 \Rightarrow \xi_{i_1 i_2 i_3}^{(1)} + \xi_{i_1 i_2 i_3}^{(2)} + \xi_{i_1 i_2 i_3}^{(3)} = 0,$$

$$X_{(i_1,i_2,i_3),(j_1,j_2,j_3),(k_1,k_2,k_3)} + X_{(j_1,j_2,j_3),(i_1,i_2,i_3),(k_1,k_2,k_3)} + X_{(i_1,i_2,i_3),(j_1,j_2,j_3),(k_1,k_2,k_3)} = 0$$

$$\Rightarrow \mathcal{E}_{i_1,i_2,i_3} x_i = 0.$$  

(3.52)

This corresponds to the elimination of the $(\frac{3}{2}, \frac{1}{2}, \frac{1}{2})$ representation leaving us only with three tensors in the $(\frac{1}{2}, \frac{1}{2}, \frac{1}{2})$ representation. Explicitly the linearly constrained embedding tensor reads as follows:

$$\Theta_{(i_1,i_2,i_3)}^A = \{(s_{x_1})_{i_1}^{j_1} \xi_{i_1 i_2 i_3}^{(1)} (s_{x_2})_{i_2}^{j_2} \xi_{i_1 i_2 i_3}^{(2)} (s_{x_3})_{i_3}^{j_3} \xi_{i_1 i_2 i_3}^{(3)}\},$$

(3.53)

and the additional linear constraint reduces further the independent tensors in the $(\frac{1}{2}, \frac{1}{2}, \frac{1}{2})$ to two, since we get condition

$$\xi_{i_1 i_2 i_3}^{(1)} + \xi_{i_1 i_2 i_3}^{(2)} + \xi_{i_1 i_2 i_3}^{(3)} = 0.$$  

(3.54)

The quadratic condition for $\Theta_M^A$ has the form (3.43), which applied to our solution of the linear constraint takes the following appearance:

$$e^{ij_1} e^{j_2} e^{i_3} \Theta_{(i_1,i_2,i_3)}^A \Theta_{(j_1,j_2,j_3)}^B = 0.$$  

(3.55)

By means of a MATHEMATICA computer code we were able to find 36 solutions to this equation, all of which correspond to non-semisimple gauge groups. We do not display them here, since, in Section 3.4.2 we show how to classify the orbits into which such solutions are organized and it will be sufficient to consider only one representative for each orbit.

### 3.4.1. The Special Geometry of the STU model

In Eq. (3.14) we derived the transformation from the Calabi–Vesentini coordinates {$S, y_{1,2}$} to a triplet of complex coordinates $z_{1,2,3}$ parameterizing the three identical copies of the coset manifold $\text{SU}(2,\mathbb{R})/\text{SO}(2)$ which compose this special instance of special Kähler manifold. Indeed setting:

$$i e^{b} + b \equiv z_1, \quad i e^{b_1} + b_1 \equiv z_2, \quad i e^{b_2} + b_2 \equiv z_3$$

(3.56)

the transformation (3.14) can be rewritten as follows:

$$S = z_1, \quad y_1 = \frac{i(z_2z_3 + 1)}{(z_2 + i)(z_3 + i)}, \quad y_2 = \frac{i(z_2 - z_3)}{(z_2 + i)(z_3 + i)}.$$  

(3.57)

In the sequel we will adopt the symmetric renaming of variables

$$z_i = i e^{b_i} + b_i.$$  

(3.58)
Applying the transformation (3.57) of its arguments to the Calabi–Vesentini holomorphic section $\Omega_{CV}$, we find:

$$
\Omega_{CV}(z) = f(z) \begin{pmatrix}
\frac{z_1}{\sqrt{2}} + \frac{z_2}{\sqrt{2}} \\
\frac{z_2 z_3}{\sqrt{2}} - \frac{1}{\sqrt{2}} \\
-\frac{z_2 z_3}{\sqrt{2}} - \frac{1}{\sqrt{2}} \\
\frac{z_1}{\sqrt{2}} - \frac{z_3}{\sqrt{2}} \\
\frac{z_1 z_2}{\sqrt{2}} + \frac{z_3}{\sqrt{2}} \\
\frac{z_1 z_2 z_3}{\sqrt{2}} - \frac{1}{\sqrt{2}} \\
-\frac{z_1 z_2 z_3}{\sqrt{2}} + \frac{1}{\sqrt{2}} \\
\frac{z_1 z_2 z_3}{\sqrt{2}} - \frac{z_1}{\sqrt{2}} \\
\frac{z_1 z_2}{\sqrt{2}} - \frac{z_3}{\sqrt{2}} \\
\frac{z_1}{\sqrt{2}} - \frac{z_2}{\sqrt{2}}
\end{pmatrix},
$$

(3.59)

$$
f(z) = \frac{i\sqrt{2}}{(z_2 + i)(z_3 + i)}.
$$

(3.60)

As it is well known the overall holomorphic factor $f(z)$ in front of the section has no consequences on the determination of the Kähler metric and simply it adds the real part of a holomorphic function to the Kähler potential. Similarly, at the level of ungauged supergravity, the symplectic frame plays no role on the Lagrangian and we are free to perform any desired symplectic rotation on the section, the preserved symplectic metric being the following one:

$$
\mathbb{C} = \begin{pmatrix}
0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 1 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 1 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 1 \\
-1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & -1 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & -1 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & -1 & 0 & 0 & 0 & 0
\end{pmatrix}.
$$

(3.61)

On the other hand at the level of gauge supergravity the choice of the symplectic frame is physically relevant. The Calabi–Vesentini frame is that one where the SO($2, 2$) isometries of the manifold are all linearly realized on the electric vector field strengths, while the SL($2, \mathbb{R}$) factor acts as a group of electric/magnetic duality transformations. For this reason the CV frame was chosen in paper [26], since in such a frame it was easy to single out the non-compact gauge group SO($2, 1$). On the other the so named special coordinate frame which admits a description in terms of a prepotential, is that one where the three group factors SL($2, \mathbb{R}$) are all on the same footing and the $\mathbf{W}$-representation is identified as the $(2, 2, 2) \sim (\frac{1}{2}, \frac{1}{2}, \frac{1}{2})$.

The philosophy underlying the embedding tensor approach to gaugings is that the embedding tensor already contains all possible symplectic frame choices, since it transforms as a good tensor under the symplectic group. Hence we can choose any preferred symplectic frame to start with.

In view of these considerations we introduce the following symplectic matrix:
\[ S = \begin{pmatrix} 0 & 0 & \frac{1}{\sqrt{2}} & \frac{1}{\sqrt{2}} & 0 & 0 & 0 & 0 \\ -\frac{1}{\sqrt{2}} & 0 & 0 & 0 & 0 & \frac{1}{\sqrt{2}} & 0 & 0 \\ -\frac{1}{\sqrt{2}} & 0 & 0 & 0 & \frac{1}{\sqrt{2}} & 0 & 0 & 0 \\ 0 & 0 & \frac{1}{\sqrt{2}} & -\frac{1}{\sqrt{2}} & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 & 0 & \frac{1}{\sqrt{2}} & \frac{1}{\sqrt{2}} & 0 \\ 0 & -\frac{1}{\sqrt{2}} & 0 & 0 & -\frac{1}{\sqrt{2}} & 0 & 0 & 0 \\ 0 & \frac{1}{\sqrt{2}} & 0 & 0 & -\frac{1}{\sqrt{2}} & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 & 0 & \frac{1}{\sqrt{2}} & -\frac{1}{\sqrt{2}} & 0 \end{pmatrix}, \] (3.62)

\[ C = S^T \odot S \] (3.63)

and we introduce the symplectic section in the special coordinate frame by setting

\[ \Omega_{SF}(z) = \left( \begin{array}{c} 1 \\ z_1 \\ z_2 \\ z_3 \\ -z_1 z_2 z_3 \\ z_2 z_3 \\ z_1 z_3 \\ z_1 z_2 \end{array} \right) = \left( \begin{array}{c} X^A(z) \\ F_A(z) \end{array} \right). \] (3.64)

Note that this frame admits a prepotential description. Introducing the following holomorphic prepotential:

\[ \mathcal{F}(z) = z_1^1 z_2^2 z_3^3 = stu. \] (3.65)

The symplectic section \( \Omega_{SF} \) can be written as follows:

\[ \Omega_{SF}(z) = \left( 1, z_m, -\mathcal{F}(z), \frac{\partial \mathcal{F}(z)}{\partial z_m} \right)_{m=1,2,3}. \] (3.66)

It is useful to work also with real fields \((\phi^r) = (b_m, h_m)\), defined as in Eq. (3.58). The Kähler potential is expressed as follows:

\[ \mathcal{K}(z^m, \bar{z}) = -\log(-i \Omega \odot \bar{\Omega}) = -\log[-i (X^A \bar{F}_A - F_A X^A)] \] (3.67)

and, in real coordinates we have:

\[ e^{-\mathcal{K}} = 8e^{b_1 + b_2 + b_3}. \] (3.68)

We also introduce the covariantly holomorphic symplectic section \( V = e^{\mathcal{K}/2} \Omega_{SF} \) satisfying the condition:

\[ \nabla_{\bar{a}} V \equiv \left( \partial_{\bar{a}} - \frac{1}{2} \partial_{\bar{a}} \mathcal{K} \right) V = 0, \] (3.69)

and its covariant derivatives:
\( U_m = (U_m^M) \equiv \nabla_m V = \left( \partial_m + \frac{1}{2} \partial_m K \right) V. \)  

(3.70)

The following properties hold:

\[
V \bar{C} \bar{V} = i, \quad U_m \bar{C} \bar{V} = i, \quad U_m \nabla \bar{V} = 0, \quad U_m \bar{U}_{\bar{n}} = -i \sigma_{m \bar{n}}.
\]  

(3.71)

If \( E_m^I, I = 1, \ldots, 3, \) is the complex vielbein matrix of the manifold, \( g_{m \bar{n}} = \sum_I E_m^I \bar{E}_{\bar{n}}^I, \) and \( E_I^m \) its inverse, we introduce the quantities \( U_I \equiv E_I^m U_m, \) in terms of which the following 8 \times 8 matrix \( \hat{L}_4 = (\hat{L}_4^M)_N \) is defined:

\[
\hat{L}_4(z, \bar{z}) = (V, \bar{U}_I, \bar{V}, U_I) \mathcal{C} = \sqrt{2} \begin{pmatrix} \text{Re}(V), \text{Re}(U_I), -i \text{Im}(V), \text{Im}(U_I) \end{pmatrix}, \quad (3.72)
\]

where \( \mathcal{C} \) is the Cayley matrix. By virtue of Eqs. (3.71), the matrix \( \hat{L}_4 \) is symplectic:

\[
\hat{L}_4^T \mathcal{C} \hat{L}_4 = \mathcal{C}.
\]

In order to find the coset representative \( L \) as an \( \text{Sp}(8, \mathbb{R}) \) matrix in the solvable gauge, and the symplectic representation of the isometry generators \( t_A \) in the special coordinate basis, we proceed as follows. We construct a symplectic matrix \( L \) which coincides with the identity at the origin where \( \phi_r \equiv 0 \Leftrightarrow h_m = b_m = 0: \)

\[
L(\phi_r) = \hat{L}_4(\phi_r) \hat{L}_4(\phi_r \equiv 0)^{-1}. \quad (3.73)
\]

The following property holds:

\[
V(\phi_r) = L(\phi_r) V(\phi_r \equiv 0). \quad (3.74)
\]

The matrix \( L \) is the coset representative in the solvable gauge. To show this we compute the following generators:

\[
\begin{align*}
\mathfrak{h}_m &= \frac{\partial L}{\partial h_m} \bigg|_{\phi_r = 0}, \quad \mathfrak{a}_m &= \frac{\partial L}{\partial b_m} \bigg|_{\phi_r = 0}.
\end{align*}
\]  

(3.75)

These generators close a solvable Lie algebra \( \text{Solv} \) which is the Borel subalgebra of \( \mathfrak{g}_{SK}. \) The above construction is general and applies to any symmetric Special Kähler manifold. In our case \( \text{Solv} = \text{Solv}^{(1)}_2 \oplus \text{Solv}^{(3)}_2, \) where

\[
\text{Solv}^{(m)}_2 = \{ \mathfrak{h}_m, \mathfrak{a}_m \}, \quad [\mathfrak{h}_m, \mathfrak{h}_n] = \delta_{mn} \mathfrak{a}_n, \quad \left[ \text{Solv}^{(m)}_2, \text{Solv}^{(n)}_2 \right] = 0. \quad (3.76)
\]

One can verify that

\[
\mathbb{L}(\mathfrak{h}_m, \mathfrak{b}_m) = \mathbb{L}_{\text{axion}}(\mathfrak{b}_m) \mathbb{L}_{\text{dilaton}}(\mathfrak{h}_m) = e^{\mathfrak{b}_m \mathfrak{a}_n} e^{\mathfrak{a}_n \mathfrak{h}_m}. \quad (3.77)
\]

Each \( \mathfrak{sl}(2) \) algebra is spanned by \( \{ \mathfrak{h}_m, \mathfrak{a}_m, \mathfrak{a}_m^T \}, \) \( [\mathfrak{a}_m, \mathfrak{a}_n^T] = 2 \delta_{mn} \mathfrak{h}_n. \) The explicit matrix representations of these generators is:

\[
\begin{align*}
\mathfrak{h}_1 &= \begin{pmatrix} -\frac{1}{2} & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & \frac{1}{2} & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & -\frac{1}{2} & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & -\frac{1}{2} & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & \frac{1}{2} & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & -\frac{1}{2} & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & \frac{1}{2} & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & \frac{1}{2} \end{pmatrix},
\end{align*}
\]

\[
\begin{align*}
\mathfrak{a}_1 &= \begin{pmatrix} 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & -1 \\
0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 \end{pmatrix}.
\end{align*}
\]
The axionic and dilatonic parts of the coset representative in (3.77) have the following matrix form:

\[
\begin{align*}
I_{\text{axion}}(b_m) &= e^{b_m a_m} = \\
&= \begin{pmatrix}
1 & 0 & 0 & 0 & 0 & 0 & 0 \\
b_1 & 1 & 0 & 0 & 0 & 0 & 0 \\
b_2 & 0 & 1 & 0 & 0 & 0 & 0 \\
b_3 & 0 & 0 & 1 & 0 & 0 & 0 \\
-b_1 b_2 b_3 & -b_2 b_3 & -b_1 b_3 & -b_1 b_2 & 1 & -b_1 & -b_2 & -b_3 \\
b_2 b_3 & 0 & b_3 & b_2 & 0 & 1 & 0 & 0 \\
b_1 b_3 & b_3 & 0 & b_1 & 0 & 0 & 1 & 0 \\
b_1 b_2 & b_2 & b_1 & 0 & 0 & 0 & 0 & 1
\end{pmatrix}, \\
I_{\text{dilaton}}(b_m) &= e^{b_m h_m} = \text{diag}(e^{-\frac{b_1}{2}}, e^{-\frac{b_2}{2}}, e^{-\frac{b_3}{2}}, e^{-\frac{b_1 - b_2}{2}}, e^{-\frac{b_1 - b_3}{2}}, e^{-\frac{b_2 - b_3}{2}}, e^{-\frac{b_1 - b_2 + b_3}{2}}, e^{-\frac{b_1 + b_2 - b_3}{2}}, e^{-\frac{b_2 - b_1 + b_3}{2}}, e^{-\frac{b_3 - b_1 - b_2}{2}}).
\end{align*}
\]

To make contact with the discussion about the embedding tensor provided in the previous section, we define the transformation from the basis \((i_1, i_2, i_3)\) and the special coordinate symplectic frame. We start with an ordering of the independent components of a vector \(W_{i_1 i_2 i_3}\), which defines a symplectic basis to be dubbed “old”:

\[
W^{\text{old}} = (W^{\text{old}}_M) = (W_{1,1,1}, W_{1,1,2}, W_{1,2,1}, W_{1,2,2}, W_{2,1,1}, W_{2,1,2}, W_{2,2,1}, W_{2,2,2}).
\]

The new special coordinate basis is related to the old one by an orthogonal transformation \(O\):

\[
W^{s,c}_M = O_M W^{\text{old}}_M, \quad O = \frac{1}{2\sqrt{2}} \begin{pmatrix}
1 & -1 & -1 & 1 & -1 & 1 & 1 & -1 \\
-1 & 1 & 1 & -1 & 1 & 1 & -1 & -1 \\
1 & -1 & -1 & 1 & 1 & -1 & -1 & -1 \\
-1 & -1 & 1 & 1 & 1 & 1 & 1 & 1 \\
1 & 1 & 1 & 1 & 1 & 1 & 1 & 1 \\
1 & 1 & -1 & -1 & -1 & -1 & -1 & -1 \\
-1 & -1 & -1 & -1 & -1 & -1 & -1 & -1 \\
1 & -1 & -1 & -1 & -1 & -1 & -1 & -1
\end{pmatrix}.
\]
The $\mathfrak{s}(2)^3$ generators $t_A$ in the old basis read:

\[
(t_{x_1})_{j_1,j_2,j_3} = (s_{x_1})_{j_1} \delta_{j_2}^{k_3} \delta_{j_3}^{k_1}, \quad (t_{x_2})_{j_1,j_2,j_3} = (s_{x_2})_{j_2} \delta_{j_1}^{k_1} \delta_{j_3}^{k_3},
\]

\[
(t_{x_3})_{j_1,j_2,j_3} = (s_{x_3})_{j_3} \delta_{j_1}^{k_3} \delta_{j_2}^{k_1}.
\] (3.82)

In the new basis their representation is deduced from their relation to the $\text{Solv}$ generators and their transpose:

\[
t_1 = 2h_m, \quad t_2 = a_m - a^T_m, \quad t_3 = -a_m - a^T_m.
\] (3.83)

The commutation relations among them read:

\[
[t_{x_m}, t_{x_n}] = -2\delta_{mn}\epsilon_{xyzt}t_{z_n},
\] (3.84)

where the adjoint index is raised with $\eta_{xy} = \text{diag}(+1,-1,+1)$.

3.4.1.1. The Killing vectors

A standard procedure in coset geometry allows to compute the Killing vectors $\{k_A\} = \{k_m, k^m\}$:

\[
k_{1m} = -2(\partial b_m + b_m \partial b_m), \quad k_{2m} = 2b_m \partial b_m + (b_m^2 - e^{2b_m} + 1) \partial b_m,
\]

\[
k_{3m} = -2b_m \partial b_m - (b_m^2 - e^{2b_m} - 1) \partial b_m.
\] (3.85)

For the purpose of computing the scalar potential, it is convenient to compute the holomorphic Killing vectors $k^m, k^\bar{m}$. To this end we solve the equation:

\[
\delta_\alpha \Omega(z) N = -\Omega(z)^M t_{\alpha M} N = k^m_\alpha \partial_m \Omega(z) + \ell_\alpha \Omega(z),
\] (3.86)

and find:

\[
k_{1m} = -2z^m \partial_m, \quad k_{2m} = (1 + (z_m) \partial_m, \quad k_{3m} = (1 - (z_m) \partial_m.
\] (3.87)

These are conveniently expressed in terms of a holomorphic prepotential $\mathcal{P}_\alpha(z)$:

\[
\mathcal{P}_\alpha = -\overline{\mathcal{P}}_{\bar{M}} t_{\alpha M} N \mathcal{C}_{NL} V^L,
\]

\[
\mathcal{P}_{1m} = -i \frac{z^m + \bar{z}^m}{z^m - \bar{z}^m}, \quad \mathcal{P}_{2m} = i \frac{1 + |z^m|^2}{z^m - \bar{z}^m}, \quad \mathcal{P}_{2m} = i \frac{1 - |z^m|^2}{z^m - \bar{z}^m},
\] (3.88)

the relation being:

\[
k^m_\alpha = -ig^\bar{m}_n \partial_n \mathcal{P}_\alpha.
\] (3.89)

3.4.2. The gaugings with no Fayet–Iliopoulos terms

We first consider the case of no Fayet–Iliopoulos terms, namely ($\Theta_{M}^\alpha = 0$). We can use the global symmetry $G$ of the theory to simplify our analysis. Indeed the field equations and Bianchi identities are invariant if we transform the field and embedding tensors at the same time. This is in particular true for the scalar potential $V(\phi, \Theta)$:

\[
\forall g \in G: V(\phi, \Theta) = V(g \circ \phi, g \circ \Theta),
\] (3.90)

where $(g \circ \phi)^\gamma$ are the scalar fields obtained from $\phi^\gamma$ by the action of the isometry $g$, and $(g \circ \Theta)_{M}^\alpha$ is the $g$-transformed embedding tensor. Notice that we can have other formal symmetries of the potential which are not in USK. Consider for instance the symplectic transformation:
\[ S = \text{diag}(1, \varepsilon_m, 1, \varepsilon_m), \]  
\[ \varepsilon_m = \pm 1, \quad \varepsilon_1 \varepsilon_2 \varepsilon_3 = 1. \]

These transformations correspond to the isometries \( \varepsilon^m \to \varepsilon_m \varepsilon^m \), which however do not preserve the physical domain defined by the upper half plane for each complex coordinate: \( \text{Im}(\varepsilon^m) > 0 \). Therefore embedding tensors connected by such transformations are to be regarded as physically inequivalent.

We have shown in Section 3.4 that the embedding tensor, solution to the linear constraints and in the absence of Fayet–Iliopoulos terms, is parameterized by two independent tensors \( \xi(2), \xi(3) \) in the \( (\frac{1}{2}, \frac{1}{2}, \frac{1}{2}) \) of \( U_{SK} \). These are then subject to the quadratic constraints that restrict the \( U_{SK} \)-orbits of these two quantities. We can think of acting by means of \( U_{SK} \) on \( \xi(2) \), so as to make it the simplest possible. By virtue of Eq. (3.90) this will not change the physics of the gauged model (vacua, spectra, interactions), but just make their analysis simpler.

Let us recall that the \( U_{SK} \)-orbits of a single object, say \( \xi(2)^M \), in the \( W = (\frac{1}{2}, \frac{1}{2}, \frac{1}{2}) \) representation are described by a quartic invariant \( I_4(\xi(2)) \), defined as:

\[ I_4(\xi(2)) = -\frac{2}{3} t_{AMN} t_{APQ} \xi(2)^M \xi(2)^N \xi(2)^P \xi(2)^Q. \]  

A very important observation is that by definition the \( W \) representation is that of the electromagnetic-charges of a black-hole solution of ungauged supergravity. Hence the components of the \( \xi(2) \)-tensor could be identified with the charges \( Q \) of such a black-hole and the classification of the orbits of \( U_{SK} \) in the representation \( W \) coincides with the classification of black-hole solutions. The quartic invariant is just the same that in the black-hole case determines the area of the horizon. Here we make the first contact with the profound relation that links the black-hole potentials with the gauging potentials. The orbits in the \( (\frac{1}{2}, \frac{1}{2}, \frac{1}{2}) \)-representation are classified as follows [36]:

(i) Regular, \( I_4 > 0 \), and there exists a \( \mathbb{Z}_3 \)-centralizer;
(ii) Regular, \( I_4 > 0 \), no \( \mathbb{Z}_3 \)-centralizer;
(iii) Regular, \( I_4 < 0 \);
(iv) Light-like, \( I_4 = 0, \partial_M I_4 \neq 0 \);
(v) Critical, \( I_4 = 0, \partial_M I_4 = 0, t_A M N \partial_M \partial_N I_4 \neq 0 \);
(vi) Doubly critical, \( I_4 = 0, \partial_M \partial_N I_4 = 0, t_A M N \partial_M \partial_N I_4 = 0 \),

where \( \partial_M \equiv \partial/\partial \xi^{(2)M} \). The quadratic constraints (3.44) restrict \( \xi(2) \) (and \( \xi(3) \)) to be either in the critical or in the doubly-critical orbit. Let us analyze the two cases separately.

\( \xi(2) \) Critical. The quadratic constraints imply \( \xi(3) = 0 \) and thus the embedding tensor is parameterized by \( \xi(1) = \xi(2) \), namely the diagonal of the first two \( \text{SL}(2, \mathbb{R}) \) groups in \( G_{SK} \). We can choose a representative of the orbit in the form:

\[ \xi(2) = g(0, 1, c, 0, 0, 0, 0), \]  

The scalar potential reads:

\[ \text{(3.93)} \]
\[ V = V_{\text{gaugino},1} = g^2 e^{-\theta_1 - \theta_2 - \theta_3} \left( (b_1 + c b_2)^2 + (e^{\theta_1} - c e^{\theta_2})^2 \right). \]  

(3.94)

The truncation to the dilatons \((b_m = 0)\) is a consistent one:

\[
\frac{\partial V}{\partial b_m} \bigg|_{b_m=0} = 0, \tag{3.95}
\]

and

\[
V|_{b_m=0} = g^2 \left( e^{-\frac{1}{2}(b_1 + b_2 + b_3)} - ce^{-\frac{1}{2}(b_1 - b_2 + b_3)} \right)^2. \tag{3.96}
\]

The above potential has an extremum if \(c > 0\), for \(e^{\theta_1} = c e^{\theta_2}\), while it is runaway if \(c < 0\). The sign of \(c\) is changed by a transformation of the kind (3.91) with \(\varepsilon_1 = -\varepsilon_2 = -\varepsilon_3 = 1\). For the reason outlined above, in passing from a negative to a positive \(c\), the critical point of the potential moves to the unphysical domain \((\text{Im}(z^2) < 0)\). The gauging for \(c = -1\) coincides with the one considered in [26], in the absence of Fayet–Iliopoulos terms, with potential (compare with Eq. (3.12):

\[
V_{CV} = \frac{e_0^2}{2 \text{Im}(S)} P_2^+(y), \tag{3.97}
\]

where \(P_2^+(y) = 1 - 2y_0 \bar{y}_0 \pm 2y_1 \bar{y}_1 + y^2 \bar{y}^2\) and \(y^2 = y_0^2 + y_1^2\). The two potentials are connected by the transformation relations between the Calabi–Vesentini and the special coordinates spelled out in Eq. (3.57) and by setting \(e_0 = 2\sqrt{2}g\).

\(\xi^{(2)}\) Doubly-critical. We can choose a representative of the orbit in the form:

\[
\xi^{(2)} = g(1, 0, 0, 0, 0, 0, 0, 0). \tag{3.98}
\]

In this case \(\xi^{(3)}\) is non-vanishing and has the form:

\[
\xi^{(3)} = g'(1, 0, 0, 0, 0, 0, 0, 0). \tag{3.99}
\]

The gauging is electric \((\Theta^A = 0)\) and the gauge generators \(X_A = (X_0, X_m), m = 1, 2, 3\), satisfy the following commutation relations:

\[
[X_0, X_m] = M_m^n X_n, \quad M_m^n = \text{diag}( -2(g + g'), 2g, 2g'), \tag{3.100}
\]

all other commutators being zero. This gauging originates from a Scherk–Schwarz reduction from \(D = 5\), in which the semisimple global symmetry generator defining the reduction is the 2-parameter combination \(M_m^n\) of the \(\text{so}(1, 1)^2\) global symmetry generators of the \(D = 5\) parent theory.

The scalar potential is axion-independent and reads:

\[
V = (g^2 + gg' + g'^2) e^{-\theta_1 - \theta_2 - \theta_3}. \tag{3.101}
\]

This potential is trivially integrable since it contains only one exponential of a single scalar field combination.

3.4.3. Adding \(U(1)\) Fayet–Iliopoulos terms

Let us now consider adding a component of the embedding tensor along one generator of the \(\text{SO}(3)\) global symmetry group: \(\theta_M = \theta_M^{a=1}\). The constraints on \(\theta_M\) are (3.43) and (3.45), which read:
While the constraints (3.44) on $\Theta^A_M$ are just the same as before and induce the same restrictions on the orbits of $\xi^{(2)}, \xi^{(3)}$. Clearly if $X P_N\theta_Q = 0$, namely $\Theta^A_M = 0$, no SK isometries are gauged and there are no constraints on $\theta_M$. We shall consider this case separately.

The potential reads

$$V = V_{\text{gaugino}, 1} + V_{\text{gaugino}, 2} + V_{\text{gravitino}},$$

where $V_{\text{gaugino}, 1}$ was constructed in the various cases in the previous section, while:

$$V_{\text{gaugino}, 2} + V_{\text{gravitino}} = (g_{\bar{M}n}D_{\bar{M}}V^M D_n V^N - 3V^M V^N)\theta_M\theta_N,$$

has just the form of an $\mathcal{N} = 1$ potential generated by a superpotential:

$$W_h = \theta_M \Omega^M_{SF},$$

as discussed later in Eq. (4.7). It is interesting to rewrite the above contribution to the potential in terms of quantities which are familiar in the context of black holes in supergravity. We use the property:

$$g_{\bar{M}n}D_{\bar{M}}V^{(M} D_n V^{N)} = -\frac{1}{2} \mathcal{M}^{-1MN},$$

where $\mathcal{M}_{MN}$ is the symplectic, symmetric, negative-definite matrix defined later in Eq. (4.30) in terms of the $\mathcal{N}_{\Lambda \Sigma}(z, \bar{z})$ matrix which appears in the $D = 4$ Lagrangian (see Eq. (4.24)). Let us now define the complex quantity $Z = V^M \theta_M$. The FI contribution to the scalar potential (3.106) can be recast in the form:

$$V_{\text{gaugino}, 2} + V_{\text{gravitino}} = -\frac{1}{2} \theta_M \mathcal{M}^{-1MN} \theta_N - 4|Z|^2 = V_{BH} - 4|Z|^2.$$

The first term has the same form as the (positive-definite) effective potential for a static black hole with charges $Q^M = C^{MN}\theta_N$, while the second one is the squared modulus of the black hole central charge. Notice that we can also write

$$V_{BH} = -\frac{1}{2} \theta_M \mathcal{M}^{-1MN} \theta_N = |Z|^2 + g_{\bar{M}n}D_{\bar{M}}Z D_n Z > 0.$$

Let us now study the full scalar potential in the relevant cases.

$\xi^{(2)}$ Critical. In this case, choosing

$$\xi^{(2)} = g(0, 1, c, 0, 0, 0, 0, 0),$$

we find for $\theta_M$ the following general solution to the quadratic constraints:

$$\theta_M = \left(0, \frac{f_1}{c}, f_1, 0, 0, f_2, \frac{f_2}{c}, 0\right),$$

where $f_1, f_2$ are constants.

The scalar potential reads:

$$V = g^2 e^{-b_1-b_2-b_3}((b_1 + c b_2)^2 + (e^{b_1} - ce^{b_2})^2) + \frac{e^{-b_3}}{c}[(f_1 + f_2 b_3)^2 + f_2^2 e^{2b_3}].$$
The above potential (if $g > 0$) has an extremum only for $c < 0$, $f_2 > 0$ and:

$$h_1 = h_2 + \log(-c), \quad h_3 = -\log\left(-\frac{2cg}{f_2}\right), \quad y_3 = -\frac{f_1}{f_2}, \quad y_1 = -cy_2.$$  \hfill (3.112)

The potential at the extremum is

$$\mathcal{V}_0 = 4gf_2 > 0,$$  \hfill (3.113)

while the squared scalar mass matrix reads:

$$\left(\partial_r \partial_s \mathcal{V}^{rt}\right)|_0 = \text{diag}(2, 2, 1, 1, 0, 0) \times \mathcal{V}_0.$$  \hfill (3.114)

In this way we retrieve the stable dS vacuum of [26], discussed in Section 3.1, the two parameters $f_1$, $f_2$ being related to $e_1$ and the de Roo–Wagemann’s angle.

$\xi^{(2)}$ Doubly-critical. In this case the constraints on $\theta_M$ impose:

$$(g + g')\theta_1 = 0, \quad g\theta_2 = 0, \quad g'\theta_3 = 0, \quad g\theta^0 = g'\theta^0 = 0,$$

$$(g\theta^1 = g'\theta^1 = 0, \quad g\theta^2 = g'\theta^2 = 0, \quad g\theta^3 = g'\theta^3 = 0.$$  \hfill (3.115)

Under these conditions, unless $g = g' = 0$, which is the case we shall consider next, the FI contribution to the scalar potential vanishes.

Case $\Theta_M^A = 0$. Pure Fayet–Iliopoulos gauging. In this case, we can act on $\theta_M$ by means of USK and reduce it the theta vector to its canonical normal form:

$$\theta_M = (0, f_1, f_2, f_3, f_0^0, 0, 0, 0).$$  \hfill (3.116)

The scalar potential reads:

$$\mathcal{V} = \mathcal{V}_{\text{gaugino}}, + \mathcal{V}_{\text{gravitino}} = -\sum_{m=1}^{3} e^{-\beta_m} \left(f_m f^0 (b_m^2 + e^{2\beta_m}) + f_n f_p\right),$$  \hfill (3.117)

where $n \neq p \neq m$. The truncation to the dilatons is consistent and we find:

$$\mathcal{V}|_{b_m=0} = -\sum_{m=1}^{3} \left(f_m f^0 e^{\beta_m} + f_n f_p e^{-\beta_m}\right),$$  \hfill (3.118)

which is extremized with respect to the dilatons by setting

$$e^{2\beta_m} = \frac{f_n f_p}{f_m f^0},$$  \hfill (3.119)

and the potential at the extremum reads:

$$\mathcal{V}_0 = -6\epsilon \sqrt{f^0 f_1 f_2 f_3},$$  \hfill (3.120)

where $\epsilon = \pm 1$. This extremum only exists if $f^0 f_1 f_2 f_3 > 0$. This implies that $\theta_M$ should be either in the orbit (i) ($\epsilon = 1$ in the above expression for $\mathcal{V}_0$) or in the orbit (ii) ($\epsilon = -1$). Using the analogy between $\theta_M$ and black hole charges, these two orbits correspond to BPS and non-BPS with $\mathcal{J}_4 > 0$ black holes. The extremum condition for $\mathcal{V}_\text{BH}$ fixes the scalar fields at the horizon according to the attractor behavior. Now the potential has an additional term $-4|Z|^2$ which, however, for the orbits (i), (ii), has the same extrema as $\mathcal{V}_\text{BH}$ since its derivative with respect
to $z^m$ is $-4D_m Z \bar{Z}$ which vanishes for the (i) orbit since at the extremum of $V_{BH}$ (BPS black hole horizon) $D_m Z = 0$, and for the (ii) orbit since at the extremum of $V_{BH}$ (black hole horizon) $Z = 0$.

We conclude that in the “BPS” orbit (i) the extremum corresponds to an AdS-vacuum where the scalar mass spectrum reads as follows:

$$\left.\left(\partial_r \partial_s V_{gst}\right)\right|_0 = \text{diag}\left(\frac{2}{3}, \frac{2}{3}, \frac{2}{3}, \frac{2}{3}, \frac{2}{3}, \frac{2}{3}\right) \times V_0 < 0. \quad (3.121)$$

These models have provided a useful supergravity framework where to study black hole solutions in anti-de Sitter spacetime [38].

In the “non-BPS” orbit (ii) the potential has a de Sitter extremum which, however, is not stable, having tachyonic directions:

$$\left.\left(\partial_r \partial_s V_{gst}\right)\right|_0 = \text{diag}\left(-\frac{2}{3}, -\frac{2}{3}, \frac{2}{3}, \frac{2}{3}, \frac{2}{3}, \frac{2}{3}\right) \times V_0. \quad (3.122)$$

We shall not consider this case in what follows.

3.5. Conclusions on the one-field cosmologies that can be derived from the gaugings of the $\mathcal{N} = 2$ STU model

Let us summarize the results of the above systematic discussion. From the gaugings of the $\mathcal{N} = 2$ STU model one can obtain following dilatonic potentials:

(A) **Critical orbit without FI terms.** We have the potential:

$$V = g^2 \left(e^{-\frac{1}{2}(\ell_1 + \ell_2 + \ell_3)} + e^{-\frac{1}{2}(\ell_1 - \ell_2 + \ell_3)}\right)^2. \quad (3.123)$$

In this case we have a consistent truncation to one dilaton by setting: $\ell_1 = \ell_2 = \ell \in \mathbb{R}$, since the derivatives of the potential with respect to $\ell_{1,2}$ vanish on such a line. The residual one dilaton potential is:

$$V = 4g^2 e^{-\ell_3} \quad (3.124)$$

which upon use of the translation rule (2.2) yields

$$\mathcal{V} = 12g^2 e^{-\frac{x}{3}}. \quad (3.125)$$

The above potential is trivially integrable, being a pure less than critical exponential.

(B) **Doubly critical orbit without FI terms.** We have the potential:

$$V = \text{const} e^{-\ell_1 - \ell_2 - \ell_3}. \quad (3.126)$$

Introducing the following field redefinitions:

$$\phi_1 = \ell_1 + \ell_2 + \ell_3, \quad \phi_2 = \ell_2 - \ell_3, \quad \phi_3 = -2\ell_1 + \ell_2 + \ell_3 \quad (3.127)$$

the kinetic term:

$$\text{kin} = \frac{1}{2} (\dot{\phi}_1 + \dot{\phi}_2 + \dot{\phi}_3) \quad (3.128)$$

is transformed into:

$$\text{kin} = \frac{1}{6} \dot{\phi}_1 + \frac{1}{4} \dot{\phi}_2 + \frac{1}{12} \dot{\phi}_3 \quad (3.129)$$
while the potential (3.126) depends only on $\phi_1$. Hence we can consistently truncate to one field by setting $\phi_2 = \phi_3 = \text{const}$ and upon use of the translation rule (2.2) we obtain a trivially integrable over critical exponential potential:

$$V = \text{const} \, e^{-\sqrt{3} \phi}. \quad (3.130)$$

(C) **Critical orbit with FI terms.** This case leads to the potential (3.111) which, as we showed, reproduces the potential (3.17) of the $\mathfrak{so}(2, 1) \times \mathfrak{u}(1)$ gauging extensively discussed in Section 3.1. Such a potential admits a stable de Sitter vacuum and a consistent one dilaton truncation to a model with a cosh potential which is not integrable, since the intrinsic index $\omega$ does not much any one of the three integrable cases.

(D) **Doubly critical orbit with FI terms.** Upon a constant shift of the dilatons in Eq. (3.118) this gauging leads to the following negative potential:

$$V = -2 \sqrt{f^0 f_1 f_2 f_3} \sum_{i=1}^{3} \cosh[h_i] \quad (3.131)$$

that has a stable anti-de Sitter extremum. We have a consistent truncation to one-field by setting to zero any two of the three dilatons. Upon use of the translation rule (2.2) we find the potential:

$$V = -\text{const} \left( 2 + \cosh \left[ \frac{\phi}{\sqrt{3}} \right] \right) \quad (3.132)$$

which does not fit into any one of the integrable series of Tables 1 and 2.

Hence apart from pure exponentials without critical points no integrable models can be fitted into any gauging of the $\mathcal{N} = 2$ STU model.

4. $\mathcal{N} = 1$ models with a superpotential

Let us now turn to consider the case of $\mathcal{N} = 1$ supergravity coupled to Wess–Zumino multiplets [24]. Following the notations of [32], the general bosonic Lagrangian of this class of models is

$$\mathcal{L}_{\text{SUGRA}} = \sqrt{-g} \left[ \mathcal{R}[g] + 2 g_{ij;}^\text{HK} \partial_i z^\mu \partial^\mu z^{i;} - 2 V(z, \bar{z}) \right], \quad (4.1)$$

where the scalar metric is Kähler (the scalar manifold must be Hodge–Kähler)

$$g_{ij;} = \partial_i \partial_j \mathcal{K}, \quad (4.2)$$

$$\mathcal{K} = \bar{\mathcal{K}} = \text{Kähler potential} \quad (4.3)$$

and the potential is

$$V = 4 e^2 \exp[\mathcal{K}] \left( g^{ij;} \mathcal{D}_i W_h(z) \mathcal{D}_j \bar{W}_h(\bar{z}) - 3 |W_h(z)|^2 \right), \quad (4.4)$$

where the superpotential $W_h(z)$ is a holomorphic function. Furthermore

---

5 Observe that here we consider only the graviton multiplet coupled to Wess–Zumino multiplets. There are no gauge multiplets and no $D$-terms. The embedding mechanisms discussed in [13] is lost a priori from the beginning.
\[ \mathcal{D}_i W = \partial_i W + \partial_i \mathcal{K} \mathcal{W}, \]
\[ \mathcal{D}_{j^*} \mathcal{W} = \partial_{j^*} \mathcal{W} + \partial_{j^*} \mathcal{K} \mathcal{W} \] (4.5)

are usually referred to as Kähler covariant derivatives. They arise since \( W_h(z) \), rather than a function, is actually a holomorphic section of the line bundle \( \mathcal{L} \to \mathcal{M}_K \) over the Kähler manifold whose first Chern class is the Kähler class, as required by the definition of Hodge–Kähler manifolds. In other words, \( c_1(\mathcal{L}) = [\mathcal{K}] \), the latter being the Kähler two-form. The fiber metric on this line bundle is \( h = \exp[\mathcal{K}] \), so that a generic section \( W(z, \bar{z}) \) of \( \mathcal{L} \) (not necessarily holomorphic) admits the invariant norm
\[ \| W \|_2 \equiv W \mathcal{W} \exp[\mathcal{K}]. \] (4.6)

A generic gauge transformation of the line bundle takes the form
\[ W'(z, \bar{z}) = \exp \left[ \frac{1}{2} f(z) \right] \times W(z, \bar{z}), \]
\[ \mathcal{W}'(z, \bar{z}) = \exp \left[ \frac{1}{2} \bar{f}(z) \right] \times \mathcal{W}(z, \bar{z}) \] (4.7)

where \( f(z) \) is a holomorphic complex function. Under the gauge transformation (4.7), the fiber metric changes according to
\[ \mathcal{K}'(z, \bar{z}) = -\mathcal{K}'(z, \bar{z}) + \Re f(z) \] (4.8)
while the norm (4.6) stays invariant. It is important to stress that the same Kähler metric \( g_{ij^*} = \partial_i \partial_{j^*} \mathcal{K}' \) would be obtained by the same token from \( \mathcal{K}' \). All transition functions from one local trivialization of the line bundle to another one are of the form (4.7) and (4.8), with an appropriate \( f(z) \). The fiber metric introduces a canonical connection \( \theta = h^{-1} \partial h \) leading to the covariant derivatives (4.5). In covariant notation, the potential (4.4) takes the form
\[ V = 4e^{2} \left( \| \mathcal{D} W \|_2^2 - 3 \| W \|_2^2 \right) \] (4.9)
where by definition
\[ \| \mathcal{D} W \|_2^2 = g^{ij^*} \mathcal{D}_i W \mathcal{D}_{j^*} \mathcal{W} \exp[\mathcal{K}], \]
\[ \| W \|_2^2 = W \mathcal{W} \exp[\mathcal{K}]. \] (4.10)

Let us now consider the notion of covariantly holomorphic section, defined by the condition
\[ \mathcal{D}_{j^*} W = 0. \] (4.11)
From any covariantly holomorphic section, one can retrieve a holomorphic one by setting
\[ W_h(z) = \exp \left[ -\frac{1}{2} \mathcal{K} \right] W \quad \Rightarrow \quad \partial_{j^*} W_h = 0. \] (4.12)

By hypothesis the superpotential \( W \) that appears in the potential (4.9) is covariantly constant. The compact notation (4.9) is very instructive since it stresses that the scalar potential results from the difference of two positive definite terms originating from two different contributions. The first contribution is the absolute square of the auxiliary fields appearing in the supersymmetry transformations of the spin \( \frac{1}{2} \)-fermions (the chiralinos belonging to Wess–Zumino multiplets), while the second is the square of the auxiliary field appearing in the supersymmetry transformation of the spin \( \frac{3}{2} \)-gravitino. Indeed
\[ \delta_{\text{SUSY}} \chi^i = i \partial_\mu z^i \gamma^\mu \epsilon^\cdot + \mathcal{H}^i \epsilon^\cdot, \quad (4.13) \]

\[ \delta_{\text{SUSY}} \Psi^\cdot = D_\mu \epsilon^\cdot + S \gamma^\mu \epsilon^\cdot, \quad (4.14) \]

where \( \epsilon^\cdot, \epsilon^\cdot_\cdot \) denote the two chiral projections of the supersymmetry parameter and the scalar field dependent auxiliary fields are

\[ S = ie \sqrt{||W||^2} = ie \sqrt{|W_h|^2} \exp \left[ \frac{1}{2} \mathcal{K} \right], \]

\[ \mathcal{H}^i = 2eg^{ij} D_j W \exp \left[ \frac{1}{2} \mathcal{K} \right]. \quad (4.15) \]

This structure of the potential shows that any de Sitter vacuum characterized by a potential \( V(z_0) \) that is positive at the extremum necessarily breaks supersymmetry since this implies that the chiralino auxiliary fields are different from zero in the vacuum \( \langle \mathcal{H}^i \rangle = \mathcal{H}^i (z_0) \neq 0 \).

Let us also stress that the parameter \( e \) appearing in the potential is just a dimensionful parameter which fixes the scale of all the masses generated by the gauging, i.e. by the introduction of a superpotential.

### 4.1. One-field models

In this general framework the simplest possibility is a model with one scalar multiplet assigned to the homogeneous Kähler manifold

\[ \mathcal{M}_K = \frac{\text{SU}(1,1)}{\text{U}(1)} \quad (4.16) \]

and Kähler potential

\[ \mathcal{K} = - \log[(z - \bar{z})^q], \quad (4.17) \]

which leads to the Kähler metric

\[ g_{\bar{z}z} = -\frac{q}{(z - \bar{z})^2}, \quad (4.18) \]

where \( q \) is an integer number. Its favorite value, \( q = 3 \), corresponds to the \( \mathcal{N} = 1 \) truncation of the \( \mathcal{N} = 2 \) model \( S^3 \) that, on its turn arises from the STU model discussed in the previous section upon identification of the three scalar multiplets \( S, T \) and \( U \). Alternatively, the case \( q = 1 \) corresponds to the \( \mathcal{N} = 1 \) truncation of an \( \mathcal{N} = 2 \) theory with vanishing Yukawa couplings. Because of their \( \mathcal{N} = 2 \) origin, both instances of the familiar Poincaré–Lobachevsky plane are not only Hodge–Kähler but actually special Kähler manifolds.

In the notation of [40], the holomorphic symplectic section governing this geometry is given by the four-component vector

\[ \Omega = \{-\sqrt{3} \bar{z}^2, z^3, \sqrt{3} \bar{z}, 1\}, \quad (4.19) \]

which transforms in the spin \( j = \frac{3}{2} \) of the \( \text{SL}(2, \mathbb{R}) \sim \text{SU}(1,1) \) group that happens to be four-dimensional symplectic

\[ \text{SL}(2, \mathbb{R}) \ni \begin{pmatrix} a & b \\ c & d \end{pmatrix} \implies \begin{pmatrix} da^2 + 2bca & -\sqrt{3}a^2 c & -cb^2 - 2adb & -\sqrt{3} b^2 d \\ -\sqrt{3} a^2 b & a^3 & \sqrt{3} ab^2 & b^3 \\ -bc^2 - 2adc & \sqrt{3} ac^2 & ad^2 + 2bcd & \sqrt{3} bd^2 \\ -\sqrt{3} c^2 d & c^3 & \sqrt{3} cd^2 & d^3 \end{pmatrix} \in \text{Sp}(4, \mathbb{R}) \quad (4.20) \]
where the preserved symplectic metric is
\[
C = \begin{pmatrix}
0 & 0 & 1 & 0 \\
0 & 0 & 0 & 1 \\
-1 & 0 & 0 & 0 \\
0 & -1 & 0 & 0
\end{pmatrix}.
\] (4.21)

According to the general set up of Special Geometry (for a recent review see [39]), the Kähler potential (4.17) is retrieved letting
\[
K(z, \bar{z}) = - \log[-i\Omega C \Omega].
\] (4.22)

Independently of the special structure that is essential for \(\mathcal{N} = 2\) supersymmetry, at the \(\mathcal{N} = 1\) level one can consider general superpotentials that are consistent with the Hodge–Kähler structure, provided they are holomorphic, namely provided they can be expanded in a power series of the unique complex field
\[
W_h(z) = \sum_{n \in \mathbb{N}} c_n z^n,
\] (4.23)
where the \(c_n\) are complex coefficients. The sum over \(n\) extends to a finite or infinite subset of the natural numbers \(\mathbb{N}\), while rational or irrational powers leading to cuts are excluded in order to obtain properly transforming sections of the Hodge line bundle.

Notwithstanding this wider choice available at the \(\mathcal{N} = 1\) level, it is interesting to note that in discussing black-hole solutions of the corresponding \(\mathcal{N} = 2\) model one is lead to an effective sigma-model whose Lagrangian resembles closely the effective Lagrangian of the cosmological sigma-model and displays a potential that is also built in terms of a superpotential, although the latter is more restricted. The comparison between cosmological and black-hole constructions provides inspiring hints on the choice of appropriate superpotentials. Let us briefly see how this works.

4.2. Cosmological versus black-hole potentials

The common starting point for black-hole and cosmological solutions is the general form of the bosonic portion of the four-dimensional supergravity, which takes the form (for a recent review see Chapter 8, Vol. 2 in [39] and all references therein)
\[
L^{(4)} = \sqrt{|\text{det} g|} \left[ \mathcal{R}[g] + \frac{1}{2} \partial_\mu \phi^a \partial^\mu \phi^b g_{ab}(\phi) + 2 \text{Im} \mathcal{N}_A \mathcal{S}(\phi) F^A_{\mu \nu} F^{\mu \nu} \right] + \text{Re} \mathcal{N}_A \mathcal{S}(\phi) F^A_{\mu \nu} F^{\mu \nu} \epsilon^{\mu \nu \rho \sigma},
\] (4.24)
where \(F^A_{\mu \nu} \equiv (\partial_\mu A^A_\nu - \partial_\nu A^A_\mu) / 2\) are the field strengths of the vector fields, \(\phi^a\) denotes the collection of \(n_s\) scalar fields parameterizing the scalar manifold \(\mathcal{M}_{\text{scalar}}^{D=4}\), with \(g_{ab}(\phi)\) its metric and the field-dependent complex matrix \(\mathcal{N}_A \mathcal{S}(\phi)\) is fully determined by constraints imposed by duality symmetries. In addition, the scalar potential \(V(\phi)\) is determined by the appropriate gauging procedures, while \(e\) is the gauge coupling constant, which vanishes in ungauged supergravity.

Although the discussion can be extended also to higher \(\mathcal{N}\), for simplicity we focus on the \(\mathcal{N} = 2, 1\) cases, where the real scalar fields are grouped in complex combinations \(z^i\) and their kinetic term becomes
\[
\frac{1}{2} \partial_\mu \phi^a \partial^\mu \phi^b g_{ab}(\phi) \mapsto 2 g_{ij} (z, \bar{z}) \partial_\mu z^i \partial^\mu \bar{z}^j.
\] (4.25)
In the case of extremal black-hole solutions of ungauged supergravity \( (e = 0) \), the four-dimensional metric is of the form

\[
ds_{\text{BH}}^2 = \exp[U(\tau)] \, dt^2 - \exp[-U(\tau)] \, dx^i \otimes dx^j \delta_{ij}
\]

(4.26)

where \( \tau = -(\sum_{i=1}^{3} x_i^2)^{-\frac{1}{2}} \) is the reciprocal of the radial distance,\(^6\) one is lead to the effective Euclidian \( \sigma \)-model (for a recent review see Chapter 9, Vol. 2 in [39] and all references therein)

\[
S_{\text{BH}} \equiv \int \mathcal{L}_{\text{BH}}(\tau) \, d\tau,
\]

\[
\mathcal{L}_{\text{BH}}(\tau) = \frac{1}{4} \left( \frac{dU}{d\tau} \right)^2 + g_{ij} \, \frac{dz^i}{d\tau} \frac{dz^j}{d\tau} + e^U V_{\text{BH}}(z, \bar{z}, Q).
\]

(4.27)

The geodesic potential \( V_{\text{BH}}(z, \bar{z}, Q) \) is defined by\(^7\)

\[
V_{\text{BH}}(z, \bar{z}, Q) = -\frac{1}{4} Q' \mathcal{M}_4^{-1}(\bar{N}) Q.
\]

(4.28)

Here \( Q \) is the vector of electric and magnetic charges of the hole, which transforms in the same representation of the Kähler isometry group \( G \) as the symplectic section of Special Geometry. In the \( S^3 \) case \( G = \text{SL}(2, \mathbb{R}) \) and the four charges of the hole

\[
Q = \{ p_1, p_2, q_1, q_2 \}
\]

(4.29)

transform by means of the matrix (4.20). The \( (2n + 2) \times (2n + 2) \) matrix \( \mathcal{M}_4^{-1} \) appearing in Eq. (4.28) is given in terms of the \( (n + 1) \times (n + 1) \) matrix \( \mathcal{N}_{A \Sigma}(\phi) \) that appears in the 4D Lagrangian. In detail,

\[
\mathcal{M}_4^{-1} = \left( \begin{array}{cc}
\text{Im} \mathcal{N} + \text{Re} \mathcal{N} \text{Im} \mathcal{N}^{-1} \text{Re} \mathcal{N} & -\text{Re} \mathcal{N} \text{Im} \mathcal{N}^{-1} \\
-\text{Im} \mathcal{N}^{-1} \text{Re} \mathcal{N} & \text{Im} \mathcal{N}^{-1} \text{Re} \mathcal{N}
\end{array} \right),
\]

(4.30)

where \( n \) is the number of vector multiplets coupled to supergravity.

Starting instead from the spatially flat cosmological metric

\[
ds_{\text{Cosm}}^2 = \exp[3A(t)] \, dt^2 - \exp[2A(t)] \, dx^i \otimes dx^j \delta_{ij}
\]

(4.31)

which, in the language of the preceding sections, corresponds to the gauge \( B = 3A \), one is led to the effective sigma model

\[
S_{\text{Cosm}} \equiv \int \mathcal{L}_{\text{Cosm}}(t) \, dt,
\]

\[
\mathcal{L}_{\text{Cosm}}(\tau) = -\frac{3}{2} \left( \frac{dA}{dt} \right)^2 + g_{ij} \frac{dz^i}{dt} \frac{dz^j}{dt} - e^{6A} V_{\text{Cosm}}(z, \bar{z}),
\]

(4.32)

where \( V_{\text{Cosm}}(z, \bar{z}) = e^2 V(\phi) \) is the scalar potential produced by gauging that, in an \( \mathcal{N} = 1 \) theory, or in an \( \mathcal{N} = 2 \) one with only abelian gauge groups (Fayet–Iliopoulos terms), admits the representation in terms of a holomorphic superpotential recalled in Eq. (4.4). The similarity between the cosmological and black-hole cases becomes striking if one recalls that the black-hole geodesic potential (4.28) admits the alternative representation

\(^6\) Not to be mistaken for the parametric time variable used in the rest of the present paper.

\(^7\) Notice that the potential \( V_{\text{BH}} \) enters the effective Lagrangian with an unusual plus sign. This is related to the fact that the variable \( \tau \) is (only in this section) a spatial coordinate.
\[ V_{\text{BH}}(z, \bar{z}, Q) = \frac{1}{2}(|Z|^2 + g^{ij} D_i Z D_j \bar{Z}). \] (4.33)

Here \( Z \) denotes the field-dependent central charge of the supersymmetry algebra

\[ Z \equiv \exp \left[ \frac{1}{2} K(z, z) \right] Q^T C \Omega(z), \] (4.34)

\( \Omega(z) \) denotes the holomorphic symplectic section of special Kähler geometry (that of Eq. (3.3) for the STU model, or that of Eq. (4.19) for the \( S^3 \) model) and \( K(z, z) \) denotes the Kähler potential. Introducing the black-hole holomorphic superpotential

\[ W_{\text{BH}}(z) = Q^T C \Omega(z). \] (4.35)

Eq. (4.33) for the geodesic potential can be recast in the form

\[ V_{\text{BH}}(z, \bar{z}, Q) = \frac{1}{2} \exp \left[ K(z, z) \right] \left( g^{ij} D_i W_{\text{BH}} D_j \bar{W}_{\text{BH}} + |W_{\text{BH}}|^2 \right) \] (4.36)

which is almost identical to Eq. (4.4) yielding the cosmological potential, up to a crucial change of sign and coefficient. The coefficient \(-3\) of the second term becomes \(+1\), and in this fashion the black hole potential is strictly positive definite since it is the sum of two squares. Yet the entire discussion suggests that black-hole superpotentials, that are group theoretically classified by the available \( G \)-orbits of charge vectors \( Q \), form a good class of superpotentials also for Gauged Supergravity models. Indeed we already saw, by means of the systematic analysis of the STU model, that black-hole superpotentials encode and exhaust the available abelian gaugings for \( N = 2 \) supergravity theories.

### 4.3. Cosmological potentials from the \( S^3 \) model

Relying on the preceding discussion, let us consider the abelian gaugings of the \( S^3 \) model provided by the superpotential

\[ W_Q(z) = QC \Omega = -q_2 z^3 + \sqrt{3} q_1 z^2 + \sqrt{3} p_1 z + p_2, \] (4.37)

which happens to be the most general third-order polynomial. Let us stress that in multi-field models based on larger special Kähler homogeneous manifolds \( G/H \), despite the existence of many coordinates \( z_i \), the order of the superpotential will stay three since this is the polynomial order of the symplectic section for all such special geometries. Inserting (4.37) into (4.4) yields the four-parameter potential

\[ V(z, \bar{z}, Q) = -i (2p_1^2 + ((z + \bar{z})q_1 + 2\sqrt{3} z \bar{z} q_2)p_1 + 2z \bar{z} q_1^2 + p_2(3(z + \bar{z}) q_2 - 2\sqrt{3} q_1)) \frac{z - \bar{z}}{z - \bar{z}} \] (4.38)

in which one can decompose \( z \) into its real and imaginary parts according to

\[ z = i e^h + b. \] (4.39)

Not every choice of the charge vector \( Q \) allows for a consistent truncation to a vanishing axion, guaranteed by the condition

\[ \partial_b V(z, \bar{z}, Q)|_{b=0} = 0 \] (4.40)

and yet there is a representative with such a property for every \( \text{SL}(2, \mathbb{R}) \)-orbit in the \( j = \frac{3}{2} \) representation. These orbits are [36] (see [40] for the notations)
1. The very small orbit with a parabolic stability group \( O_1 = \{ p_1 \to 0, p_2 \to 0, q_1 \to 0, q_2 \to q \} \).
2. The small orbit with no stability group \( O_2 = \{ p_1 \to \sqrt{3} p, p_2 \to 0, q_1 \to 0, q_2 \to 0 \} \).
3. The large orbit with positive quartic-invariant (i.e. \( p_2 q_1 < 0 \) regular BPS in black hole constructions) \( O_3 = \{ p_1 \to 0, p_2 \to p, q_1 \to -\sqrt{3} q, q_2 \to 0 \} \).
4. The large orbit with negative quartic-invariant (i.e. \( p_2 q_1 > 0 \) regular non-BPS in black hole constructions) \( O_4 = \{ p_1 \to 0, p_2 \to p, q_1 \to \sqrt{3} q, q_2 \to 0 \} \).

and the following superpotentials and potentials:

\( O_1 \) The superpotential is purely cubic \( W = -q z^3 \) and the potential vanishes
\[
V = 0.
\] (4.41)
This is an instance of flat potentials \([22]\). Namely supersymmetry is broken by the presence of non-vanishing auxiliary fields, yet the vacuum energy is exactly zero and the ground state is Minkowski space.

\( O_2 \) The superpotential is linear \( W = 3 p z \) and the consistent truncation to zero axion yields a pure exponential
\[
V = -3 e^{-h} p^2.
\] (4.42)
This potential is trivially integrable.

\( O_3 \) The superpotential is quadratic \( W = p - 3 q z^2 \) and the consistent truncation to zero axion yields the following potential
\[
V = -3 e^{-h} q (p + e^{2h} q) \simeq -3 q^2 \cosh h.
\] (4.43)
The last form of the potential can be always achieved by means of a constant shift of the scalar field \( h \leftrightarrow h + \text{const} \). In this case the intrinsic index is:
\[
\omega = \frac{1}{3}
\] (4.44)
since the kinetic term of the \( S^3 \)-model corresponds to \( q = 3 \). It is different from the value \( \omega = 1 \) which is obtained from the non-abelian \( \mathfrak{so}(1, 2) \)-gauging of the same model, yet it is still different from either one of the integrable indices: \( \omega \neq \sqrt{3} \) and \( \omega \neq \frac{2}{\sqrt{3}} \). This result confirms what we already learned. Consistent one-field truncations of Gauged Supergravity easily yield cosmological models of the cosh-type, yet non-integrable ones. It is interesting to note that the cosh case is in correspondence with the regular BPS black holes.

\( O_4 \) The superpotential is quadratic \( W = p + 3 q z^2 \) but with a different relative sign between the constant and quadratic terms. The consistent truncation to zero axion yields the following potential
\[
V = 3 e^{-h} p q - 3 e^{h} q^2 \simeq -3 q^2 \sinh h.
\] (4.45)
As in the previous case the last form of the potential can always be achieved by means of a constant shift of the scalar field. It is interesting to note that the sinh case of potential is in correspondence with the regular non-BPS black holes. Once again the index \( \omega \) is not a critical one for integrability.
5. The supersymmetric integrable model with one multiplet

The unique integrable model that so far we have been able to fit into the considered supersymmetric framework with just one multiplet (the Kähler metric is fixed once for all to the choice (4.16) belongs to the series $I_2$ of Table 1 and occurs for the under-critical value $\gamma = \frac{2}{3}$. Before proceeding further into the analysis of this particular integrable model it is just appropriate to stress that in a couple of separate publications Sagnotti and collaborators [16,19] have also shown that the phenomenon of climbing scalars, displayed by all of the integrable models we were able to classify, has the potential ability to explain the oscillations in the low angular momentum part of the CMB spectrum, apparently observed by PLANCK. In his recent talk given at the Dubna SQS2013 workshop, our coauthor Sagnotti has also shown a best fit to the PLANCK data for the low $\ell$ part of the spectrum, by using precisely the series of integrable potentials $I_2$\(^8\)

\[ V(\phi) = a \exp[2\sqrt{3}\gamma \phi] + b \exp[\sqrt{3}(\gamma + 1)\phi] \]  

(5.1)

with the particularly nice value $\gamma = -\frac{7}{6}$ (see [13] for details about the $D$-map insertion into supergravity). Here the different subcritical value $\gamma = \frac{2}{3}$ is select by supergravity when we try to realize the integral model through a superpotential ($F$-embedding). Indeed this potential can be obtained from the $S^3$-model with a carefully calibrated and unique superpotential that now we describe. We immediately anticipate that such a superpotential is not of the form discussed in the previous section and therefore strictly corresponds to an $\mathcal{N} = 1$ theory and not to an abelian gauging of the $\mathcal{N} = 2$ model. Technically, the difficulty met when trying to fit an integrable case into supergravity coupled just to one multiplet resides in the following. If the superpotential involves powers only up to the cubic order, as pertains to the construction via symplectic sections, the dilaton truncation can contain at most two types of exponentials, one positive and one negative, so that one can reach either cosh $p\phi$ or sinh $p\phi$ models. Yet the obtained index $p$ is always 1, different from the $p = 3, 2$ required by integrability. In order to get higher values of $p$, one would need higher powers $z^n$ in the superpotential, but as the degree of $W(z)$ increases one is confronted with new problems: one can generate higher exponentials $\exp[p\phi]$ but only positive ones, while negative exponents are bounded from below, so that the list ends with $\exp[-3\phi]$. On the other hand, together with the highest positive exponential $\exp[p_{\text{max}}\phi]$, also subleading ones for $0 < p < p_{\text{max}}$ appear and cannot be eliminated by a choice of coefficients in $W(z)$. As a result, the possible match with integrable models of the cosh, and sinh type is ruled out, as the match with the sporadic potentials of Table 2, all of which have the property of being symmetric in positive and negative exponentials. One is thus left with the two series $I_2$ and [9] of Table 1. The last is easily ruled out, since the exponents $6\gamma$ and $\frac{6}{\gamma}$ can be simultaneously integer only for $\gamma = 1, 2, 3$, and no superpotential produces these values without producing other exponentials with intermediate subleading exponents. The hunting ground is thus restricted to the series $I_2$ of Table 1 (the cosh models have already been discussed), where one is to spot a combination of powers in $W(z)$ that gives rise to only two exponents in the potential, whose indices should be related by the very restricted relation defining the series. A careful and systematic analysis led us to the unique solution provided by the following superpotential:

\[ W_{\text{int}} = \lambda z^4 + i\kappa z^3, \]  

(5.2)

---

\(^8\) In comparing the following equation with the table of paper [1], please note the coefficient $\sqrt{3}$ appearing in the exponents that has been introduced to convert the unconventional normalization of the field $\psi$ used there to the canonical normalization of the field $\phi$ used here.
where \( \lambda \) and \( \kappa \) are real constants. Performing the construction of the scalar potential one is led to

\[
V_{\text{int}}(z, \bar{z}) = \frac{z^2 \bar{z}^2 \lambda (3z^2 \kappa + 4i \bar{z}z^2 \lambda - 4i \bar{z}^2 \lambda z + 3\bar{z}^2 \kappa)}{3(z - \bar{z})^2}.
\]

(5.3)

To study the extrema of the above potential and for convenience in the further development of the integration it is useful to change parametrization, reabsorbing the overall coupling constant \( \lambda \) into a rescaling of the space–time coordinates and setting:

\[
\lambda = \frac{6}{\sqrt{5}}, \quad \kappa = \frac{2 \omega}{\sqrt{5}}.
\]

(5.4)

In this way the potential (5.3) becomes

\[
V_{\text{int}}(z, \bar{z}) = \frac{12z^2 \bar{z}^2 ((4i \bar{z} + \omega)z^2 - 4i \bar{z}^2 z + \bar{z}^2 \omega)}{5(z - \bar{z})^2}.
\]

(5.5)

Next let us consider the derivative of the potential with respect to the complex field \( z \):

\[
\partial_z V_{\text{int}} = \frac{24z^2 \bar{z}^2 ((4i \bar{z} + \omega)z^3 + 2\bar{z}(-5i \bar{z} - \omega)z^2 + 6i \bar{z}^3 z - \bar{z}^3 \omega)}{5(z - \bar{z})^3}
\]

\[
= \frac{6}{5} be^{-2h}(b^2 + e^{2h})(3(4e^h - \omega)b^2 + e^{2h}(\omega + 12e^h))
\]

\[
+ i \left( -\frac{6}{5} e^{-3h}(b^2 + e^{2h})((\omega - 2e^h)b^4 + e^{2h}(8e^h - \omega)b^2 + 2e^{4h}(\omega + 5e^h)) \right).
\]

(5.6)

In Eqs. (5.7), (5.8) we have separated the real and imaginary parts of the potential derivative after replacing the field \( z \) with its standard parametrization in terms of a dilaton and an axion:

\[
z = i \exp[h] + b.
\]

(5.9)

In order to get a true extremum both the real and imaginary part of the derivative should vanish for appropriate values of \( b \) and \( h \). We begin with considering the zeros of the real part (5.7) in the axion \( b \). It is immediately evident that there are three of them:

\[
b = 0, \quad b = \pm \frac{\sqrt{-e^{2h} \omega - 12e^{3h}}}{\sqrt{3} \sqrt{4e^h - \omega}}.
\]

(5.10)

The first zero in (5.10) is always available. The other two can occur only if \( \omega \) is negative and

\[-e^{2h} \omega - 12e^{3h} > 0.\]

(5.11)

If we choose the first zero (truncation of the axion) and we insert it into the imaginary part of the derivative we get:

\[
-\frac{12}{5} e^{3h} \omega - 12e^{4h} = 0 \quad \Rightarrow \quad h = \left\{ \begin{array}{ll}
- \log(-\frac{5}{\omega}) & \text{if } \omega < 0, \\
-\infty & \text{always.}
\end{array} \right.
\]

(5.12)

In the case the second and third zeros displayed in (5.10) are permissible (\( \omega < 0 \)), substituting their values in the imaginary part of the derivative (5.8) we obtain the condition

\[
-\frac{128e^{3h}(2e^h - \omega)\omega^3}{45(4e^h - \omega)^3} = 0 \quad \Rightarrow \quad h = -\infty.
\]

(5.13)
Indeed for $\omega < 0$, no other solution of the above equation are available. We conclude that if $\omega > 0$ the only extremum of the potential is at $\tilde{h} = -\infty$ where the potential vanishes so that such an extremum corresponds to a Minkowski vacuum. If $\omega < 0$ we have instead an additional extremum at:

$$z_0 = i\frac{|\omega|}{5}$$

(5.14)

where the potential takes the following negative value:

$$V_{\text{int}}(z_0) = \frac{6\omega^5}{15 \times 625} < 0.$$  

(5.15)

Hence the extremum (5.14) defines an anti-de Sitter vacuum. We can wonder whether such an AdS vacuum is either supersymmetric or stable. The first possibility can be immediately ruled out by computing the derivative of the superpotential at the extremum:

$$\partial_z W_{\text{int}}(z) \bigg|_{z = z_0} = -\frac{6i\omega^3}{125\sqrt{5}} \neq 0.$$  

(5.16)

Since $\partial_z W_{\text{int}}(z)$ does not vanish at the extremum, the auxiliary field of the chiralino is different from zero and supersymmetry is broken. In order to investigate stability of the AdS vacuum we have to consider the Breitenlohner–Freedman bound [41] which, in the normalizations of [32] (see p. 462 of Vol. I) is given by:

$$\lambda_i \geq \frac{3}{4} V_{\text{int}}(z_0),$$  

(5.17)

where by $\lambda_i$ we have denoted the Hessian of the potential $\partial_i \partial_j V_{\text{int}}$ calculated at the extremum (5.14). Using $h, b$ as field basis we immediately obtain:

$$\partial_i \partial_j V_{\text{int}}|_{z = z_0} = \begin{pmatrix} -\frac{24\omega^5}{3125} & 0 \\ 0 & -\frac{84\omega^3}{625} \end{pmatrix}$$  

(5.18)

from which the two eigenvalues are immediately read off and seen to be both positive. Hence the Breitenlohner–Freedman bound is certainly satisfied and we can conclude that for $\omega < 0$ we have two vacua, a Minkowski vacuum at infinity and a stable, non-supersymmetric AdS vacuum at the extremum (5.14). For $\omega > 0$, instead we have only the Minkowski vacuum at infinity.

5.1. Truncation to zero axion

The potential (5.5) of the considered supergravity model can be consistently truncated to a vanishing value of the axion $b$, since its derivative with respect to $b$ vanishes at $b = 0$. Imposing such a truncation from (5.5) we obtain the following dilatonic potential:

$$V_{\text{int}} = \frac{6}{5} e^{4h} (\omega + 4e^b)$$  

(5.19)

which by means of the replacement (2.2) is mapped into the case $\gamma = \frac{2}{3}$ of the series $I_2$ of integrable potentials listed in Table 1. According to the previous analysis of extrema of the full theory, we see that, depending on the sign of the parameter $\omega$, this potential is either monotonic or it has a minimum (see Fig. 2) The important thing to note is that when it exists, the extremum of the potential is always at a negative value of the potential. It corresponds to the stable AdS
In this figure we display the behavior of the supersymmetric integrable potential (5.19) for the two choices $\omega = 1$ (above) and $\omega = -1$ (below).

vacuum discussed in the previous subsection. As its well known the AdS has no parametrization in terms of spatially flat constant time slices. Hence if we assume, to begin with, a spatially flat ansatz for the metric, as we do in Eq. (1.1), no solution of the Friedman equations can stabilize the scalar field at the AdS extremum. Indeed the exact solutions produced by the available general integral show that the scalar field always flows to infinity at the beginning and end of cosmic time.

5.2. Explicit integration of the supersymmetric integrable model

In order to integrate the field equation of this model, it is convenient to write down the explicit form of the Lagrangian which has the following form

$$\mathcal{L}_{int} = e^{3A-B} \left( -\frac{3}{2} A'^2 + \frac{3}{2} h'^2 - \frac{6}{5} e^{2B+4h} (\omega + 4e^h) \right)$$

(5.20)

and following the strategy outlined in [1], one can move on to two new functions $U(\tau)$ and $V(\tau)$

$$A(\tau) = \frac{1}{5} \log(U(\tau)) + \log(V(\tau)),$$

$$B(\tau) = 2 \log(V(\tau)) - \frac{2}{5} \log(U(\tau)),$$
\[ h(\tau) = \frac{1}{5} \left( \log(U(\tau)) - 5 \log(V(\tau)) \right). \quad (5.21) \]

Inserting the transformation (5.21) into the Lagrangian (5.20) this becomes
\[ \mathcal{L}_{\text{int}} = -4U(\tau)^{6/5} - \omega V(\tau)U(\tau) - U'(\tau)V'(\tau) \]
while the Hamiltonian constraint takes the form
\[ H = 4(U(\tau)^{6/5} + \omega V(\tau)U(\tau) - U'(\tau)V'(\tau) = 0. \quad (5.23) \]
The field equations associated with (5.22) have the following triangular form:
\[ \omega U(\tau) - U''(\tau) = 0, \]
\[ \omega V(\tau) - V''(\tau) + \frac{24}{5} \sqrt[6]{U(\tau)} = 0 \quad (5.24) \]
and can be integrated by means of trigonometric or hyperbolic functions depending on the sign of \( \omega \).

5.3. Trigonometric solutions in the potential with AdS extremum

If we pose \( \omega = -v^2 \) the first equation becomes the equation of the standard harmonic oscillator and we have:
\[ U(\tau) = a \cos(v \tau) + b \sin(v \tau), \quad (5.25) \]
\[ V(\tau) = 4 \cot \left( v \tau + \tan^{-1} \left( \frac{a}{b} \right) \right) \]
\[ \times {}_2F_1 \left( \frac{1}{2}, \frac{9}{10}; \frac{3}{2}; \cos^2 \left( v \tau + \tan^{-1} \left( \frac{a}{b} \right) \right) \right) \sin^2 \left( v \tau + \tan^{-1} \left( \frac{a}{b} \right) \right)^{9/10} \]
\[ \times (b \cos(v \tau) - a \sin(v \tau)) + 5(\cos(v \tau) \sqrt[5]{(a \cos(v \tau) + b \sin(v \tau))^4 + 4a}) \]
\[ + \sin(v \tau) \left( b \cos(v \tau) + a \sin(v \tau) \right)^{4/5} \]
\[ \times (5v^2(a \cos(v \tau) + b \sin(v \tau))^{-4/5}), \quad (5.26) \]
where \( {}_2F_1 \) denotes a hypergeometric function of the specified indices. The parameters \( a, b, c, d \) are four integration constants, on which the Hamiltonian constraint imposes the condition
\[ (bc + ad) = 0. \quad (5.27) \]
We solve the constraint by setting \( d = -\rho a, c = \rho b \). In this way we obtain an explicit general integral depending on three parameters \( (a, b, \rho) \). The explicit form of the solution for the scale factor, for the \( \exp[B] \) function and for the scalar field \( h(\tau) \) are given below.
\[ a(\tau; a, b, \rho, v) = \frac{3(\tau; a, b, \rho, v)}{5v^2(a \cos(v \tau) + b \sin(v \tau))^{3/5}}, \]
\[ 3(\tau; a, b, \rho, v) = 5(\sin(v \tau) (4b - av^2 \rho(a \cos(v \tau) + b \sin(v \tau))^{4/5}) \]
\[ + \cos(v \tau) (b \rho(a \cos(v \tau) + b \sin(v \tau))^{4/5} v^2 + 4a)) \]
\[ + 4 \cot \left( v \tau + \tan^{-1} \left( \frac{a}{b} \right) \right) {}_2F_1 \left( \frac{1}{2}, \frac{9}{10}; \frac{3}{2}; \cos^2 \left( v \tau + \tan^{-1} \left( \frac{a}{b} \right) \right) \right) \]
\[(b \cos(\nu \tau) - a \sin(\nu \tau)) \sin^2 \left(\nu \tau + \tan^{-1} \left(\frac{a}{b}\right)\right)^{9/10},\]

\[\exp[\mathcal{B}(\tau; a, b, \rho, \nu)] = \frac{(\mathfrak{I}(\tau; a, b, \rho, \nu))^2}{25\nu^4(a \cos(\nu \tau) + b \sin(\nu \tau))^2}, \quad (5.28)\]

\[\mathfrak{h}(\tau; a, b, \rho, \nu) = \log \left[\frac{5\nu^2(a \cos(\nu \tau) + b \sin(\nu \tau))^{7/5}}{\mathfrak{I}(\tau; a, b, \rho, \nu)}\right]. \quad (5.29)\]

From the explicit form of the solution the structure of its time development is not immediately evident. Yet it is clear that it must be periodic, since all addends are constructed in terms of trigonometric functions with the same frequency \(\nu\). Therefore we are lead to suspect that the scalar field will go to infinity and the scale factor to zero in a periodic fashion. In other words we expect solutions with a Big Bang and a Big Crunch. This expectation is sustained by the general arguments of paper [1]. Indeed the considered scalar potential has an absolute minimum, yet this minimum is at a negative value, so that in the phase portrait of the equivalent first system there is no fixed point and under these conditions the only possible solutions are blow-up solutions, physically corresponding to Big Bang/Big Crunch Universes.

### 5.3.1. Structure of the moduli space of the general integral

In order to understand the actual form and the behavior of these type of solutions it is convenient to investigate first the physical interpretation of the three integration constants \(a, b, \rho\) that we have introduced and reduce the general integral to a simpler canonical form.

An a priori observation valid for all the solutions of Friedman equations is that the effective parameter labeling such solutions is only one, two parameters being accounted for by the uninteresting overall scale of the scale-factor and by the equally uninteresting possibility of shifting the parametric time \(\tau\) by a constant. What has to be done case by case is to work out those combinations of the parameters that can be disposed of by the above mentioned symmetries and single out the unique meaningful deformation parameter.

In the present case we begin by noting that all functions in the solution depend on \(\tau\) only through the combination \(\nu \tau\). Hence the frequency parameter \(\nu\) can be reabsorbed by rescaling the parametric time:

\[\nu \tau = \tau'. \quad (5.30)\]

In other words, without loss of generality we can set \(\nu = 1\). Secondly let us consider the following rescaling of the solution parameters:

\[a \mapsto \lambda a, \quad b \mapsto \lambda b, \quad \rho \mapsto \lambda^{-4/5} \rho \quad (5.31)\]

under such a transformation we have:

\[a(\tau; \lambda a, \lambda b, \lambda^{-4/5} \rho, 1) = \lambda^{2/5}a(\tau; a, b, \rho, 1),\]

\[\exp[\mathcal{B}(\tau; \lambda a, \lambda b, \lambda^{-4/5} \rho, 1)] = \exp[\mathcal{B}(\tau; a, b, \rho, 1)].\]

\[\mathfrak{h}(\tau; \lambda a, \lambda b, \lambda^{-4/5} \rho, 1) = \mathfrak{h}(\tau; a, b, \rho, 1). \quad (5.32)\]

From this we deduce that a suitable combination of the parameters \(a, b, \rho\) is just the overall scale of the scale factor, as we announced. Using the symmetry of Eq. (5.32) we could for instance fix the gauge where either one of the three parameters \(a, b, \rho\) is 1. Yet we can do much better if we realize that the ratio \(a/b\) actually amounts to a shift of the parametric time. Indeed by means of several analytic manipulations we can prove the following identities:
\[
a(a, b, \rho, 1) = \frac{4}{5} \left( \sqrt{b^2 + a^2 \cos(\tau)} \right)^{2/5} \left( \cos^2(\tau)^{9/10} F_1 \left( \frac{1}{2}, \frac{9}{10}; \frac{3}{2}; \sin^2(\tau) \right) \tan^2(\tau) + 5 \right) \\
- \rho \left( \sqrt{b^2 + a^2 \cos(\tau)} \right)^{6/5} \tan(\tau),
\]
\[
\exp \left[ B(a, b, \rho, 1) \right] = \frac{1}{25} \left( 4 \cos^2(\tau)^{9/10} F_1 \left( \frac{1}{2}, \frac{9}{10}; \frac{3}{2}; \sin^2(\tau) \right) \tan^2(\tau) \\
- 5\rho \left( \sqrt{b^2 + a^2 \cos(\tau)} \right)^{4/5} \tan(\tau) + 20 \right)^2,
\]
\[
h(a, b, \rho, 1) = - \log \left( \frac{4}{5} \cos^2(\tau)^{9/10} F_1 \left( \frac{1}{2}, \frac{9}{10}; \frac{3}{2}; \sin^2(\tau) \right) \tan^2(\tau) \\
- \rho \left( \sqrt{b^2 + a^2 \cos(\tau)} \right)^{4/5} \tan(\tau) + 4 \right). \tag{5.33}
\]

In this way we realize that after the shift \( \tau \mapsto \tau + \arctan \left( \frac{b}{a} \right) \) the solution functions depend only on the two parameters \( \sqrt{b^2 + a^2} \) and \( \rho \). Furthermore the explicit result (5.33) suggests that we redefine these latter as follows:
\[
\sqrt{b^2 + a^2} = \Lambda^\frac{5}{2}, \quad \rho = \frac{Y}{\Lambda^2} \tag{5.34}
\]
so that we obtain:
\[
a(\tau, \Lambda, Y) = \Lambda a(\tau, Y) \\
= \Lambda \left[ \frac{4}{5} \cos^\frac{3}{2}(\tau) \left( \cos(\tau)^{9/5} \right) F_1 \left( \frac{1}{2}, \frac{9}{10}; \frac{3}{2}; \sin^2(\tau) \right) \tan^2(\tau) + 5 \right] - Y \cos^\frac{1}{2}(\tau) \sin(\tau),
\]
\[
\exp[B(\tau, Y)] = \frac{1}{25} \left( 4 \cos^2(\tau)^{9/10} F_1 \left( \frac{1}{2}, \frac{9}{10}; \frac{3}{2}; \sin^2(\tau) \right) \tan^2(\tau) - 5Y \sin(\tau) \cos^\frac{1}{2}(\tau) + 20 \right)^2,
\]
\[
h(\tau, Y) = - \log \left( \frac{4}{5} \cos^2(\tau)^{9/10} F_1 \left( \frac{1}{2}, \frac{9}{10}; \frac{3}{2}; \sin^2(\tau) \right) \tan^2(\tau) - \frac{Y \sin(\tau)}{\cos^\frac{1}{2}(\tau)} + 4 \right). \tag{5.35}
\]

which puts into evidence the only relevant deformation parameter, namely \( Y \). We can get some understanding of the meaning of this latter by plotting the solution functions for various values of \( Y \). We begin by analyzing the simplest and most symmetrical solution at \( Y = 0 \).

### 5.3.2. The simplest trigonometric solution at \( Y = 0 \)

In Fig. 3 we display the behavior of the real and imaginary parts of three main functions composing the solution in the case \( Y = 0 \). From these plots it is evident that a physical solution of scalar-matter coupled gravity exists only in those windows where the three functions are simultaneously real, for instance in the interval \([-\pi/2, \pi/2]\). In such an interval of the parametric time \( \tau \) the scale factor goes from a zero to another zero so that the cosmic evolution should correspond
Fig. 3. Plots of the real and imaginary parts of the scale factor, $\exp[B]$ factor and scalar field for the case of parameter $Y = 0$. In the three diagrams the solid line represents the real part, while the dashed line represents the imaginary part. We note the periodicity of all the functions and the windows where the three of them are simultaneously real. Taking as basic reference window the interval $[-\frac{\pi}{2}, \frac{\pi}{2}]$ we note the reflection symmetry of the plots with respect to the point $\tau = 0$. 
Fig. 4. In this figure we display the behavior of the cosmic time with respect to the parametric time for the trigonometric type of solution of the supersymmetric cosmological model with parameter $Y = 0$.

to a Universe that starts with a Big Bang and finishes its life collapsing into a Big Crunch. To put such a conclusion on a firm ground we have actually to verify that both zeros of the scale factor do indeed correspond to a true space-like singularity and this can only be done by considering the intrinsic components of the curvature two-form showing that they all blow up to infinity in the initial and final point. This we will do shortly. First let us verify analytically the limit of the scale factor and of the scalar field in the initial and final point of the reality domain of the solution. We find:

$$\lim_{\tau \to \pm \frac{\pi}{2}} a(\tau; 0) = 0,$$

$$\lim_{\tau \to \pm \frac{\pi}{2}} h(\tau, 0) = -\infty,$$

$$\lim_{\tau \to \pm \frac{\pi}{2}} \exp[B(\tau, 0)] = +\infty.$$  (5.36)

This means that a life-cycle of this Universe is contained in the following finite interval of parametric time $[-\frac{\pi}{2}, \frac{\pi}{2}]$, which by the $\exp[B(\tau; 1, 1, 0, 1)]$ function is monotonically mapped into a finite interval of cosmic time. Indeed defining:

$$T_c(\tau) = \int_{-\frac{\pi}{2}}^{\tau} \exp[B(t; 1, 1, 0, 1)] dt$$  (5.37)

we find:

$$T_c \left( -\frac{\pi}{2} \right) = 0, \quad T_c \left( \frac{\pi}{2} \right) = 84.7046$$  (5.38)

the plot of $T_c(\tau)$ being displayed in Fig. 4. Due to these properties of the cosmic time function we do not lose any essential information by plotting the solution in parametric rather than in cosmic time. The essential difference between this case and the case of positive potentials with positive extrema discussed in [1] is best appreciated by considering the phase-portrait of the this solution presented in Fig. 5. The absence of a fixed point implies the structure of a blow-up solution with a Big Bang and a Big Crunch which is displayed in Fig. 6. As we already emphasized the interpretation of Big Bang and Big Crunch is suggested by the plots, yet it has to be verified by an appropriate study of the curvature singularity.
Fig. 5. In this figure we display the phase portrait of the solution defined by $Y = 0$. The axes are the scalar field $\Phi \equiv h$ and its derivative with respect to the cosmic time $V \equiv \partial_\tau h$. The extremum of the potential is at $\Phi_0 = -\log[5]$. It is reached by the solution however with a non-vanishing velocity. The field also reaches vanishing velocity yet not an extremum of the potential. Hence there is no fixed point and the trajectory is from infinity to infinity with no fixed point.

Fig. 6. In this figure we display plots (in parametric time) of the scale factor (solid line) and of the scalar field (dashed line) for the trigonometric type at $Y = 0$. It is evident that in a finite time the Universe undergoes a Big Bang, a decelerated expansion and then a Big Crunch. At the same time the scalar field climbs from $-\infty$ to a maximum and then descends again to $-\infty$.

5.3.3. The curvature two-form and its singularities

Throughout this paper we consider metrics of the form (1.1). It is important to calculate the explicit general form of the curvature 2-form associated with such metrics. To this effect we introduce the vielbein:

$$E^1 = \exp[B(\tau)] d\tau, \quad E^i = \exp[A(\tau)] dx^i \quad (i = 2, 3, 4)$$

and we obtain the following result for the matrix-valued curvature two-form:
having denoted by $\omega^{AB}$ the Levi-Civita spin connection defined by:

$$dE^A + \omega^{AB} \wedge E^C \eta_{BC} = 0$$

and having introduced the following notation:

$$a \equiv \exp[A(\tau)] = a(\tau), \quad b \equiv \exp[B(\tau)].$$

If the functions $a, b$ are specialized to the form (5.35), we obtain some rather formidable, yet fully explicit analytic expressions that correspond to the intrinsic components of the Riemann tensor for the solution under consideration. We can calculate the limit of the curvature two-form when $\tau$ approaches a zero of the scale-factor. In the case $Y = 0$, the only zeros are at $\tau = \pm \pi/2$ and we find:

$$\lim_{\tau \to \pm \pi/2} R^{AB} = \left(\begin{array}{cccc}
0 & \infty E^1 \wedge E^2 & \infty E^1 \wedge E^3 & \infty E^1 \wedge E^4 \\
-\infty E^1 \wedge E^2 & 0 & -\infty E^2 \wedge E^3 & -\infty E^2 \wedge E^4 \\
-\infty E^1 \wedge E^3 & \infty E^2 \wedge E^3 & 0 & -\infty E^3 \wedge E^4 \\
-\infty E^1 \wedge E^4 & \infty E^2 \wedge E^4 & \infty E^3 \wedge E^4 & 0
\end{array}\right).$$

Hence the Riemann tensor diverges in all directions and both the initial and final zero of the scale factor correspond to true singularities confirming their interpretation as the Big Bang and Big Crunch points. Actually we can make the statement even more precise. In the case of the $Y = 0$ solution we can calculate the asymptotic expansion of the curvature components in the neighborhood of the two singularities and we find:

$$R^{AB} \tau \to \pm \pi/2 \approx \frac{1}{(\mp \pi/2 + \tau)^{b/5}} \frac{25\Gamma(\frac{3}{2})^4}{8\pi^2 \Gamma(\frac{1}{10})^4} \times \left(\begin{array}{cccc}
0 & E^1 \wedge E^2 & E^1 \wedge E^3 & E^1 \wedge E^4 \\
-E^1 \wedge E^2 & 0 & -\frac{1}{2} E^2 \wedge E^3 & -\frac{1}{2} E^2 \wedge E^4 \\
-E^1 \wedge E^3 & \frac{1}{2} E^2 \wedge E^3 & 0 & -\frac{1}{2} E^3 \wedge E^4 \\
-E^1 \wedge E^4 & \frac{1}{2} E^2 \wedge E^4 & \frac{1}{2} E^3 \wedge E^4 & 0
\end{array}\right).$$

Hence all the components of the intrinsic curvature tensor have the same degree of divergence which is identically at the Big Bang and at the Big Crunch. This reflects the already noted $\mathbb{Z}_2$ symmetry of the solution.

### 5.3.4. $Y$-deformed solutions

We established that the actual moduli space of the trigonometric solutions is provided by the deformation parameter $Y$. It is interesting to explore the quality of the solutions that this latter parameterizes. The first fundamental question is whether all solutions have a Big Bang and a Big Crunch or other behaviors are possible. Periodicity of the solution functions guarantees that,
Fig. 7. Plot of the function \( f(\tau) \). A zero of the scale factor occurs when \( f(\tau_0) = Y \). Hence for all those values of \( Y \) that are never attained by the function \( f(\tau) \) in the range \([−\frac{\pi}{2}, \frac{\pi}{2}]\) there is no early Big Crunch. The two straight asymptotic lines are at the values \( \pm Y_0 = \pm \frac{4}{5} \sqrt{\pi} \Gamma\left(\frac{1}{10}\right) / \Gamma\left(\frac{2}{5}\right) \).

In any case, the scale factor has zeros at \( \tau = \pm \frac{\pi}{2} + n \times \pi \) for \( n \in \mathbb{Z} \), yet there is also another possibility which has to be taken into account: an additional zero might or might not occur in the interval \([−\frac{\pi}{2}, \frac{\pi}{2}]\). This depends on the value of \( Y \). Given the form (5.35) of the solution, a zero of the scale factor can occur at a value \( \tau_0 \) which satisfies the equation:

\[
Y = \frac{4}{5} \cos^\frac{3}{10} (\tau_0) \csc(\tau_0) \left( \cos^2(\tau_0)^{9/10} \right)_{2F1} \left( \frac{1}{2}, \frac{9}{10}; \frac{3}{2}; \sin^2(\tau_0) \right) \tan^2(\tau_0) + 5 \equiv \hat{f}(\tau_0).
\]

The plot of the function \( \hat{f}(\tau) \) defined above is displayed in Fig. 7. We see that for

\[
|Y| \leq Y_0 \equiv \frac{4 \sqrt{\pi} \Gamma\left(\frac{1}{10}\right)}{\Gamma\left(\frac{2}{5}\right)}
\]

the candidate Big Bang and Big Crunch are at \( \tau = \pm \frac{\pi}{2} \), while for \(|Y| > Y_0\) the candidate Big Bang is at \( \tau = −\frac{\pi}{2} \), while the Big Crunch occurs earlier at:

\[
\tau_0 = \hat{f}^{-1}(Y).
\]

It is reasonable to expect a significantly different structure of the solution in the two cases.

**\( Y \neq 0 \) but less than critical.** In Fig. 8 we display the plot of the real and imaginary parts for the three functions composing the solutions for the less than critical case \( Y = 1 \). As we see the shape of the plots is no longer symmetric and imaginary parts are developed also by the scalar field and by the \( \exp[B]-factor \), yet the candidate Big Crunch occurs once again at the parametric time \( \tau = \frac{\pi}{2} \). Furthermore the scalar field, after climbing to some finite value, drops again to \( -\infty \) at the end of the life cycle of this Universe. The no longer symmetric phase-portrait of this solution is displayed in Fig. 9. The verification that the zeros of the scale factor are indeed singularities is done by inspecting the divergences of the curvature two-form in \( \pm \frac{\pi}{2} \). In this case it is much more difficult to calculate the coefficients of the asymptotic expansion, yet is sufficiently easy to determine the divergence order of the various components. We find:
Fig. 8. Plots of the real and imaginary parts of the scale factor, of the \( \exp[B]\)-factor and of the scalar field for the case of parameter \( Y = 1 \), which is non-zero but smaller than the critical value \( Y_0 \). In the three diagrams the solid line represents the real part, while the dashed line represents the imaginary part. The interval in which the three function are simultaneously real is still \([-\frac{\pi}{2}, \frac{\pi}{2}]\) as in the \( Y = 0 \) case. Yet the shape of the plots is no longer symmetric and also the \( \exp[B] \) factor and the scalar field start developing imaginary parts that instead are identically zero over the full range of \( \tau \) when \( Y = 0 \).
Fig. 9. In this figure we display the phase portrait of the solution defined by $Y = 1$. The axes are the scalar field $\Phi \equiv h$ and its derivative with respect to the cosmic time $V \equiv \partial T_c h$. The extremum of the potential is at $\Phi_0 = -\log[5]$. It is reached by the solution however with a non-vanishing velocity. The field also reaches vanishing velocity yet not at the extremum of the potential. Hence there is no fixed point and the trajectory is infinite with no fixed point. The symmetric shape of the $Y = 0$ trajectory is lost.

At $Y \neq 0$, differently from the $Y = 0$ case there are two different velocities of approach to infinity for the curvature components. Half of them go faster and half of them go slower. Yet the relevant point is that all of them blow up and certify that we are in presence of a true singularity, both at the beginning and at the end of time. The overall shape of the solution for the scalar field and for the scale factor is displayed in Fig. 10.

Overcritical $Y > Y_0$. When the parameter $Y$ is over critical we have a new zero of the scale factor which occurs at some $\tau_0$ in the fundamental interval $[-\pi/2, \pi/2]$. A practical way to deal with this type of solutions is to invert the procedure and use $\tau_0$ as parameter by setting $Y = f(\tau_0)$. As an illustrative example we choose $\tau_0 = \pi/6$ and we get:

$$Y_\bullet = f\left(\frac{\pi}{6}\right) = 4 \times 2^{4/5} \sqrt{3} + \frac{2}{5} \cdot 2 \cdot \frac{9}{10} \cdot \frac{3}{2} \cdot \frac{1}{4}. \quad (5.49)$$

The behavior of the real and imaginary parts of the three functions composing the solution for $Y = Y_\bullet$ is displayed in Fig. 11. The new character of the solution is immediately evident from such plots. The earlier zero of the scale factor is in correspondence with a divergence of the scalar field that now climbs from $-\infty$ to $+\infty$ as it is displayed in Fig. 12. The structure of the phase-portrait changes significantly with respect to the subcritical cases and it is displayed in Fig. 13. At the Big Bang point $\tau = -\pi/2$ the curvature components diverge just in the same way as in Eq. (5.48) namely the fastest approach to infinity is $O(1/(\tau+\pi/2)^{14/5})$. Instead at the new Big Crunch point $\tau = \pi/3$ the fastest diverging components of the curvature tensor have a
much stronger singularity, namely they diverge as $O\left(\frac{1}{(\tau + \pi/2)^7}\right)$. This further shows the clear-cut separation between less than critical and over critical solutions of the trigonometric type.

Apart from this finer structure the above detailed analysis has explicitly demonstrated the main point which we want to stress since it is somehow new in General Relativity. Notwithstanding the spatial flatness of the metric and notwithstanding the positive asymptotic behavior of the potential $V(h)$ that goes to $+\infty$ for large values of the scalar field $h$, the presence of a negative extremum of $V(h)$, (does not matter whether maximum or minimum) always implies a collapse of the Universe at a finite value of the cosmic or parametric time.

The Big Crunch collapse is the typical destiny of a closed Universe with positive spatial curvature. Therefore one is naturally led to inquiry whether such Universes as those discussed above have just the same causal structure as a closed Universe. To give an answer to such a question we consider the Particle and Event Horizons.

5.3.5. Particle and event horizons

Two important concept in Cosmology are those of Particle and Event Horizons. Given a metric of the form (1.1) let us rewrite it in polar coordinates:

$$ds^2 = \exp[B(\tau)] d\tau^2 - a^2(\tau) (dr^2 + r^2 d\Omega^2)$$

(5.50)

where, as usual, $d\Omega^2$ denotes the metric on a two-sphere, and let us consider the radial light-like geodesics defined by the equation:

$$0 = \exp[2B(\tau)] d\tau^2 - a^2(\tau) dr^2 \rightarrow \int_0^R dr = \int_{-\pi/2}^T d\tau \frac{\exp[B(\tau)]}{a(\tau)}.$$  (5.51)

From (5.51) it follows that at any parametric time $T$ after the Big Bang, the remotest radial coordinate from which we can receive a signal is given by:

$$R(T) = \int_{-\pi/2}^T d\tau \frac{\exp[B(\tau)]}{a(\tau)}.$$  (5.52)
Fig. 11. Plots of the real and imaginary parts of the scale factor, of the $\exp[\mathcal{B}]$-factor and of the scalar field for the case of parameter $Y = Y_\star$, which is overcritical $Y_\star > Y_0$. In the three diagrams the solid line represents the real part, while the dashed line represents the imaginary part. The interval in which the three functions are simultaneously real is now reduced to $[-\frac{\pi}{2}, \frac{\pi}{2}]$. In the real range the scalar fields climbs from $-\infty$ to $+\infty$. 
Fig. 12. In this figure we display plots (in parametric time) of the scale factor (solid line) and of the scalar field (dashed line) for the supercritical trigonometric type of solutions at $Y = Y_\bullet$. It is evident that in the finite parametric time interval $[-\frac{\pi}{2}, \frac{\pi}{3}]$ the Universe undergoes a Big Bang, a decelerated expansion and then a Big Crunch. At the same time the scalar field climbs from $-\infty$ to $+\infty$.

Fig. 13. In this figure we display the phase portrait of the solution defined by the supercritical value $Y = Y_\bullet$. The axes are the scalar field $\Phi \equiv h$ and its derivative with respect to the cosmic time $V \equiv \partial_\tau h$. The extremum of the potential is at $\Phi_0 = -\log[5]$. It is reached by the solution however with a non-vanishing velocity. The field also reaches vanishing velocity yet differently from the previous case before rather than after the extremum. This allows for the continuous climbing of the field up to $+\infty$.

and in any case the maximal physical value of such a radial coordinate is given by:

$$r_{\text{max}} = \int_{-\frac{\pi}{2}}^{T_{\text{max}}} d\tau \frac{\exp[B(\tau)]}{a(\tau)}$$  \hspace{1cm} (5.53)

where $T_{\text{max}}$ is the Big Crunch parametric time. It is therefore convenient to measure radial coordinates $r$ in fractions of this maximal one and measure the scale factor in fraction of the maximal one attained during time evolution:

$$a_{\text{max}} \equiv a(\hat{\tau}), \text{ where } \partial_\tau a(\tau)|_{\tau = \hat{\tau}} = 0.$$ \hspace{1cm} (5.54)

Setting:
we conclude that the farthest distance from which an observer can receive a signal at any instant of time is:

\[
\mathcal{P}(T) = \frac{a(T)}{a_{\text{max}}} \int_{\frac{-\pi}{2}}^{\pi} d\tau \frac{\exp[B(\tau)]}{a(\tau)}. \tag{5.56}
\]

By definition this distance is the Particle Horizon and defines the portion of Space that is visible by any Observer living at time \(T\). On the other hand the Event Horizon is the boundary of the Physical Space from which no signal will ever reach an Observer living at time \(T\) at any time of his future. In full analogy with Eq. (5.56) the Event Horizon is defined by:

\[
\mathcal{E}(T) = \frac{a(T)}{a_{\text{max}}} \int_{\frac{-\pi}{2}}^{\pi} d\tau \frac{\exp[B(\tau)]}{a(\tau)}. \tag{5.57}
\]

It is well known, (see for instance [39]) that in a matter dominated, closed Universe the Particle Horizon and the Event Horizon exactly coincide. This means that in such a Universe, the portion of space that is invisible to an observer living at time \(T\) will remain invisible to him also at all later times. Furthermore in such a Universe the Particle/Event Horizon contracts to zero exactly at the moment when the Universe reaches its maximum extension. An observer living at that time is completely blind and will stay blind all the rest of his life. In the Universes we have considered in this section things go quite differently since the Particle and the Event Horizon do not coincide and actually have a somehow opposite behavior. Plots of the Particle and Event horizon are shown in Fig. 14 for the three solutions of trigonometric type we have been considering. In all cases the Particle Horizon does not coincide with the Event Horizon. At the beginning the latter is larger than the former, which means that there is a portion of the invisible Universe which will reveal itself to the a given observer in his future. Then the Particle Horizon grows and the Event Horizon rapidly decreases. This means that as time goes on larger and larger becomes the visible Universe but also larger and larger the portion that will never reveal itself to an observer living at that time. In all cases before the end of its life the entire Universe becomes visible to an observer living at that time. This happens at relatively early times for supercritical solutions with \(Y > Y_0\).

This rather intricate structure is quite different from that of a Closed Universe with positive spatial curvature. Contrary to generally accepted lore, these solutions of Einstein–Klein–Gordon equations show that spatial flatness of the Universe should not lead us to automatically exclude the possibility of a final collapse into a singularity. We think that this is an important warning and for this reason we analyzed the case in-depth. A second motivation for our detailed analysis was the confirmation of the climbing, descending mechanism that strictly correlates the space–time singularities with the divergences of the scalar field that can only start and end up at infinity. Without a positive extremum the scalar field cannot stop at a fixed point and space–time has no other option than exploding and then collapsing.

In the next section we consider the much simpler and rather smoothly behaved hyperbolic solutions that occur when the potential has no extremum.
Fig. 14. The three plots in this figure respectively refer to the solutions $Y = 0$, $Y = 1$ and $Y = Y_\bullet$. In each plot the solid line represents the Scale Factor, the dashed line with longer dashes represents the Particle Horizon, while the dashed line with shorter and denser dashes represents the Event Horizon. In all cases the Event Horizon goes to zero faster than the Particle Horizon and invisible portion of the Universe becomes visible to the same observer at later times.
5.4. Hyperbolic solutions in the run away potential without finite extrema

When we choose $\omega = v^2$ the potential has no minimum, the solution of Eqs. (5.24) drastically simplifies and it is provided in terms of exponential functions. Explicitly we obtain:

$$U(\tau) = ae^{v\tau} + be^{-v\tau}, \quad (5.58)$$

$$V(\tau) = \left(e^{-v\tau} \left(e^{2v\tau}a + b\right) \left(c\left(e^{2v\tau}v^2 + dv^2 - 4e^{v\tau} \sqrt{e^{v\tau}a + be^{-v\tau}}\right)ight.ight.$$  

$$\left.\times \left(e^{v\tau}b + be^{-v\tau}(b - ae^{2v\tau})\right)\right) + 4F1\left(\begin{array}{ccc}
\frac{2}{5}, \frac{4}{5}, \frac{7}{5}; - \frac{ae^{2v\tau}}{b} \end{array}\right). \quad (5.59)$$

where $a, b, c, d$ are integration constants. Once inserted in the formula (5.21) for the physical fields the solution (5.58) produces a solution of the original equation upon implementation of the same constraint (5.27) as in the trigonometric case that we can solve with the same position, namely by setting $d = -pa$, $c = \rho b$. The final form of the solution depending on three parameters is the following one:

$$\Psi(\tau; a, b, \rho, v) = e^{-v\tau}\sqrt[e^{v\tau}a + be^{-v\tau}]\left(e^{2v\tau}a + b\right)\left(be^{2v\tau}v^2 - \rho v^2ight)$$  

$$- 4e^{v\tau}\sqrt[e^{v\tau}a + be^{-v\tau}]\times e^{v\tau}\sqrt[e^{v\tau}a + be^{-v\tau}]\left(b - ae^{2v\tau}\right)$$  

$$\times \left(e^{2v\tau}a + b\right)\left(1 + \frac{e^{2v\tau}a + b}{b}\right)^{4/5}2F1\left(\begin{array}{ccc}
\frac{2}{5}, \frac{4}{5}, \frac{7}{5}; - \frac{ae^{2v\tau}}{b} \end{array}\right), \quad (5.60)$$

$$a(\tau; a, b, \rho, v) = \frac{\Psi(\tau; a, b, \rho, v)}{(e^{2v\tau}a + b)v^2},$$

$$\exp[B](\tau; a, b, \rho, v) = \left(\frac{\Psi(\tau; a, b, \rho, v)}{e^{v\tau}a + be^{-v\tau})^{4/5}(e^{2v\tau}a + b)^2v^2\right),$n

$$\mathfrak{h}(\tau; a, b, \rho, v) = \log\left[\frac{(e^{2v\tau}a + be^{-v\tau})^{2/5}(e^{2v\tau}a + b)v^2}{\Psi(\tau; a, b, \rho, v)}\right]. \quad (5.61)$$

5.4.1. The simplest hyperbolic solution

The simplest solution of the hyperbolic type is obtained for the choice $a = 0$, $c = 0$, $b = 1$, $\rho = 1$, since in this case the hypergeometric function disappears and we simply get:

$$A(t, \nu) \equiv a\left(t - \frac{5 \log(\nu^2)}{6\nu}; 0, 1, 1, \nu\right) = \frac{5^{2/3}e^{-\nu^2/6}(1 + \frac{6\nu}{\nu})}{\nu^{4/3}},$$

$$\exp[B(t, \nu)] \equiv \exp\left[B\left(t - \frac{5 \log(\nu^2)}{6\nu}; 0, 1, 1, \nu\right)\right] = \frac{25(-1 + \frac{6\nu}{\nu})^2}{\nu^4},$$

$$\mathfrak{h}(t, \nu) \equiv \mathfrak{h}\left(t - \frac{5 \log(\nu^2)}{6\nu}; a, b, \rho, \nu\right) = \log\left[\frac{1}{5(-1 + \frac{6\nu}{\nu})} + 2\log(\nu)\right]. \quad (5.62)$$

The shift in the parametric time variable $\tau \rightarrow t - \frac{5 \log(\nu^2)}{6\nu}$ has been specifically arranged in such a way that $t = 0$ is a zero of the scale factor, namely corresponds to the Big Bang. Furthermore,
in this case, which involves only elementary transcendental functions, the relation between para-
metric and cosmic time can be explicitly evaluated. We have:

\[ T_c(t) \equiv \int_0^t dx \exp [B(x, \nu)] = \frac{25t}{\nu^4} - \frac{125e^{6\nu}}{3\nu^5} + \frac{125e^{12\nu}}{12\nu^5} + \frac{125}{4\nu^5}. \] (5.63)

This allows for a simple evaluation of the asymptotic behavior of both the scale factor and the
scalar field for asymptotic very late and very early times. We calculate the limit:

\[ \lim_{t \to \infty} \frac{\log[a(t, \nu)]}{\log[T_c(t)]} = \frac{1}{3}. \] (5.64)

This means that at late times, independently from the parameter \( \nu \) the scale factor behaves like
the cubic root of the cosmic time.

\[ a(T_c, \nu) \underset{T_c \to \infty}{\simeq} \text{const} \times T_c^{\frac{1}{3}}. \]

This corresponds to an equation of state of type (1.7) with \( w = 1 \). In view of Eqs. (1.5) this
means that at late times the predominant contribution to the energy density is the kinetic one, the
potential energy being negligible. Such a conclusion can be matched with the information on the
asymptotic behavior of the scalar field for late times. This latter can be worked in the following
way. As \( t \to \infty \) (for \( \nu > 0 \)) we have:

\[ T_c \underset{t \to \infty}{\simeq} \frac{125e^{12\nu}}{12\nu^5}, \] (5.65)

while for the scalar field we get:

\[ h(t, \nu) \underset{t \to \infty}{\simeq} -\frac{6t\nu}{5}. \] (5.66)

Combining the two results we get:

\[ h \underset{T_c \to \infty}{\simeq} -\frac{1}{2} \log[T_c] \] (5.67)

namely, the scalar field goes logarithmically to \(-\infty\) when plotted against cosmic time. Obviously
the value of the potential (5.19) at \( h = -\infty \) is zero and this explains the asymptotic dominance
of the kinetic energy.

To work out the behavior at very early times it is more complicated, yet we can predict it
by inspecting the behavior of the energy density and of the pressure. Inserting the form of the
solution and of the potential in Eq. (1.5) we obtain the parametric time behavior of the energy
density and of the pressure\(^9\):

\[ \rho = \frac{3\nu^8(-4\nu^2 + 2e^{6\nu}(2\nu^2 + 5) + e^{12\nu}(3\nu^2 - 5) - 5)}{15625(-1 + e^{6\nu})^6}, \] (5.68)

\[ p = \frac{3\nu^8(4\nu^2 - 2e^{6\nu}(2\nu^2 + 5) + e^{12\nu}(3\nu^2 + 5) + 5)}{15625(-1 + e^{6\nu})^6}. \] (5.69)

\(^9\) Note that to obtain this result we have calculated the derivative of the scalar field with respect to the cosmic time and
not with respect to the parametric time.
Expanding both functions in power series for $t \sim 0$ we obtain:

$$\rho \sim 0 \sim \nu^{-4} \frac{1728 t^6}{1728 t^6} - \frac{v^5}{2592 t^5} + O\left(\frac{1}{t^4}\right), \quad (5.70)$$

$$p \sim 0 \sim \nu^{-4} \frac{1728 t^6}{1728 t^6} - \frac{13 v^5}{12960 t^5} + O\left(\frac{1}{t^4}\right). \quad (5.71)$$

Both the pressure and the energy density diverge as $1/t^6$ plus subleading singularities; the identity of the coefficient in the leading pole of both expansions implies that also at very early times the effective equation of state is

$$p = \rho \Leftrightarrow w = 1 \quad (5.72)$$

which implies the following behavior for the scale factor:

$$a(T_c, v) \underset{T_c \rightarrow 0}{\sim} \text{const} \times T_c^{\frac{1}{3}}. \quad (5.73)$$

With same technique we can also work out the asymptotic behavior of the scalar field in the origin of time:

$$h \underset{T_c \rightarrow 0}{\sim} -\frac{1}{3} \log[T_c]. \quad (5.73)$$

In Fig. 15 we present the plots of an explicit example of such solutions. Finally in Fig. 16 we present the phase portrait for this type of solutions of the hyperbolic type. We compare the phase portrait of the simple solutions we have analyzed above with the phase portrait of the generic solutions that involve also the hypergeometric term. The quality of the picture is essentially the same, yet there is an important critical difference concerning the asymptotic behavior.

### 5.4.2. Hyperbolic solutions displaying the second asymptotic behavior for late times

According to the analysis of [1] and of the previous literature [16–18], when the potential is of the exponential type considered in this paper, namely

$$V(\phi) = \sum_{k=1}^{n} V_{0k} e^{2 \gamma k \phi}, \quad \gamma_k > \gamma_{k+1} \quad (5.74)$$

there are two possible different asymptotic behaviors of the scale factor in the vicinity of a Big Bang or of a Big Crunch. One behavior is universal and it is the one we have met in the previous example of the simplest hyperbolic solution, namely:

$$a(T_c) \sim T_c^{\frac{1}{3}} \Leftrightarrow w = 1 \quad \text{kinetic asymptopia}. \quad (5.75)$$

The universal asymptopia is the only one available both at the beginning and at the end of time when the solution for the scalar field is climbing. As already stressed it corresponds to a complete dominance of the kinetic energy of the scalar field with respect to its potential one. On the other hand if the scalar is descending there are a priori two possible asymptopia available: in addition to the universal kinetic one (5.75), there is also the following additional one:

$$a(T_c) \sim T_c^{\frac{1}{3} \gamma_{dom}} \Leftrightarrow w = 2 \gamma_{dom}^{-2} - 1 \quad \text{potential asymptopia}, \quad (5.76)$$

where $\gamma_{dom}$ is the coefficient of the dominant exponential appearing in the potential, once this latter is written in its normal form (5.74) by means of the replacement (2.2). For descending
Fig. 15. Here we present the behavior of the scale factor and of the scalar field for the simplest of the hyperbolic type solutions ($\omega = \nu^2$) of the cosmological model based on the potential of Eq. (5.19). The analytic form of the solution is given in Eq. (5.62). For the plot we have chosen $\nu = \frac{1}{4}$. In the first graph, describing the scale factor, the solid line is the actual solution while the dashed curves are of the form $\alpha_1 T^{1/3}_{c}$ with two different coefficient $\alpha_1 = \frac{32}{101} T^{1/3}$ and $\alpha_2 = \frac{35}{101} T^{1/3}$. The first curve is tangential to the solution at $T_c \to 0$ while the second is tangential to the solution at $T_c \to \infty$. The same style of presentation is adopted in the second picture. Here we plot the scalar field against the logarithm of the cosmic time. The two dashed straight lines represent the curves $-\frac{1}{3} \log[T_c]$, and $-\frac{1}{2} \log[T_c]$. The first is tangential to the solution at $T_c \to 0$, the second is tangential to the solution at $T_c \to \infty$.

Scalars that tend to $-\infty$, the dominant exponential is that with the smaller positive $\gamma_k$. In our case $\gamma_{dom} = \frac{3}{2}$ so that the second asymptopia available in the case of descending solutions, is:

$$a(T_c) \sim T_c^{\frac{3}{2}} \quad \Leftrightarrow \quad w = -\frac{1}{9} \quad \text{potential asymptopia.} \quad (5.77)$$

The simplest asymptotic solution described in the previous section does not take advantage of the second possibility for its asymptotic behavior: both at very early and at very late times the scale factors goes as $T_c^{\frac{3}{2}}$. This is no longer the case for the solutions with parameter $\alpha \neq 0$ where the contribution from the hypergeometric function is switched on. Indeed we have verified that for all such solutions the scale factor goes as $T_c^{\frac{1}{3}}$ near the initial singularity but diverges as $T_c^{\frac{3}{2}}$. 
Fig. 16. In this picture we present some trajectories making up the phase portrait for the solutions of hyperbolic type presented in this section. In the upper figure we use only the simple solution involving elementary functions. In the lower figure we utilize also solutions involving the hypergeometric term. The quality of the portraits is just the same.

for late times. An example of such a behavior is provided in the plots of Fig. 17. The effective equation of state at late times corresponds to a negative pressure although to a very weak one \( w = -\frac{1}{9} \). This negative pressure is provided by a small dominance of the potential energy over the kinetic one and causes an indefinite expansion of the Universe slightly stronger than that of a matter dominated Universe \( (T_c)^{3/2} \) yet very far from an exponential one.

6. Analysis of the cosh-model

We come finally to one of the main questions posed in this paper, namely how much valuable are the exact solutions of integrable one-field cosmologies as simulations or approximations of the solutions of the physical cosmological models, in particular of those provided by consistent one-field truncations of supergravity theories. These latter, as we have emphasized, are mostly not integrable, as we have already seen and as we are going to show further in the last section. In other words the present section is devoted to assess the viability of the minimalist approach.

Assuming that cosmological models derived from supergravity are not integrable, we still would like to ascertain whether integrable potentials that are similar to actual physical ones have solutions that simulate in a reasonable way the solutions of the supergravity derived models.

In this contest a primary example is provided by the cosh model already introduced in previous sections and defined by the Lagrangian (3.28) that depends on the three parameters \( q, p \) and \( \mu^2 \). As we already pointed out, the family of cosh-models defined by the Lagrangian (3.28) with the potential

\[
V(\phi) = -\mu^2 \cosh[p\phi], \quad \mu^2 = \text{either positive or negative}
\]
Fig. 17. In this picture we present the plots of the scale factor (first plot) and of the scalar field (second plot) in the hyperbolic solution with parameters \( a = -1, b = \rho = v = 1 \). The asymptotic behavior for late cosmic times of the scale factor is visualized by the third plot of the logarithm of the function divided by the logarithm of its argument. It is quite evident that this ratio goes rapidly to \( 0.75 = \frac{3}{4} \).

includes two integrable cases when \( \frac{p}{\sqrt{3q}} = 1 \) or \( \frac{p}{\sqrt{3q}} = \frac{2}{3} \). At the same time the case \( p = 1, q = 1 \), which by no means is integrable, corresponds to the consistent one-field truncation of the STU model. Furthermore, other instances of the same Lagrangian are expected to appear in the consistent truncation of other Gauged Supergravity models.

Hence the model (3.28) is a perfect testing ground for the questions we have posed.

An important conclusion that was reached in section [1] is that the qualitative behavior of solutions is dictated by the type of critical points possessed by the equivalent first order dynamical system that, on its turn is just dictated by the properties of extrema of the potential. Hence the first question that arises in connection with our model (3.28) is whether its critical points are always of the same type or fall into different classes depending on the parameters \( p, q \).

To this effect we begin by summarizing some of the results of sections [1] in a language less mathematically oriented and closer to the jargon of the supergravity community. Furthermore in such a summary we utilize the customary physical normalizations of Friedman equations, recalled in Eqs. (1.3), rather than the normalizations and notations of [1] that are less familiar to cosmologists.

6.1. Summary of the mathematical results on the structure of Friedman equations and on the qualitative description of their solutions

Choosing the standard gauge \( B = 0 \) and considering the standard form (1.3) of Friedman equations, the main crucial observation that was put forward in [1], is that the logarithm of the scale factor \( A(t) \) is a cyclic variable since it appears only through the Hubble function, namely covered by a derivative. The next crucial observation was that the second order differential system (1.3) can be rewritten in two different ways as a system of first order ordinary differential equations for two variables. These rewritings were named irreducible subsystems and it was ad-
vocated that each of them, when solved, generates solutions of the initial second order system (1.3). Adopting such a language allowed for the use of some powerful theorems that can predict the general qualitative behavior of the solutions of Friedman equations once the potential $V(\phi)$ is specified.

The two subsystems are hereby rewritten in the standard notation of General Relativity and Cosmology:

**Subsystem I.** The first subsystem uses as independent variables the scalar field $\phi(t)$ and its time derivative $\dot{\phi}(t)$, which is renamed $v(t)$. Hence one writes:

$$
\dot{\phi} = v, \\
\dot{v} = -3\sigma v \sqrt{\frac{1}{3} v^2 + \frac{2}{3} V(\phi)} - V'(\phi),
$$

(6.2)

where $\sigma = \pm 1$ takes into account the two branches of the square root in solving the quadratic equation for the Hubble function. One has to exclude those branches of the solutions of (6.2) that satisfy the following conditions:

$$
\frac{1}{3} v^2 + \frac{2}{3} V(\phi) = 0 \quad \text{if} \quad V(\phi) \neq 0, \\
\frac{1}{3} v^2 + \frac{2}{3} V(\phi) < 0.
$$

(6.3)

The remaining solutions are named admissible. When an admissible solution of Eqs. (6.2) is given, namely when the functions $\phi(t)$ and $v(t)$ have been determined, the Hubble function $H(t)$ is immediately obtained,

$$
H(t) = \sigma \sqrt{\frac{1}{3} v^2(t) + \frac{2}{3} V(\phi(t))}
$$

(6.4)

and, by means of a further integration, one obtains also the scale factor:

$$
a(t) = \exp \left[ \int H(t) \, dt \right].
$$

(6.5)

**Subsystem II.** The second subsystem uses as independent variables the scalar field and the Hubble function. So doing we are lead to the following first order equations:

$$
\dot{\phi} = \sigma (3H^2 - 2V(\phi))^\frac{1}{2}, \quad \sigma = \pm 1, \\
\dot{H} = -(3H^2 - 2V(\phi)).
$$

(6.6)

As in the first instance of irreducible subsystem also here we have to exclude inadmissible branches of solutions to Eqs. (6.6), namely those that satisfy the following conditions:

$$
3H^2 - 2V(\phi) = 0 \quad \text{if} \quad V'(\phi) \neq 0, \\
3H^2 - 2V(\phi) = 0 < 0.
$$

(6.7)

It is easily verified that all equations of the original Friedman system (1.3) follow from either one of the subsystems (6.2) and (6.6).

In mathematical language, both subsystems (6.2) and (6.6) are nonlinear autonomous first-order ordinary differential equations over a two-dimensional Euclidean plane, namely either $\mathbb{R}^2 \ni (\phi, v)$ or $\mathbb{R}^2 \ni (\phi, H)$. 
The mathematical theory of planar dynamical systems is highly developed and allows for a qualitative analysis of both the local and the global behavior of their phase portraits, namely of their trajectories (also named orbits). According to such a theory a generic planar system is very regular: it admits only a few different types of trajectories and limit sets. Explicitly we can have:

(a) periodic orbits, named also cycles,
(b) heteroclinic orbits that connect two different critical points of the system,
(c) homoclinic orbits that start from a critical point and return to it at the end of time,
(d) trajectories that connect the point at infinity of $\mathbb{R}^2$ with a fixed point.

As a result no planar dynamical system can be chaotic. This property distinguishes planar systems very strongly from dynamical systems in dimensions higher than two, where various chaotic regimes are generically allowed.

From these considerations it follows that one-field cosmologies are not chaotic and one can obtain a qualitative understanding of their solutions. In the integrable case the solutions can also be worked out analytically. The relevant point is that such analytical solutions can be taken as trustworthy models of the behavior of the solutions also for entire classes of potentials whose integrable representatives occur only at very special values of their parameters.

6.1.1. Subsystem I: Qualitative analysis

To illustrate in a concrete manner these general ideas we choose to work with the subsystem I. Fixed points of the subsystem (6.2) are defined by the following equations

$$v_0 = 0, \quad V'(\phi_0) = 0$$

(6.8)

and are admissible if they satisfy the condition:

$$V(\phi_0) \geq 0.$$  

(6.9)

In plain physical words a fixed point of this dynamical system is just a vacuum solution of scalar coupled gravity, namely a constant configuration of the scalar field that is an extremum of the potential. At the same time the space–time metric is either the Minkowski metric if $V(\phi_0) = 0$ or the de Sitter metric if $V(\phi_0) = \Lambda > 0$.

From the dynamical system point of view, if the condition (6.9) is not fulfilled, then the subsystem does not possess fixed points at all, i.e. all its phase space points are regular. Without fixed points a nonlinear system admits only monotonic solutions, which can also blow up in a finite time. From the physical point of view an extremum of the potential at a negative value of the potential corresponds to an anti-de Sitter space, yet anti-de Sitter, differently from de Sitter admits no representation in terms of flat constant time slices, which is our initial assumption. Hence the only consequence can be a blowing up solution with a Big Bang followed by a Big Crunch.

6.1.1.1. Linearization of the first order system in a neighborhood of a fixed point

Let us consider the linearization of the first order system in a neighborhood of the fixed point by setting:

$$\phi = \phi_0 + \Delta \phi,$$

$$v = \Delta v.$$  

(6.10)

To first order in the deviations we obtain
\[ \Delta \dot{\phi} = \Delta v, \]
\[ \Delta \dot{v} = -\sigma \sqrt{6V(\phi_0)} \Delta v - V''(\phi_0) \Delta \phi + \text{h.o.t.} \]  
(6.11)

where the abbreviation h.o.t. means higher order terms.

As explained in [1], the eigenvalues of the linearization matrix:
\[ M = \begin{bmatrix} 0 & 1 \\ -V'' & -\sigma \sqrt{6V_0} \end{bmatrix} \]  
(6.12)

namely:
\[ \lambda_{\pm} = \frac{1}{\sqrt{2}} \left( -\sigma \sqrt{3V_0} \pm \sqrt{3V_0 - 2V''} \right) \]  
(6.13)

characterize the type of the corresponding critical points and consequently define the phase portrait of the linearization. The main theorem that allows to predict the qualitative behavior of the solution states the following. In the case the fixed point is hyperbolic, namely both eigenvalues have a non-vanishing real part, (i.e. \( \text{Re}(\lambda_{\pm}) \neq 0 \)) the phase portraits of the nonlinear system and of its linearization are diffeomorphic in a finite neighborhood of the hyperbolic fixed point. Hence the analysis of the linearization gives a valuable information about the phase portrait of the original nonlinear system we are interested in.

We have in particular the following classification of non-degenerate fixed points for which, by definition, the linearization matrix has no zero eigenvalue:

### 6.1.1.2. Classification of fixed point types

(a) **Saddle.** When the two real eigenvalues have opposite sign \( \lambda_+ > 0, \lambda_- < 0 \) or \( \lambda_+ < 0, \lambda_- > 0 \).
(b) **Node.** When the two real eigenvalues have the same sign \( \lambda_+ > 0, \lambda_- > 0 \) or \( \lambda_+ < 0, \lambda_- < 0 \).
(c) **Improper node.** When the two eigenvalues coincide \( \lambda_+ = \lambda_- \), yet the linearization matrix (6.12) is not diagonal.
(d) **Degenerate node.** When the linearization matrix (6.12) is proportional to identity.
(e) **Focus.** When the two eigenvalues have non-vanishing both the imaginary part and the real part and are complex conjugate to each other \( \lambda_{\pm} = x \pm iy \).
(f) **Center.** When the two eigenvalues are purely imaginary and conjugate to each other.

For each of these fixed point types the trajectories have a distinct behavior that we are going to illustrate by means of our concrete example, namely the **Cosh model** analyzed in the present section.

Furthermore we should recall the result that the subsystem (6.2) has no periodic trajectories according to Dulac’s criterion since:
\[ \frac{\partial \dot{\phi}}{\partial \phi} + \frac{\partial \dot{v}}{\partial v} = -2\sigma \frac{v^2 + V(\phi)}{\sqrt{\frac{1}{3}v^2 + \frac{2}{3}V(\phi)}} \]  
(6.14)

does not change sign over the whole two-dimensional plane. Thus, we are led to the conclusion that this subsystem can have only fixed points (i.e. vacuum solutions) as well as heteroclinic/homoclinic orbits, orbits connecting infinity with a fixed point or orbits connecting infinity with infinity in the case of fixed points of the saddle type.
As we know the case \( p = 1, q = 1 \) is the one which appears in the non-abelian gauging of the STU model and the case \( p = 1, q = 3 \) can be obtained in the \( S^3 \) model by means of an abelian gauging.

Although the \( \mathcal{N} = 2 \) case is not integrable, yet it belongs to same subclass (Node) as the integrable case \( \frac{p}{\sqrt{3q}} = \frac{2}{3} \). Hence we can probably learn about its behavior from an analysis of the integrable case close to it.

The natural question which we have posed, namely how much do the solutions of the physically relevant cosh-models depart from the exact solutions of the integrable members of the same family can now be partially answered. As we just stressed, the physically relevant \( \mathcal{N} = 2 \) case is of the Node type so that any integrable model with the same type of fixed point would just capture all the features of the physical STU model. The other integrable case \( \frac{p}{\sqrt{3q}} = 1 \) is instead of the Focus type, therefore it has a little bit more of structure with respect to the STU-model. Any other supersymmetric one-field model with a Cosh potential of the focus type, although not integrable might be well described by the \( \frac{p}{\sqrt{3q}} = 1 \) integrable member of the family (3.28).

### 6.2. Normal form of the Cosh-model

In this spirit let us first show how the Cosh-model can be put into a normal form, displaying a unique parameter \( \omega \) whose value will determine the type of fixed point and, for two special choices, yield two integrable models.

To this effect we introduce the following rescaled fields and variables:

\[
\begin{align*}
\h(t) &= \frac{\phi(\tau)}{\sqrt{q}}, \\
t &= \frac{\sqrt{2} \tau}{\mu}, \\
A(t) &= A(\tau), \\
\omega &= \frac{p}{\sqrt{3q}}.
\end{align*}
\] (6.15)

In terms of these new items, the effective Lagrangian (3.28) becomes

\[
L = e^{3A(\tau) - B(\tau)} \left\{ -\frac{3}{2} A'(\tau)^2 + \frac{1}{2} \phi'(\tau)^2 \mp 2e^{2B(\tau)} \cosh[\omega \phi(\tau)] \right\},
\] (6.16)

where the sign choice distinguishes two drastically different systems. The first choice yields a positive definite potential with an absolute minimum that allows for a stable de Sitter vacuum, while the second yields a potential unbounded from below with an absolute maximum. If we choose the gauge \( B = 0 \), the field equations of this system, including the Hamiltonian constraint can be written as the following three Friedman equations:

\[
\begin{align*}
\frac{a'(\tau)^2}{a(\tau)^2} - \frac{1}{3} \phi'(\tau)^2 \mp \frac{4}{3} \cosh[\omega \phi(\tau)] &= 0, \\
\frac{2}{3} \phi'(\tau)^2 \mp \frac{4}{3} \cosh[\omega \phi(\tau)] + \frac{a''(\tau)}{a(\tau)} &= 0, \\
\pm 2\omega \sinh[\omega \phi(\tau)] + \frac{3a'(\tau) \phi'(\tau)}{a(\tau)} + \phi''(\tau) &= 0
\end{align*}
\] (6.17)

where \( a(\tau) = \exp[A(\tau)] \).

### 6.3. The general integral for the case \( \omega = \sqrt{3} \)

In the integrable case \( \omega = \sqrt{3} \), by means of the integrating transformation described in [1] we obtain the following general solution of Eqs. (6.17) depending on three parameters, the scale \( \lambda \)
and the two angles $\psi$ and $\theta$, which applies to the case of the positive potential (upper choice in Eq. (6.16)):

$$a_+(\tau) = \sqrt[3]{\lambda} \cos(\psi) \cosh(\sqrt{3}\tau) + \lambda \sinh(\sqrt{3}\tau) \cos^2(\theta - \sqrt{3}\tau) \sin^2(\psi),$$

$$\phi_+(\tau) = \frac{1}{\sqrt{3}} \log \left[ \frac{\cos(\psi) \cosh(\sqrt{3}\tau) - \cos(\theta - \sqrt{3}\tau) \sin(\psi) + \sinh(\sqrt{3}\tau)}{\cos(\psi) \cosh(\sqrt{3}\tau) + \cos(\theta - \sqrt{3}\tau) \sin(\psi) + \sinh(\sqrt{3}\tau)} \right].$$

(6.18)

(6.19)

For the negative potential (lower choice in Eq. (6.16)) we find instead:

$$a_-(\tau) = \sqrt[3]{\lambda^2} \left( \cosh^2(\psi) \sin^2(\theta - \sqrt{3}\tau) - \sin(\sqrt{3}\tau) \cos(\sqrt{3}\tau) \sin(\psi) \right)^2,$$

$$\phi_-(\tau) = \frac{1}{\sqrt{3}} \log \left[ \frac{\sin(\sqrt{3}\tau) + \cosh(\psi) \sinh(\theta - \sqrt{3}\tau) - \cos(\sqrt{3}\tau) \sin(\psi)}{\sin(\sqrt{3}\tau) - \cosh(\psi) \sinh(\theta - \sqrt{3}\tau) - \cos(\sqrt{3}\tau) \sin(\psi)} \right].$$

(6.20)

(6.21)

One important observation is the following. In the case of the positive potential, by choosing the parameters $\psi = 0$, $\theta = 0$ we obtain the very simple solution:

$$a_0(\tau) = \exp \left[ \frac{2\tau}{3^{2/3}} \right], \quad \phi = 0 \Rightarrow \hbar = 0.$$ 

(6.22)

This is the de Sitter solution where the scalar field is stationary at its absolute minimum and the scale factor grows exponentially. Such a solution is ruled out in the case of the negative potential which does not allow for any static scalar field solution.

6.3.1. The second Hamiltonian structure

We have explicitly integrated the integrable cosmological model and we have obtained its general integral. We can address the question why is it integrable? The answer is that it admits not just one rather two functionally independent conserved Hamiltonians. Examining their structure is worth doing since it provides hints about the underlying properties of the field theory that might be responsible for the emergence of integrability at the cosmological level. Consider then the Lagrangian and the Hamiltonian for the model under consideration:

$$\mathcal{L}_0 = \exp[3A] \left\{ \frac{q}{2} \dot{h}^2 - \frac{3}{2} \dot{A}^2 - \mu^2 \cosh[3h] \right\},$$

$$\mathcal{H}_0 = \exp[3A] \left\{ \frac{q}{2} \dot{h}^2 - \frac{3}{2} \dot{A}^2 + \mu^2 \cosh[3h] \right\}.$$ 

(6.23)

By means of direct evaluation we can easily check that the following two functionals:

$$\mathcal{H}_1 = -\frac{1}{2} e^{3A} \left( 2\mu^2 \cosh^2 \left( \frac{3h}{2} \right) + 3\dot{A}^2 \cosh^2 \left( \frac{3h}{2} \right) + 3 \sinh^2 \left( \frac{3h}{2} \right) \dot{h}^2 \right. \left. + 3 \sinh(3h) \dot{A} \dot{h} \right),$$

(6.24)

$$\mathcal{H}_2 = \frac{1}{2} e^{3A} \left( -2\mu^2 \sinh^2 \left( \frac{3h}{2} \right) + 3\dot{A}^2 \sinh^2 \left( \frac{3h}{2} \right) + 3 \cosh^2 \left( \frac{3h}{2} \right) \dot{h}^2 + 3 \sinh(3h) \dot{A} \dot{h} \right)$$

(6.25)

satisfy the following conditions:
\[ \mathcal{H}_1 + \mathcal{H}_2 = \mathcal{H}_0, \]  
(6.26)

\[ \frac{d}{d\tau} \mathcal{H}_{1,2} = 0 \quad \text{upon use of field equations from Lagrangian.} \]  
(6.27)

Hence \( \mathcal{H}_{1,2} \) are the two conserved Hamiltonians that guarantee the integrability of the system. As we know, the actual solution of Friedman equations is obtained by enforcing also the constraint \( \mathcal{H}_0 = 0 \) so that for our general solution (6.19)–(6.21) we have: \( \mathcal{H}_1 = -\mathcal{H}_2. \)

### 6.4. The general integral in the case \( \omega = \frac{2}{\sqrt{3}} \)

In this case as in the previous one the solution can be obtained by means of the same substitution described [1], yet with respect to the previous case there is one relevant difference. In this case the gauge \( B = 0 \) cannot be chosen and there is a difference between the cosmic time \( t_c \) and the parametric time \( \tau. \) The form of the space–time metric is the following:

\[ ds^2 = \exp[-2A(\tau)]d\tau^2 - \exp[2A(\tau)]d\vec{x}^2 \]  
(6.28)

corresponding to the gauge \( B(\tau) = -A(\tau). \) In principle the general integral depends on three integration constants, but in this case one of them can be immediately reabsorbed into a shift of the parametric time and thus we are left only with two relevant constants.

At the end of the computations the general integral can be written in a very suggestive and elegant form in terms of the four roots of a quartic polynomial. Precisely we have:

\[ A(\tau) = \log \left[ \frac{2\sqrt[3]{(\tau - \lambda_1)(\tau - \lambda_2)(\tau - \lambda_3)(\tau - \lambda_4)}}{\sqrt{3}} \right], \]  
(6.29)

\[ \phi(\tau) = \frac{1}{4} \sqrt{3} \log \left( \frac{(\tau - \lambda_3)(\tau - \lambda_4)}{(\tau - \lambda_1)(\tau - \lambda_2)} \right), \]  
(6.30)

\[ B(\tau) = -A(\tau). \]  
(6.31)

One important caveat, however, is the following. Let us name \( \text{Re} \lambda_{\text{min}}, \text{Re} \lambda_{\text{max}} \) the smallest and the largest of the real parts of the four roots. The functions (6.31) provide an exact solution of the Friedman equations under two conditions:

(A) The parametric time \( \tau \) is either in the range \( \text{Re} \lambda_{\text{max}} \leq \tau \leq +\infty \) or in the range \( \text{Re} \lambda_{\text{min}} \geq \tau \geq -\infty. \)

(B) The four roots satisfy the following constraint:

\[ 2\lambda_1\lambda_2 - \lambda_3\lambda_2 - \lambda_4\lambda_2 - \lambda_1\lambda_3 - \lambda_1\lambda_4 + 2\lambda_3\lambda_4 = 0. \]  
(6.33)

When \( \tau \) is in the range \( \text{Re} \lambda_{\text{min}} \leq \tau \leq \text{Re} \lambda_{\text{max}}, \) the expressions (6.31) do not satisfy Friedman equations.

Solving (6.33) explicitly and replacing such solution back into the expression (6.31) of the fields we obtain the general integral depending on three parameters. As we already stressed one of these three parameters can always be reabsorbed by means of a shift of the parametric time coordinate \( \tau, \) as it is evident from the form (6.31). A convenient way of taking into account these gauge fixing is provided by solving the constraint (6.33) in terms of two real parameters \((\alpha, \beta),\) as it follows:
\[
\lambda_1 = \alpha - \sqrt{2}\sqrt{\alpha^2 - \beta}, \\
\lambda_2 = \alpha + \sqrt{2}\sqrt{\alpha^2 - \beta}, \\
\lambda_3 = -\alpha - \sqrt{2}\sqrt{\alpha^2 + \beta}, \\
\lambda_4 = \sqrt{2}\sqrt{\alpha^2 + \beta - \alpha},
\]

which greatly facilitates the discussion since depending on whether \(|\beta| < \alpha^2\) or \(|\beta| > \alpha^2\), we either have four real roots and hence four zeros of the scale factor or two real roots and two complex conjugate ones. If \(|\beta| = \alpha^2\) we have only real roots but three of them coincide and the fourth is different. Finally if both \(\alpha\) and \(\beta\) vanish the four roots coincide and the corresponding solution degenerates into the de Sitter solution. Let us substitute explicitly the values (6.34) into (6.31) and obtain the general integral in the form:

\[
a(\tau, \alpha, \beta) \equiv \exp\left[ A(\tau, \alpha, \beta) \right] = 2^{4/3} \sqrt{\alpha^4 - 6\tau^2\alpha^2 + 8\beta\tau\alpha + \tau^4 - 4\beta^2}, \\
\phi(\tau, \alpha, \beta) = \frac{1}{4} \sqrt{3} \log \left( \frac{\alpha^2 - 2\tau\alpha - \tau^2 + 2\beta}{\alpha^2 + 2\tau\alpha - \tau^2 - 2\beta} \right), \\
\exp\left[ B(\tau, \alpha, \beta) \right] = \frac{\sqrt{3}}{2^{4/3} \sqrt{\alpha^4 - 6\tau^2\alpha^2 + 8\beta\tau\alpha + \tau^4 - 4\beta^2}},
\]

which will be very useful in the discussion of the of the properties of solutions in the vicinity of the fixed point which for this case happen to be of the Node-type.

**6.5. Discussion of the fixed points**

Inserting the potential \(2 \cosh(\omega \phi)\) into the formulae (6.12) and (6.13) for the linearization matrix and for its eigenvalues we immediately obtain:

\[
\mathcal{M} = \begin{bmatrix} 0 & 1 \\ -2\omega^2 & \sqrt{12} \end{bmatrix} \Rightarrow \lambda_{\pm} = -\sqrt{3} \pm \sqrt{3 - 2\omega^2}.
\]

From this we learn that for \(0 < \omega < \sqrt{\frac{3}{2}}\) the fixed point at \(\phi = 0\) is of the Node type. For \(\omega = \sqrt{\frac{3}{2}}\) the fixed point is exactly of the improper node type, while for \(\omega > \sqrt{\frac{3}{2}}\) it is always of the focus type. What this implies for the behavior of the scalar field we will shortly see.

One important point to stress concerns the initial conditions that have always got to be of the Big Bang type (initial singularity). The exact solutions of the integrable case are very instructive in this respect. Fixing the Big Bang initial condition \(a(0) = 0\) at a finite reference time \((\tau = 0)\) from (6.19) we obtain a relation between the angle \(\psi\) and the angle \(\theta\)

\[
\cos(\psi) - \cos^2(\theta) \sin^2(\psi) = 0
\]

which inserted back into the formula (6.19) for the scalar field implies \(\phi(0) = \infty\). This is not a peculiarity of the integrable model, rather it is a general fact. The zeros at a finite time of the scale factor are always in correspondence with a singularity of the scalar field. Hence the only initial condition that can be fixed independently at the Big Bang is the initial velocity of the field \(\dot{\phi}(0)\).
Emerging from the initial singularity at $\pm \infty$ the scalar field can flow to its fixed point value, namely zero (if the potential is positive) and the way it does so depends on the fixed-point type. However it can also happen that before reaching the fixed point the scalar field goes again to $\pm \infty$. In this case we have a blow-up solution where, at a finite time the scale factor goes again to zero and we have a Big Crunch. This happens only in the case the extremal point of the potential (either minimum or maximum) is negative. If the extremal point of the potential is positive we always have an asymptotic de Sitter destiny of the considered Universe.

6.6. Behavior of solutions in the neighborhood of a Node critical point

The best way to discuss the quality of solutions is by means of the so-called phase portrait of the dynamical system, where we plot the trajectories of the solution in the plane $\{\phi, v\}$. As we have already emphasized there are solutions that go to the fixed point $\{\phi, v\} = \{0, 0\}$ and solutions that never reach it. The solutions that reach the fixed point have a universal type of behavior in its neighborhood which we now describe for the case of the Node. We illustrate such a behavior by means of the exact solutions of the integrable case $\omega = \frac{2}{\sqrt{3}}$. The two eigenvectors corresponding to the two eigenvalues in Eq. (6.38) are:

$$v_\pm = \begin{cases} \{\frac{\sqrt{3} - 2\omega^2 - \sqrt{3}}{2\omega^2}, 1\} & \text{if } \omega = \frac{2}{\sqrt{3}} \\ \{\frac{\sqrt{3} - 2\omega^2 + \sqrt{3}}{2\omega^2}, 1\} & \text{if } \omega = \frac{1}{\sqrt{3}} \end{cases}$$

According to theory all solutions of the differential system approach the fixed point along the eigenvector vector $\left\{\frac{\sqrt{3} - 2\omega^2 + \sqrt{3}}{2\omega^2}, 1\right\}$ corresponding to the eigenvalue of largest absolute value, except a unique, exceptional one, named the separatrix that approaches the fixed point along the other eigenvector $\left\{-\frac{\sqrt{3}}{2}, 1\right\}$. This can be checked analytically computing a limit. Let us define the function:

$$T(\tau, \alpha, \beta) = \frac{\partial_\tau \phi(\tau, \alpha, \beta) a(\tau, \alpha, \beta)}{\phi(\tau, \alpha, \beta)} = \frac{\partial_t \phi(t_e)}{\phi(t_e)}$$

which, due to the metric (6.28), represents the ratio of the logarithmic derivative of the scalar field with respect to the cosmic time. By explicit calculation we find:

$$\lim_{\tau \to \pm \infty} T(\tau, \alpha, \beta) = \pm \frac{2}{\sqrt{3}}$$

$$\lim_{\tau \to \pm \infty} T(\tau, \alpha, 0) = \pm \frac{2}{\sqrt{3}}$$

$$\lim_{\tau \to \pm \infty} T(\tau, 0, \beta) = \pm \frac{4}{\sqrt{3}}$$

We conclude that the separatrix solution is given by the choice $\alpha = 0$, all the other solutions being regular. An inspiring view of the phase portrait is given in Fig. 18.

6.7. Analysis of the separatrix solution

Because its exceptional status, the separatrix solution is worth to be analyzed in some detail. Explicitly we have:
Fig. 18. In this figure are shown the \((\phi, v)\) trajectories corresponding to six solutions of the \(\omega = \frac{2}{\sqrt{3}} \) Cosh-model with parameters \((\alpha, \beta) = (0, \frac{1}{2}), (2, 0), (5, 0), (3, -3), (2, 5), (1, -18)\). The case \((0, \frac{1}{2})\) corresponds to the separatrix which approaches the fixed point along the tangent vector \(\left\{ -\sqrt{3}\frac{4}{3}, 1 \right\} \). All the other solutions approach the fixed point along the tangent vector \(\left\{ -\sqrt{3}, 1 \right\} \). The two dashed lines represent these two tangent vectors.

\[
a(\tau) = \frac{2\sqrt{\frac{4}{3}} - 1}{\sqrt{3}}, \quad \phi(\tau) = \frac{1}{4\sqrt{3}} \log \left( \frac{\tau^2 - 1}{\tau^2 + 1} \right), \quad e^B = \frac{\sqrt{3}}{2\sqrt{\frac{4}{3}} - 1}.
\]

(6.45)

(6.46)

(6.47)

The roots of \(a(\tau)\) are \pm i and \pm 1. Hence the solution exists and is real only for \(|\tau| > 1\). So we have two real identical branches of the solution in the range \([-\infty, -1]\) and in the range \([1, +\infty]\). The cosmic time can be explicitly integrated and admits the following expression in terms of a generalized hypergeometric function:

\[
\int_1^T \frac{\sqrt{3}}{2\sqrt{\frac{4}{3}} - 1} \, d\tau = \frac{1}{2} \sqrt{3} \left( -\frac{3F_2(1, 1, \frac{5}{3}; 2, 2; \frac{1}{T^4}}{16T^4} + \log(T) + \frac{1}{8} (-\pi + \log(64)) \right).
\]

(6.48)

The behavior of the scale factor and of the scalar field for the separatrix solution are displayed in Fig. 19.

6.8. Analysis of a solution with four real roots

Next we analyze the exact solution with parameters \((\alpha, \beta) = (2, 0)\). In this case the solution has the following form:

\[
a(\tau) = 2\sqrt{\frac{4}{3}} - 1, \quad \phi(\tau) = \frac{1}{4\sqrt{3}} \log \left( \frac{\tau^2 - 1}{\tau^2 + 1} \right), \quad e^B = \frac{\sqrt{3}}{2\sqrt{\frac{4}{3}} - 1}.
\]
Fig. 19. Behavior of the scale factor and of the scalar field in the case of the separatrix solution. As one can see we have a climbing scalar that from minus infinity goes asymptotically to its extremum value $\phi = 0$, while the scale factor has an asymptotic exponential behavior as in most of the other regular solutions that go the fixed point, the specialty of this solution is visible only in the phase portrait.

\[
a(\tau) = \frac{2 \sqrt{\tau^4 - 24\tau^2 + 16}}{\sqrt{3}}, \tag{6.49}
\]

\[
\phi(\tau) = \frac{1}{4} \sqrt{3} \log \left( \frac{\tau^2 + 4\tau - 4}{\tau^2 - 4\tau - 4} \right), \tag{6.50}
\]

\[
e^B = \frac{\sqrt{3}}{2 \sqrt{\tau^4 - 24\tau^2 + 16}} \tag{6.51}
\]

and the scale factor admit four real roots, namely:

\[
\lambda_{1,2,3,4} = \{2 + 2\sqrt{2}, -2 + 2\sqrt{2}, 2 - 2\sqrt{2}, -2 - 2\sqrt{2}\}. \tag{6.52}
\]

The scalar factor is real in the intervals $[-\infty, \lambda_4], [\lambda_3, \lambda_2]$ and $[\lambda_1, +\infty]$. In the two identical branches (it suffices to change $\tau \leftrightarrow -\tau$) $[-\infty, \lambda_4]$ and $[\lambda_1, +\infty]$, the solution reaches the fixed point and it is asymptotically de Sitter (exponential increase of the scale factor). On the other hand in the branch $[\lambda_3, \lambda_2]$, we might think that we have a solution that begins with a Big Bang and ends up with Big Crunch (the Universe collapses) in a finite amount of cosmic time. Yet this branch of the functions (6.49) simply does not satisfy Friedman equations and such an embarrassing solution does not exist! The two phase trajectory composing the solution and the fake solution are displayed in Fig. 20. The physical branch of this solution which is connected with the fixed point and has an asymptotically de Sitter behavior is displayed in Fig. 21.
Fig. 20. In this figure we display the two trajectories in $(\phi, v)$-plane corresponding to the solution of parameter $(\alpha, \beta) = (2, 0)$. One branch connects infinity with infinity and never reaches the fixed point. This branch is actually fake since in this range of the parametric time Friedman equations are not satisfied! The other physical branch which satisfies Friedman equation, connects infinity with the fixed point which is reached along the universal tangent vector of all solutions except the separatrix.

6.9. Numerical simulations

In this section our goal is to explore the phase portrait and the behavior of the solutions of Friedman equations for a few different values of $\omega$ which corresponds to different fixed point types. We consider the following four cases:

\[ \omega = \begin{array}{l}
1 \\
\frac{2}{\sqrt{3}} \\
\sqrt{3} \\
3 \\
\end{array} \]

Node
Node & integrable
focus & integrable
focus

and we make a comparison of their behavior by solving numerically the Friedman equations with the same initial conditions in the four cases. We cannot choose exactly $a(0) = 0$ since this corresponds to a singularity, so we just choose $a(0)$ quite small and $\phi(0)$ quite large. A precise way of choosing the initial conditions to be applied to all four cases can be provided by the analytic solution determined by the integrable cases. This time we use the solution Eq. (6.19) of the integral model $\omega = \sqrt{3}$ characterized by parameters:

\[ \lambda = 1, \quad \theta = \pi, \quad \xi = \frac{\pi}{6} \]

we obtain:

\[ a(0) = 2^{-1/3}, \quad \phi(0) = \frac{\log \left( \frac{1+\sqrt{3}}{1-\sqrt{3}} \right)}{\sqrt{3}}, \quad \dot{\phi}(0) = -2 \]

which will be the initial values for the integration programme in all four cases (6.53). The result is provided in a series of figures.
We begin with the node case $\omega = 1$ whose behavior is displayed in Figs. 22 and 23. We continue with the integrable case $\omega = \frac{3}{\sqrt{3}}$ whose behavior is displayed in Figs. 24 and 25. Next we consider the integrable case $\omega = \sqrt{3}$ whose behavior, already of the focus type, is displayed in Figs. 26 and 27. Finally we consider the case $\omega = 3$ which is not integrable, yet shares with the integrable case $\omega = \sqrt{3}$ the quality of its fixed point, namely the focus type. The behavior is displayed in Fig. 28 and in Fig. 29. The conclusion of this comparison is that what we can learn from the integrable models is the behavior of solutions that share the same type of fixed point. Also the asymptotic behavior for very late times of both the scalar field and of the scale factor is captured by the integrable case and is shared by all other members of the family. The structure of the scalar field behavior at finite times is rather strongly dependent from the value of $\omega$. For sufficiently large $\omega$ we get the focus case and the scalar oscillates around the fixed point. The larger is $\omega$ the more the trajectory winds around the fixed point. This winding corresponds to oscillations of the scalar field which in Cosmology are potentially interesting since they might be at the heart of the reheating mechanism after inflation.

The structure of possible trajectories is summarized in Fig. 30 that plots together the behavior of the scalar field for the various considered values of $\omega$ and so does for the scale factors.
Fig. 22. Setting the initial conditions Eq. (6.55), in this figure we display the evolution of the scale factor $a(\tau)$ and of the scalar field $\phi(\tau)$ for the case $\omega = 1$, where the fixed point is of the node type. As we see we just have a descending scalar that goes smoothly to the fixed value with no oscillations.

Fig. 23. In this figure we display the trajectory determined by the initial conditions of Eq. (6.55) in the phase space $(\phi, v)$ for the case $\omega = 1$. This trajectory goes from infinity to the fixed point without winding around it (node) and following the direction fixed by the linearization matrix.
Fig. 24. Setting the initial conditions Eq. (6.55), in this figure we display the evolution of the scale factor $a(\tau)$ and of the scalar field $\phi(\tau)$ for the case $\omega = \frac{2}{\sqrt{3}}$, which is actually integrable. Its analytic behavior was extensively analyzed before.

Fig. 25. In this figure we display the trajectory determined by the initial conditions of Eq. (6.55) in the phase space $(\phi, v)$ for the case $\omega = \frac{2}{\sqrt{3}}$. This trajectory goes from infinity to the fixed point without winding around it: Node.
Fig. 26. Setting the initial conditions Eq. (6.55), in this figure we display the evolution of the scale factor \(a(\tau)\) and of the scalar field \(\phi(\tau)\) for the case \(\omega = \sqrt{3}\), where the fixed point is of the focus type. This is an integrable case for which we possess the analytical solutions. As we see we have a descending scalar that goes to a negative-valued minimum and then climbs again to its fixed point value at zero.

Fig. 27. In this figure we display the trajectory determined by the initial conditions of Eq. (6.55) in the phase space \((\phi, v)\) for the case \(\omega = \sqrt{3}\). This trajectory goes from infinity to its fixed point winding a little bit around it (focus).
Fig. 28. Setting the initial conditions Eq. (6.55), in this figure we display the evolution of the scale factor $a(\tau)$ and of the scalar field $\phi(\tau)$ for the case $\omega = 3$, where the fixed point is of the focus type. As we see we have a descending scalar that goes to a negative-valued passing through various oscillations.

Fig. 29. In this figure we display the trajectory determined by the initial conditions of Eq. (6.55) in the phase space $(h = \phi, v)$ for the case $\omega = 3$. This trajectory goes from infinity to its fixed point winding few times around it (focus).
7. A brief scan of $\mathcal{N} = 1$ superpotentials from flux compactifications

A lot of progress has been made since the year 2000 in the context of flux compactifications of string theory with the aim of obtaining four-dimensional effective theories with phenomenologically desirable features, among which, after the discovery of the Universe late-time acceleration, with high relevance, ranks the need to find de Sitter vacua.

A vast literature in this field deals with the search of flux backgrounds that are compatible with minimal $\mathcal{N} = 1$ supersymmetry in $D = 4$ and relies on the mechanism, firstly discovered in [23] of inducing an effective $\mathcal{N} = 1$ superpotential from fluxes. This token has been extensively utilized in those compactifications that give rise to an STU-model as low energy description [42–44,46,45]. The final outcome of these procedures, which take their motivation in orbifold and orientifold compactifications [14], is just an explicit expression for the superpotential $W(S, T, U)$, which can be used to calculate the scalar potential, whose extrema and consistent one-field truncations can be studied in a systematic way. The marvelous thing of this description is that each coefficient in the development of $W(S, T, U)$ in the constituent fields, $S, T, U$ has a direct interpretation in terms of fluxes and, in an appropriate basis, it admits only quantized integer values.
Recently a series of papers has appeared \cite{48,47,49} aiming at a systematic charting of the landscape of these superpotentials and of the extrema of their corresponding scalar potentials. In the present section we briefly consider such a landscape with the aim to single out consistent one-dilaton truncations and work out the corresponding potential $V(\phi)$, to be compared with our list of integrable ones. As a result we obtain several examples of one field potentials, all falling into the general family of linear combinations of exponentials that we consider, yet very seldom satisfying the severe relations on the exponents and the coefficients that are required for integrability. Although in this run we identified only one new integrable model, the lesson that we learn from it is that, by considering truncations of multi-field supergravities, the variety of possible outcomes is significantly enlarged. Indeed what matters are the intrinsic indices $\omega_i = p_i / \sqrt{3q}$, the numbers $p_i$ being the exponent coefficients of the field $\phi$ that is kept in the truncation, while $q$ is the coefficient of the kinetic term of the latter. In view of this, when $\phi$ happens to be a linear combination of several other dilatons, the coefficients of the linear combination play a role, both in generating a variety of $p_i$'s and in giving rise to non-standard $q$'s, the final result being difficult to be predicted a priori. On the other hand, the coefficients of the appropriate linear combination that can constitute a consistent truncation, are searched for by diagonalizing the mass matrix of the theory in the vicinity of an extremum. Indeed the mass eigenstates around an exact vacuum of the theory constitute a natural basis where some fields can be put consistently to zero with the exception of one which survives.

We plan to use the landscape chartered by Dibitetto et al. as a mean to illustrate the above ideas.

7.1. The STU playing ground

Originally, in orbifold compactifications on $T^6 / \mathbb{Z}_2 \times \mathbb{Z}_2$, one arrives at seven complex moduli fields:

$$S, T_1, T_2, T_3, U_1, U_2, U_3$$

(7.1)

the first being related to the original dilaton and Kalb–Ramond field of 10-dimensional supergravity, the remaining six being the appropriate complexifications of the six radii of $T^6$. Embedding these fields into $\mathcal{N} = 4$ supergravity, which is a necessary intermediate step when supersymmetry is halved by the orbifold projection, Dibitetto et al. \cite{47} have been able to reduce the playing ground for the search of flux induced superpotentials to three fields:

$$S, T = T_1 = T_2 = T_3, U = U_1 = U_2 = U_3.$$  

(7.2)

This is done by consistently truncating $\mathcal{N} = 4$ supergravity to the singlets with respect to an appropriately chosen global $\text{SO}(3)$ symmetry group. This truncation breaks supersymmetry to $\mathcal{N} = 1$. As a result of the identifications in Eq. (7.2) the Kähler potential for the residual fields takes the following form:

$$\mathcal{K} = - \log[-i(S - \bar{S})] - \log[i(T - \bar{T})^3] - \log[i(U - \bar{U})^3].$$

(7.3)

Next the authors of \cite{47} have identified a list of 32 monomials out of which the superpotential can be constructed as a linear combination. Assigning a real coefficient $\lambda_1, \ldots, 16$ to the even powers and an imaginary coefficient $i\mu_1, \ldots, 16$ to the odd ones, one guarantees a priori that the truncation to zero axions for all the three fields is consistent. With this proviso, the most general $\mathcal{N} = 1$ superpotential considered by Dibitetto et al. is the following one:
The interesting point is that each of $\lambda$ and each of $\mu$ has a precise interpretation in terms of 10-dimensional fluxes of various type.

A completely general approach to the study of vacua and possible one-field truncations would consist of the following precise algorithm:

**Truncation Charting Algorithm from Flux Superpotentials (TCAFS)**

1. Calculate from $W_{\text{gen}}$ the scalar potential $V_{\text{gen}}(\lambda, \mu, h_{1,2,3})$ depending on three dilatons, three axions and 32 real parameters $\{\lambda, \mu\}$.
2. Consistently truncate this potential to zero axions $\hat{V}_{\text{gen}}(\lambda, \mu, h_{1,2,3}) = V_{\text{gen}}(\lambda, \mu, h_{1,2,3}, 0)$.
3. Calculate the three derivatives of the potential with respect to the remaining dilatons $\partial h_i \hat{V}$ and impose that they are zero at $h_{1,2,3} = 0$:

$$\partial h_i \hat{V}(\lambda, \mu, h_{1,2,3})|_{h_{1,2,3}=0} = 0.$$  \hfill (7.5)

These conditions impose that the base point of the manifold $S = T = U = i$ should be an extremum of the potential. This choice corresponds to no loss of generality, since the dilatons are defined up to a translation and any extremum can be mapped into the reference point $S = T = U = i$ at the prize of rescaling some of the coefficients $\{\lambda, \mu\}$. Hence, as long as we keep $\{\lambda, \mu\}$ general we do not loose anything by deciding a priori where the extremum should be located. In this way Eq. (7.5) become a set of three algebraic equations of higher order for the coefficients $\{\lambda, \mu\}$.

4. Solve, if possible, the algebraic equations (7.5). In principle this results into a set of $m$ solutions:

$$\lambda_i = \lambda_i^{(\alpha)}, \quad \mu_i = \mu_i^{(\alpha)}, \quad \alpha = 1, \ldots, m,$$  \hfill (7.6)

5. Replace one by one the solutions (7.6) into $\hat{V}_{\text{gen}}(\lambda, \mu, h_{1,2,3})$, obtaining $m$ potentials of the three dilaton fields:

$$V^{(\alpha)}(h_1, h_2, h_3) \equiv \hat{V}_{\text{gen}}(\lambda^{(\alpha)}, \mu^{(\alpha)}, h_{1,2,3}), \quad \alpha = 1, \ldots, m$$  \hfill (7.7)

which, by construction, have an extremum in $h_{1,2,3} = 0$. Verify whether each extremum corresponds to Minkowski ($V^{(\alpha)}(0) = 0$), de Sitter ($V^{(\alpha)}(0) > 0$) or anti-de Sitter ($V^{(\alpha)}(0) < 0$) space.

6. For each potential $V^{(\alpha)}(h_1, h_2, h_3)$ calculate the mass-matrix in the extremum:

$$M_{ij}^{(\alpha)} = \partial_i \partial_j V^{(\alpha)}(h_1, h_2, h_3)|_{h_{1,2,3}=0}, \quad \alpha = 1, \ldots, m$$  \hfill (7.8)

and the corresponding eigenvalues $\Lambda_i^{(\alpha)} (I = 1, 2, 3)$ and eigenvectors $\nu_i^{(\alpha)}$. 

\[ W_{\text{gen}} = \lambda_1 + U^2 \lambda_2 + SU \lambda_3 + SU^3 \lambda_4 + TU \lambda_5 + TU^3 \lambda_6 + ST \lambda_7 + STU^2 \lambda_8 + T^3 U^3 \lambda_9 + T^3 U \lambda_{10} + ST^3 U^2 \lambda_{11} + ST^3 \lambda_{12} + T^2 U^2 \lambda_{13} + T^2 \lambda_{14} + ST^2 U^3 \lambda_{15} + ST^2 U \lambda_{16} + U \mu_1 + iU^3 \mu_2 + iSU^2 \mu_3 + iSU^2 \mu_4 + iTU^2 \mu_5 + iT U^2 \mu_6 + iSTU \mu_7 + iSTU^3 \mu_8 + iT^3 U^2 \mu_9 + iT^3 \mu_{10} + iST^3 U^3 \mu_{11} + iST^3 U \mu_{12} + iT^2 U^3 \mu_{13} + iT^2 U \mu_{14} + iST^2 U^2 \mu_{15} + iST^2 \mu_{16}. \]  \hfill (7.4)
7. From the eigenvalues $\Lambda^{(\alpha)}_I$ ($I = 1, 2, 3$) we learn about stability or instability of the corresponding vacuum. By means of the corresponding eigenvectors introduce a new basis of three fields well-adapted to the potential $V^{(\alpha)}$:

$$\phi^{(\alpha)}_I \equiv v^{(\alpha)}_I \cdot \vec{h}, \quad \alpha = 1, \ldots, m, \; I = 1, 2, 3. \quad (7.9)$$

8. Transform the potential and the kinetic term to the new well adapted basis and inspect if truncation to any of the $\phi_{1,2,3}$, by setting to zero the other two is consistent.

9. In case of positive answer to the previous question calculate the effective coefficient $q$ in the kinetic term and by means of the transformation (2.2) produce a potential $V(\varphi)$ to be compared with the list of integrable ones.

The problem with the above algorithm is simply computational. Using all of the 32 terms in the superpotential and truncating to zero axions we are left with a three-dilaton potential that contains 480 terms and the three algebraic equations (7.5) in 32 unknowns are too large to be solved by standard codes in MATHEMATICA. Some strategy to reduce the parameter space has to be found. What we were able to do with ease was just to test the TCAFS on some reduced space suggested by the special superpotentials reviewed in [47]. We did not assume the coefficients presented in that paper, we simply restricted the parameter space to that spanned by the monomials included in each of these superpotentials and by running the TACFS algorithm on such parameter space we retrieved exactly the same results presented in [47]. For all these extrema we have also calculated the mass-matrix and we have found some consistent one field truncations for which we could determine the corresponding one-field potential. As anticipated they all fall in the family of linear combination of exponentials and although none coincides with an integrable one we start seeing new powers and new structures that in the one-field constructions were absent.

Finally with some ingenuity we were able to derive new interesting instances of superpotentials that lead to interesting one-field truncations. In one case we obtained a new instance of a supersymmetric integrable cosmological model.

### 7.2. Locally geometric flux induced superpotentials

In [47], the authors consider a particular superpotential that is denominated *locally geometric* since its origin is claimed to arise from a combination of geometric type IIB fluxes with non-geometric ones, the resulting composition still admitting a *locally geometric* description. Leaving aside the discussion of its ten-dimensional origin in type IIB or type IIA compactifications the above mentioned superpotential has the following form:

$$W_{\text{locgeo}} = \lambda_1 + SU^3\lambda_4 + TU\lambda_5 + STU^2\lambda_8 \quad (7.10)$$

and corresponds to a truncation of the 32-dimensional parameter space to a four-dimensional one spanned by $\{\lambda_1, \lambda_4, \lambda_5, \lambda_8\}$. Using the standard parametrization of the fields:

$$S = i \exp[h_1] + b_1, \quad T = i \exp[h_2] + b_2, \quad U = i \exp[h_3] + b_3 \quad (7.11)$$

and implementing the first two steps of the TCAFS we obtain the following potential:

$$V = \frac{1}{64} e^{-h_1-3h_2-3h_3} \lambda_1^2 - \frac{1}{32} e^{-3h_2} \lambda_1 \lambda_4 + \frac{1}{64} e^{h_1-3h_2+3h_3} \lambda_4^2 + \frac{1}{32} e^{-2h_2+h_3} \lambda_5 \lambda_8 - \frac{1}{192} e^{-h_1-h_2-h_3} \lambda_5^2 - \frac{1}{32} e^{-2h_2-h_3} \lambda_1 \lambda_8 + \frac{1}{32} e^{-h_2} \lambda_5 \lambda_8 - \frac{1}{192} e^{h_1-h_2+h_3} \lambda_8^2. \quad (7.12)$$
Implementing next the steps 3 and 4 of the TCAFS we obtain five non-vanishing solutions for the \( \lambda \)-coefficients that can be displayed by writing the corresponding superpotentials:

\[
W^{\text{locgeo}}_{(2)} = 1 + 3TU + 3STU^2 - SU^3, \tag{7.13}
\]

\[
W^{\text{locgeo}}_{(3a)} = 5 - 9TU + 3STU^2 - SU^3, \tag{7.14}
\]

\[
W^{\text{locgeo}}_{(\text{Mink})} = 1 + SU^3, \tag{7.15}
\]

\[
W^{\text{locgeo}}_{(1)} = -1 - 3TU + 3STU^2 - SU^3, \tag{7.16}
\]

\[
W^{\text{locgeo}}_{(3b)} = \frac{1}{3} - TU + 3STU^2 + \frac{5SU^3}{3}. \tag{7.17}
\]

The names given to these solutions are taken from the nomenclature utilized in Table 5 of [47] since the superpotentials we found exactly correspond to those considered there, up to a multiplicative overall constant in the last case. The values in the extremum of the corresponding scalar potentials are:

\[
\left\{ V^{\text{locgeo}}_{(2)}(0), V^{\text{locgeo}}_{(3a)}(0), V^{\text{locgeo}}_{(\text{Mink})}(0), V^{\text{locgeo}}_{(1)}(0), V^{\text{locgeo}}_{(3b)}(0) \right\}
\]

\[
= \left\{ \frac{1}{16}, \frac{15}{16}, 0, -\frac{3}{16}, -\frac{5}{48} \right\}. \tag{7.18}
\]

Hence we conclude that we have one de Sitter vacuum, one Minkowski vacuum and three anti-de Sitter vacua. This is just the same result found by the authors of [47]. The corresponding dilaton potentials that give rise to such vacua at their extremum have the following explicit form:

\[
V^{\text{locgeo}}_{(2)}(\vec{h}) = -\left( -\frac{1}{32}e^{-3h_2} - \frac{9e^{-h_2}}{32} - \frac{1}{64}e^{-h_1-3h_2-3h_3} + \frac{3}{32}e^{-2h_2-h_3} \right.
\]

\[
+ \frac{3}{64}e^{-h_1-h_2-h_3} \right) + \frac{3}{32}e^{-2h_2+h_3} + \frac{3}{64}e^{h_1-h_2+h_3} - \frac{1}{64}e^{h_1-3h_2+3h_3}, \tag{7.19}
\]

\[
V^{\text{locgeo}}_{(3a)}(\vec{h}) = -\left( -\frac{5}{32}e^{-3h_2} + \frac{27e^{-h_2}}{32} - \frac{25}{64}e^{-h_1-3h_2-3h_3} + \frac{15}{32}e^{-2h_2-h_3} \right.
\]

\[
+ \frac{27}{64}e^{-h_1-h_2-h_3} - \frac{9}{32}e^{h_3-2h_2} + \frac{3}{64}e^{h_1-h_2+h_3} - \frac{1}{64}e^{h_1-3h_2+3h_3} \right), \tag{7.20}
\]

\[
V^{\text{locgeo}}_{(\text{Mink})}(\vec{h}) = \left( -\frac{1}{32}e^{-3h_2} - \frac{1}{64}e^{-h_1-3h_2-3h_3} - \frac{1}{64}e^{h_1-3h_2+3h_3} \right), \tag{7.21}
\]

\[
V^{\text{locgeo}}_{(1)}(\vec{h}) = -\left( -\frac{1}{32}e^{-3h_2} + \frac{9e^{-h_2}}{32} - \frac{1}{64}e^{-h_1-3h_2-3h_3} - \frac{3}{32}e^{-2h_2-h_3} \right.
\]

\[
+ \frac{3}{64}e^{-h_1-h_2-h_3} - \frac{3}{32}e^{h_3-2h_2} + \frac{3}{64}e^{h_1-h_2+h_3} - \frac{1}{64}e^{h_1-3h_2+3h_3}, \tag{7.22}
\]

\[
V^{\text{locgeo}}_{(3b)}(\vec{h}) = -\left( -\frac{5}{288}e^{-3h_2} + \frac{3e^{-h_2}}{32} - \frac{1}{576}e^{-h_1-3h_2-3h_3} - \frac{1}{32}e^{-2h_2-h_3} \right.
\]

\[
+ \frac{1}{192}e^{-h_1-h_2-h_3} + \frac{5}{96}e^{h_3-2h_2} + \frac{3}{64}e^{h_1-h_2+h_3} - \frac{25}{576}e^{h_1-3h_2+3h_3} \right). \tag{7.23}
\]

We continue the development of the TCAFS algorithm, case by case.
7.2.1. The dS potential

Calculating the mass matrix in the extremum of the potential $V_{(2)}^\text{locgeo}(\vec{h})$ we obtain:

$$M_{\text{mass}} = \begin{pmatrix} -\frac{1}{16} & 0 & 0 \\ 0 & 0 & 0 \\ 0 & 0 & 0 \end{pmatrix}.$$  \hspace{1cm} (7.24)

Hence we have one negative and two null eigenvalues which means that this dS vacuum is unstable. Since the mass matrix is diagonal the charge eigenstates, namely the fields coincide with the mass eigenstates and we can explore if there are consistent truncations. By direct evaluation of the derivatives we find that there are two consistent one-field truncations:

**A-truncation.** $h_1 = h_2 = 0$. In this case the residual potential is:

$$V = \frac{1}{16} e^{-3h_3} (1 - 3e^{h_3} + 3e^{2h_3}).$$

Since the kinetic term of $h_3$ has a factor $q = 3$, by means of the substitution (2.2), we obtain:

$$V(\varphi) = \frac{e^{-\varphi}}{16} - \frac{3}{16} e^{-2\varphi/3} + \frac{3 e^{-\varphi/3}}{16}$$ \hspace{1cm} (7.25)

which is not any of the integrable potentials but belongs to the same class.

**B-truncation.** $h_2 = h_3 = 0$. In this case the residual potential is:

$$V = -\frac{1}{32} (-4 + e^{-h_1} + e^{h_1}).$$

Since the kinetic term of the $h_1$ field has a factor $q = 1$, by means of the substitution (2.2), we obtain:

$$V(\varphi) = \frac{1}{16} \left( 2 - \text{Cosh} \left[ \frac{\varphi}{\sqrt{3}} \right] \right)$$ \hspace{1cm} (7.26)

which also belongs to the class of exponential potentials here considered but does not fit into any integrable series or sporadic case.

7.2.2. The AdS potential

Calculating the mass matrix in the extremum of the potential $V_{(3a)}^\text{locgeo}(\vec{h})$ we obtain:

$$M_{\text{mass}} = \begin{pmatrix} -\frac{1}{16} & -\frac{3}{4} & -\frac{3}{4} \\ -\frac{3}{4} & -3 & -\frac{3}{2} \\ -\frac{3}{4} & -\frac{3}{2} & -3 \end{pmatrix}.$$  \hspace{1cm} (7.27)

The eigenvalues of this mass-matrix are:

$$\text{Eigenval} = \left\{ \frac{1}{32} (-71 - \sqrt{6481}) > 0, \frac{3}{2} > 0, -\frac{1}{32} (-71 + \sqrt{6481}) < 0 \right\}$$ \hspace{1cm} (7.28)

showing that this anti-de Sitter vacuum is stable since all the eigenvalues satisfy the Breitenlohner–Freedman bound $\lambda_i > -\frac{45}{48}$. The corresponding mass eigenstates are the following fields
\[ \{ \phi_1, \phi_2, \phi_3 \} = \left\{ \frac{1}{24} (-73 + \sqrt{6481})h_1 + h_2 + h_3, -h_2 + h_3, \right. \\
\left. -\frac{1}{24} (73 + \sqrt{6481})h_1 + h_2 + h_3 \right\}. \tag{7.29} \]

Calculating the derivatives we verify that there is no consistent truncation of this potential.

### 7.2.3. The AdS potential 3b
Calculating the mass matrix in the extremum of the potential \( V_{\text{loc geo}}^{3b}(\vec{h}) \) we obtain:

\[ M_{\text{mass}} = -\begin{pmatrix} \frac{1}{144} & \frac{1}{12} & -\frac{1}{12} \\ \frac{1}{12} & -\frac{1}{5} & \frac{1}{6} \\ -\frac{1}{12} & \frac{1}{6} & -\frac{1}{3} \end{pmatrix}. \tag{7.30} \]

The eigenvalues of this mass-matrix are:

\[ \text{Eigenval} = \begin{cases} \frac{1}{288} (71 + \sqrt{6481}), & > 0 \\ \frac{1}{6}, & > 0 \\ \frac{1}{288} (71 - \sqrt{6481}), & < 0 \end{cases} \tag{7.31} \]

showing that also this anti-de Sitter vacuum is stable since all eigenvalues satisfy the Breitenlohner–Freedman bound \( \lambda_i > -\frac{5}{64} \). The corresponding mass eigenstates are the following fields

\[ \{ \phi_1, \phi_2, \phi_3 \} = \left\{ \frac{1}{24} (-73 + \sqrt{6481})h_1 - h_2 + h_3, h_2 + h_3, \right. \\
\left. -\frac{1}{24} (73 + \sqrt{6481})h_1 - h_2 + h_3 \right\}. \tag{7.32} \]

Calculating the derivatives we verify that there is no consistent truncation of this potential.

### 7.2.4. The AdS potential 1
Calculating the mass matrix in the extremum of the potential \( V_{\text{loc geo}}^{1}(\vec{h}) \) we obtain:

\[ M_{\text{mass}} = -\begin{pmatrix} \frac{1}{16} & 0 & 0 \\ 0 & -\frac{3}{8} & 0 \\ 0 & 0 & -\frac{3}{8} \end{pmatrix}. \tag{7.33} \]

which is diagonal and the eigenvalues are immediately read off:

\[ \text{Eigenval} = \begin{cases} -\frac{1}{16}, & < 0 \\ \frac{3}{8}, & < 0 \\ \frac{3}{8}, & > 0 \end{cases} \tag{7.34} \]

showing that also this anti-de Sitter vacuum is stable. Indeed also in this case the Breitenlohner–Freedman bound is satisfied \( \lambda_i > -\frac{5}{64} \). The mass eigenstates coincide in this case with the charge eigenstates and we have two consistent one-field truncations:

**A-truncation.** \( h_1 = h_3 = 0 \). In this case the residual potential is:

\[ V = \frac{3}{16} e^{-2h_2} - \frac{3 e^{-h_2}}{8}. \]
Since the kinetic term of $h_2$ has a factor $q = 3$, by means of the substitution (2.2), we obtain:

$$V(\varphi) = \frac{3}{16} e^{-2\varphi/3} - \frac{3}{8} e^{-\varphi/3}$$

(7.35)

which is not any of the integrable potentials but belongs to the same class.

**B-truncation.** $h_2 = h_3 = 0$. In this case the residual potential is:

$$V = -\frac{1}{32} (4 + e^{-h_1} + e^{h_1}).$$

Since the kinetic term of the $h_1$ field has a factor $q = 1$, by means of the substitution (2.2), we obtain:

$$V(\varphi) = -\frac{1}{16} \left( \cosh \left( \frac{\varphi}{\sqrt{3}} \right) + 2 \right)$$

(7.36)

which also belongs to the class of exponential potentials here considered but does not fit into any integrable series or sporadic case.

### 7.3. Another more complex example

Inspired by the superpotential presented in Eq. (5.12) of [47] we have also considered the following extension of the superpotential (7.10):

$$\hat{W}_{\text{locgeo}} = \lambda_1 + SU^3 \lambda_4 + TU \lambda_5 + STU^2 \lambda_8$$

$$+ T^3 U^3 \lambda_9 + ST^3 \lambda_{12} + T^2 U^2 \lambda_{13} + ST^2 U \lambda_{16}$$

(7.37)

which leads to a potential with 30 terms and 26 different type of exponentials. Finding all the roots of the equations that determine the existence of an extremum turned out to be difficult, yet apart from the already known solution of the previous sections we were able to find by trial and error another solution corresponding to the following superpotential which depends on the overall parameter $\lambda_4$:

$$\hat{W}_0 = -TU \lambda_4 - ST^2 U \lambda_4 - 2STU^2 \lambda_4 + 2T^2 U^2 \lambda_4 + SU^3 \lambda_4 + T^3 U^3 \lambda_4.$$  

(7.38)

The corresponding scalar potential, which for brevity we do not write, can be consistently truncated to the dilatons by setting all the axions to zero and by construction it has an extremum in $S = T = U = i$ where it takes the positive value $\frac{\lambda_4^2}{12}$. Hence, choosing the superpotential (7.38) we find a dS vacuum. Calculating the mass matrix in this extremum we find:

$$M = -\frac{\lambda_4^2}{24} \begin{pmatrix} 1 & 3 & 0 \\ 3 & -1 & 0 \\ 0 & 0 & -4 \end{pmatrix}.$$  

(7.39)

It is convenient to choose $\lambda_4 = \sqrt{24}$ and in this way the eigenvalues of $M$ have the following simple form:

$$\text{Eigenvalues}[M] \equiv \Lambda_i = \{4, \sqrt{10}, -\sqrt{10}\}.$$  

(7.40)

The presence of a negative one among the eigenvalues (7.40) shows that the constructed dS-vacuum is unstable. The eigenstates corresponding to the above eigenvalues are the following fields:
\{ \phi_1, \phi_2, \phi_3, \} = \left\{ h_3, \left( \frac{1}{3} - \frac{\sqrt{10}}{3} \right) h_1 + h_2, \left( \frac{1}{3} + \frac{\sqrt{10}}{3} \right) h_1 + h_2 \right\}. \tag{7.41}

By calculating the derivatives we find that setting \( \phi_2 = \phi_3 = 0 \) is a consistent truncation. The surviving potential has the form:

\[
V = 1 + \frac{e^{-\phi_1}}{2} + 2e^{\phi_1} - 3e^{2\phi_1} + \frac{3e^{3\phi_1}}{2}
\tag{7.42}
\]

while the kinetic term of the field \( \phi_1 \) has \( q = 3 \). By means of the transformation (2.2) the potential (7.42) is mapped into:

\[
V(\varphi) = 1 + \frac{e^{-\varphi/3}}{2} + 2e^{\varphi/3} - 3e^{2\varphi/3} + \frac{3e^{\varphi}}{2}
\tag{7.43}
\]

which is a combination of four different exponentials but does not fit into any of the integrable cases listed by us in Tables 1 and 2.

We might find still more examples, yet we think that those provided already illustrate the variety of one-field multi exponential potentials one can obtain by consistent truncations of Gauged Supergravity. In this large variety identifying, if any, combinations that perfectly match one of the integrable cases is quite difficult, in want of some strategy able to orient a priori our choices, yet with some art we were able to single out at least one of them.

7.4. A new integrable model embedded into supergravity

Working in a reduced parameter space that was determined with some inspired guessing we found the following very simple superpotential:

\[
W_{\text{integ}} = (iT^3 + 1)(SU^3 - 1)
\tag{7.44}
\]

which inserted into the formula for the scalar potential, upon consistent truncation to no axions produces the following dilatonic potential:

\[
V_{\text{dil}}(\vec{h}) = \frac{5}{32} + \frac{1}{32} e^{-3h_2} + \frac{e^{3h_2}}{32} - \frac{1}{64} e^{-h_1 - 3h_3} + \frac{1}{64} e^{-h_1 - h_2 - 3h_3}

+ \frac{1}{64} e^{-h_1 + 3h_2 - 3h_3} - \frac{1}{64} e^{h_1 + 3h_3} + \frac{1}{64} e^{-3h_2 + 3h_3} + \frac{1}{64} e^{h_1 + 3h_2 + 3h_3}. \tag{7.45}
\]

Performing the following field redefinition:

\[
h_1 \rightarrow \sqrt{3}\phi_2, \quad h_2 \rightarrow \frac{\phi_1}{\sqrt{3}}, \quad h_3 \rightarrow \phi_3
\tag{7.46}
\]

which is a rotation that preserves the form of the dilaton kinetic term:

\[
\frac{1}{2} h_1^2 + \frac{3}{2} h_2^2 + \frac{3}{2} h_3^2 \rightarrow \frac{1}{2} \phi_1^2 + \frac{3}{2} \phi_2^2 + \frac{3}{2} \phi_3^2
\tag{7.47}
\]

the dilaton potential (7.45) transforms into

\[
V_{\text{dil}}(\vec{\phi}) = \frac{5}{32} + \frac{1}{32} e^{-\sqrt{3}\phi_1} + \frac{1}{32} e^{\sqrt{3}\phi_1} - \frac{1}{64} e^{-\sqrt{3}\phi_2 - 3\phi_3} + \frac{1}{64} e^{-\sqrt{3}\phi_1 - \sqrt{3}\phi_2 - 3\phi_3}

+ \frac{1}{64} e^{\sqrt{3}\phi_1 - \sqrt{3}\phi_2 - 3\phi_3} - \frac{1}{64} e^{\sqrt{3}\phi_2 + 3\phi_3} + \frac{1}{64} e^{-\sqrt{3}\phi_1 + \sqrt{3}\phi_2 + 3\phi_3}

+ \frac{1}{64} e^{\sqrt{3}\phi_1 + \sqrt{3}\phi_2 + 3\phi_3}. \tag{7.48}
\]
The above potential has a de Sitter extremum at \( \phi_{1,2,3} = 0 \):

\[
\left. \frac{\partial}{\partial \phi_{1,2,3}} V_{\text{dil}}(\vec{\phi}) \right|_{\phi_{1,2,3}=0} = 0, \quad V_{\text{dil}}(\vec{\phi})|_{\phi_{1,2,3}=0} = \frac{1}{4} > 0. \tag{7.49}
\]

This dS vacuum is stable since the mass matrix:

\[
\text{Mass}^2 \equiv \left. \frac{\partial^2}{\partial \phi_i \partial \phi_j} V_{\text{dil}}(\vec{\phi}) \right|_{\phi_{1,2,3}=0} = \begin{pmatrix}
\frac{3}{8} & 0 & 0 \\
0 & \frac{3}{32} & \frac{3\sqrt{3}}{32} \\
0 & \frac{3\sqrt{3}}{32} & \frac{9}{32}
\end{pmatrix} \tag{7.50}
\]

has two positive and one null eigenvalue:

\[
\text{Eigenvalues \text{Mass}^2} = \left\{ \frac{3}{8}, \frac{3}{8}, 0 \right\}. \tag{7.51}
\]

The corresponding mass eigenstates are the following fields:

\[
\{\omega_1, \omega_2, \omega_3\} = \left\{ \frac{\phi_2}{\sqrt{3}} + \phi_3, \phi_1, \phi_3 - \sqrt{3}\phi_2 \right\} \tag{7.52}
\]

and transformed to the \( \omega_i \) basis the dilatonic potential (7.48) becomes:

\[
V_{\text{dil}}(\vec{\omega}) = \frac{5}{32} - \frac{1}{64} e^{-3\omega_1} - \frac{e^{3\omega_1}}{64} + \frac{1}{32} e^{-\sqrt{3}\omega_2} + \frac{1}{32} e^{\sqrt{3}\omega_2} + \frac{1}{64} e^{-3\omega_1 - \sqrt{3}\omega_2} + \frac{1}{64} e^{3\omega_1 - \sqrt{3}\omega_2} + \frac{1}{64} e^{\sqrt{3}\omega_2 - 3\omega_1} + \frac{1}{64} e^{3\omega_1 + \sqrt{3}\omega_2} \tag{7.53}
\]

which is explicitly independent from the field \( \omega_3 \), (the massless field) and can be consistently truncated to either one of the two massive modes \( \omega_{1,2} \). In terms of the mass-eigenstates the kinetic term has the following form:

\[
\text{kin} = \frac{9\omega_1^2}{8} + \frac{\omega_2^2}{2} + \frac{3\omega_3^2}{8}. \tag{7.54}
\]

In view of this and using the translation rule (2.2), the two potentials that we obtain from the two consistent truncations are:

\[
3V_{\text{dil}}(\vec{\omega})|_{\omega_1=0; \omega_2=\frac{\sqrt{3}}{3}} = \frac{3}{8} (\cosh(\varphi) + 1). \tag{7.55}
\]

\[
3V_{\text{dil}}(\vec{\omega})|_{\omega_1=\frac{2\varphi}{\sqrt{3}}; \omega_2=0} = \frac{3}{32} \left( \cosh\left(\frac{2\varphi}{\sqrt{3}}\right) + 7 \right). \tag{7.56}
\]

The potential (7.55) fits into the integrable series \( I_1 \) of Table 1 with \( C_{11} = C_{22} = C_{12} \), while the potential (7.56) associated with the second consistent truncation does not fit into any integrable series.

The lesson that we learn from this example is very much illuminating in order to appreciate the significance and the role of integrable models in the framework of supergravity. Leaving aside the axions, that should be in any case taken into account, but that we assume to be stabilized at their vanishing values, a generic solution of the supergravity field equations would involve a scalar field moving, in the current example, on the two-dimensional surface of the potential displayed in Fig. 31. Certainly the two field model is not integrable since it admits non-integrable reductions so that deriving generic solutions cannot be attained. Yet the existence of an integrable one-field
reduction implies that we can work out some special exact solution of the multi-field theory by using the integrability of a particular reduction. Conceptually this means that the emphasis on general integrals usually attached to integrability is to be dismissed in this case. The general integral of the reduction is in any case a set of particular solutions of the complete physical theory which does not capture the full extensions of initial conditions. The correct attitude is to consider the whole machinery of integrability as an algorithm to construct particular solutions of Einstein equations that might be more or less relevant depending on their properties.

8. Summary and conclusions

In the present paper we have addressed two questions

(A) Whether any of the integrable one-field cosmological models classified in [1] can be fitted into Gauged Supergravity (in the $\mathcal{N} = 1$ case we have restricted ourselves to $F$-term supergravities, where the contribution of the vector multiplets to the scalar potential is negligible) based on constant curvature scalar manifolds $G/H$ as consistent truncations to one-field of appropriate multi-field models.

(B) Whether the solutions of integrable one-field cosmological models can be used as a handy simulation of the behavior of unknown exact cosmological solutions of Friedman equations.

Both questions have received a positive answer.

As for question (A) we have shown that the embedding of integrable one-field cosmologies into Gauged Supergravity is rather difficult but not impossible. Indeed we were able to identify two independent examples of integrable potentials that can be embedded into $\mathcal{N} = 1$ supergravity, gauged by means of suitable and particularly nice superpotentials:

Model One. The first integrable supersymmetric model corresponds to $\mathcal{N} = 1$ supergravity coupled to a single Wess–Zumino multiplet with the Kähler Geometry of $SU(1,1)/U(1)$ and the following quartic superpotential:
\[ W_{\text{int1}}(z) = \frac{2}{\sqrt{5}} \left( 3z^4 + i\omega z^3 \right) \] (8.1)

and gives rise, after truncation to the dilaton, to the integrable potential of series \( I_2 \) in Table 1 with \( \gamma = \frac{2}{3} \).

**Model Two.** It is obtained within the STU-model with a pure \( \mathcal{N} = 1 \) Kähler structure (the special Kähler structure is violated) and a very specific superpotential:

\[ W_{\text{int2}}(S, T, U) = (iT^3 + 1)(SU^3 - 1) \] (8.2)

whose interpretation within flux compactification is an interesting issue to be pursued further. A consistent truncation of this \( \mathcal{N} = 1 \) model yields the integrable potential \( I_1 \) of Table 1.

An important additional result of the present paper is the complete classification of all possible \( \mathcal{N} = 2 \) gaugings of the STU model that was presented in Section 3. This classification was performed within the embedding tensor formalism by means of which we reduced the enumeration of non-abelian gaugings to the enumeration of admissible G-orbits in the \( \mathbf{W} \)-representation, the same which black-hole charges are assigned to. In each admissible orbit we have the choice of switching on the Fayet–Iliopoulos terms or keeping them zero. This yields two different gaugings for each admissible orbit. Finally one can consider pure abelian gaugings that are once again in correspondence with the G-orbits in the \( \mathbf{W} \)-representation. This classification provided two results. On one hand we verified that the only (stable) de Sitter vacuum is the one that was obtained several years ago in [26]. On the other hand we might conclude that no integrable model can be embedded in any of these gauged models.

Provisionally it follows that the very few examples of integrable cosmologies admitting a supergravity embedding are found within the \( \mathcal{N} = 1 \) framework with \( F \)-term gauging. On the other hand, utilizing the \( D \)-terms and the axial symmetric Kählerian manifolds that are in the image of the \( D \)-map, infinite series of integrable cosmological models can be embedded into \( \mathcal{N} = 1 \) supergravity [13]. We plan to pursue further the classification of all the gaugings for the \( \mathcal{N} = 2 \) models of Table 3 in order to ascertain whether these conclusion hold true in all cases or whether there are new integrable truncations [33].

As for question (B) we have addressed it concretely in the case of the cosh-model which emerges in many one-field truncations but it is integrable only for a few distinct values of the parameters. We have shown that if the non-integrable and the integrable considered models share the same type of fixed point (for instance node), then the solutions of the integrable case capture all the features of the solutions to the non-integrable model and are actually numerically very close to them. Such a demonstration is forcibly only qualitative and can be just appreciated by looking at the plots. A precise algorithm to estimate the error is so far absent.

The detailed analysis, presented in Section 5, of the space of solutions to the supersymmetric integrable **Model One** revealed a new interesting phenomenon that, to the best of our knowledge, was so far undiscovered in General Relativity. As we stressed in paper [1], whenever the scalar potential has an extremum at negative values the solutions of Friedman equation describe a Universe that, notwithstanding its spatial flatness, ends its life into a Big Crunch like a closed Universe with positive spatial curvature. This is already an interesting novelty but an even more striking one was discovered in our analysis of Section 5.3.5. The causal structure of these spatially flat collapsing Universes is significantly different from that of the closed Universe, since here the particle and event horizons do not coincide and have an interesting evolution during
the Universe life cycle. We think that the possible cosmological implications of this mechanism should be attentively considered.

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